## Some Examples of PDE's

## PDE Examples

Example 36.1 (Traffic Equation). Consider cars travelling on a straight road, i.e. $\mathbb{R}$ and let $u(t, x)$ denote the density of cars on the road at time $t$ and space $x$ and $v(t, x)$ be the velocity of the cars at $(t, x)$. Then for $J=[a, b] \subset \mathbb{R}$, $N_{J}(t):=\int_{a}^{b} u(t, x) d x$ is the number of cars in the set $J$ at time $t$. We must have

$$
\begin{aligned}
\int_{a}^{b} \dot{u}(t, x) d x & =\dot{N}_{J}(t)=u(t, a) v(t, a)-u(t, b) v(t, b) \\
& =-\int_{a}^{b} \frac{\partial}{\partial x}[u(t, x) v(t, x)] d x
\end{aligned}
$$

Since this holds for all intervals $[a, b]$, we must have

$$
\dot{u}(t, x) d x=-\frac{\partial}{\partial x}[u(t, x) v(t, x)]
$$

To make life more interesting, we may imagine that $v(t, x)=-F\left(u(t, x), u_{x}(t, x)\right)$, in which case we get an equation of the form

$$
\frac{\partial}{\partial t} u=\frac{\partial}{\partial x} G\left(u, u_{x}\right) \text { where } G\left(u, u_{x}\right)=-u(t, x) F\left(u(t, x), u_{x}(t, x)\right) \text {. }
$$

A simple model might be that there is a constant maximum speed, $v_{m}$ and maximum density $u_{m}$, and the traffic interpolates linearly between 0 (when $\left.u=u_{m}\right)$ to $v_{m}$ when $(u=0)$, i.e. $v=v_{m}\left(1-u / u_{m}\right)$ in which case we get

$$
\frac{\partial}{\partial t} u=-v_{m} \frac{\partial}{\partial x}\left(u\left(1-u / u_{m}\right)\right)
$$

Example 36.2 (Burger's Equation). Suppose we have a stream of particles travelling on $\mathbb{R}$, each of which has its own constant velocity and let $u(t, x)$ denote the velocity of the particle at $x$ at time $t$. Let $x(t)$ denote the trajectory of the particle which is at $x_{0}$ at time $t_{0}$. We have $C=\dot{x}(t)=u(t, x(t))$. Differentiating this equation in $t$ at $t=t_{0}$ implies

$$
0=\left.\left[u_{t}(t, x(t))+u_{x}(t, x(t)) \dot{x}(t)\right]\right|_{t=t_{0}}=u_{t}\left(t_{0}, x_{0}\right)+u_{x}\left(t_{0}, x_{0}\right) u\left(t_{0}, x_{0}\right)
$$

which leads to Burger's equation

$$
0=u_{t}+u u_{x}
$$

Example 36.3 (Minimal surface Equation). (Review Dominated convergence theorem and differentiation under the integral sign.) Let $D \subset \mathbb{R}^{2}$ be a bounded region with reasonable boundary, $u_{0}: \partial D \rightarrow \mathbb{R}$ be a given function. We wish to find the function $u: D \rightarrow \mathbb{R}$ such that $u=u_{0}$ on $\partial D$ and the graph of $u$, $\Gamma(u)$ has least area. Recall that the area of $\Gamma(u)$ is given by

$$
A(u)=\operatorname{Area}(\Gamma(u))=\int_{D} \sqrt{1+|\nabla u|^{2}} d x
$$

Assuming $u$ is a minimizer, let $v \in C^{1}(D)$ such that $v=0$ on $\partial D$, then

$$
\begin{aligned}
0 & =\left.\frac{d}{d s}\right|_{0} A(u+s v)=\left.\frac{d}{d s}\right|_{0} \int_{D} \sqrt{1+|\nabla(u+s v)|^{2}} d x \\
& =\left.\int_{D} \frac{d}{d s}\right|_{0} \sqrt{1+|\nabla(u+s v)|^{2}} d x \\
& =\int_{D} \frac{1}{\sqrt{1+|\nabla u|^{2}}} \nabla u \cdot \nabla v d x \\
& =-\int_{D} \nabla \cdot\left(\frac{1}{\sqrt{1+|\nabla u|^{2}}} \nabla u\right) v d x
\end{aligned}
$$

from which it follows that

$$
\nabla \cdot\left(\frac{1}{\sqrt{1+|\nabla u|^{2}}} \nabla u\right)=0
$$

Example 36.4 (Heat or Diffusion Equation). Suppose that $\Omega \subset \mathbb{R}^{n}$ is a region of space filled with a material, $\rho(x)$ is the density of the material at $x \in \Omega$ and $c(x)$ is the heat capacity. Let $u(t, x)$ denote the temperature at time $t \in[0, \infty)$ at the spatial point $x \in \Omega$. Now suppose that $B \subset \mathbb{R}^{n}$ is a "little" volume in $\mathbb{R}^{n}, \partial B$ is the boundary of $B$, and $E_{B}(t)$ is the heat energy contained in the volume $B$ at time $t$. Then

$$
E_{B}(t)=\int_{B} \rho(x) c(x) u(t, x) d x
$$

So on one hand,

$$
\begin{equation*}
\dot{E}_{B}(t)=\int_{B} \rho(x) c(x) \dot{u}(t, x) d x \tag{36.1}
\end{equation*}
$$



Fig. 36.1. A test volume $B$ in $\Omega$.
while on the other hand,

$$
\begin{equation*}
\dot{E}_{B}(t)=\int_{\partial B}\langle G(x) \nabla u(t, x), n(x)\rangle d \sigma(x), \tag{36.2}
\end{equation*}
$$

where $G(x)$ is a $n \times n$-positive definite matrix representing the conduction properties of the material, $n(x)$ is the outward pointing normal to $B$ at $x \in$ $\partial B$, and $d \sigma$ denotes surface measure on $\partial B$. (We are using $\langle\cdot, \cdot\rangle$ to denote the standard dot product on $\mathbb{R}^{n}$.)

In order to see that we have the sign correct in (36.2), suppose that $x \in \partial B$ and $\nabla u(x) \cdot n(x)>0$, then the temperature for points near $x$ outside of $B$ are hotter than those points near $x$ inside of $B$ and hence contribute to a increase in the heat energy inside of $B$. (If we get the wrong sign, then the resulting equation will have the property that heat flows from cold to hot!)

Comparing Eqs. (36.1) to (36.2) after an application of the divergence theorem shows that

$$
\begin{equation*}
\int_{B} \rho(x) c(x) \dot{u}(t, x) d x=\int_{B} \nabla \cdot(G(\cdot) \nabla u(t, \cdot))(x) d x \tag{36.3}
\end{equation*}
$$

Since this holds for all volumes $B \subset \Omega$, we conclude that the temperature functions should satisfy the following partial differential equation.

$$
\begin{equation*}
\rho(x) c(x) \dot{u}(t, x)=\nabla \cdot(G(\cdot) \nabla u(t, \cdot))(x) \tag{36.4}
\end{equation*}
$$

or equivalently that

$$
\begin{equation*}
\dot{u}(t, x)=\frac{1}{\rho(x) c(x)} \nabla \cdot(G(x) \nabla u(t, x)) \tag{36.5}
\end{equation*}
$$

Setting $g^{i j}(x):=G_{i j}(x) /(\rho(x) c(x))$ and

$$
z^{j}(x):=\sum_{i=1}^{n} \partial\left(G_{i j}(x) /(\rho(x) c(x))\right) / \partial x^{i}
$$

the above equation may be written as:

$$
\begin{equation*}
\dot{u}(t, x)=L u(t, x), \tag{36.6}
\end{equation*}
$$

where

$$
\begin{equation*}
(L f)(x)=\sum_{i, j} g^{i j}(x) \frac{\partial^{2}}{\partial x^{i} \partial x^{j}} f(x)+\sum_{j} z^{j}(x) \frac{\partial}{\partial x^{j}} f(x) \tag{36.7}
\end{equation*}
$$

The operator $L$ is a prototypical example of a second order "elliptic" differential operator.

Example 36.5 (Laplace and Poisson Equations). Laplaces Equation is of the form $L u=0$ and solutions may represent the steady state temperature distribution for the heat equation. Equations like $\Delta u=-\rho$ appear in electrostatics for example, where $u$ is the electric potential and $\rho$ is the charge distribution.

Example 36.6 (Shrodinger Equation and Quantum Mechanics).

$$
i \frac{\partial}{\partial t} \psi(t, x)=-\frac{\Delta}{2} \psi(t, x)+V(x) \psi(t, x) \text { with }\|\psi(\cdot, 0)\|_{2}=1
$$

Interpretation,

$$
\int_{A}|\psi(t, x)|^{2} d t=\text { the probability of finding the particle in } A \text { at time } t .
$$

(Notice similarities to the heat equation.)
Example 36.7 (Wave Equation). Suppose that we have a stretched string supported at $x=0$ and $x=L$ and $y=0$. Suppose that the string only undergoes vertical motion (pretty bad assumption). Let $u(t, x)$ and $T(t, x)$ denote the height and tension of the string at $(t, x), \rho_{0}(x)$ denote the density in equilibrium and $T_{0}$ be the equilibrium string tension. Let $J=[x, x+\Delta x] \subset[0, L]$,


Fig. 36.2. A piece of displace string
is the momentum of the piece of string above $J$. (Notice that $\rho_{0}(x) d x$ is the weight of the string above $x$.) Newton's equations state

$$
\frac{d M_{J}(t)}{d t}=\int_{J} u_{t t}(t, x) \rho_{0}(x) d x=\text { Force on String. }
$$

Since the string is to only undergo vertical motion we require

$$
T(t, x+\Delta x) \cos \left(\alpha_{x+\Delta x}\right)-T(t, x) \cos \left(\alpha_{x}\right)=0
$$

for all $\Delta x$ and therefore that $T(t, x) \cos \left(\alpha_{x}\right)=T_{0}$, i.e.

$$
T(t, x)=\frac{T_{0}}{\cos \left(\alpha_{x}\right)}
$$

The vertical tension component is given by

$$
\begin{aligned}
T(t, x+\Delta x) \sin \left(\alpha_{x+\Delta x}\right)-T(t, x) \sin \left(\alpha_{x}\right) & =T_{0}\left[\frac{\sin \left(\alpha_{x+\Delta x}\right)}{\sin \left(\alpha_{x+\Delta x}\right)}-\frac{\sin \left(\alpha_{x}\right)}{\cos \left(\alpha_{x}\right)}\right] \\
& =T_{0}\left[u_{x}(t, x+\Delta x)-u_{x}(t, x)\right]
\end{aligned}
$$

Finally there may be a component due to gravity and air resistance, say

$$
\text { gravity }=-\int_{J} \rho_{0}(x) d x \text { and resistance }=-\int_{J} k(x) u_{t}(t, x) d x
$$

So Newton's equations become

$$
\begin{aligned}
\int_{x}^{x+\Delta x} u_{t t}(t, x) \rho_{0}(x) d x & =T_{0}\left[u_{x}(t, x+\Delta x)-u_{x}(t, x)\right] \\
& -\int_{x}^{x+\Delta x} \rho_{0}(x) d x-\int_{x}^{x+\Delta x} k(x) u_{t}(t, x) d x
\end{aligned}
$$

and differentiating this in $\Delta x$ gives

$$
u_{t t}(t, x) \rho_{0}(x)=u_{x x}(t, x)-\rho_{0}(x)-k(x) u_{t}(t, x)
$$

or equivalently that

$$
\begin{equation*}
u_{t t}(t, x)=\frac{1}{\rho_{0}(x)} u_{x x}(t, x)-1-\frac{k(x)}{\rho_{0}(x)} u_{t}(t, x) \tag{36.8}
\end{equation*}
$$

Example 36.8 (Maxwell Equations in Free Space).

$$
\begin{aligned}
\frac{\partial \mathbf{E}}{\partial t} & =\nabla \times \mathbf{B} \\
\frac{\partial \mathbf{B}}{\partial t} & =-\nabla \times \mathbf{E} \\
\nabla \cdot \mathbf{E} & =\nabla \cdot \mathbf{B}=0 .
\end{aligned}
$$

Notice that

$$
\frac{\partial^{2} \mathbf{E}}{\partial t^{2}}=\nabla \times \frac{\partial \mathbf{B}}{\partial t}=-\nabla \times(\nabla \times \mathbf{E})=\Delta \mathbf{E}-\nabla(\nabla \cdot \mathbf{E})=\Delta \mathbf{E}
$$

and similarly, $\frac{\partial^{2} \mathbf{B}}{\partial t^{2}}=\Delta \mathbf{B}$ so that all the components of the electromagnetic fields satisfy the wave equation.

Example 36.9 (Navier - Stokes). Here $u(t, x)$ denotes the velocity of a fluid ad $(t, x), p(t, x)$ is the pressure. The Navier - Stokes equations state,

$$
\begin{align*}
\frac{\partial u}{\partial t}+\partial_{u} u & =\nu \Delta u-\nabla p+f \text { with } u(0, x)=u_{0}(x)  \tag{36.9}\\
\nabla \cdot u & =0 \text { (incompressibility) } \tag{36.10}
\end{align*}
$$

where $f$ are the components of a given external force and $u_{0}$ is a given divergence free vector field, $\nu$ is the viscosity constant. The Euler equations are found by taking $\nu=0$. Equation (36.9) is Newton's law of motion again, $F=m a$. See http://www.claymath.org for more information on this Million Dollar problem.

### 36.1 Some More Geometric Examples

Example 36.10 (Einstein Equations). Einstein's equations from general relativity are

$$
\operatorname{Ric}_{g}-\frac{1}{2} g S_{g}=T
$$

where $T$ is the stress energy tensor.
Example 36.11 (Yamabe Problem). Does there exists a metric $g_{1}=u^{4 /(n-2)} g_{0}$ in the conformal class of $g_{0}$ so that $g_{1}$ has constant scalar curvature. This is equivalent to solving

$$
-\gamma \Delta_{g_{0}} u+S_{g_{0}} u=k u^{\alpha}
$$

where $\gamma=4 \frac{n-1}{n-2}, \alpha=\frac{n+2}{n-2}, k$ is a constant and $S_{g_{0}}$ is the scalar curvature of $g_{0}$.
Example 36.12 (Ricci Flow). Hamilton introduced the Ricci - flow,

$$
\frac{\partial g}{\partial t}=\operatorname{Ric}_{g}
$$

as another method to create "good" metrics on manifolds. This is a possible solution to the 3 dimensional Poincaré conjecture, again go to the Clay math web site for this problem.

## First Order Quasi-Linear Scalar PDE

### 37.1 Linear Evolution Equations

Consider the first order partial differential equation

$$
\begin{equation*}
\partial_{t} u(t, x)=\sum_{i=1}^{n} a_{i}(x) \partial_{i} u(t, x) \text { with } u(0, x)=f(x) \tag{37.1}
\end{equation*}
$$

where $x \in \mathbb{R}^{n}$ and $a_{i}(x)$ are smooth functions on $\mathbb{R}^{n}$. Let $A(x)=$ $\left(a_{1}(x), \ldots, a_{n}(x)\right)$ and for $u \in C^{1}\left(\mathbb{R}^{n}, \mathbb{C}\right)$, let

$$
\tilde{A} u(x):=\left.\frac{d}{d t}\right|_{0} u(x+t Z(x))=\nabla u(x) \cdot A(x)=\sum_{i=1}^{n} a_{i}(x) \partial_{i} u(x)
$$

i.e. $\tilde{A}(x)$ is the first order differential operator, $\tilde{A}(x)=\sum_{i=1}^{n} a_{i}(x) \partial_{i}$. With this notation we may write Eq. (37.1) as

$$
\partial_{t} u=\tilde{A} u \text { with } u(0, \cdot)=f .
$$

The following lemma contains the key observation needed to solve Eq. (37.2).
Lemma 37.1. Let $A$ and $\tilde{A}$ be as above and $f \in C^{1}\left(\mathbb{R}^{n}, \mathbb{R}\right)$, then

$$
\begin{equation*}
\frac{d}{d t} f \circ e^{t A}(x)=\tilde{A} f \circ e^{t A}(x)=\tilde{A}\left(f \circ e^{t A}\right)(x) . \tag{37.3}
\end{equation*}
$$

Proof. By definition,

$$
\frac{d}{d t} e^{t A}(x)=A\left(e^{t A}(x)\right)
$$

and so by the chain rule

$$
\frac{d}{d t} f \circ e^{t A}(x)=\nabla f\left(e^{t A}(x)\right) \cdot A\left(e^{t A}(x)\right)=\tilde{A} f\left(e^{t A}(x)\right)
$$

which proves the first Equality in Eq. (37.3). For the second we will need to use the following two facts: 1) $e^{(t+s) A}=e^{t A} \circ e^{s Z}$ and 2) $e^{t A}(x)$ is smooth in $x$. Assuming this we find
$\frac{d}{d t} f \circ e^{t A}(x)=\left.\frac{d}{d s}\right|_{0} f \circ e^{(t+s) A}(x)=\left.\frac{d}{d s}\right|_{0}\left[f \circ e^{t A} \circ e^{s Z}(x)\right]=\tilde{A}\left(f \circ e^{t A}\right)(x)$
which is the second equality in Eq. (37.3).
Theorem 37.2. The function $u \in C^{1}(\mathcal{D}(A), \mathbb{R})$ defined by

$$
\begin{equation*}
u(t, x):=f\left(e^{t A}(x)\right) \tag{37.4}
\end{equation*}
$$

solves $E q$. (37.2). Moreover this is the unique function defined on $\mathcal{D}(A)$ which solves Eq. (37.3).

Proof. Suppose that $u \in C^{1}(\mathcal{D}(A), \mathbb{R})$ solves Eq. (37.2), then

$$
\frac{d}{d t} u\left(t, e^{-t A}(x)\right)=u_{t}\left(t, e^{-t A}(x)\right)-\tilde{A} u\left(t, e^{-t A}(x)\right)=0
$$

and hence

$$
u\left(t, e^{-t A}(x)\right)=u(0, x)=f(x)
$$

Let $\left(t_{0}, x_{0}\right) \in \mathcal{D}(A)$ and apply the previous computations with $x=e^{t A}\left(x_{0}\right)$ to find $u\left(t_{0}, x\right)=f\left(e^{t A}\left(x_{0}\right)\right)$. This proves the uniqueness assertion. The verification that $u$ defined in Eq. (37.4) solves Eq. (37.2) is simply the second equality in Eq. (37.3).

Notation 37.3 Let $e^{t \tilde{A}} f(x)=u(t, x)$ where $u$ solves Eq. (37.2), i.e.

$$
e^{t \tilde{A}} f(x)=f\left(e^{t A}(x)\right)
$$

The differential operator $\tilde{A}: C^{1}\left(\mathbb{R}^{n}, \mathbb{R}\right) \rightarrow \mathbb{C}\left(\mathbb{R}^{n}, \mathbb{R}\right)$ is no longer bounded so it is not possible in general to conclude

$$
\begin{equation*}
e^{t \tilde{A}} f=\sum_{n=0}^{\infty} \frac{t^{n}}{n!} \tilde{A}^{n} f \tag{37.5}
\end{equation*}
$$

Indeed, to make sense out of the right side of Eq. (37.5) we must know $f$ is infinitely differentiable and that the sum is convergent. This is typically not the case. because if $f$ is only $C^{1}$. However there is still some truth to Eq. (37.5). For example if $f \in C^{k}\left(\mathbb{R}^{n}, \mathbb{R}\right)$, then by Taylor's theorem with remainder,

$$
e^{t \tilde{A}} f-\sum_{n=0}^{k} \frac{t^{n}}{n!} \tilde{A}^{n} f=o\left(t^{k}\right)
$$

by which I mean, for any $x \in \mathbb{R}^{n}$,

$$
t^{-k}\left[e^{t \tilde{A}} f(x)-\sum_{n=0}^{k} \frac{t^{n}}{n!} \tilde{A}^{n} f(x)\right] \rightarrow 0 \text { as } t \rightarrow 0
$$

Example 37.4. Suppose $n=1$ and $A(x)=1, \tilde{A}(x)=\partial_{x}$ then $e^{t A}(x)=x+t$ and hence

$$
e^{t \partial_{x}} f(x)=f(x+t)
$$

It is interesting to notice that

$$
e^{t \partial_{x}} f(x)=\sum_{n=0}^{\infty} \frac{t^{n}}{n!} f^{(n)}(x)
$$

is simply the Taylor series expansion of $f(x+t)$ centered at $x$. This series converges to the correct answer (i.e. $f(x+t)$ ) iff $f$ is "real analytic." For more details see the Cauchy - Kovalevskaya Theorem in Section 39.

Example 37.5. Suppose $n=1$ and $A(x)=x^{2}, \tilde{A}(x)=x^{2} \partial_{x}$ then $e^{t A}(x)=$ $\frac{x}{1-t x}$ on $\mathcal{D}(A)=\{(t, x): 1-t x>0\}$ and hence $e^{t \tilde{A}} f(x)=f\left(\frac{x}{1-t x}\right)=u(t, x)$ on $\mathcal{D}(A)$, where

$$
\begin{equation*}
u_{t}=x^{2} u_{x} . \tag{37.6}
\end{equation*}
$$

It may or may not be possible to extend this solution, $u(t, x)$, to a $C^{1}$ solution on all $\mathbb{R}^{2}$. For example if $\lim _{x \rightarrow \infty} f(x)$ does not exist, then $\lim _{t \uparrow x} u(t, x)$ does not exist for any $x>0$ and so $u$ can not be the restriction of $C^{1}-$ function on $\mathbb{R}^{2}$. On the other hand if there are constants $c_{ \pm}$and $M>0$ such that $f(x)=c_{+}$for $x>M$ and $f(x)=c_{-}$for $x<-M$, then we may extend $u$ to all $\mathbb{R}^{2}$ by defining

$$
u(t, x)=\left\{\begin{array}{l}
c_{+} \text {if } x>0 \text { and } t>1 / x \\
c_{-} \text {if } x<0 \text { and } t<1 / x
\end{array}\right.
$$

It is interesting to notice that $x(t)=1 / t$ solves $\dot{x}(t)=-x^{2}(t)=-A(x(t))$, so any solution $u \in C^{1}\left(\mathbb{R}^{2}, \mathbb{R}\right)$ to Eq. (37.6) satisfies $\frac{d}{d t} u(t, 1 / t)=0$, i.e. $u$ must be constant on the curves $x=1 / t$ for $t>0$ and $x=1 / t$ for $t<0$. See Example 37.13 below for a more detailed study of Eq. (37.6).

Example 37.6. Suppose $n=2$.

1. If $A(x, y)=(-y, x)$, i.e. $A\binom{x}{y}=\left(\begin{array}{cc}0 & -1 \\ 1 & 0\end{array}\right)\binom{x}{y}$ then

$$
e^{t A}\binom{x}{y}=\left(\begin{array}{cc}
\cos t-\sin t \\
\sin t & \cos t
\end{array}\right)\binom{x}{y}
$$

and hence

$$
e^{t \tilde{A}} f(x, y)=f(x \cos t-y \sin t, y \cos t+x \sin t)
$$

2. If $A(x, y)=(x, y)$, i.e. $A\binom{x}{y}=\left(\begin{array}{ll}1 & 0 \\ 0 & 1\end{array}\right)\binom{x}{y}$ then

$$
e^{t A}\binom{x}{y}=\left(\begin{array}{cc}
e^{t} & 0 \\
0 & e^{t}
\end{array}\right)\binom{x}{y}
$$

and hence

$$
e^{t \tilde{A}} f(x, y)=f\left(x e^{t}, y e^{t}\right)
$$

Theorem 37.7. Given $A \in C^{1}\left(\mathbb{R}^{n}, \mathbb{R}^{n}\right)$ and $h \in C^{1}\left(\mathbb{R} \times \mathbb{R}^{n}, \mathbb{R}\right)$.

1. (Duhamel's Principle) The unique solution $u \in C^{1}(\mathcal{D}(A), \mathbb{R})$ to

$$
\begin{equation*}
u_{t}=\tilde{A} u+h \text { with } u(0, \cdot)=f \tag{37.7}
\end{equation*}
$$

is given by

$$
u(t, \cdot)=e^{t \tilde{A}} f+\int_{0}^{t} e^{(t-\tau) \tilde{A}} h(\tau, \cdot) d \tau
$$

or more explicitly,

$$
\begin{equation*}
u(t, x):=f\left(e^{t A}(x)\right)+\int_{0}^{t} h\left(\tau, e^{(t-\tau) A}(x)\right) d \tau \tag{37.8}
\end{equation*}
$$

2. The unique solution $u \in C^{1}(\mathcal{D}(A), \mathbb{R})$ to

$$
u_{t}=\tilde{A} u+h u \text { with } u(0, \cdot)=f
$$

is given by

$$
\begin{equation*}
u(t, \cdot)=e^{\int_{0}^{t} h\left(\tau, e^{(t-\tau) A}(x)\right) d \tau} f\left(e^{t A}(x)\right) \tag{37.10}
\end{equation*}
$$

which we abbreviate as

$$
\begin{equation*}
e^{t\left(\tilde{A}+M_{h}\right)} f(x)=e^{\int_{0}^{t} h\left(\tau, e^{(t-\tau) A}(x)\right) d \tau} f\left(e^{t A}(x)\right) . \tag{37.11}
\end{equation*}
$$

Proof. We will verify the uniqueness assertions, leaving the routine check the Eqs. (37.8) and (37.9) solve the desired PDE's to the reader. Assuming $u$ solves Eq. (37.7), we find

$$
\begin{aligned}
\frac{d}{d t}\left[e^{-t \tilde{A}} u(t, \cdot)\right](x) & =\frac{d}{d t} u\left(t, e^{-t A}(x)\right)=\left(u_{t}-\tilde{A} u\right)\left(t, e^{-t A}(x)\right) \\
& =h\left(t, e^{-t A}(x)\right)
\end{aligned}
$$

and therefore

$$
\left[e^{-t \tilde{A}} u(t, \cdot)\right](x)=u\left(t, e^{-t A}(x)\right)=f(x)+\int_{0}^{t} h\left(\tau, e^{-\tau A}(x)\right) d \tau
$$

and so replacing $x$ by $e^{t A}(x)$ in this equation implies

$$
u(t, x)=f\left(e^{t A}(x)\right)+\int_{0}^{t} h\left(\tau, e^{(t-\tau) A}(x)\right) d \tau
$$

Similarly if $u$ solves Eq. (37.9), we find with $z(t):=\left[e^{-t \tilde{A}} u(t, \cdot)\right](x)=$ $u\left(t, e^{-t A}(x)\right)$ that

$$
\begin{aligned}
\dot{z}(t) & =\frac{d}{d t} u\left(t, e^{-t A}(x)\right)=\left(u_{t}-\tilde{A} u\right)\left(t, e^{-t A}(x)\right) \\
& =h\left(t, e^{-t A}(x)\right) u\left(t, e^{-t A}(x)\right)=h\left(t, e^{-t A}(x)\right) z(t)
\end{aligned}
$$

Solving this equation for $z(t)$ then implies

$$
u\left(t, e^{-t A}(x)\right)=z(t)=e^{\int_{0}^{t} h\left(\tau, e^{-\tau A}(x)\right) d \tau} z(0)=e^{\int_{0}^{t} h\left(\tau, e^{-\tau A}(x)\right) d \tau} f(x)
$$

Replacing $x$ by $e^{t A}(x)$ in this equation implies

$$
u(t, x)=e^{\int_{0}^{t} h\left(\tau, e^{(t-\tau) A}(x)\right) d \tau} f\left(e^{t A}(x)\right)
$$

Remark 37.8. It is interesting to observe the key point to getting the simple expression in Eq. (37.11) is the fact that

$$
e^{t \tilde{A}}(f g)=(f g) \circ e^{t A}=\left(f \circ e^{t A}\right) \cdot\left(g \circ e^{t A}\right)=e^{t \tilde{A}} f \cdot e^{t \tilde{A}} g
$$

That is to say $e^{t \tilde{A}}$ is an algebra homomorphism on functions. This property does not happen for any other type of differential operator. Indeed, if $L$ is some operator on functions such that $e^{t L}(f g)=e^{t L} f \cdot e^{t L} g$, then differentiating at $t=0$ implies

$$
L(f g)=L f \cdot g+f \cdot L g
$$

i.e. $L$ satisfies the product rule. One learns in differential geometry that this property implies $L$ must be a vector field.

Let us now use this result to find the solution to the wave equation

$$
\begin{equation*}
u_{t t}=u_{x x} \text { with } u(0, \cdot)=f \text { and } u_{t}(0, \cdot)=g \tag{37.12}
\end{equation*}
$$

To this end, let us notice the $u_{t t}=u_{x x}$ may be written as

$$
\left(\partial_{t}-\partial_{x}\right)\left(\partial_{t}+\partial_{x}\right) u=0
$$

and therefore noting that

$$
\left.\left(\partial_{t}+\partial_{x}\right) u(t, x)\right|_{t=0}=g(x)+f^{\prime}(x)
$$

we have

$$
\left(\partial_{t}+\partial_{x}\right) u(t, x)=e^{t \partial_{x}}\left(g+f^{\prime}\right)(x)=\left(g+f^{\prime}\right)(x+t)
$$

The solution to this equation is then a consequence of Duhamel' s Principle which gives

$$
\begin{aligned}
u(t, x) & =e^{-t \partial_{x}} f(x)+\int_{0}^{t} e^{-(t-\tau) \partial_{x}}\left(g+f^{\prime}\right)(x+\tau) d \tau \\
& =f(x-t)+\int_{0}^{t}\left(g+f^{\prime}\right)(x+\tau-(t-\tau)) d \tau \\
& =f(x-t)+\int_{0}^{t}\left(g+f^{\prime}\right)(x+2 \tau-t) d \tau \\
& =f(x-t)+\int_{0}^{t} g(x+2 \tau-t) d \tau+\left.\frac{1}{2} f(x+2 \tau-t)\right|_{\tau=0} ^{\tau=t} \\
& =\frac{1}{2}[f(x+t)+f(x-t)]+\frac{1}{2} \int_{-t}^{t} g(x+s) d s
\end{aligned}
$$

The following theorem summarizes what we have proved.
Theorem 37.9. If $f \in C^{2}(\mathbb{R}, \mathbb{R})$ and $g \in C^{1}(\mathbb{R}, \mathbb{R})$, then Eq. (37.12) has a unique solution given by

$$
\begin{equation*}
u(t, x)=\frac{1}{2}[f(x+t)+f(x-t)]+\frac{1}{2} \int_{-t}^{t} g(x+s) d s \tag{37.13}
\end{equation*}
$$

Proof. We have already proved uniqueness above. The proof that $u$ defined in Eq. (37.13) solves the wave equation is a routine computation. Perhaps the most instructive way to verify that $u$ solves $u_{t t}=u_{x x}$ is to observe, letting $y=x+s$, that

$$
\begin{aligned}
\int_{-t}^{t} g(x+s) d s & =\int_{x-t}^{x+t} g(y) d y=\int_{0}^{x+t} g(y) d y+\int_{x-t}^{0} g(y) d y \\
& =\int_{0}^{x+t} g(y) d y-\int_{0}^{x-t} g(y) d y
\end{aligned}
$$

From this observation it follows that

$$
u(t, x)=F(x+t)+G(x-t)
$$

where

$$
F(x)=\frac{1}{2}\left(f(x)+\int_{0}^{x} g(y) d y\right) \text { and } G(x)=\frac{1}{2}\left(f(x)-\int_{0}^{x} g(y) d y\right)
$$

Now clearly $F$ and $G$ are $C^{2}-$ functions and

$$
\left(\partial_{t}-\partial_{x}\right) F(x+t)=0 \text { and }\left(\partial_{t}+\partial_{x}\right) G(x-t)=0
$$

so that

$$
\left(\partial_{t}^{2}-\partial_{x}^{2}\right) u(t, x)=\left(\partial_{t}-\partial_{x}\right)\left(\partial_{t}+\partial_{x}\right)(F(x+t)+G(x-t))=0
$$

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Now let us formally apply Exercise 37.45 to the wave equation $u_{t t}=u_{x x}$, in which case we should let $A^{2}=-\partial_{x}^{2}$, and hence $A=\sqrt{-\partial_{x}^{2}}$. Evidently we should take

$$
\begin{aligned}
& \cos \left(t \sqrt{-\partial_{x}^{2}}\right) f(x)=\frac{1}{2}[f(x+t)+f(x-t)] \text { and } \\
& \frac{\sin \left(t \sqrt{-\partial_{x}^{2}}\right)}{\sqrt{-\partial_{x}^{2}}} g(x)=\frac{1}{2} \int_{-t}^{t} g(x+s) d s=\frac{1}{2} \int_{x-t}^{x+t} g(y) d y
\end{aligned}
$$

Thus with these definitions, we can try to solve the equation

$$
\begin{equation*}
u_{t t}=u_{x x}+h \text { with } u(0, \cdot)=f \text { and } u_{t}(0, \cdot)=g \tag{37.14}
\end{equation*}
$$

by a formal application of Exercise 37.43. According to Eq. (37.73) we should have

$$
u(t, \cdot)=\cos (t A) f+\frac{\sin (t A)}{A} g+\int_{0}^{t} \frac{\sin ((t-\tau) A)}{A} h(\tau, \cdot) d \tau
$$

i.e.
$u(t, x)=\frac{1}{2}[f(x+t)+f(x-t)]+\frac{1}{2} \int_{-t}^{t} g(x+s) d s+\frac{1}{2} \int_{0}^{t} d \tau \int_{x-t+\tau}^{x+t-\tau} d y h(\tau, y)$.
(37.15)

An alternative way to get to this equation is to rewrite Eq. (37.14) in first order (in time) form by introducing $v=u_{t}$ to find

$$
\begin{align*}
\frac{\partial}{\partial t}\binom{u}{v} & =A\binom{u}{v}+\binom{0}{h} \text { with } \\
\binom{u}{v} & =\binom{f}{g} \text { at } t=0 \tag{37.16}
\end{align*}
$$

where

$$
A:=\left(\begin{array}{cc}
0 & 1 \\
\frac{\partial^{2}}{\partial x^{2}} & 0
\end{array}\right)
$$

A restatement of Theorem 37.9 is simply
$e^{t A}\binom{f}{g}(x)=\binom{u(t, x)}{u_{t}(t, x)}=\frac{1}{2}\binom{f(x+t)+f(x-t)+\int_{-t}^{t} g(x+s) d s}{f^{\prime}(x+t)-f^{\prime}(x-t)+g(x+t)+g(x-t)}$.
According to Du hamel's principle the solution to Eq. (37.16) is given by

$$
\binom{u(t, \cdot)}{u_{t}(t, \cdot)}=e^{t A}\binom{f}{g}+\int_{0}^{t} e^{(t-\tau) A}\binom{0}{h(\tau, \cdot)} d \tau
$$

The first component of the last term is given by

$$
\frac{1}{2} \int_{0}^{t}\left[\int_{\tau-t}^{t-\tau} h(\tau, x+s) d s\right] d \tau=\frac{1}{2} \int_{0}^{t}\left[\int_{x-t+\tau}^{x+t-\tau} h(\tau, y) d y\right] d \tau
$$

which reproduces Eq. (37.15).
To check Eq. (37.15), it suffices to assume $f=g=0$ so that

$$
u(t, x)=\frac{1}{2} \int_{0}^{t} d \tau \int_{x-t+\tau}^{x+t-\tau} d y h(\tau, y)
$$

Now

$$
\begin{aligned}
u_{t} & =\frac{1}{2} \int_{0}^{t}[h(\tau, x+t-\tau)+h(\tau, x-t+\tau)] d \tau, \\
u_{t t} & =\frac{1}{2} \int_{0}^{t}\left[h_{x}(\tau, x+t-\tau)-h_{x}(\tau, x-t+\tau)\right] d \tau+h(t, x) \\
u_{x}(t, x) & =\frac{1}{2} \int_{0}^{t} d \tau[h(\tau, x+t-\tau)-h(\tau, x-t+\tau)] \text { and } \\
u_{x x}(t, x) & =\frac{1}{2} \int_{0}^{t} d \tau\left[h_{x}(\tau, x+t-\tau)-h_{x}(\tau, x-t+\tau)\right]
\end{aligned}
$$

so that $u_{t t}-u_{x x}=h$ and $u(0, x)=u_{t}(0, x)=0$. We have proved the following theorem.

Theorem 37.10. If $f \in C^{2}(\mathbb{R}, \mathbb{R})$ and $g \in C^{1}(\mathbb{R}, \mathbb{R})$, and $h \in C\left(\mathbb{R}^{2}, \mathbb{R}\right)$ such that $h_{x}$ exists and $h_{x} \in C\left(\mathbb{R}^{2}, \mathbb{R}\right)$, then Eq. (37.14) has a unique solution $u(t, x)$ given by Eq. (37.14).

Proof. The only thing left to prove is the uniqueness assertion. For this suppose that $v$ is another solution, then $(u-v)$ solves the wave equation (37.12) with $f=g=0$ and hence by the uniqueness assertion in Theorem $37.9, u-v \equiv 0$.

### 37.1.1 A 1-dimensional wave equation with non-constant

 coefficientsTheorem 37.11. Let $c(x)>0$ be a smooth function and $\tilde{C}=c(x) \partial_{x}$ and $f, g \in C^{2}(\mathbb{R})$. Then the unique solution to the wave equation

$$
\begin{equation*}
u_{t t}=\tilde{C}^{2} u=c u_{x x}+c^{\prime} u_{x} \text { with } u(0, \cdot)=f \text { and } u_{t}(0, \cdot)=g \tag{37.17}
\end{equation*}
$$

is

$$
\begin{equation*}
u(t, x)=\frac{1}{2}\left[f\left(e^{-t C}(x)\right)+f\left(e^{t C}(x)\right)\right]+\frac{1}{2} \int_{-t}^{t} g\left(e^{s C}(x)\right) d s \tag{37.18}
\end{equation*}
$$

defined for $(t, x) \in \mathcal{D}(C) \cap \mathcal{D}(-C)$.
Proof. (Uniqueness) If $u$ is a $C^{2}-$ solution of Eq. (37.17), then

$$
\left(\partial_{t}-\tilde{C}\right)\left(\partial_{t}+\tilde{C}\right) u=0
$$

$$
\left.\left(\partial_{t}+\tilde{C}\right) u(t, x)\right|_{t=0}=g(x)+\tilde{C} f(x) .
$$

Therefore

$$
\begin{equation*}
\left(\partial_{t}+\tilde{C}\right) u(t, x)=e^{t \tilde{C}}\left(g+f^{\prime}\right)(x)=\left(g+f^{\prime}\right)\left(e^{t C}(x)\right) \tag{37.19}
\end{equation*}
$$

which has solution given by Duhamel's Principle as

$$
\begin{aligned}
u(t, x) & =e^{-t \grave{A}} f(x)+\int_{0}^{t} e^{-(t-\tau) \tilde{C}}(g+\tilde{C} f)\left(e^{\tau C}(x)\right) d \tau \\
& =f\left(e^{-t C}(x)\right)+\int_{0}^{t}(g+\tilde{C} f)\left(e^{(2 \tau-t) C}(x)\right) d \tau \\
& =f\left(e^{-t C}(x)\right)+\frac{1}{2} \int_{-t}^{t}(g+\tilde{C} f)\left(e^{s C}(x)\right) d s \\
& =f\left(e^{-t C}(x)\right)+\frac{1}{2} \int_{-t}^{t} g\left(e^{s C}(x)\right) d s+\frac{1}{2} \int_{-t}^{t} \frac{d}{d s} f\left(e^{s C}(x)\right) d s \\
& =\frac{1}{2}\left[f\left(e^{-t C}(x)\right)+f\left(e^{t C}(x)\right)\right]+\frac{1}{2} \int_{-t}^{t} g\left(e^{s C}(x)\right) d s .
\end{aligned}
$$

(Existence.) Let $y=e^{s C}(x)$ so $d y=c\left(e^{s C}(x)\right) d s=c(y) d s$ in the integral in Eq. (37.18), then

$$
\begin{aligned}
\int_{-t}^{t} g\left(e^{s C}(x)\right) d s & =\int_{e^{-t C}(x)}^{e^{t C}(x)} g(y) \frac{d y}{c(y)} \\
& =\int_{0}^{e^{t C}(x)} g(y) \frac{d y}{c(y)}+\int_{e^{-t C}(x)}^{0} g(y) \frac{d y}{c(y)} \\
& =\int_{0}^{e^{t C}(x)} g(y) \frac{d y}{c(y)}-\int_{0}^{e^{-t C}(x)} g(y) \frac{d y}{c(y)}
\end{aligned}
$$

From this observation, it follows follows that

$$
u(t, x)=F\left(e^{t C}(x)\right)+G\left(e^{-t C}(x)\right)
$$

where

$$
\begin{equation*}
F(x)=\frac{1}{2}\left(f(x)+\int_{0}^{x} g(y) \frac{d y}{c(y)}\right) \text { and } G(x)=\frac{1}{2}\left(f(x)-\int_{0}^{x} g(y) \frac{d y}{c(y)}\right) . \tag{37.20}
\end{equation*}
$$

Now clearly $F$ and $G$ are $C^{2}$ - functions and

$$
\begin{equation*}
\left(\partial_{t}-\tilde{C}\right) F\left(e^{t C}(x)\right)=0 \text { and }\left(\partial_{t}+\tilde{C}\right) G\left(e^{-t C}(x)\right)=0 \tag{37.21}
\end{equation*}
$$

$$
\left(\partial_{t}^{2}-\tilde{C}^{2}\right) u(t, x)=\left(\partial_{t}-\tilde{C}\right)\left(\partial_{t}+\tilde{C}\right)\left[F\left(e^{t C}(x)\right)+G\left(e^{-t C}(x)\right)\right]=0
$$

By Du hamel's principle, we can similarly solve

$$
u_{t t}=\tilde{C}^{2} u+h \text { with } u(0, \cdot)=0 \text { and } u_{t}(0, \cdot)=0
$$

Corollary 37.12. The solution to Eq. (37.19) is

$$
\begin{aligned}
u(t, x) & =\frac{1}{2} \int_{0}^{t}\left(\begin{array}{c}
\text { Solution to Eq. (37.17) } \\
\text { at time } t-\tau \\
\text { with } f=0 \text { and } g=h(\tau, \cdot)
\end{array}\right) d \tau \\
& =\frac{1}{2} \int_{0}^{t} d \tau \int_{\tau-t}^{t-\tau} h\left(\tau, e^{s C}(x)\right) d s
\end{aligned}
$$

Proof. This is simply a matter of computing a number of derivatives:

$$
\begin{aligned}
u_{t} & =\frac{1}{2} \int_{0}^{t} d \tau\left[h\left(\tau, e^{(t-\tau) C}(x)\right)+h\left(\tau, e^{(\tau-t) C}(x)\right)\right] \\
u_{t t} & =h(t, x)+\frac{1}{2} \int_{0}^{t} d \tau\left[\tilde{C} h\left(\tau, e^{(t-\tau) C}(x)\right)-\tilde{C} h\left(\tau, e^{(\tau-t) C}(x)\right)\right] \\
\tilde{C} u & =\frac{1}{2} \int_{0}^{t} d \tau \int_{\tau-t}^{t-\tau} \tilde{C} h\left(\tau, e^{s C}(x)\right) d s=\frac{1}{2} \int_{0}^{t} d \tau \int_{\tau-t}^{t-\tau} \frac{d}{d s} h\left(\tau, e^{s C}(x)\right) d s \\
& =\frac{1}{2} \int_{0}^{t} d \tau\left[h\left(\tau, e^{(t-\tau) C}(x)\right)-h\left(\tau, e^{(\tau-t) C}(x)\right)\right] \text { and } \\
\tilde{C}^{2} u & =\frac{1}{2} \int_{0}^{t} d \tau\left[\tilde{C} h\left(\tau, e^{(t-\tau) C}(x)\right)-\tilde{C} h\left(\tau, e^{(\tau-t) C}(x)\right)\right] .
\end{aligned}
$$

Subtracting the second and last equation then shows $u_{t t}=\tilde{A}^{2} u+h$ and it is clear that $u(0, \cdot)=0$ and $u_{t}(0, \cdot)=0$.

### 37.2 General Linear First Order PDE

In this section we consider the following PDE,

$$
\sum_{i=1}^{n} a_{i}(x) \partial_{i} u(x)=c(x) u(x)
$$

where $a_{i}(x)$ and $c(x)$ are given functions. As above Eq. (37.20) may be written simply as

$$
\tilde{A} u(x)=c(x) u(x) .
$$

The key observation to solving Eq. (37.21) is that the chain rule implies
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which we will write briefly as

$$
\frac{d}{d s} u \circ e^{s A}=\tilde{A} u \circ e^{s A}
$$

Combining Eqs. (37.21) and (37.22) implies

$$
\frac{d}{d s} u\left(e^{s A}(x)\right)=c\left(e^{s A}(x)\right) u\left(e^{s A}(x)\right)
$$

which then gives

$$
\begin{equation*}
u\left(e^{s A}(x)\right)=e^{\int_{0}^{s} c\left(e^{\sigma A}(x)\right) d \sigma} u(x) \tag{37.23}
\end{equation*}
$$

Equation (37.22) shows that the values of $u$ solving Eq. (37.21) along any flow line of $A$, are completely determined by the value of $u$ at any point on this flow line. Hence we can expect to construct solutions to Eq. (37.21) by specifying $u$ arbitrarily on any surface $\Sigma$ which crosses the flow lines of $A$ transversely, see Figure 37.1 below.


Fig. 37.1. The flow lines of $A$ through a non-characteristic surface $\Sigma$.

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$$
u_{t}=x^{2} u_{x} \text { with } u(\lambda, \lambda)=f(\lambda)
$$

for some $f \in C^{1}(\mathbb{R}, \mathbb{R})$. The characteristic equations are given by

$$
\begin{equation*}
t^{\prime}(s)=1 \text { and } x^{\prime}(s)=-x^{2}(s) \tag{37.24}
\end{equation*}
$$

and the flow lines of this equations must live on the solution curves to $\frac{d x}{d t}=$ $-x^{2}$, i.e. on curves of the form $x(t)=\frac{1}{t-C}$ for $C \in \mathbb{R}$ and $x=0$, see Figure 37.13.


Any solution to $u_{t}=x^{2} u_{x}$ must be constant on these characteristic curves. Notice the line $x=t$ crosses each characteristic curve exactly once while the line $t=0$ crosses some but not all of the characteristic curves.
Solving Eqs. (37.24) with $t(0)=\lambda=x(0)$ gives

$$
\begin{equation*}
t(s)=s+\lambda \text { and } x(s)=\frac{\lambda}{1+s \lambda} \tag{37.25}
\end{equation*}
$$

and hence

$$
u\left(s+\lambda, \frac{\lambda}{1+s \lambda}\right)=f(\lambda) \text { for all } \lambda \text { and } s>-1 / \lambda
$$

(for a plot of some of the integral curves of Eq. (37.24).) Let

$$
\begin{equation*}
(t, x)=\left(s+\lambda, \frac{\lambda}{1+s \lambda}\right) \tag{37.26}
\end{equation*}
$$

and solve for $\lambda$ :

$$
x=\frac{\lambda}{1+(t-\lambda) \lambda} \text { or } x \lambda^{2}-(x t-1) \lambda-x=0
$$

which gives

$$
\begin{equation*}
\lambda=\frac{(x t-1) \pm \sqrt{(x t-1)^{2}+4 x^{2}}}{2 x} \tag{37.27}
\end{equation*}
$$

Now to determine the sign, notice that when $s=0$ in Eq. (37.26) we have $t=\lambda=x$. So taking $t=x$ in the right side of Eq. (37.27) implies

$$
\begin{aligned}
\frac{\left(x^{2}-1\right) \pm \sqrt{\left(x^{2}-1\right)^{2}+4 x^{2}}}{2 x} & =\frac{\left(x^{2}-1\right) \pm\left(x^{2}+1\right)}{2 x} \\
& =\left\{\begin{array}{cc}
x & \text { with }+ \\
-2 / x & \text { with }-
\end{array}\right.
\end{aligned}
$$

Therefore, we must take the plus sing in Eq. (37.27) to find

$$
\lambda=\frac{(x t-1)+\sqrt{(x t-1)^{2}+4 x^{2}}}{2 x}
$$

and hence

$$
\begin{equation*}
u(t, x)=f\left(\frac{(x t-1)+\sqrt{(x t-1)^{2}+4 x^{2}}}{2 x}\right) . \tag{37.28}
\end{equation*}
$$

When $x$ is small,

$$
\lambda=\frac{(x t-1)+(1-x t) \sqrt{1+\frac{4 x^{2}}{(x t-1)^{2}}}}{2 x} \cong \frac{(1-x t) \frac{2 x^{2}}{(x t-1)^{2}}}{2 x}=\frac{x}{1-x t}
$$

so that

$$
u(t, x) \cong f\left(\frac{x}{1-x t}\right) \text { when } x \text { is small. }
$$

Thus we see that $u(t, 0)=f(0)$ and $u(t, x)$ is $C^{1}$ if $f$ is $C^{1}$. Equation (37.28) sets up a one to one correspondence between solution $u$ to $u_{t}=x^{2} u_{x}$ and $f \in C^{1}(\mathbb{R}, \mathbb{R})$.

Example 37.14. To solve

$$
\begin{equation*}
x u_{x}+y u_{y}=\lambda x y u \text { with } u=f \text { on } S^{1}, \tag{37.29}
\end{equation*}
$$

let $A(x, y)=(x, y)=x \partial_{x}+y \partial_{y}$. The equations for $(x(s), y(s)):=e^{s A}(x, y)$ are

$$
x^{\prime}(s)=x(s) \text { and } y^{\prime}(s)=y(s)
$$

from which we learn

$$
\begin{equation*}
e^{s A}(x, y)=e^{s}(x, y) \tag{37.32}
\end{equation*}
$$

Then by Eq. (37.23),

$$
\begin{equation*}
u\left(e^{s}(x, y)\right)=e^{\lambda \int_{0}^{s} e^{2 \sigma} x y d \sigma} u(x, y)=e^{\frac{\lambda}{2}\left(e^{2 s}-1\right) x y} u(x, y) \tag{37.33}
\end{equation*}
$$

$$
u(x, y)=e^{\frac{\lambda}{2}\left(1-e^{-2 s}\right) x y} u\left(e^{-s}(x, y)\right)
$$

and then choosing $s$ so that

$$
1=\left\|e^{-s}(x, y)\right\|^{2}=e^{-2 s}\left(x^{2}+y^{2}\right)
$$

i.e. so that $s=\frac{1}{2} \ln \left(x^{2}+y^{2}\right)$. We then find

$$
u(x, y)=\exp \left(\frac{\lambda}{2}\left(1-\frac{1}{x^{2}+y^{2}}\right) x y\right) f\left(\frac{(x, y)}{\sqrt{x^{2}+y^{2}}}\right)
$$

Notice that this solution always has a singularity at $(x, y)=(0,0)$ unless $f$ is constant.


Characteristic curves for Eq. (37.29) along with the plot of $S^{1}$.

## Example 37.15. The PDE,

$$
\begin{equation*}
e^{x} u_{x}+u_{y}=u \text { with } u(x, 0)=g(x) \tag{37.30}
\end{equation*}
$$

has characteristic curves determined by $x^{\prime}:=e^{x}$ and $y^{\prime}:=1$ and along these curves solutions $u$ to Eq. (37.30) satisfy

$$
\begin{equation*}
\frac{d}{d s} u(x, y)=u(x, y) \tag{37.31}
\end{equation*}
$$

Solving these "characteristic equations" gives

$$
-e^{-x(s)}+e^{-x_{0}}=\int_{0}^{s} e^{-x} x^{\prime} d s=\int_{0}^{s} 1 d s=s
$$

so that

$$
x(s)=-\ln \left(e^{-x_{0}}-s\right) \text { and } y(s)=y_{0}+s .
$$

From Eqs. (37.32) and (37.33) one shows

$$
y(s)=y_{0}+e^{-x_{0}}-e^{-x(s)}
$$

so the "characteristic curves" are contained in the graphs of the functions

$$
y=C-e^{-x} \text { for some constant } C
$$



Some characteristic curves for Eq. (37.30). Notice that the line $y=0$ intersects some but not all of the characteristic curves. Therefore Eq. (37.30) does not uniquely determine a function $u$ defined on all of $\mathbb{R}^{2}$. On the otherhand if the intial condition were $u(0, y)=g(y)$ the method would produce an everywhere defined solution.

Since the initial condition is at $y=0$, set $y_{0}=0$ in Eq. (37.33) and notice from Eq.(37.31) that

$$
\begin{equation*}
u\left(-\ln \left(e^{-x_{0}}-s\right), s\right)=u(x(s), y(s))=e^{s} u\left(x_{0}, 0\right)=e^{s} g\left(x_{0}\right) \tag{37.34}
\end{equation*}
$$

Setting $(x, y)=\left(-\ln \left(e^{-x_{0}}-s\right), s\right)$ and solving for $\left(x_{0}, s\right)$ implies

$$
s=y \text { and } x_{0}=-\ln \left(e^{-x}+y\right)
$$

and using this in Eq. (37.34) then implies

$$
u(x, y)=e^{y} g\left(-\ln \left(y+e^{-x}\right)\right)
$$

This solution is only defined for $y>-e^{-x}$.
Example 37.16. In this example we will use the method of characteristics to solve the following non-linear equation,

$$
\begin{equation*}
x^{2} u_{x}+y^{2} u_{y}=u^{2} \text { with } u:=1 \text { on } y=2 x . \tag{37.35}
\end{equation*}
$$

As usual let $(x, y)$ solve the characteristic equations, $x^{\prime}=x^{2}$ and $y^{\prime}=y^{2}$ so that

$$
(x(s), y(s))=\left(\frac{x_{0}}{1-s x_{0}}, \frac{y_{0}}{1-s y_{0}}\right) .
$$

Now let $\left(x_{0}, y_{0}\right)=(\lambda, 2 \lambda)$ be a point on line $y=2 x$ and supposing $u$ solves Eq. (37.35). Then $z(s)=u(x(s), y(s))$ solves

$$
z^{\prime}=\frac{d}{d s} u(x, y)=x^{2} u_{x}+y^{2} u_{y}=u^{2}(x, y)=z^{2}
$$

with $z(0)=u(\lambda, 2 \lambda)=1$ and hence

$$
\begin{equation*}
u\left(\frac{\lambda}{1-s \lambda}, \frac{2 \lambda}{1-2 s \lambda}\right)=u(x(s), y(s))=z(s)=\frac{1}{1-s} . \tag{37.36}
\end{equation*}
$$

Let

$$
\begin{equation*}
(x, y)=\left(\frac{\lambda}{1-s \lambda}, \frac{2 \lambda}{1-2 s \lambda}\right)=\left(\frac{1}{\lambda^{-1}-s}, \frac{1}{\lambda^{-1} / 2-s}\right) \tag{37.37}
\end{equation*}
$$

and solve the resulting equations:

$$
\lambda^{-1}-s=x^{-1} \text { and } \lambda^{-1} / 2-s=y^{-1}
$$

for $s$ gives $s=x^{-1}-2 y^{-1}$ and hence

$$
\begin{equation*}
1-s=1+2 y^{-1}-x^{-1}=x^{-1} y^{-1}(x y+2 x-y) . \tag{37.38}
\end{equation*}
$$

Combining Eqs. (37.36) - (37.38) gives

$$
u(x, y)=\frac{x y}{x y+2 x-y} .
$$

Notice that the characteristic curves here lie on the trajectories determined by $\frac{d x}{x^{2}}=\frac{d y}{y^{2}}$, i.e. $y^{-1}=x^{-1}+C$ or equivalently

$$
y=\frac{x}{1+C x}
$$



## 37 First Order Quasi-Linear Scalar PDE

### 37.3 Quasi-Linear Equations

In this section we consider the following PDE,

$$
\begin{equation*}
A(x, z) \cdot \nabla_{x} u(t, x)=\sum_{i=1}^{n} a_{i}(x, u(x)) \partial_{i} u(x)=c(x, u(x)) \tag{37.39}
\end{equation*}
$$

where $a_{i}(x, z)$ and $c(x, z)$ are given functions on $(x, z) \in \mathbb{R}^{n} \times \mathbb{R}$ and $A(x, z):=$ $\left(a_{1}(x, z), \ldots, a_{1}(x, z)\right)$. Assume $u$ is a solution to Eq. (37.39) and suppose $x(s)$ solves $x^{\prime}(s)=A(x(s), u(x(s))$.Then from Eq. (37.39) we find

$$
\frac{d}{d s} u(x(s))=\sum_{i=1}^{n} a_{i}(x(s), u(x(s))) \partial_{i} u(x(s))=c(x(s), u(x(s)))
$$

see Figure 37.2 below. We have proved the following Lemma.


Fig. 37.2. Determining the values of $u$ by solving ODE's. Notice that potential problem though where the projection of characteristics cross in $x$ - space.

Lemma 37.17. Let $w=(x, z), \pi_{1}(w)=x, \pi_{2}(w)=z$ and $Y(w)=$ $(A(x, z), c(x, z))$. If $u$ is a solution to Eq. (37.39), then

$$
u\left(\pi_{1} \circ e^{s Y}\left(x_{0}, u\left(x_{0}\right)\right)=\pi_{2} \circ e^{s Y}\left(x_{0}, u\left(x_{0}\right)\right)\right.
$$

Let $\Sigma$ be a surface in $\mathbb{R}^{n}(x-$ space $)$, i.e. $\Sigma: U \subset_{o} \mathbb{R}^{n-1} \longrightarrow \mathbb{R}^{n}$ such that $\Sigma(0)=x_{0}$ and $D \Sigma(y)$ is injective for all $y \in U$. Now suppose $u_{0}: \Sigma \rightarrow \mathbb{R}$ is given we wish to solve for $u$ such that (37.39) holds and $u=u_{0}$ on $\Sigma$. Let

$$
\begin{equation*}
\phi(s, y):=\pi_{1} \circ e^{s Y}\left(\Sigma(y), u_{0}(\Sigma(y))\right) \tag{37.40}
\end{equation*}
$$

then

$$
\begin{aligned}
\frac{\partial \phi}{\partial s}(0,0) & =\pi_{1} \circ Y\left(x_{0}, u_{0}\left(x_{0}\right)\right)=A\left(x_{0}, u_{0}\left(x_{0}\right)\right) \text { and } \\
D_{y} \phi(0,0) & =D_{y} \Sigma(0)
\end{aligned}
$$

Assume $\Sigma$ is non-characteristic at $x_{0}$, that is $A\left(x_{0}, u_{0}\left(x_{0}\right)\right) \notin \operatorname{Ran} \Sigma^{\prime}(0)$ where $\Sigma^{\prime}(0): \mathbb{R}^{n-1} \rightarrow \mathbb{R}^{n}$ is defined by

$$
\Sigma^{\prime}(0) v=\partial_{v} \Sigma(0)=\left.\frac{d}{d s}\right|_{0} \Sigma(s v) \text { for all } v \in \mathbb{R}^{n-1}
$$

Then $\left(\frac{\partial \phi}{\partial s}, \frac{\partial \phi}{\partial y^{1}}, \ldots, \frac{\partial \phi}{\partial y^{n-1}}\right)$ are all linearly independent vectors at $(0,0) \in$ $\mathbb{R} \times \mathbb{R}^{n-1}$. So $\phi: \mathbb{R} \times \mathbb{R}^{n-1} \rightarrow \mathbb{R}^{n}$ has an invertible differential at $(0,0)$ and so the inverse function theorem gives the existence of open neighborhood $0 \in W \subset U$ and $0 \in J \subset \mathbb{R}$ such that $\left.\phi\right|_{J \times W}$ is a homeomorphism onto an open set $V:=\phi(J \times W) \subset \mathbb{R}^{n}$, see Figure 37.3. Because of Lemma 37.17, if


Fig. 37.3. Constructing a neighborhood of the surface $\Sigma$ near $x_{0}$ where we can solve the quasi-linear PDE.
we are going to have a $C^{1}$ - solution $u$ to Eq. (37.39) with $u=u_{0}$ on $\Sigma$ it would have to satisfy

$$
\begin{equation*}
u(x)=\pi_{2} \circ e^{s Y}\left(\Sigma(y), u_{0}(\Sigma(y))\right) \text { with }(s, y):=\phi^{-1}(x) \tag{37.41}
\end{equation*}
$$

i.e. $x=\phi(s, y)$.

Proposition 37.18. The function $u$ in Eq. (37.41) solves Eq. (37.39) on $V$ with $u=u_{0}$ on $\Sigma$.

Proof. By definition of $u$ in Eq. (37.41) and $\phi$ in Eq. (37.40),

$$
\phi^{\prime}(s, y)=\pi_{1} Y \circ e^{s Y}\left(\Sigma(y), u_{0}(\Sigma(y))\right)=A(\phi(s, y)), u(\phi(s, y))
$$

and

$$
\begin{equation*}
\frac{d}{d s} u(\phi(s, y))=\pi_{2} Y(\phi(s, y), u(\phi(s, y)))=c(\phi(s, y), u(\phi(s, y))) . \tag{37.42}
\end{equation*}
$$

On the other hand by the chain rule,

$$
\begin{align*}
\frac{d}{d s} u(\phi(s, y)) & =\nabla u(\phi(s, u)) \cdot \phi^{\prime}(s, y) \\
& =\nabla u(\phi(s, u)) \cdot A(\phi(s, y)), u(\phi(s, y)) \tag{37.43}
\end{align*}
$$

Comparing Eqs. (37.42) and (37.43) implies

$$
\nabla u(\phi(s, y)) \cdot A(\phi(s, y), u(\phi(s, y)))=c(\phi(s, y), u(\phi(s, y)))
$$

Since $\phi(J \times W)=V$, $u$ solves Eq. (37.39) on $V$. Clearly $u(\phi(0, y))=u_{0}(\Sigma(y))$ so $u=u_{0}$ on $\Sigma$.

Example 37.19 (Conservation Laws). Let $F: \mathbb{R} \rightarrow \mathbb{R}$ be a smooth function, we wish to consider the PDE for $u=u(t, x)$,

$$
\begin{equation*}
0=u_{t}+\partial_{x} F(u)=u_{t}+F^{\prime}(u) u_{x} \text { with } u(0, x)=g(x) \tag{37.44}
\end{equation*}
$$

The characteristic equations are given by

$$
\begin{equation*}
t^{\prime}(s)=1, \quad x^{\prime}(s)=F^{\prime}(z(s)) \text { and } \frac{d}{d s} z(s)=0 \tag{37.45}
\end{equation*}
$$

The solution to Eqs. (37.45) with $t(0)=0, x(0)=x$ and hence

$$
z(0)=u(t(0), x(0))=u(0, x)=g(x)
$$

are given by

$$
t(s)=s, z(s)=g(x) \text { and } x(s)=x+s F^{\prime}(g(x))
$$

So we conclude that any solution to Eq. (37.44) must satisfy,

$$
u\left(s, x+s F^{\prime}(g(x))=g(x)\right.
$$

This implies, letting $\psi_{s}(x):=x+s F^{\prime}(g(x))$, that

$$
u(t, x)=g\left(\psi_{t}^{-1}(x)\right)
$$

In order to find $\psi_{t}^{-1}$ we need to know $\psi_{t}$ is invertible, i.e. that $\psi_{t}$ is monotonic in $x$. This becomes the condition

$$
0<\psi_{t}^{\prime}(x)=1+t F^{\prime \prime}(g(x)) g^{\prime}(x)
$$

If this holds then we will have a solution.

Example 37.20 (Conservation Laws in Higher Dimensions). Let $F: \mathbb{R} \rightarrow \mathbb{R}^{n}$ be a smooth function, we wish to consider the PDE for $u=u(t, x)$,

$$
\begin{equation*}
0=u_{t}+\nabla \cdot F(u)=u_{t}+F^{\prime}(u) \cdot \nabla u \text { with } u(0, x)=g(x) . \tag{37.46}
\end{equation*}
$$

The characteristic equations are given by

$$
\begin{equation*}
t^{\prime}(s)=1, \quad x^{\prime}(s)=F^{\prime}(z(s)) \text { and } \frac{d}{d s} z(s)=0 . \tag{37.47}
\end{equation*}
$$

The solution to Eqs. (37.47) with $t(0)=0, x(0)=x$ and hence

$$
z(0)=u(t(0), x(0))=u(0, x)=g(x),
$$

are given by

$$
t(s)=s, z(s)=g(x) \text { and } x(s)=x+s F^{\prime}(g(x))
$$

So we conclude that any solution to Eq. (37.46) must satisfy,

$$
\begin{equation*}
u\left(s, x+s F^{\prime}(g(x))\right)=g(x) \tag{37.48}
\end{equation*}
$$

This implies, letting $\psi_{s}(x):=x+s F^{\prime}(g(x))$, that

$$
u(t, x)=g\left(\psi_{t}^{-1}(x)\right)
$$

In order to find $\psi_{t}^{-1}$ we need to know $\psi_{t}$ is invertible. Locally by the implicit function theorem it suffices to know,

$$
\psi_{t}^{\prime}(x) v=v+t F^{\prime \prime}(g(x)) \partial_{v} g(x)=\left[I+t F^{\prime \prime}(g(x)) \nabla g(x) \cdot\right] v
$$

is invertible. Alternatively, let $y=x+s F^{\prime}(g(x))$, (so $\left.x=y-s F^{\prime}(g(x))\right)$ in Eq. (37.48) to learn, using Eq. (37.48) which asserts $g(x)=u\left(s, x+s F^{\prime}(g(x))\right)=$ $u(s, y)$,

$$
u(s, y)=g\left(y-s F^{\prime}(g(x))\right)=g\left(y-s F^{\prime}(u(s, y))\right)
$$

This equation describes the solution $u$ implicitly.
Example 37.21 (Burger's Equation). Recall Burger's equation is the PDE,

$$
\begin{equation*}
u_{t}+u u_{x}=0 \text { with } u(0, x)=g(x) \tag{37.49}
\end{equation*}
$$

where $g$ is a given function. Also recall that if we view $u(t, x)$ as a time dependent vector field on $\mathbb{R}$ and let $x(t)$ solve

$$
\dot{x}(t)=u(t, x(t)),
$$

then

$$
\ddot{x}(t)=u_{t}+u_{x} \dot{x}=u_{t}+u_{x} u=0
$$

Therefore $x$ has constant acceleration and

$$
x(t)=x(0)+\dot{x}(0) t=x(0)+g(x(0)) t .
$$

This equation contains the same information as the characteristic equations Indeed, the characteristic equations are

$$
\begin{equation*}
t^{\prime}(s)=1, \quad x^{\prime}(s)=z(s) \text { and } z^{\prime}(s)=0 \tag{37.50}
\end{equation*}
$$

Taking initial condition $t(0)=0, x(0)=x_{0}$ and $z(0)=u\left(0, x_{0}\right)=g\left(x_{0}\right)$ we find

$$
t(s)=s, \quad z(s)=g\left(x_{0}\right) \text { and } x(s)=x_{0}+s g\left(x_{0}\right)
$$

According to Proposition 37.18, we must have

$$
\begin{equation*}
u\left(\left(s, x_{0}+s g\left(x_{0}\right)\right)=u(s, x(s))=u(0, x(0))=g\left(x_{0}\right)\right. \tag{37.51}
\end{equation*}
$$

Letting $\psi_{t}\left(x_{0}\right):=x_{0}+\operatorname{tg}\left(x_{0}\right)$, "the" solution to $(t, x)=\left(s, x_{0}+s g\left(x_{0}\right)\right)$ is given by $s=t$ and $x_{0}=\psi_{t}^{-1}(x)$. Therefore, we find from Eq. (37.51) that

$$
\begin{equation*}
u(t, x)=g\left(\psi_{t}^{-1}(x)\right) \tag{37.52}
\end{equation*}
$$

This gives the desired solution provided $\psi_{t}^{-1}$ is well defined.
Example 37.22 (Burger's Equation Continued). Continuing Example 37.21. Suppose that $g \geq 0$ is an increasing function (i.e. the faster cars start to the right), then $\psi_{t}$ is strictly increasing and for any $t \geq 0$ and therefore Eq. (37.52) gives a solution for all $t \geq 0$. For a specific example take $g(x)=\max (x, 0)$, then

$$
\psi_{t}(x)=\left\{\begin{array}{cc}
(1+t) x & \text { if } x \geq 0 \\
x & \text { if } x \leq 0
\end{array}\right.
$$

and therefore,

$$
\begin{gathered}
\psi_{t}^{-1}(x)=\left\{\begin{array}{c}
(1+t)^{-1} x \\
\text { if } x \geq 0 \\
x \quad \text { if } x \leq 0
\end{array}\right. \\
u(t, x)=g\left(\psi_{t}^{-1}(x)\right)=\left\{\begin{array}{c}
(1+t)^{-1} x \text { if } x \geq 0 \\
0 \quad \text { if } x \leq 0
\end{array}\right.
\end{gathered}
$$

Notice that $u(t, x) \rightarrow 0$ as $t \rightarrow \infty$ since all the fast cars move off to the right leaving only slower and slower cars passing $x \in \mathbb{R}$.

Example 37.23. Now suppose $g \geq 0$ and that $g^{\prime}\left(x_{0}\right)<0$ at some point $x_{0} \in$ $\mathbb{R}$, i.e. there are faster cars to the left of $x_{0}$ then there are to the right of $x_{0}$, see Figure 37.4. Without loss of generality we may assume that $x_{0}=$ 0 . The projection of a number of characteristics to the $(t, x)$ plane for this velocity profile are given in Figure 37.5 below. Since any $C^{2}-$ solution to Eq.(37.49) must be constant on these lines with the value given by the slope, it is impossible to get a $C^{2}$ - solution on all of $\mathbb{R}^{2}$ with this initial condition. Physically, there are collisions taking place which causes the formation of a shock wave in the system.


Fig. 37.4. An intial velocity profile where collisions are going to occur. This is the graph of $g(x)=1 /\left(1+(x+1)^{2}\right)$


Fig. 37.5. Crossing of projected characteristics for Burger's equation

### 37.4 Distribution Solutions for Conservation Laws

Let us again consider the conservation law in Example 37.19 above. We will now restrict our attention to non-negative times. Suppose that $u$ is a $C^{1}-$ solution to

$$
\begin{equation*}
u_{t}+(F(u))_{x}=0 \text { with } u(0, x)=g(x) \tag{37.53}
\end{equation*}
$$

and $\phi \in C_{c}^{2}([0, \infty) \times \mathbb{R})$. Then by integration by parts,

$$
\begin{aligned}
0 & =-\int_{\mathbb{R}} d x \int_{t \geq 0} d t\left(u_{t}+F(u)_{x}\right) \varphi \\
& =-\left.\int_{\mathbb{R}}[u \varphi]\right|_{t=0} ^{t=\infty} d x+\int_{\mathbb{R}} d x \int_{t \geq 0} d t\left(u \varphi_{t}+F(u) \phi_{x}\right) \\
& =\int_{\mathbb{R}} g(x) \varphi(0, x) d x+\int_{\mathbb{R}} d x \int_{t \geq 0} d t\left(u(t, x) \phi_{t}(t, x)+F(u(t, x)) \phi_{x}(t, x)\right)
\end{aligned}
$$

Definition 37.24. A bounded measurable function $u(t, x)$ is a distributional solution to Eq. (37.53) iff

$$
0=\int_{\mathbb{R}} g(x) \varphi(0, x) d x+\int_{\mathbb{R}} d x \int_{t \geq 0} d t\left(u(t, x) \phi_{t}(t, x)+F(u(t, x)) \phi_{x}(t, x)\right)
$$

for all test functions $\phi \in C_{c}^{2}(D)$ where $D=[0, \infty) \times \mathbb{R}$.
Proposition 37.25. If $u$ is a distributional solution of Eq. (37.53) and $u$ is $C^{1}$ then $u$ is a solution in the classical sense. More generally if $u \in C^{1}(R)$ where $R$ is an open region contained in $D^{0}:=(0, \infty) \times \mathbb{R}$ and

$$
\begin{equation*}
\int_{\mathbb{R}} d x \int_{t \geq 0} d t\left(u(t, x) \phi_{t}(t, x)+F(u(t, x)) \phi_{x}(t, x)\right)=0 \tag{37.54}
\end{equation*}
$$

for all $\phi \in C_{c}^{2}(R)$ then $u_{t}+(F(u))_{x}:=0$ on $R$.
Proof. Undo the integration by parts argument to show Eq. (37.54) implies

$$
\int_{R}\left(u_{t}+(F(u))_{x}\right) \varphi d x d t=0
$$

for all $\phi \in C_{c}^{1}(R)$. This then implies $u_{t}+(F(u))_{x}=0$ on $R$.
Theorem 37.26 (Rankine-Hugoniot Condition). Let $R$ be a region in $D^{0}$ and $x=c(t)$ for $t \in[a, b]$ be a $C^{1}$ curve in $R$ as pictured below in Figure 37.6.

Suppose $u \in C^{1}(R \backslash c([a, b]))$ and $u$ is bounded and has limiting values $u^{+}$ and $u^{-}$on $x=c(t)$ when approaching from above and below respectively. Then $u$ is a distributional solution of $u_{t}+(F(u))_{x}=0$ in $R$ if and only if

$$
\begin{equation*}
u_{t}+\frac{\partial}{\partial x} F(u):=0 \text { on } \quad R \backslash c([a, b]) \tag{37.55}
\end{equation*}
$$

and for all $t \in[a, b]$,

$$
\begin{equation*}
\dot{c}(t)\left[u^{+}(t, c(t))-u^{-}(t, c(t))\right]=F\left(u^{+}(t, c(t))\right)-F\left(u^{-}(t, c(t)) .\right. \tag{37.56}
\end{equation*}
$$



Fig. 37.6. A curve of discontinuities of $u$.

Proof. The fact that Equation 37.55 holds has already been proved in the previous proposition. For (37.56) let $\Omega$ be a region as pictured in Figure 37.6 above and $\phi \in C_{c}^{1}(\Omega)$. Then

$$
\begin{align*}
0 & =\int_{\Omega^{\prime}}\left(u \phi_{t}+F(u) \phi_{x}\right) d t d x \\
& =\int_{\Omega_{+}}\left(u \phi_{t}+F(u) \phi_{x}\right) d t d x+\int_{\Omega_{-}}\left(u \phi_{t}+F(u) \phi_{x}\right) d t d x \tag{37.57}
\end{align*}
$$

where

$$
\Omega_{ \pm}=\left\{(t, x) \in \Omega: \begin{array}{l}
x>c(t) \\
x<c(t)
\end{array}\right\}
$$

Now the outward normal to $\Omega_{ \pm}$along $c$ is

$$
n(t)= \pm \frac{(\dot{c}(t),-1)}{\sqrt{1+\dot{c}(t)^{2}}}
$$

and the "surface measure" along $c$ is given by $d \sigma(t)=\sqrt{1+\dot{c}(t)^{2}} d t$. Therefore

$$
n(t) d \sigma(t)= \pm(\dot{c}(t),-1) d t
$$

where the sign is chosen according to the sign in $\Omega_{ \pm}$. Hence by the divergence theorem,

$$
\begin{aligned}
\int_{\Omega_{ \pm}}\left(u \phi_{t}\right. & \left.+F(u) \phi_{x}\right) d t d x \\
& =\int_{\Omega_{ \pm}}(u, F(u)) \cdot\left(\phi_{t}, \phi_{x}\right) d t d x \\
& =\int_{\partial \Omega_{ \pm}} \phi(u, F(u)) \cdot n(t) d \sigma(t) \\
& = \pm \int_{\alpha}^{\beta} \phi(t, c(t))\left(u_{t}^{ \pm}(t, c(t)) \dot{c}(t)-F\left(u_{t}^{ \pm}(t, c(t))\right)\right) d t
\end{aligned}
$$

Putting these results into Eq. (37.57) gives

$$
\begin{aligned}
0 & =\int_{\alpha}^{\beta}\left\{\dot{c}(t)\left[u^{+}(t, c(t))-u^{-}(t, c(t))\right]\right. \\
& -\left(F\left(u^{+}(t, c(t))\right)-F\left(u^{-}(t, c(t))\right)\right\} \phi(t, c(t)) d t
\end{aligned}
$$

for all $\phi$ which implies

$$
\dot{c}(t)\left[u^{+}(t, c(t))-u^{-}(t, c(t))\right]=F\left(u^{+}(t, c(t))\right)-F\left(u^{-}(t, c(t))\right.
$$

Example 37.27. In this example we will find an integral solution to Burger's Equation, $u_{t}+u u_{x}=0$ with initial condition

$$
u(0, x)= \begin{cases}0 & x \geq 1 \\ 1-x & 0 \leq x \leq 1 \\ 1 & x \leq 0\end{cases}
$$

The characteristics are given from above by

$$
\begin{aligned}
& x(t)=\left(1-x_{0}\right) t+x_{0} \quad x_{0} \in(0,1) \\
& x(t)=x_{0}+t \text { if } x_{0} \leq 0 \text { and } \\
& x(t)=x_{0} \text { if } x_{0} \geq 1
\end{aligned}
$$



## Projected characteristics

For the region bounded determined by $t \leq x \leq 1$ and $t \leq 1$ we have $u(t, x)$ is equal to the slope of the line through $(t, x)$ and $(1,1)$, i.e.

$$
u(t, x)=\frac{x-1}{t-1}
$$

Notice that the solution is not well define in the region where characteristics cross, i.e. in the shock zone

$$
R_{2}:=\{(t, x): t \geq 1, x \geq 1 \text { and } x \leq t\}
$$

see Figure 37.7. Let us now look for a distributional solution of the equation


Fig. 37.7. The schock zone and the values of $u$ away from the shock zone.
valid for all $(x, t)$ by looking for a curve $c(t)$ in $R_{2}$ such that above $c(t), u=0$ while below $c(t), u=1$.

To this end we will employ the Rankine-Hugoniot Condition of Theorem 37.26. To do this observe that Burger's Equation may be written as $u_{t}+$ $(F(u))_{x}=0$ where $F(u)=\frac{u^{2}}{2}$. So the Jump condition is

$$
\dot{c}\left(u_{+}-u_{-}\right)=\left(F\left(u_{+}\right)-F\left(u_{-}\right)\right)
$$

that is

$$
(0-1) \dot{c}=\left(\frac{0^{2}}{2}-\frac{1^{2}}{2}\right)=-\frac{1}{2} .
$$

Hence $\dot{c}(t)=\frac{1}{2}$ and therefore $c(t)=\frac{1}{2} t+1$ for $t \geq 0$. So we find a distributional solution given by the values in shown in Figure 37.8.


Fig. 37.8. A distributional solution to Burger's equation.

### 37.5 Exercises

Exercise 37.28. For $A \in L(X)$, let

$$
\begin{equation*}
e^{t A}:=\sum_{n=0}^{\infty} \frac{t^{n}}{n!} A^{n} \tag{37.58}
\end{equation*}
$$

Show directly that:

1. $e^{t A}$ is convergent in $L(X)$ when equipped with the operator norm.
2. $e^{t A}$ is differentiable in $t$ and that $\frac{d}{d t} e^{t A}=A e^{t A}$.

Exercise 37.29. Suppose that $A \in L(X)$ and $v \in X$ is an eigenvector of $A$ with eigenvalue $\lambda$, i.e. that $A v=\lambda v$. Show $e^{t A} v=e^{t \lambda} v$. Also show that $X=\mathbb{R}^{n}$ and $A$ is a diagonalizable $n \times n$ matrix with

$$
A=S D S^{-1} \text { with } D=\operatorname{diag}\left(\lambda_{1}, \ldots, \lambda_{n}\right)
$$

then $e^{t A}=S e^{t D} S^{-1}$ where $e^{t D}=\operatorname{diag}\left(e^{t \lambda_{1}}, \ldots, e^{t \lambda_{n}}\right)$.
Exercise 37.30. Suppose that $A, B \in L(X)$ and $[A, B]:=A B-B A=0$. Show that $e^{(A+B)}=e^{A} e^{B}$.

Exercise 37.31. Suppose $A \in C(\mathbb{R}, L(X))$ satisfies $[A(t), A(s)]=0$ for all $s, t \in \mathbb{R}$. Show

$$
y(t):=e^{\left(\int_{0}^{t} A(\tau) d \tau\right)} x
$$

is the unique solution to $\dot{y}(t)=A(t) y(t)$ with $y(0)=x$.
Exercise 37.32. Compute $e^{t A}$ when

$$
A=\left(\begin{array}{cc}
0 & 1 \\
-1 & 0
\end{array}\right)
$$

and use the result to prove the formula

$$
\cos (s+t)=\cos s \cos t-\sin s \sin t
$$

Hint: Sum the series and use $e^{t A} e^{s A}=e^{(t+s) A}$.
Exercise 37.33. Compute $e^{t A}$ when

$$
A=\left(\begin{array}{lll}
0 & a & b \\
0 & 0 & c \\
0 & 0 & 0
\end{array}\right)
$$

with $a, b, c \in \mathbb{R}$. Use your result to compute $e^{t(\lambda I+A)}$ where $\lambda \in \mathbb{R}$ and $I$ is the $3 \times 3$ identity matrix. Hint: Sum the series.

Theorem 37.34. Suppose that $T_{t} \in L(X)$ for $t \geq 0$ satisfies

1. (Semi-group property.) $T_{0}=I d_{X}$ and $T_{t} T_{s}=T_{t+s}$ for all $s, t \geq 0$.
2. (Norm Continuity) $t \rightarrow T_{t}$ is continuous at 0 , i.e. $\left\|T_{t}-I\right\|_{L(X)} \rightarrow 0$ as $t \downarrow 0$.

Then there exists $A \in L(X)$ such that $T_{t}=e^{t A}$ where $e^{t A}$ is defined in Eq. (37.58).

Exercise 37.35. Prove Theorem 37.34 using the following outline.

1. First show $t \in[0, \infty) \rightarrow T_{t} \in L(X)$ is continuous.
2. For $\epsilon>0$, let $S_{\epsilon}:=\frac{1}{\epsilon} \int_{0}^{\epsilon} T_{\tau} d \tau \in L(X)$. Show $S_{\epsilon} \rightarrow I$ as $\epsilon \downarrow 0$ and conclude from this that $S_{\epsilon}$ is invertible when $\epsilon>0$ is sufficiently small. For the remainder of the proof fix such a small $\epsilon>0$.
3. Show

$$
T_{t} S_{\epsilon}=\frac{1}{\epsilon} \int_{t}^{t+\epsilon} T_{\tau} d \tau
$$

and conclude from this that

$$
\lim _{t \downarrow 0} t^{-1}\left(T_{t}-I\right) S_{\epsilon}=\frac{1}{\epsilon}\left(T_{\epsilon}-I d_{X}\right)
$$

4. Using the fact that $S_{\epsilon}$ is invertible, conclude $A=\lim _{t \downarrow 0} t^{-1}\left(T_{t}-I\right)$ exists in $L(X)$ and that

$$
A=\frac{1}{\epsilon}\left(T_{\epsilon}-I\right) S_{\epsilon}^{-1}
$$

5. Now show using the semigroup property and step 4 . that $\frac{d}{d t} T_{t}=A T_{t}$ for all $t>0$.
6. Using step 5 , show $\frac{d}{d t} e^{-t A} T_{t}=0$ for all $t>0$ and therefore $e^{-t A} T_{t}=$ $e^{-0 A} T_{0}=I$.

Exercise 37.36 (Higher Order ODE). Let $X$ be a Banach space,, $\mathcal{U} \subset_{o}$ $X^{n}$ and $f \in C(J \times \mathcal{U}, X)$ be a Locally Lipschitz function in $\mathbf{x}=\left(x_{1}, \ldots, x_{n}\right)$. Show the $n^{\text {th }}$ ordinary differential equation,

$$
\begin{align*}
y^{(n)}(t) & =f\left(t, y(t), \dot{y}(t), \ldots y^{(n-1)}(t)\right) \\
\text { with } y^{(k)}(0) & =y_{0}^{k} \text { for } k=0,1,2 \ldots, n-1 \tag{37.59}
\end{align*}
$$

where $\left(y_{0}^{0}, \ldots, y_{0}^{n-1}\right)$ is given in $\mathcal{U}$, has a unique solution for small $t \in J$. Hint: let $\mathbf{y}(t)=\left(y(t), \dot{y}(t), \ldots y^{(n-1)}(t)\right)$ and rewrite Eq. (37.59) as a first order ODE of the form

$$
\dot{\mathbf{y}}(t)=Z(t, \mathbf{y}(t)) \text { with } \mathbf{y}(0)=\left(y_{0}^{0}, \ldots, y_{0}^{n-1}\right)
$$

Exercise 37.37. Use the results of Exercises 37.33 and 37.36 to solve

$$
\ddot{y}(t)-2 \dot{y}(t)+y(t)=0 \text { with } y(0)=a \text { and } \dot{y}(0)=b .
$$

Hint: The $2 \times 2$ matrix associated to this system, $A$, has only one eigenvalue 1 and may be written as $A=I+B$ where $B^{2}=0$.

Exercise 37.38. Suppose that $A: \mathbb{R} \rightarrow L(X)$ is a continuous function and $U, V: \mathbb{R} \rightarrow L(X)$ are the unique solution to the linear differential equations

$$
\begin{equation*}
\dot{V}(t)=A(t) V(t) \text { with } V(0)=I \tag{37.60}
\end{equation*}
$$

and

$$
\begin{equation*}
\dot{U}(t)=-U(t) A(t) \text { with } U(0)=I \tag{37.61}
\end{equation*}
$$

Prove that $V(t)$ is invertible and that $V^{-1}(t)=U(t)$. Hint: 1) show $\frac{d}{d t}[U(t) V(t)]=0$ (which is sufficient if $\left.\operatorname{dim}(X)<\infty\right)$ and 2) show compute $y(t):=V(t) U(t)$ solves a linear differential ordinary differential equation that has $y \equiv 0$ as an obvious solution. Then use the uniqueness of solutions to ODEs. (The fact that $U(t)$ must be defined as in Eq. (37.61) is the content of Exercise 19.32 in the analysis notes.)

Exercise 37.39 (Duhamel's Principle I). Suppose that $A: \mathbb{R} \rightarrow L(X)$ is a continuous function and $V: \mathbb{R} \rightarrow L(X)$ is the unique solution to the linear differential equation in Eq. (37.60). Let $x \in X$ and $h \in C(\mathbb{R}, X)$ be given. Show that the unique solution to the differential equation:

$$
\begin{equation*}
\dot{y}(t)=A(t) y(t)+h(t) \text { with } y(0)=x \tag{37.62}
\end{equation*}
$$

is given by

$$
\begin{equation*}
y(t)=V(t) x+V(t) \int_{0}^{t} V(\tau)^{-1} h(\tau) d \tau \tag{37.63}
\end{equation*}
$$

Hint: compute $\frac{d}{d t}\left[V^{-1}(t) y(t)\right]$ when $y$ solves Eq. (37.62).

Exercise 37.40 (Duhamel' s Principle II). Suppose that $A: \mathbb{R} \rightarrow L(X)$ is a continuous function and $V: \mathbb{R} \rightarrow L(X)$ is the unique solution to the linear differential equation in Eq. (37.60). Let $W_{0} \in L(X)$ and $H \in C(\mathbb{R}, L(X))$ be given. Show that the unique solution to the differential equation:

$$
\begin{equation*}
\dot{W}(t)=A(t) W(t)+H(t) \text { with } W(0)=W_{0} \tag{37.64}
\end{equation*}
$$

is given by

$$
\begin{equation*}
W(t)=V(t) W_{0}+V(t) \int_{0}^{t} V(\tau)^{-1} H(\tau) d \tau \tag{37.65}
\end{equation*}
$$

Exercise 37.41 (Non-Homogeneous ODE). Suppose that $U \subset_{o} X$ is open and $Z: \mathbb{R} \times U \rightarrow X$ is a continuous function. Let $J=(a, b)$ be an interval and $t_{0} \in J$. Suppose that $y \in C^{1}(J, U)$ is a solution to the "nonhomogeneous" differential equation:

$$
\begin{equation*}
\dot{y}(t)=Z(t, y(t)) \text { with } y\left(t_{o}\right)=x \in U . \tag{37.66}
\end{equation*}
$$

Define $Y \in C^{1}\left(J-t_{0}, \mathbb{R} \times U\right)$ by $Y(t):=\left(t+t_{0}, y\left(t+t_{0}\right)\right)$. Show that $Y$ solves the "homogeneous" differential equation

$$
\begin{equation*}
\dot{Y}(t)=\tilde{A}(Y(t)) \text { with } Y(0)=\left(t_{0}, y_{0}\right) \tag{37.67}
\end{equation*}
$$

where $\tilde{A}(t, x):=(1, Z(x))$. Conversely, suppose that $Y \in C^{1}\left(J-t_{0}, \mathbb{R} \times U\right)$ is a solution to Eq. (37.67). Show that $Y(t)=\left(t+t_{0}, y\left(t+t_{0}\right)\right)$ for some $y \in$ $C^{1}(J, U)$ satisfying Eq. (37.66). (In this way the theory of non-homogeneous ode's may be reduced to the theory of homogeneous ode's.)

Exercise 37.42 (Differential Equations with Parameters). Let $W$ be another Banach space, $U \times V \subset_{o} X \times W$ and $Z \in C(U \times V, X)$ be a locally Lipschitz function on $U \times V$. For each $(x, w) \in U \times V$, let $t \in J_{x, w} \rightarrow \phi(t, x, w)$ denote the maximal solution to the ODE

$$
\begin{equation*}
\dot{y}(t)=Z(y(t), w) \text { with } y(0)=x . \tag{37.68}
\end{equation*}
$$

Prove

$$
\begin{equation*}
\mathcal{D}:=\left\{(t, x, w) \in \mathbb{R} \times U \times V: t \in J_{x, w}\right\} \tag{37.69}
\end{equation*}
$$

is open in $\mathbb{R} \times U \times V$ and $\phi$ and $\dot{\phi}$ are continuous functions on $\mathcal{D}$.
Hint: If $y(t)$ solves the differential equation in (37.68), then $v(t):=$ $(y(t), w)$ solves the differential equation,

$$
\begin{equation*}
\dot{v}(t)=\tilde{A}(v(t)) \text { with } v(0)=(x, w) \tag{37.70}
\end{equation*}
$$

where $\tilde{A}(x, w):=(Z(x, w), 0) \in X \times W$ and let $\psi(t,(x, w)):=v(t)$. Now apply the Theorem 6.21 to the differential equation (37.70).

$$
\begin{aligned}
& \cos (t A):=\sum_{n=0}^{\infty} \frac{(-1)^{2 n}}{(2 n)!} t^{2 n} A^{2 n} \text { and } \\
& \frac{\sin (t A)}{A}:=\sum_{n=0}^{\infty} \frac{(-1)^{2 n+1}}{(2 n+1)!} t^{2 n+1} A^{2 n} .
\end{aligned}
$$

Show that the unique solution $y \in C^{2}(\mathbb{R}, X)$ to

$$
\begin{equation*}
\ddot{y}(t)+A^{2} y(t)=0 \text { with } y(0)=y_{0} \text { and } \dot{y}(0)=\dot{y}_{0} \in X \tag{37.71}
\end{equation*}
$$

is given by

$$
y(t)=\cos (t A) y_{0}+\frac{\sin (t A)}{A} \dot{y}_{0}
$$

Remark 37.44. Exercise 37.43 can be done by direct verification. Alternatively and more instructively, rewrite Eq. (37.71) as a first order ODE using Exercise 37.36. In doing so you will be lead to compute $e^{t B}$ where $B \in L(X \times X)$ is given by

$$
B=\left(\begin{array}{cc}
0 & I \\
-A^{2} & 0
\end{array}\right)
$$

where we are writing elements of $X \times X$ as column vectors, $\binom{x_{1}}{x_{2}}$. You should then show

$$
e^{t B}=\left(\begin{array}{cc}
\cos (t A) & \frac{\sin (t A)}{A} \\
-A \sin (t A) & \cos (t A)
\end{array}\right)
$$

where

$$
A \sin (t A):=\sum_{n=0}^{\infty} \frac{(-1)^{2 n+1}}{(2 n+1)!} t^{2 n+1} A^{2(n+1)}
$$

Exercise 37.45 (Duhamel's Principle for the Abstract Wave Equation). Continue the notation in Exercise 37.43, but now consider the ODE,

$$
\begin{equation*}
\ddot{y}(t)+A^{2} y(t)=f(t) \text { with } y(0)=y_{0} \text { and } \dot{y}(0)=\dot{y}_{0} \in X \tag{37.72}
\end{equation*}
$$

where $f \in C(\mathbb{R}, X)$. Show the unique solution to Eq. (37.72) is given by

$$
\begin{equation*}
y(t)=\cos (t A) y_{0}+\frac{\sin (t A)}{A} \dot{y}_{0}+\int_{0}^{t} \frac{\sin ((t-\tau) A)}{A} f(\tau) d \tau \tag{37.73}
\end{equation*}
$$

Hint: Again this could be proved by direct calculation. However it is more instructive to deduce Eq. (37.73) from Exercise 37.39 and the comments in Remark 37.44.

Exercise 37.46. Number 3 on p. 163 of Evans.

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$$
\begin{aligned}
x^{\prime}(s) & =F_{p}(x(s), z(s), p(s)) \\
z^{\prime}(s) & =p(s) \cdot F_{p}(x(s), z(s), p(s)) \\
p^{\prime}(s) & =-F_{x}(x(s), z(s), p(s))-F_{z}(x(s), z(s), p(s)) p(s) .
\end{aligned}
$$

## Fully nonlinear first order PDE

In this section let $\mathcal{U} \subset_{o} \mathbb{R}^{n}$ be an open subset of $\mathbb{R}^{n}$ and $(x, z, p) \in \overline{\mathcal{U}} \times \mathbb{R}^{n} \times$ $\mathbb{R} \rightarrow F(x, z, p) \in \mathbb{R}$ be a $C^{2}$ - function. Actually to simplify notation let us suppose $\mathcal{U}=\mathbb{R}^{n}$. We are now looking for a solution $u: \mathbb{R}^{n} \rightarrow \mathbb{R}$ to the fully non-linear PDE,

$$
\begin{equation*}
F(x, u(x), \nabla u(x))=0 . \tag{38.1}
\end{equation*}
$$

As above, we "reduce" the problem to solving ODE's. To see how this might be done, suppose $u$ solves (38.1) and $x(s)$ is a curve in $\mathbb{R}^{n}$ and let

$$
z(s)=u(x(s)) \text { and } p(s)=\nabla u(x(s)) .
$$

Then

$$
\begin{align*}
& z^{\prime}(s)=\nabla u(x(s)) \cdot x^{\prime}(s)=p(s) \cdot x^{\prime}(s) \text { and }  \tag{38.2}\\
& p^{\prime}(s)=\partial_{x^{\prime}(s)} \nabla u(x(s)) . \tag{38.3}
\end{align*}
$$

We would now like to find an equation for $x(s)$ which along with the above system of equations would form and ODE for $(x(s), z(s), p(s))$. The term, $\partial_{x^{\prime}(s)} \nabla u(x(s))$, which involves two derivative of $u$ is problematic and we would like to replace it by something involving only $\nabla u$ and $u$. In order to get the desired relation, differentiate Eq. (38.1) in $x$ in the direction $v$ to find

$$
\begin{aligned}
0 & =F_{x} \cdot v+F_{z} \partial_{v} u+F_{p} \cdot \partial_{v} \nabla u=F_{x} \cdot v+F_{z} \partial_{v} u+F_{p} \cdot \nabla \partial_{v} u \\
& =F_{x} \cdot v+F_{z} \nabla u \cdot v+\left(\partial_{F_{p}} \nabla u\right) \cdot v,
\end{aligned}
$$

wherein we have used the fact that mixed partial derivative commute. This equation is equivalent to

$$
\begin{equation*}
\left.\partial_{F_{p}} \nabla u\right|_{(x, u(x), \nabla u(x))}=-\left.\left(F_{x}+F_{z} \nabla u\right)\right|_{(x, u(x), \nabla u(x))} . \tag{38.4}
\end{equation*}
$$

By requiring $x(s)$ to solve $x^{\prime}(s)=F_{p}(x(s), z(s), p(s))$, we find, using Eq. (38.4) and Eqs. (38.2) and (38.3) that $(x(s), z(s), p(s))$ solves the characteristic equations,

We will in the future simply abbreviate these equations by

$$
\begin{align*}
x^{\prime} & =F_{p} \\
z^{\prime} & =p \cdot F_{p}  \tag{38.5}\\
p^{\prime} & =-F_{x}-F_{z} p .
\end{align*}
$$

The above considerations have proved the following Lemma.

## Lemma 38.1. Let

$$
\begin{gathered}
A(x, z, p):=\left(F_{p}(x, z, p), p \cdot F_{p}(x, z, p),-F_{x}(x, z, p)-F_{z}(x, z, p) p\right), \\
\pi_{1}(x, z, p)=x \text { and } \pi_{2}(x, z, p)=z .
\end{gathered}
$$

If $u$ solves $E q$. (38.1) and $x_{0} \in U$, then

$$
\begin{align*}
e^{s A}\left(x_{0}, u\left(x_{0}\right), \nabla u\left(x_{0}\right)\right) & =(x(s), u(x(s)), \nabla u(x(s))) \text { and } \\
u(x(s)) & =\pi_{2} \circ e^{s A}\left(x_{0}, u\left(x_{0}\right), \nabla u\left(x_{0}\right)\right) \tag{38.6}
\end{align*}
$$

where $x(s)=\pi_{1} \circ e^{s A}\left(x_{0}, u\left(x_{0}\right), \nabla u\left(x_{0}\right)\right)$.
We now want to use Eq. (38.6) to produce solutions to Eq. (38.1). As in the quasi-linear case we will suppose $\Sigma: U \subset_{o} \mathbb{R}^{n-1} \longrightarrow \mathbb{R}^{n}$ is a surface, $\Sigma(0)=x_{0}, D \Sigma(y)$ is injective for all $y \in U$ and $u_{0}: \Sigma \rightarrow \mathbb{R}$ is given. We wish to solve Eq. (38.1) for $u$ with the added condition that $u(\Sigma(y))=u_{0}(y)$. In order to make use of Eq. (38.6) to do this, we first need to be able to find $\nabla u(\Sigma(y))$. The idea is to use Eq. (38.1) to determine $\nabla u(\Sigma(y))$ as a function of $\Sigma(y)$ and $u_{0}(y)$ and for this we will invoke the implicit function theorem. If $u$ is a function such that $u(\Sigma(y))=u_{0}(y)$ for $y$ near 0 and $p_{0}=\nabla u\left(x_{0}\right)$ then

$$
\partial_{v} u_{0}(0)=\left.\partial_{v} u(\Sigma(y))\right|_{y=0}=\nabla u\left(x_{0}\right) \cdot \Sigma^{\prime}(0) v=p_{0} \cdot \Sigma^{\prime}(0) v
$$

Notation 38.2 Let $\nabla_{\Sigma} u_{0}(y)$ denote the unique vector in $\mathbb{R}^{n}$ which is tangential to $\Sigma$ at $\Sigma(y)$ and such that

$$
\partial_{v} u_{0}(y)=\nabla_{\Sigma} u_{0}(y) \cdot \Sigma^{\prime}(0) v \text { for all } v \in \mathbb{R}^{n-1}
$$

Theorem 38.3. Let $F: \mathbb{R}^{n} \times \mathbb{R} \times \mathbb{R}^{n} \rightarrow \mathbb{R}$ be a $C^{2}$ function, $0 \in U \subset \subset_{o}^{n-1}$, $\Sigma: U \subset_{o} \mathbb{R}^{n-1} \xrightarrow{C^{2}} \mathbb{R}^{n}$ be an embedded submanifold, $\left(x_{0}, z_{0}, p_{0}\right) \in \Sigma \times \mathbb{R} \times \mathbb{R}^{n}$ such that $F\left(x_{0}, z_{0}, p_{0}\right)=0$ and $x_{0}=\Sigma(0), u_{0}: \Sigma \xrightarrow{C^{1}} \mathbb{R}$ such that $u_{0}\left(x_{0}\right)=$ $z_{0}, \mathbf{n}(y)$ be a normal vector to $\Sigma$ at $y$. Further assume

1. $\partial_{v} u_{0}(0)=p_{0} \cdot \Sigma^{\prime}(0) v=p_{0} \cdot \partial_{v} \Sigma(0)$ for all $v \in \mathbb{R}^{n-1}$.

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## 2. $F_{p}\left(x_{0}, y_{0}, z_{0}\right) \cdot \mathbf{n}(0) \neq 0$

Then there exists a neighborhood $V \subset \mathbb{R}^{n}$ of $x_{0}$ and a $C^{2}$-function $u: V \rightarrow$ $\mathbb{R}$ such that $u \circ \Sigma=u_{0}$ near 0 and $E q$. (38.1) holds for all $x \in V$.

Proof. Step 1. There exist a neighborhood $U_{0} \subset U$ and a function $p_{0}$ : $U_{0} \rightarrow \mathbb{R}^{n}$ such that

$$
\begin{equation*}
p_{0}(y)^{\tan }=\nabla_{\Sigma} u_{0}(y) \text { and } F\left(\Sigma(y), u\left(\Sigma((y)), p_{0}(y)\right)=0\right. \tag{38.7}
\end{equation*}
$$

for all $y \in U_{0}$, where $p_{0}(y)^{\tan }$ is component of $p_{0}(y)$ tangential to $\Sigma$. This is


Fig. 38.1. Decomposing $p$ into its normal and tangential components.
an exercise in the implicit function theorem.
Choose $\alpha_{0} \in \mathbb{R}$ such that $\nabla_{\Sigma} u_{0}(0)+\alpha_{0} \mathbf{n}(0)=p_{0}$ and define

$$
f(\alpha, y):=F\left(\Sigma(y), u_{0}(y), \nabla_{\Sigma} u_{0}(y)+\alpha \mathbf{n}(y)\right)
$$

Then

$$
\frac{\partial f}{\partial \alpha}(\alpha, 0)=F_{p}\left(x_{0}, z_{0}, \nabla_{\Sigma} u_{0}(0)+\alpha \mathbf{n}(0)\right) \cdot \mathbf{n}(0) \neq 0
$$

so by the implicit function theorem there exists $0 \in U_{0} \subset U$ and $\alpha: U_{0} \rightarrow \mathbb{R}$ such that $f(\alpha(y), y)=0$ for all $y \in U_{0}$. Now define

$$
p_{0}(y):=\nabla_{\Sigma} u_{0}(y)+\alpha(y) \mathbf{n}(y) \text { for } y \in U_{0}
$$

To simplify notation in the future we will from now on write $U$ for $U_{0}$.
Step 2. Suppose $(x, z, p)$ is a solution to $(38.5)$ such that $F(x(0), z(0), p(0))=$ 0 then

$$
\begin{equation*}
F(x(s), z(s), p(s))=0 \text { for all } t \in J \tag{38.8}
\end{equation*}
$$

$$
\begin{align*}
\frac{d}{d s} F(x(s), z(s), p(s)) & =F_{x} \cdot x^{\prime}+F_{z} z^{\prime}+F_{p} \cdot p^{\prime} \\
& =F_{x} \cdot F_{p}+F_{z}\left(p \cdot F_{p}\right)-F_{p} \cdot\left(F_{x}+F_{z} p\right)=0 \tag{38.9}
\end{align*}
$$

Step 3. (Notation). For $y \in U$ let

$$
(X(s, y), Z(s, y), P(s, y))=e^{s A}\left(\Sigma(y), u_{0}(y), p_{0}(y)\right)
$$

ie. $X(s, y), Z(s, y)$ and $P(s, y)$ solve the coupled system of O.D.E.'s:

$$
\begin{align*}
X^{\prime} & =F_{p} \text { with } X(0, y)=\Sigma(y) \\
Z^{\prime} & =P \cdot F_{p} \text { with } Z(0, y)=u_{0}(y) \\
P^{\prime} & =F_{x}-F_{z} P \text { with } P(0, y)=p_{0}(y) \tag{38.10}
\end{align*}
$$

With this notation Eq. (38.8) becomes

$$
\begin{equation*}
F(X(s, y), Z(s, y), P(s, y))=0 \text { for all } t \in J \tag{38.11}
\end{equation*}
$$

Step 4. There exists a neighborhood $0 \in U_{0} \subset U$ and $0 \in J \subset \mathbb{R}$ such that $X: J \times U_{0} \rightarrow \mathbb{R}^{n}$ is a $C^{1}$ diffeomorphism onto an open set $V:=X\left(J \times U_{0}\right) \subset$ $\mathbb{R}^{n}$ with $x_{0} \in V$. Indeed, $X(0, y)=\Sigma(y)$ so that $\left.D_{y} X(0, y)\right|_{y=0}=\Sigma^{\prime}(0)$ and hence

$$
D X(0,0)(a, v)=\frac{\partial X}{\partial s}(0,0) a+\Sigma^{\prime}(0) v=F_{p}\left(x_{0}, z_{0}, p_{0}\right) a+\Sigma^{\prime}(0) v
$$

By the assumptions, $F_{p}\left(x_{0}, z_{0}, p_{0}\right) \notin \operatorname{Ran} \Sigma^{\prime}(0)$ and $\Sigma^{\prime}(0)$ is injective, it follows that $D X(0,0)$ is invertible So the assertion is a consequence of the inverse function theorem.

Step 5. Define

$$
u(x):=Z\left(X^{-1}(x)\right)
$$

then $u$ is the desired solution. To prove this first notice that $u$ is uniquely characterized by

$$
u(X(s, y))=Z(s, y) \text { for all }(s, y) \in J_{0} \times U_{0}
$$

Because of Step 2., to finish the proof it suffices to show $\nabla u(X(s, y))=P(s, y)$.
Step 6. $\nabla u(X(s, y))=P(s, y)$. From Eq. (38.10),

$$
\begin{equation*}
P \cdot X^{\prime}=P \cdot F_{p}=Z^{\prime}=\frac{d}{d s} u(X)=\nabla u(X) \cdot X^{\prime} \tag{38.12}
\end{equation*}
$$

which shows

$$
[P-\nabla u(X)] \cdot X^{\prime}=0
$$

So to finish the proof it suffices to show

$$
[P-\nabla u(X)] \cdot \partial_{v} X=0
$$

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for all $v \in \mathbb{R}^{n-1}$ or equivalently that

$$
\begin{equation*}
P(s, y) \cdot \partial_{v} X=\nabla u(x) \cdot \partial_{v} X=\partial_{v} u(X)=\partial_{v} Z \tag{38.13}
\end{equation*}
$$

for all $v \in \mathbb{R}^{n-1}$.
To prove Eq. (38.13), fix a $y$ and let

$$
r(s):=P(s, y) \cdot \partial_{v} X(s, y)-\partial_{v} Z(s, y)
$$

Then using Eq. (38.10),

$$
\begin{align*}
r^{\prime} & =P^{\prime} \cdot \partial_{v} X+P \cdot \partial_{v} X^{\prime}-\partial_{v} Z^{\prime} \\
& =\left(-F_{x}-F_{z} P\right) \cdot \partial_{v} X+P \cdot \partial_{v} F_{p}-\partial_{v}\left(P \cdot F_{p}\right) \\
& =\left(-F_{x}-F_{z} P\right) \cdot \partial_{v} X-\left(\partial_{v} P\right) \cdot F_{p} . \tag{38.14}
\end{align*}
$$

Further, differentiating Eq. (38.11) in $y$ implies for all $v \in \mathbb{R}^{n-1}$ that

$$
\begin{equation*}
F_{x} \cdot \partial_{v} X+F_{z} \partial_{v} Z+F_{p} \cdot \partial_{v} P=0 \tag{38.15}
\end{equation*}
$$

Adding Eqs. (38.14) and (38.15) then shows

$$
r^{\prime}=-F_{z} P \cdot \partial_{v} X+F_{z} \partial_{v} Z=-F_{z} r
$$

which implies

$$
r(s)=e^{-\int_{0}^{s} F_{z}(X, Z, P)(\sigma, y) d \sigma} r(0) .
$$

This shows $r \equiv 0$ because $p_{0}(y)^{T}=\left(\nabla_{\Sigma} u_{0}\right)(\Sigma(y))$ and hence

$$
\begin{aligned}
r(0) & =p_{0}(y) \cdot \partial_{v} \Sigma(y)-\partial_{v} u_{0}(\Sigma(y)) \\
& =\left[p_{0}(y)-\nabla_{\Sigma} u_{0}(\Sigma(y))\right] \cdot \partial_{v} \Sigma(y)=0 .
\end{aligned}
$$

Example 38.4 (Quasi-Linear Case Revisited). Let us consider the quasi-linear PDE in Eq. (37.39),

$$
\begin{equation*}
A(x, z) \cdot \nabla_{x} u(x)-c(x, u(x))=0 . \tag{38.16}
\end{equation*}
$$

in light of Theorem 38.3. This may be written in the form of Eq. (38.1) by setting

$$
F(x, z, p)=A(x, z) \cdot p-c(x, z)
$$

The characteristic equations (38.5) for this $F$ are

$$
\begin{aligned}
x^{\prime} & =F_{p}=A \\
z^{\prime} & =p \cdot F_{p}=p \cdot A \\
p^{\prime} & =-F_{x}-F_{z} p=-\left(A_{x} \cdot p-c_{x}\right)-\left(A_{z} \cdot p-c_{z}\right) p
\end{aligned}
$$

Recalling that $p(s)=\nabla u(s, x)$, the $z$ equation above may expressed, by using Eq. (38.16) as

$$
z^{\prime}=p \cdot A=c
$$

Therefore the equations for $(x(s), z(s))$ may be written as

$$
x^{\prime}(s)=A(x, z) \text { and } z^{\prime}=c(x, z)
$$

and these equations may be solved without regard for the $p$-equation. This is what makes the quasi-linear PDE simpler than the fully non-linear PDE.

### 38.1 An Introduction to Hamilton Jacobi Equations

A Hamilton Jacobi Equation is a first order PDE of the form

$$
\begin{equation*}
\frac{\partial S}{\partial t}(t, x)+H\left(x, \nabla_{x} S(t, x)\right)=0 \text { with } S(0, x)=g(x) \tag{38.17}
\end{equation*}
$$

where $H: \mathbb{R}^{n} \times \mathbb{R}^{n} \rightarrow \mathbb{R}$ and $g: \mathbb{R}^{n} \rightarrow \mathbb{R}$ are given functions. In this section we are going to study the connections of this equation to the Euler Lagrange equations of classical mechanics.

### 38.1.1 Solving the Hamilton Jacobi Equation (38.17) by characteristics

Now let us solve the Hamilton Jacobi Equation (38.17) using the method of characteristics. In order to do this let

$$
\left(p_{0}, p\right)=\left(\frac{\partial S}{\partial t}, \nabla_{x} S(t, x)\right) \text { and } F(t, x, z, p):=p_{0}+H(x, p)
$$

Then Eq. (38.17) becomes

$$
0=F\left(t, x, S, \frac{\partial S}{\partial t}, \nabla_{x} S\right)
$$

Hence the characteristic equations are given by

$$
\begin{gathered}
\frac{d}{d s}(t(s), x(s))=F_{\left(p_{0}, p\right)}=\left(1, \nabla_{p} H(x(s), p(s))\right. \\
\frac{d}{d s}\left(p_{0}, p\right)(s)=-F_{(t, x)}-F_{z}\left(p_{0}, p\right)=-F_{(t, x)}=\left(0,-\nabla_{x} H(x(s), p(s))\right) \\
z^{\prime}(s)=\left(p_{0}, p\right) \cdot F_{\left(p_{0}, p\right)}=p_{0}(s)+p(s) \cdot \nabla_{p} H(x(s), p(s))
\end{gathered}
$$

and
Solving the $t$ equation with $t(0)=0$ gives $t=s$ and so we identify $t$ and $s$ and our equations become

$$
\begin{align*}
\dot{x}(t) & =\nabla_{p} H(x(t), p(t))  \tag{38.18}\\
\dot{p}(t) & =-\nabla_{\mathbf{x}} H(x(t), p(t))  \tag{38.19}\\
\frac{d}{d t}\left[\frac{\partial S}{\partial t}(t, x(t))\right] & =\frac{d}{d t} p_{0}(t)=0 \text { and } \\
\frac{d}{d t}[S(t, x(t))] & =\frac{d}{d t} z(t)=\frac{\partial S}{\partial t}(t, x(t))+p(t) \cdot \nabla_{p} H(x(t), p(t)) \\
& =-H(x(t), p(t))+p(t) \cdot \nabla_{p} H(x(t), p(t))
\end{align*}
$$

Hence we have proved the following proposition.
Proposition 38.5. If $S$ solves the Hamilton Jacobi Equation Eq. 38.17 and $(x(t), p(t))$ are solutions to the Hamilton Equations (38.18) and (38.19) (see also Eq. (38.29) below) with $p(0)=\left(\nabla_{x} g\right)(x(0))$ then

$$
S(T, x(T))=g(x(0))+\int_{0}^{T}\left[p(t) \cdot \nabla_{p} H(x(t), p(t))-H(x(t), p(t))\right] d t
$$

In particular if $(T, x) \in \mathbb{R} \times \mathbb{R}^{n}$ then

$$
\begin{equation*}
S(T, x)=g(x(0))+\int_{0}^{T}\left[p(t) \cdot \nabla_{p} H(x(t), p(t))-H(x(t), p(t))\right] d t \tag{38.20}
\end{equation*}
$$

provided $(x, p)$ is a solution to Hamilton Equations (38.18) and (38.19) satisfying the boundary condition $x(T)=x$ and $p(0)=\left(\nabla_{x} g\right)(x(0))$.
Remark 38.6. Let $X\left(t, x_{0}, p_{0}\right)=x(t)$ and $P\left(t, x_{0}, p_{0}\right)=p(t)$ where $(x(t), p(t))$ satisfies Hamilton Equations (38.18) and (38.19) with $(x(0), p(0))=\left(x_{0}, p_{0}\right)$ and let $\Psi(t, x):=(t, X(t, x, \nabla g(x))$. Then $\Psi(0, x)=(0, x)$ so

$$
\partial_{v} \Psi(0,0)=(0, v) \text { and } \partial_{t} \Psi(0,0)=\left(1, \nabla_{p} H(x, \nabla g(x))\right)
$$

from which it follows that $\Psi^{\prime}(0,0)$ is invertible. Therefore given $a \in \mathbb{R}^{n}$, the exists $\epsilon>0$ such that $\Psi^{-1}(t, x)$ is well defined for $|t|<\epsilon$ and $|x-a|<\epsilon$. Writing $\Psi^{-1}(T, x)=\left(T, x_{0}(T, x)\right)$ we then have that

$$
(x(t), p(t)):=\left(X \left(t, x_{0}(T, x), \nabla g\left(x_{0}(T, x)\right), P\left(t, x_{0}, \nabla g\left(x_{0}(T, x)\right)\right)\right.\right.
$$

solves Hamilton Equations (38.18) and (38.19) satisfies the boundary condition $x(T)=x$ and $p(0)=\left(\nabla_{x} g\right)(x(0))$.

### 38.1.2 The connection with the Euler Lagrange Equations

Our next goal is to express the solution $S(T, x)$ in Eq. (38.20) solely in terms of the path $x(t)$. For this we digress a bit to Lagrangian mechanics and the notion of the "classical action."

Definition 38.7. Let $T>0, L: \mathbb{R}^{n} \times \mathbb{R}^{n} \longrightarrow \mathbb{R}$ be a smooth"Lagrangian" and $g: \mathbb{R}^{n} \rightarrow \mathbb{R}$ be a smooth function. The $g$ - weighted action $I_{T}^{g}(q)$ of a function $q \in C^{2}\left([0, T], \mathbb{R}^{n}\right)$ is defined to be

$$
I_{T}^{g}(q)=g(q(0))+\int_{0}^{T} L(q(t), \dot{q}(t)) d t
$$

When $g=0$ we will simply write $I_{T}$ for $I_{0, T}$.
We are now going to study the function $S(T, x)$ of "least action,"

$$
\begin{aligned}
S(T, x) & :=\inf \left\{I_{T}^{g}(q): q \in C^{2}([0, T]) \text { with } q(T)=x\right\} \\
& =\inf \left\{g(q(0))+\int_{0}^{T} L(q(t), \dot{q}(t)) d t: q \in C^{2}([0, T]) \text { with } q(T)=x\right\}
\end{aligned}
$$

The next proposition records the differential of $I_{T}^{g}(q)$.
Proposition 38.8. Let $L \in C^{\infty}\left(\mathbb{R}^{n} \times \mathbb{R}^{n}, \mathbb{R}\right)$ be a smooth Lagrangian, then for $q \in C^{2}\left([0, T], \mathbb{R}^{n}\right)$ and $h \in C^{1}\left([0, T], \mathbb{R}^{n}\right)$

$$
\begin{align*}
D I_{T}^{g}(q) h & =\left[\left(\nabla g(q)-D_{2} L(q, \dot{q})\right) \cdot h\right]_{t=0}+\left[D_{2} L(q, \dot{q}) \cdot h\right]_{t=T} \\
& +\int_{0}^{T}\left(D_{1} L(q, \dot{q})-\frac{d}{d t} D_{2} L(q, \dot{q})\right) h d t \tag{38.22}
\end{align*}
$$

Proof. By differentiating past the integral,

$$
\begin{aligned}
\partial_{h} I_{T}(q) & =\left.\frac{d}{d s}\right|_{0} I_{T}(q+s h)=\left.\int_{0}^{T} \frac{d}{d s}\right|_{0} L(q(t)+s h(t), \dot{q}(t)+s \dot{h}(t)) d t \\
& =\int_{0}^{T}\left(D_{1} L(q, \dot{q}) h+D_{2} L(q, \dot{q}) \dot{h}\right) d t \\
& =\int_{0}^{T}\left(D_{1} L(q, \dot{q})-\frac{d}{d t} D_{2} L(q, \dot{q})\right) h d t+\left.D_{2} L(q, \dot{q}) h\right|_{0} ^{T}
\end{aligned}
$$

This completes the proof since $I_{T}^{g}(q)=g(q(0))+I_{T}(q)$ and $\partial_{h}[g(q(0))]=$ $\nabla g(q(0)) \cdot h(0)$.

Definition 38.9. A function $q \in C^{2}\left([0, T], \mathbb{R}^{n}\right)$ is said to solve the Euler Lagrange equation for $L$ if $q$ solves

$$
\begin{equation*}
D_{1} L(q, \dot{q})-\frac{d}{d t}\left[D_{2} L(q, \dot{q})\right]=0 \tag{38.23}
\end{equation*}
$$

This is equivalently to $q$ satisfying $D I_{T}(q) h=0$ for all $h \in C^{1}\left([0, T], \mathbb{R}^{n}\right)$ which vanish on $\partial[0, T]=\{0, T\}$.

Let us note that the Euler Lagrange equations may be written as:

$$
D_{1} L(q, \dot{q})=D_{1} D_{2} L(q, \dot{q}) \dot{q}+D_{2}^{2} L(q, \dot{q}) \ddot{q}
$$

Corollary 38.10. Any minimizer $q$ (or more generally critical point) of $I_{T}^{g}(\cdot)$ must satisfy the Euler Lagrange Eq. (38.23) with the boundary conditions

$$
\begin{equation*}
q(T)=x \text { and } \nabla g(q(0))=\nabla_{\dot{q}} L(q(0), q(0))=D_{2} L(q(0), q(0)) . \tag{38.24}
\end{equation*}
$$

Proof. The corollary is a consequence Proposition 38.8 and the first derivative test which implies $D I_{T}^{g}(q) h=0$ for all $h \in C^{1}\left([0, T], \mathbb{R}^{n}\right)$ such that $h(T)=0$.

Example 38.11. Let $U \in C^{\infty}\left(\mathbb{R}^{n}, \mathbb{R}\right), m>0$ and $L(q, v)=\frac{1}{2} m|v|^{2}-U(q)$. Then

$$
D_{1} L(q, v)=-\nabla U(q) \text { and } D_{2} L(q, v)=m v
$$

and the Euler Lagrange equations become

$$
-\nabla U(q)=\frac{d}{d t}[m \dot{q}]=m \ddot{q}
$$

which are Newton's equations of motion for a particle of mass $m$ subject to a force $-\nabla U$. In particular if $U=0$, then $q(t)=q(0)+t \dot{q}(0)$.

The following assumption on $L$ will be assumed for the rest of this section.
Assumption 1 We assume $\left[D_{2}^{2} L(q, v)\right]^{-1}$ exists for all $(q, v) \in \mathbb{R}^{n} \times \mathbb{R}^{n}$ and $v \rightarrow D_{2} L(q, v)$ is invertible for all $q \in \mathbb{R}^{n}$.
Notation 38.12 For $q, p \in \mathbb{R}^{n}$ let

$$
\begin{equation*}
V(q, p):=\left[D_{2} L(q, \cdot)\right]^{-1}(p) \tag{38.25}
\end{equation*}
$$

Equivalently, $V(q, p)$ is the unique element of $\mathbb{R}^{n}$ such that

$$
\begin{equation*}
D_{2} L(q, V(q, p))=p . \tag{38.26}
\end{equation*}
$$

Remark 38.13. The function $V: \mathbb{R}^{n} \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ is smooth in $(q, p)$. This is a consequence of the implicit function theorem applied to $\Psi(q, v):=$ $\left(q, D_{2} L(q, v)\right)$.

Under Assumption 1, Eq. (38.23) may be written as

$$
\begin{equation*}
\ddot{q}=F(q, \dot{q}) \tag{38.27}
\end{equation*}
$$

where

$$
F(q, \dot{q})=D_{2}^{2} L(q, \dot{q})^{-1}\left\{D_{1} L(q, \dot{q})-D\left(D_{2} L(q, \dot{q}) \dot{q}\right\}\right.
$$

Definition 38.14 (Legendre Transform). Let $L \in C^{\infty}\left(\mathbb{R}^{n} \times \mathbb{R}^{n}, \mathbb{R}\right)$ be a function satisfying Assumption 1. The Legendre transform $L^{*} \in C^{\infty}\left(\mathbb{R}^{n} \times\right.$ $\left.\mathbb{R}^{n}, \mathbb{R}\right)$ is defined by

$$
L^{*}(x, p):=p \cdot v-L(x, v) \text { where } p=\nabla_{v} L(x, v)
$$

i.e.

$$
\begin{equation*}
L^{*}(x, p)=p \cdot V(x, p)-L(x, V(x, p)) . \tag{38.28}
\end{equation*}
$$

Proposition 38.15. Let $H(x, p):=L^{*}(x, p), q \in C^{2}\left([0, T], \mathbb{R}^{n}\right)$ and $p(t):=$ $L_{v}(q(t), \dot{q}(t))$. Then

1. $H \in C^{\infty}\left(\mathbb{R}^{n} \times \mathbb{R}^{n}, \mathbb{R}\right)$ and

$$
H_{x}(x, p)=-L_{x}(x, V(x, p)) \text { and } H_{p}(x, p)=V(x, p) . .
$$

2. $H$ satisfies Assumption 1 and $H^{*}=L$, i.e. $\left(L^{*}\right)^{*}=L$.
3. The path $q \in C^{2}\left([0, T], \mathbb{R}^{n}\right)$ solves the Euler Lagrange Eq. (38.23) then $(q(t), p(t))$ satisfies Hamilton's Equations:

$$
\begin{align*}
\dot{q}(t) & =H_{p}(q(t), p(t)) \\
\dot{p}(t) & =-H_{x}(q(t), p(t)) . \tag{38.29}
\end{align*}
$$

4. Conversely if $(q, p)$ solves Hamilton's equations (38.29) then $q$ solves the Euler Lagrange Eq. (38.23) and

$$
\begin{equation*}
\frac{d}{d t} H(q(t), p(t))=0 . \tag{38.30}
\end{equation*}
$$

Proof. The smoothness of $H$ follows by Remark 38.13.

1. Using Eq. (38.28) and Eq. (38.26)

$$
\begin{aligned}
H_{x}(x, p) & =p \cdot V_{x}(x, p)-L_{x}(x, V(x, p))-L_{v}(x, V(x, p)) V_{x}(x, p) \\
& =p \cdot V_{x}(x, p)-L_{x}(x, V(x, p))-p \cdot V_{x}(x, p) \\
& =-L_{x}(x, V(x, p))
\end{aligned}
$$

and similarly,

$$
\begin{aligned}
H_{p}(x, p) & =V(x, p)+p \cdot V_{p}(x, p)-L_{v}(x, V(x, p)) V_{p}(x, p) \\
& =V(x, p)+p \cdot V_{p}(x, p)-p \cdot V_{p}(x, p)=V(x, p) .
\end{aligned}
$$

2. Since $H_{p}(x, p)=V(x, p)=\left[L_{v}(x, \cdot)\right]^{-1}(p)$ and by Remark 38.13, $p \rightarrow$ $V(x, p)$ is smooth with a smooth inverse $L_{v}(x, \cdot)$, it follows that $H$ satisfies Assumption 1. Letting $p=L_{v}(x, v)$ in Eq. (38.28) shows

$$
\begin{aligned}
H\left(x, L_{v}(x, v)\right) & =L_{v}(x, v) \cdot V\left(x, L_{v}(x, v)\right)-L\left(x, V\left(x, L_{v}(x, v)\right)\right) \\
& =L_{v}(x, v) \cdot v-L(x, v)
\end{aligned}
$$

and using this and the definition of $H^{*}$ we find

$$
\begin{aligned}
H^{*}(x, v) & =v \cdot\left[H_{p}(x, \cdot)\right]^{-1}(v)-H\left(x,\left[H_{p}(x, \cdot)\right]^{-1}(v)\right) \\
& =v \cdot L_{v}(x, v)-H\left(x, L_{v}(x, v)\right)=L(x, v) .
\end{aligned}
$$

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$$
\begin{equation*}
q_{x, T}(t):=\phi_{t}\left(x_{0}(T, x), V\left(x_{0}(T, x), \nabla g\left(x_{0}(T, x)\right)\right)\right. \tag{38.32}
\end{equation*}
$$

so that $q_{x, T}$ solves the Euler Lagrange equations, $q_{x, T}(T)=x, q_{x, T}(0)=$ $x_{0}(T, x)$ and $\dot{q}_{x, T}(0)=V\left(x_{0}(T, x), \nabla g\left(x_{0}(T, x)\right)\right.$ or equivalently

$$
\partial_{v} L\left(q_{x, T}(0), \dot{q}_{x, T}(0)\right)=\nabla g\left(x_{0}(T, x)\right)
$$

Then the function

$$
\begin{equation*}
S(T, x):=I_{T}^{g}\left(q_{x, T}\right)=g\left(q_{x, T}(0)\right)+\int_{0}^{T} L\left(q_{x, T}(t), \dot{q}_{x, T}(t)\right) d t \tag{38.33}
\end{equation*}
$$

solves the Hamilton Jacobi Equation (38.17).
Conjecture 38.19. For general $g$ and $L$ convex in $v$, the function

$$
S(t, x)=\inf _{q \in C^{2}\left([0, t], \mathbb{R}^{n}\right)}\left\{g(q(0))+\int_{0}^{t} L(q(\tau), \dot{q}(\tau)) d \tau: q(t)=x\right\}
$$

is a distributional solution to the Hamilton Jacobi Equation Eq. 38.17. See Evans to learn more about this conjecture.

Proof. We will give two proofs of this Theorem.
First Proof. One need only observe that the theorem is a consequence of Definition 38.14 and Proposition 38.15 and 38.5 .

Second Direct Proof. By the fundamental theorem of calculus and differentiating past the integral,

$$
\begin{align*}
\frac{\partial S(T, x)}{\partial T} & =\nabla g\left(x_{0}(T, x)\right) \cdot \frac{\partial}{\partial T} x_{0}(T, x)+L\left(q_{x, T}(T), \dot{q}_{x, T}(T)\right) \\
& \left.+\int_{0}^{T} \frac{\partial}{\partial T} L\left(q_{x, T}(t), \dot{q}_{x, T}(t)\right)\right) d t \\
& =\nabla g\left(x_{0}(T, x)\right) \cdot \frac{\partial}{\partial T} x_{0}(T, x)+L\left(q_{x, T}(T), \dot{q}_{x, T}(T)\right) \\
& +D I_{T}\left(q_{x, T}\right)\left[\frac{\partial}{\partial T} q_{x, T}\right] \\
& =L\left(q_{x, T}(T), \dot{q}_{x, T}(T)\right)+D I_{T}^{g}\left(q_{x, T}\right)\left[\frac{\partial}{\partial T} q_{x, T}\right] \tag{38.34}
\end{align*}
$$

Using Proposition 38.8 and the fact that $q_{x, T}$ satisfies the Euler Lagrange equations and the boundary conditions in Corollary 38.10 we find

$$
\begin{equation*}
D I_{T}^{g}\left(q_{x, T}\right)\left[\frac{\partial}{\partial T} q_{x, T}\right]=\left.\left(D_{2} L\left(q_{x, T}(t), \dot{q}_{x, T}(t)\right) \frac{\partial}{\partial T} q_{x, T}(t)\right)\right|_{t=T} \tag{38.35}
\end{equation*}
$$

Furthermore differentiating the identity, $q_{x, T}(T)=x$, in $T$ implies

$$
\begin{equation*}
0=\frac{d}{d T} x=\frac{d}{d T} q_{x, T}(T)=\dot{q}_{x, T}(T)+\left.\frac{d}{d T} q_{x, T}(t)\right|_{t=T} \tag{38.36}
\end{equation*}
$$

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Combining Eqs. (38.34) - (38.36) gives

$$
\begin{equation*}
\left.\frac{\partial S(T, x)}{\partial T}=L\left(x, \dot{q}_{x, T}(T)\right)\right)-D_{2} L\left(x, \dot{q}_{x, T}(T)\right) \dot{q}_{x, T}(T) \tag{38.37}
\end{equation*}
$$

Similarly for $v \in \mathbb{R}^{n}$,

$$
\begin{aligned}
\partial_{v} S(T, x) & =\partial_{v} I_{T}^{g}\left(q_{x, T}\right)=D I_{T}^{g}\left(\left(q_{x, T}\right)\right)\left[\partial_{v} q_{x, T}\right] \\
& =D_{2} L\left(q_{x, T}(T), \dot{q}_{x, T}(T)\right) \partial_{v} q_{x, T}(T)=D_{2} L\left(x, \dot{q}_{x, T}(T)\right) v
\end{aligned}
$$

wherein the last equality we have use $q_{x, T}(T)=x$. This last equation is equivalent to

$$
D_{2} L\left(x, \dot{q}_{x, T}(T)\right)=\nabla_{x} S(T, x)
$$

from which it follows that

$$
\begin{equation*}
\dot{q}_{x, T}(T)=V\left(x, \nabla_{x} S(T, x)\right) \tag{38.38}
\end{equation*}
$$

Combining Eqs. (38.37) and (38.38) and the definition of $H$, shows

$$
\begin{aligned}
\frac{\partial S(T, x)}{\partial T} & =L\left(x, V\left(x, \nabla_{x} S(T, x)\right)\right)-D_{2} L\left(x, \dot{q}_{x, T}(T)\right) V\left(x, \nabla_{x} S(T, x)\right) \\
& =-H\left(x, \nabla_{x} S(T, x)\right)
\end{aligned}
$$

Remark 38.20. The hypothesis of Theorem 38.18 may always be satisfied locally, for let $\psi: \mathbb{R} \times \mathbb{R}^{n} \rightarrow \mathbb{R} \times \mathbb{R}^{n}$ be given by $\psi(t, y):=\left(t, \phi_{t}(y, V(y, \nabla g(y))\right.$. Then $\psi(0, y):=(0, y)$ and so

$$
\dot{\psi}(0, y)=(1, *) \text { and } \psi_{y}(0, y)=i d_{\mathbb{R}^{n}}
$$

from which it follows that $\psi^{\prime}(0, y)^{-1}$ exists for all $y \in \mathbb{R}^{n}$. So the inverse function theorem guarantees for each $a \in \mathbb{R}^{n}$ that there exists an open interval $J \subset \mathbb{R}$ with $0 \in J$ and $a \in U \subset_{o} \mathbb{R}^{n}$ and a smooth function $x_{0}: J \times U \rightarrow \mathbb{R}^{n}$ such that

$$
\psi\left(T, x_{0}(T, x)\right)=\left(T, x_{0}(T, x)\right) \text { for } T \in J \text { and } x \in U
$$

i.e.

$$
\phi_{T}\left(x_{0}(T, x), V\left(x_{0}(T, x), \nabla g\left(x_{0}(T, x)\right)\right)=x\right.
$$

### 38.2 Geometric meaning of the Legendre Transform

Let $V$ be a finite dimensional real vector space and $f: V \rightarrow \mathbb{R}$ be a strictly convex function. Then the function $f^{*}: V^{*} \rightarrow \mathbb{R}$ defined by

$$
\begin{equation*}
f^{*}(\alpha)=\sup _{v \in V}(\alpha(v)-f(v)) \tag{38.39}
\end{equation*}
$$

is called the Legendre transform of $f$. Now suppose the supremum on the right side of Eq. (38.39) is obtained at a point $v \in V$, see Figure 38.2 below. Eq. (38.39) may be rewritten as $f^{*}(\alpha) \geq \alpha(\cdot)-f(\cdot)$ with equality at $v$ or equivalently that

$$
-f^{*}(\alpha)+\alpha(\cdot) \leq f(\cdot) \text { with equality at some point } v \in V
$$

Geometrically, the graph of $\alpha \in V^{*}$ defines a hyperplane which if translate by $-f^{*}(\alpha)$ just touches the graph of $f$ at one point, say $v$, see Figure 38.2. At the point of contact, $v, \alpha$ and $f$ must have the same tangent plane and


Fig. 38.2. Legendre Transform of $f$.
since $\alpha$ is linear this means that $f^{\prime}(v)=\alpha$. Therefore the Legendre transform $f^{*}: V^{*} \rightarrow \mathbb{R}$ of $f$ may be given explicitly by

$$
f^{*}(\alpha)=\alpha(v)-f(v) \text { with } v \text { such that } f^{\prime}(v)=\alpha
$$

## Cauchy - Kovalevskaya Theorem

As a warm up we will start with the corresponding result for ordinary differential equations.

Theorem 39.1 (ODE Version of Cauchy - Kovalevskaya, I.). Suppose $a>0$ and $f:(-a, a) \rightarrow \mathbb{R}$ is real analytic near 0 and $u(t)$ is the unique solution to the $O D E$

$$
\begin{equation*}
\dot{u}(t)=f(u(t)) \text { with } u(0)=0 . \tag{39.1}
\end{equation*}
$$

Then $u$ is also real analytic near 0 .
We will give four proofs. However it is the last proof that the reader should focus on for understanding the PDE version of Theorem 39.1.

Proof. (First Proof.)If $f(0)=0$, then $u(t)=0$ for all $t$ is the unique solution to Eq. (39.1) which is clearly analytic. So we may now assume that $f(0) \neq 0$. Let $G(z):=\int_{0}^{z} \frac{1}{f(u)} d u$, another real analytic function near 0 . Then as usual we have

$$
\frac{d}{d t} G(u(t))=\frac{1}{f(u(t))} \dot{u}(t)=1
$$

and hence $G(u(t))=t$. We then have $u(t)=G^{-1}(t)$ which is real analytic near $t=0$ since $G^{\prime}(0)=\frac{1}{f(0)} \neq 0$.

Proof. (Second Proof.) For $z \in \mathbb{C}$ let $u_{z}(t)$ denote the solution to the ODE

$$
\begin{equation*}
\dot{u}_{z}(t)=z f\left(u_{z}(t)\right) \text { with } u_{z}(0)=0 . \tag{39.2}
\end{equation*}
$$

Notice that if $u(t)$ is analytic, then $t \rightarrow u(t z)$ satisfies the same equation as $u_{z}$. Since $G(z, u)=z f(u)$ is holomorphic in $z$ and $u$, it follows that $u_{z}$ in Eq. (39.2) depends holomorphically on $z$ as can be seen by showing $\bar{\partial}_{z} u_{z}=0$, i.e. showing $z \rightarrow u_{z}$ satisfies the Cauchy Riemann equations. Therefore if $\epsilon>0$ is chosen small enough such that Eq. (39.2) has a solution for $|t|<\epsilon$ and $|z|<2$, then

$$
\begin{equation*}
u(t)=u_{1}(t)=\left.\sum_{n=0}^{\infty} \frac{1^{n}}{n!} \partial_{z}^{n} u_{z}(t)\right|_{z=0} . \tag{39.3}
\end{equation*}
$$

Now when $z \in \mathbb{R}, u_{z}(t)=u(t z)$ and therefore

$$
\left.\partial_{z}^{n} u_{z}(t)\right|_{z=0}=\left.\partial_{z}^{n} u(t z)\right|_{z=0}=u^{(n)}(0) t^{n} .
$$

Putting this back in Eq. (39.3) shows

$$
u(t)=\sum_{n=0}^{\infty} \frac{1}{n!} u^{(n)}(0) t^{n}
$$

which shows $u(t)$ is analytic for $t$ near 0 .
Proof. (Third Proof.) Go back to the original proof of existence of solutions, but now replace $t$ by $z \in \mathbb{C}$ and $\int_{0}^{t} f(u(\tau)) d \tau$ by $\int_{0}^{z} f(u(\xi)) d \xi=$ $\int_{0}^{1} f(u(t z)) z d t$. Then the usual Picard iterates proof work in the class of holomorphic functions to give a holomorphic function $u(z)$ solving Eq. (39.1).

Proof. (Fourth Proof: Method of Majorants) Suppose for the moment we have an analytic solution to Eq. (39.1). Then by repeatedly differentiating Eq. (39.1) we learn

$$
\begin{aligned}
\ddot{u}(t) & =f^{\prime}(u(t)) \dot{u}(t)=f^{\prime}(u(t)) f(u(t)) \\
u^{(3)}(t) & =f^{\prime \prime}(u(t)) f^{2}(u(t))+\left[f^{\prime}(u(t))\right]^{2} f(u(t)) \\
& \vdots \\
u^{(n)}(t) & =p_{n}\left(f(u(t)), \ldots, f^{(n-1)}(u(t))\right)
\end{aligned}
$$

where $p_{n}$ is a polynomial in $n$ variables with all non-negative integer coefficients. The first few polynomials are $p_{1}(x)=x, p_{2}(x, y)=x y, p_{3}(x, y, z)=$ $x^{2} z+x y^{2}$. Notice that these polynomials are universal, i.e. are independent of the function $f$ and

$$
\begin{aligned}
\left|u^{(n)}(0)\right| & =\left|p_{n}\left(f(0), \ldots, f^{(n-1)}(0)\right)\right| \\
& \leq p_{n}\left(|f(0)|, \ldots,\left|f^{(n-1)}(0)\right|\right) \leq p_{n}\left(g(0), \ldots, g^{(n-1)}(0)\right)
\end{aligned}
$$

where $g$ is any analytic function such that $\left|f^{(k)}(0)\right| \leq g^{(k)}(0)$ for all $k \in \mathbb{Z}_{+}$. (We will abbreviate this last condition as $f \ll g$.) Now suppose that $v(t)$ is a solution to

$$
\begin{equation*}
\dot{v}(t)=g(v(t)) \text { with } v(0)=0 \tag{39.4}
\end{equation*}
$$

then we know from above that

$$
v^{(n)}(0)=p_{n}\left(g(0), \ldots, g^{(n-1)}(0)\right) \geq\left|u^{(n)}(0)\right| \text { for all } n .
$$

Hence if knew that $v$ were analytic with radius of convergence larger that some $\rho>0$, then by comparison we would find

$$
\sum_{n=0}^{\infty} \frac{1}{n!}\left|u^{(n)}(0)\right| \rho^{n} \leq \sum_{n=0}^{\infty} \frac{1}{n!} v^{(n)}(0) \rho^{n}<\infty
$$

and this would show

$$
u(t):=\sum_{n=0}^{\infty} \frac{1}{n!} p_{n}\left(f(0), \ldots, f^{(n-1)}(0)\right) t^{n}
$$

is a well defined analytic function for $|t|<\rho$.
I now claim that $u(t)$ solves Eq. (39.1). Indeed, both sides of Eq. (39.1) are analytic in $t$, so it suffices to show the derivatives of each side of Eq. (39.1) agree at $t=0$. For example $\dot{u}(0)=f(0), \ddot{u}(0)=\left.\frac{d}{d t}\right|_{0} f(u(t))$, etc. However this is the case by the very definition of $u^{(n)}(0)$ for all $n$.

So to finish the proof, it suffices to find an analytic function $g$ such that $\left|f^{(k)}(0)\right| \leq g^{(k)}(0)$ for all $k \in \mathbb{Z}_{+}$and for which we know the solution to Eq. (39.4) is analytic about $t=0$. To this end, suppose that the power series expansion for $f(t)$ at $t=0$ has radius of convergence larger than $r>0$, then $\sum_{n=0}^{\infty} \frac{1}{n!} f^{(n)}(0) r^{n}$ is convergent and in particular,

$$
C:=\max _{n}\left|\frac{1}{n!} f^{(n)}(0) r^{n}\right|<\infty
$$

from which we conclude

$$
\max _{n}\left|\frac{1}{n!} f^{(n)}(0)\right| \leq C r^{-n}
$$

Let

$$
g(u):=\sum_{n=0}^{\infty} C r^{-n} u^{n}=C \frac{1}{1-u / r}=C \frac{r}{r-u}
$$

Then clearly $f \ll g$. To conclude the proof, we will explicitly solve Eq. (39.4) with this function $g(t)$,

$$
\dot{v}(t)=C \frac{r}{r-v(t)} \text { with } v(0)=0
$$

By the usual separation of variables methods we find $r v(t)-\frac{1}{2} v^{2}(t)=C r t$, i.e.

$$
2 C r t-2 r v(t)+v^{2}(t)=0
$$

which has solutions, $v(t)=r \pm \sqrt{r^{2}-2 C r t}$. We must take the negative sign to get the correct initial condition, so that

$$
\begin{equation*}
v(t)=r-\sqrt{r^{2}-2 C r t}=r-r \sqrt{1-2 C t / r} \tag{39.5}
\end{equation*}
$$

which is real analytic for $|t|<\rho:=r / C$.
Let us now Jazz up this theorem to that case of a system of ordinary differential equations. For this we will need the following lemma.

Lemma 39.2. Suppose $h:(-a, a)^{d} \rightarrow \mathbb{R}^{d}$ is real analytic near $0 \in(-a, a)^{d}$, then

$$
h \ll \frac{C r}{r-z_{1}-\cdots-z_{d}}
$$

for some constants $C$ and $r$.
Proof. By definition, there exists $\rho>0$ such that

$$
h(z)=\sum_{\alpha} h_{\alpha} z^{\alpha} \text { for }|z|<\rho
$$

where $h_{\alpha}=\frac{1}{\alpha!} \partial^{\alpha} h(0)$. Taking $z=r(1,1, \ldots, 1)$ with $r<\rho$ implies there exists $C<\infty$ such that $\left|h_{\alpha}\right| r^{|\alpha|} \leq C$ for all $\alpha$, i.e.

$$
\left|h_{\alpha}\right| \leq C r^{-|\alpha|} \leq C \frac{|\alpha|!}{\alpha!} r^{-|\alpha|}
$$

This completes the proof since

$$
\begin{aligned}
\sum_{\alpha} C \frac{|\alpha|!}{\alpha!} r^{-|\alpha|} z^{\alpha} & =C \sum_{n=0}^{\infty} \sum_{|\alpha|=n} \frac{|\alpha|!}{\alpha!}\left(\frac{z}{r}\right)^{\alpha}=C \sum_{n=0}^{\infty}\left(\frac{z_{1}+\cdots+z_{d}}{r}\right)^{n} \\
& =C \frac{1}{1-\left(\frac{z_{1}+\cdots+z_{d}}{r}\right)}=\frac{C r}{r-z_{1}-\cdots-z_{d}}
\end{aligned}
$$

all of which is valid provided $|z|:=\left|z_{1}\right|+\cdots+\left|z_{d}\right|<r$.
Theorem 39.3 (ODE Version of Cauchy - Kovalevskaya, II.). Suppose $a>0$ and $f:(-a, a)^{d} \rightarrow \mathbb{R}^{d}$ be real analytic near $0 \in(-a, a)^{d}$ and $u(t)$ is the unique solution to the $O D E$

$$
\begin{equation*}
\dot{u}(t)=f(u(t)) \text { with } u(0)=0 . \tag{39.6}
\end{equation*}
$$

Then $u$ is also real analytic near 0 .
Proof. All but the first proof of Theorem 39.1 may be adapted to the cover this case. The only proof which perhaps needs a little more comment is the fourth proof. By Lemma 39.2, we can find $C, r>0$ such that

$$
f_{j}(z) \ll g_{j}(z):=\frac{C r}{r-z_{1}-\cdots-z_{d}}
$$

for all $j$. Let $v(t)$ denote the solution to the ODE,

$$
\begin{equation*}
\dot{v}(t)=g(v(t))=\frac{C r}{r-v_{1}(t)-\cdots-v_{d}(t)}(1,1, \ldots, 1) \tag{39.7}
\end{equation*}
$$

with $v(0)=0$. By symmetry, $v_{j}(t)=v_{1}(t)=: w(t)$ for each $j$ so Eq. (39.7) implies

$$
\dot{w}(t)=\frac{C r}{r-d w(t)}=\frac{C(r / d)}{(r / d)-w(t)} \text { with } w(0)=0 .
$$

We have already solved this equation (see Eq. (39.5) with $r$ replaced by $r / d$ ) to find

$$
\begin{equation*}
w(t)=r / d-\sqrt{r^{2} / d^{2}-2 C r t / d}=r / d(1-\sqrt{1-2 C d t / r}) . \tag{39.8}
\end{equation*}
$$

Thus $v(t)=w(t)(1,1, \ldots, 1)$ is a real analytic function which is convergent for $|t|<r /(2 C d)$.

Now suppose that $u$ is a real analytic solution to Eq. (39.6). Then by repeatedly differentiating Eq. (39.6) we learn

$$
\begin{align*}
\ddot{u}_{j}(t) & =\partial_{i} f_{j}(u(t)) \dot{u}_{i}(t)=\partial_{i} f_{j}(u(t)) f_{i}(u(t)) \\
u_{j}^{(3)}(t) & =\partial_{k} \partial_{i} f_{j}(u(t)) \dot{u}_{k}(t) \dot{u}_{i}(t)+\partial_{i} f_{j}(u(t)) \ddot{u}_{i}(t) \\
& \vdots  \tag{39.9}\\
u_{j}^{(n)}(t) & =p_{n}\left(\left\{\partial^{\alpha} f_{j}(u(t))\right\}_{|\alpha|<n},\left\{u_{i}^{(k)}(t)\right\}_{k<n, 1 \leq i \leq d}\right)
\end{align*}
$$

where $p_{n}$ is a polynomial with all non-negative integer coefficients. We now define $u_{j}^{(n)}(0)$ inductively so that

$$
u_{j}^{(n)}(0)=p_{n}\left(\left\{\partial^{\alpha} f_{j}(u(0))\right\}_{|\alpha|<n},\left\{u_{i}^{(k)}(0)\right\}_{k<n, 1 \leq i \leq d}\right)
$$

for all $n$ and $j$ and we will attempt to define

$$
\begin{equation*}
u(t)=\sum_{n=0}^{\infty} \frac{1}{n!} u^{(n)}(0) t^{n} \tag{39.10}
\end{equation*}
$$

To see this sum is convergent we make use of the fact that the polynomials $p_{n}$ are universal i.e. are independent of the function $f_{j}$ ) and have non-negative coefficients so that by induction

$$
\begin{aligned}
\left|u_{j}^{(n)}(0)\right| & \leq p_{n}\left(\left\{\left|\partial^{\alpha} f_{j}(u(0))\right|\right\}_{|\alpha|<n},\left\{\left|u_{i}^{(k)}(0)\right|\right\}_{k<n, 1 \leq i \leq d}\right) \\
& \leq p_{n}\left(\left\{\partial^{\alpha} g_{j}(u(0))\right\}_{|\alpha|<n},\left\{v_{i}^{(k)}(0)\right\}_{k<n, 1 \leq i \leq d}\right)=v_{j}^{(n)}(0)
\end{aligned}
$$

Notice the when $n=0$ that $\left|u_{j}(0)\right|=0=v_{j}(0) .{ }^{1}$ Thus we have shown $u \ll v$ and so by comparison the sum in Eq. (39.10) is convergent for $t$ near 0 . As before $u(t)$ solves Eq. (39.6) since both functions $\dot{u}(t)$ and $f(u(t))$ are analytic functions of $t$ which have common values for all derivatives in $t$ at $t=0$.

[^0]
### 39.1 PDE Cauchy Kovalevskaya Theorem

In this section we will consider the following general quasi-linear system of partial differential equations

$$
\begin{equation*}
\sum_{|\alpha|=k} a_{\alpha}\left(x, J^{k-1} u\right) \partial_{x}^{\alpha} u(x)+c\left(x, J^{k-1} u\right)=0 \tag{39.11}
\end{equation*}
$$

where

$$
J^{l} u(x)=\left(u(x), D u(x), D^{2} u(x), \ldots, D^{l} u(x)\right)
$$

is the " $l$ - jet" of $u$. Here $u: \mathbb{R}^{n} \rightarrow \mathbb{R}^{m}$ and $a_{\alpha}\left(J^{k-1} u, x\right)$ is an $m \times m$ matrix. As usual we will want to give boundary data on some hypersurface $\Sigma \subset \mathbb{R}^{n}$. Let $\nu$ denote a smooth vector field along $\Sigma$ such that $v(x) \notin T_{x} \Sigma\left(T_{x} \Sigma\right.$ is the tangent space to $\Sigma$ at $x$ ) for $x \in \Sigma$. For example we might take $\nu(x)$ to be orthogonal to $T_{x} \Sigma$ for all $x \in \Sigma$. To hope to get a unique solution to Eq. (39.11) we will further assume there are smooth functions $g_{l}$ on $\Sigma$ for $l=0, \ldots, k-1$ and we will require

$$
\begin{equation*}
D^{l} u(x)(v(x), \ldots, v(x))=g_{l}(x) \text { for } x \in \Sigma \text { and } l=0, \ldots, k-1 . \tag{39.12}
\end{equation*}
$$

Proposition 39.4. Given a smooth function $u$ on a neighborhood of $\Sigma$ satisfying Eq. (39.12), we may calculate $D^{l} u(x)$ for $x \in \Sigma$ and $l<k$ in terms of the functions $g_{l}$ and there tangential derivatives.

Proof. Let us begin by choosing a coordinate system $y$ on $\mathbb{R}^{n}$ such that $\Sigma \cap D(y)=\left\{y_{n}=0\right\}$ and let us extend $\nu$ to a neighborhood of $\Sigma$ by requiring $\frac{\partial \nu}{\partial y_{n}}=0$. To complete the proof, we are going to show by induction on $k$ that we may compute

$$
\left(\frac{\partial}{\partial y}\right)^{\alpha} u(x) \text { for all } x \in \Sigma \text { and }|\alpha|<k
$$

from Eq. (39.12).
The claim is clear when $k=1$, since $u=g_{0}$ on $\Sigma$. Now suppose that $k=2$ and let $\nu_{i}=\nu_{i}\left(y_{1}, \ldots, y_{n-1}\right)$ such that

$$
\nu=\sum_{i=1}^{n} \nu_{i} \frac{\partial}{\partial y_{i}} \text { in a neighborhood of } \Sigma .
$$

Then

$$
g_{1}=(D u) \nu=\nu u=\sum_{i=1}^{n} \nu_{i} \frac{\partial u}{\partial y_{i}}=\sum_{i<n} \nu_{i} \frac{\partial g_{0}}{\partial y_{i}}+\nu_{n} \frac{\partial u}{\partial y_{n}} .
$$

Since $\nu$ is not tangential to $\Sigma=\left\{y_{n}=0\right\}$, it follows that $\nu_{n} \neq 0$ and hence

$$
\begin{equation*}
\frac{\partial u}{\partial y_{n}}=\frac{1}{\nu_{n}}\left(g_{1}-\sum_{i<n} \nu_{i} \frac{\partial g_{0}}{\partial y_{i}}\right) \text { on } \Sigma \tag{39.13}
\end{equation*}
$$

39 Cauchy - Kovalevskaya Theorem

For $k=3$, first observe from the equality $u=g_{0}$ on $\Sigma$ and Eq. (39.13) we may compute all derivatives of $u$ of the form $\frac{\partial^{\alpha} u}{\partial y^{\alpha}}$ on $\Sigma$ provided $\alpha_{n} \leq 1$. From Eq. (39.12) for $l=2$, we have

$$
\begin{aligned}
g_{2} & =\left(D^{2} u\right)(v, v)=v^{2} u+\text { l.o.ts. } \\
& =\sum \nu_{j} \frac{\partial}{\partial y_{j}}\left(\nu_{i} \frac{\partial u}{\partial y_{i}}\right)+\text { l.o.ts. }=\nu_{n}^{2} \frac{\partial^{2} u}{\partial y_{n}^{2}}+\text { l.o.ts. }
\end{aligned}
$$

where l.o.ts. denotes terms involving $\frac{\partial^{\alpha} u}{\partial y^{\alpha}}$ with $\alpha_{n} \leq 1$. From this result, it follows that we may compute $\frac{\partial^{2} u}{\partial y_{n}^{2}}$ in terms of derivatives of $g_{0}, g_{1}$ and $g_{2}$. The reader is asked to finish the full inductive argument of the proof.

Remark 39.5. The above argument shows that from Eq. (39.12) we may compute $\frac{\partial^{\alpha} u}{\partial y^{\alpha}}$ for any $\alpha$ such that $\alpha_{n}<k$.

To study Eq. (39.11) in more detail, let us rewrite Eq. (39.11) in the $y$ coordinates. Using the product and the chain rule repeatedly Eq. (39.11) may be written as

$$
\begin{equation*}
\sum_{|\alpha|=k} b_{\alpha}\left(y, J^{k-1} u\right) \partial_{y}^{\alpha} u(y)+c\left(y, J^{k-1} u\right)=0 \tag{39.14}
\end{equation*}
$$

where

$$
J^{l} u(y)=\left(u(y), D u(y), D^{2} u(y), \ldots, D^{l} u(y)\right)
$$

We will be especially concerned with the $b_{(0,0, \ldots, 0, k)}$ coefficient which can be determined as follows:

$$
\begin{aligned}
\sum_{|\alpha|=k} a_{\alpha}\left(\frac{\partial}{\partial x}\right)^{\alpha} & =\sum_{|\alpha|=k} a_{\alpha}\left(\sum_{j=1}^{n} \frac{\partial y_{j}}{\partial x} \frac{\partial}{\partial y_{j}}\right)^{\alpha}=\sum_{|\alpha|=k} a_{\alpha}\left(\frac{\partial y_{n}}{\partial x} \frac{\partial}{\partial y_{n}}\right)^{\alpha}+\text { l.o.ts. } \\
& =\sum_{|\alpha|=k} a_{\alpha}\left(\frac{\partial y_{n}}{\partial x}\right)^{\alpha}\left(\frac{\partial}{\partial y_{n}}\right)^{k}+\text { l.o.ts. }
\end{aligned}
$$

where l.o.ts. now denotes terms involving $\frac{\partial^{\alpha} u}{\partial y^{\alpha}}$ with $\alpha_{n}<k$. From this equation we learn that

$$
b_{(0,0, \ldots, 0, k)}\left(y, J^{k-1} u\right)=\sum_{|\alpha|=k} a_{\alpha}\left(\frac{\partial y_{n}}{\partial x}\right)^{\alpha}=\sum_{|\alpha|=k} a_{\alpha}\left(d y_{n}\left(\frac{\partial}{\partial x}\right)\right)^{\alpha} .
$$

Definition 39.6. We will say that boundary data $\left(\Sigma, g_{0}, \ldots, g_{k-1}\right)$ is noncharacteristic for Eq. (39.11) at $x \in \Sigma$ if

$$
b_{(0,0, \ldots, 0, k)}\left(y, J^{k-1} u\right)=\sum_{|\alpha|=k} a_{\alpha}\left(x, J^{k-1} u(x)\right)\left(d y_{n}\left(\frac{\partial}{\partial x}\right)\right)^{\alpha}
$$

is invertible at $x$.

Notice that this condition is independent of the choice of coordinate system $y$. To see this, for $\xi \in\left(\mathbb{R}^{n}\right)^{*}$ let

$$
\sigma(\xi)=\sum_{|\alpha|=k} a_{\alpha}\left(x, J^{k-1} u(x)\right)\left(\xi\left(\frac{\partial}{\partial x}\right)\right)^{\alpha}
$$

which is $k$ - linear form on $\left(\mathbb{R}^{n}\right)^{*}$. This form is coordinate independent since if $f$ is a smooth function such that $f(x)=0$ and $d f_{x}=\xi$, then

$$
\sigma(\xi)=\left.\frac{1}{k!} \sum_{|\alpha|=k} a_{\alpha}\left(x, J^{k-1} u(x)\right)\left(\frac{\partial}{\partial x}\right)^{\alpha} f^{k}\right|_{x}
$$

Noting that

$$
b_{(0,0, \ldots, 0, k)}\left(y, J^{k-1} u\right)=\sigma\left(d y_{n}\right)
$$

our non-characteristic condition becomes, $\sigma\left(d y_{n}\right)$ is invertible. Finally $d y_{n}$ is the unique element $\xi$ of $\left(\mathbb{R}^{n}\right)^{*} \backslash\{0\}$ up to scaling such that $\left.\xi\right|_{T_{x} \Sigma} \equiv 0$. So the non-characteristic condition may be written invariantly as $\sigma(\xi)$ is invertible for all (or any) $\xi \in\left(\mathbb{R}^{n}\right)^{*} \backslash\{0\}$ such that $\left.\xi\right|_{T_{x} \Sigma} \equiv 0$.

Assuming the given boundary data is non-characteristic, Eq. (39.11) may be put into "standard form,"

$$
\begin{equation*}
\sum_{|\alpha|=k} b_{\alpha}\left(y, J^{k-1} u\right) \partial_{y}^{\alpha} u(y)+c\left(y, J^{k-1} u\right)=0 \tag{39.15}
\end{equation*}
$$

with

$$
\frac{\partial^{l} u}{\partial y_{n}^{l}}=g_{l} \text { on } y_{n}=0 \text { for } l<k
$$

where $b_{(0,0, \ldots, 0, k)}\left(y, J^{k-1} u\right)=I d$ - matrix and

$$
J^{l} u(y)=\left(u(y), D u(y), D^{2} u(y), \ldots, D^{l} u(y)\right) .
$$

By adding new dependent variables and possible a new independent variable for $y_{n}$ one may reduce the problem to solving the system in Eq. (39.20) below. The resulting theorem may be stated as follows.

Theorem 39.7 (Cauchy Kovalevskaya). Suppose all the coefficients in Eq. (39.11) are real analytic and the boundary data in Eq. (39.12) are also real analytic and non-characteristic near some point $a \in \Sigma$. Then there is a unique real analytic solution to Eqs. (39.11) and (39.12). (The boundary data in Eq. (39.12) is said to be real analytic if there exists coordinates $y$ as above which are real analytic and the functions $\nu$ and $g_{l}$ for $l=0, \ldots, k-1$ are real analytic functions in the $y$ - coordinate system.)

Example 39.8. Suppose $a, b, C, r$ are positive constants. We wish to show the solution to the quasi-linear PDE

$$
\begin{equation*}
w_{t}=\frac{C r}{r-y-a w}\left[b w_{y}+1\right] \text { with } w(0, y)=0 \tag{39.16}
\end{equation*}
$$

is real analytic near $(t, y)=(0,0)$. To do this we will solve the equation using the method of characteristics. Let $g(y, z):=\frac{C r}{r-y-a z}$, then the characteristic equations are

$$
\begin{aligned}
t^{\prime} & =0 \text { with } t(0)=0 \\
y^{\prime} & =-b g(y, z) \text { with } y(0)=y_{0} \text { and } \\
z^{\prime} & =g(y, z) \text { with } z(0)=0
\end{aligned}
$$

From these equations we see that we may identify $t$ with $s$ and that $y+b z=y_{0}$. Thus $z(t)=w(t, y(t))$ satisfies

$$
\begin{aligned}
\dot{z} & =g\left(y_{0}-b z, z\right)=\frac{C r}{r-y_{0}+b z-a z} \\
& =\frac{C r}{r-y_{0}+(b-a) z} \text { with } z(0)=0 .
\end{aligned}
$$

Integrating this equation gives

$$
\begin{aligned}
C r t & =\int_{0}^{t}\left(r-y_{0}+(b-a) z(\tau)\right) \dot{z}(\tau) d \tau=\left(r-y_{0}\right) z-\frac{1}{2}(a-b) z^{2} \\
& =(r-y-b z) z-\frac{1}{2}(a-b) z^{2}=(r-y) z-\frac{1}{2}(a+b) z^{2}
\end{aligned}
$$

i.e.

$$
\frac{1}{2}(a+b) z^{2}-(r-y) z+C r t=0
$$

The quadratic formula gives

$$
w(t, y)=\frac{1}{a+b}\left[(r-y) \pm \sqrt{(r-y)^{2}-2(a+b) C r t}\right]
$$

and using $w(0, y)=0$ we conclude

$$
\begin{equation*}
w(t, y)=\frac{1}{a+b}\left[(r-y)-\sqrt{(r-y)^{2}-2(a+b) C r t}\right] . \tag{39.17}
\end{equation*}
$$

Notice the $w$ is real analytic for $(t, y)$ near $(0,0)$.
In general we could use the method of characteristics and ODE properties (as in Example 39.8) to show

$$
u_{t}=a(x, u) u_{x}+b(x, u) \text { with } u(0, x)=g(x)
$$

has local real analytic solutions if $a, b$ and $g$ are real analytic. The method would also work for the fully non-linear case as well. However, the method of characteristics fails for systems while the method we will present here works in this generality.

Exercise 39.9. Verify $w$ in Eq. (39.17) solves Eq. (39.16).
Solution 39.10 (39.9). Let $\rho:=\sqrt{(r-y)^{2}-2(a+b) C r t}$, then

$$
\begin{gathered}
w(t, y)=\frac{1}{a+b}[r-y-\rho]=\frac{r-y}{a+b}-\frac{1}{a+b} \rho, \\
w_{t}=C r / \rho, \rho=r-y-(a+b) w \text { and } \\
b w_{y}+1=\frac{b}{a+b}[-1+(r-y) / \rho]+1=\frac{1}{a+b}[a+b(r-y) / \rho] .
\end{gathered}
$$

Hence

$$
\begin{aligned}
\frac{b w_{y}+1}{w_{t}} & =\frac{1}{(a+b) C r}[\rho a+b(r-y)] \\
& =\frac{1}{(a+b) C r}[(r-y-(a+b) w) a+b(r-y)] \\
& =\frac{1}{C r}[r-y-a w]
\end{aligned}
$$

as desired.
Example 39.11. Now let us solve for

$$
v(t, x)=\left(v^{1}, \ldots, v^{m}\right)\left(t, x_{1}, \ldots, x_{n}\right)
$$

where $v$ satisfies

$$
v_{t}^{j}=\frac{C r}{r-x_{1}-\cdots-x_{n}-\sum_{k=1}^{m} v^{k}}\left[1+\sum_{i=1}^{n} \sum_{k=1}^{m} \partial_{i} v^{k}\right] \text { with } v(0, x)=0 .
$$

By symmetry, $v^{j}=v^{1}=: w(t, y)$ for all $j$ where $y=x_{1}+\cdots+x_{n}$. Since $\partial_{i} v^{j}=w_{y}$, the above equations all may be written as

$$
w_{t}=\frac{C r}{r-y-m w}\left[m n w_{y}+1\right] \text { with } w(0, y)=0
$$

Therefore from Example 39.8 with $a=m$ and $b=m n$, we find

$$
\begin{equation*}
w(t, y)=\frac{1}{m(n+1)}\left[(r-y)-\sqrt{(r-y)^{2}-2 m(n+1) C r t}\right] . \tag{39.18}
\end{equation*}
$$

and hence that

$$
\begin{equation*}
v(t, x)=w\left(t, x_{1}+\cdots+x_{n}\right)(1,1,1, \ldots, 1) \in \mathbb{R}^{m} \tag{39.19}
\end{equation*}
$$

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### 39.2 Proof of Theorem 39.7

As is outlined in Evans, Theorem 39.7 may be reduced to the following theorem.

Theorem 39.12. Let $(t, x, z)=\left(t, x_{1}, \ldots, x_{n}, z_{1}, \ldots, z_{m}\right) \in \mathbb{R} \times \mathbb{R}^{n} \times \mathbb{R}^{m}$ and assume $(t, x, z) \rightarrow B_{j}(t, x, z) \in\{m \times m$-matrices $\}($ for $j=1, \ldots, n)$ and $(t, x, z) \rightarrow c(t, x, z) \in \mathbb{R}^{m}$ are real analytic functions near $(0,0,0) \in$ $\mathbb{R} \times \mathbb{R}^{n} \times \mathbb{R}^{m}$ and $x \rightarrow f(x) \in \mathbb{R}^{m}$ is real analytic near $0 \in \mathbb{R}^{n}$. Then there exists, in a neighborhood of $(t, x)=(0,0) \in \mathbb{R} \times \mathbb{R}^{n}$, a unique real analytic solution $u(t, x) \in \mathbb{R}^{m}$ to the quasi-linear system

$$
\begin{equation*}
u_{t}(t, x)=\sum_{j=1}^{n} B_{j}(t, x, u(t, x)) \partial_{j} u(t, x)+c(t, x, u(t, x)) \text { with } u(0, x)=f(x) . \tag{39.20}
\end{equation*}
$$

Proof. (Sketch.)
Step 0. By replacing $u(t, x)$ by $u(t, x)-f(x)$, we may assume $f \equiv 0$. By letting $u^{m+1}(t, x)=t$ if necessary, we may assume $B_{j}$ and $c$ do not depend on $t$. With these reductions we are left to solve

$$
\begin{equation*}
u_{t}(t, x)=\sum_{j=1}^{n} B_{j}(x, u(t, x)) \partial_{j} u(t, x)+c(x, u(t, x)) \text { with } u(0, x)=0 \tag{39.21}
\end{equation*}
$$

Step 1. Let

$$
g(x, z):=\frac{C r}{r-x_{1}-\cdots-x_{n}-z_{1}-\cdots-z_{m}}
$$

where $C$ and $r$ are positive constants such that

$$
\left(B_{j}\right)_{k l} \ll g \text { and } c_{k} \ll g
$$

for all $k, l, j$. For this choice of $C$ and $r$, let $v$ denote the solution constructed in Example 39.11 above.

Step 2. By repeatedly differentiating Eq. (39.20), show that if $u$ solves Eq. (39.20) then $\partial_{x}^{\alpha} \partial_{t}^{k} u^{j}(0,0)$ is a universal polynomial in the derivatives $\left\{\partial_{t}^{l} \partial_{x}^{\alpha}\right\}_{\alpha, l<k}$ of the entries of $B_{j}$ and $c$ and $u$ with all coefficients being nonnegative. Use this fact and induction to conclude

$$
\left|\partial_{x}^{\alpha} \partial_{t}^{k} u^{j}(0,0)\right| \leq \partial_{x}^{\alpha} \partial_{t}^{k} v^{j}(0,0) \text { for all } \alpha, k \text { and } l .
$$

Step 3. Use the computation in Step 2. to define $\partial_{x}^{\alpha} \partial_{t}^{k} u^{j}(0,0)$ for all $\alpha$ and $k$ and then defined

$$
\begin{equation*}
u(t, x):=\sum_{\alpha, k} \frac{\partial_{x}^{\alpha} \partial_{t}^{k} u(0,0)}{\alpha!k!} t^{k} x^{\alpha} \tag{39.22}
\end{equation*}
$$

Because of step 2. and Example 39.11, this series is convergent for $(t, x)$ sufficiently close to zero.

Step 4. The function $u$ defined in Step 3. solves Eq. (39.20) because both

$$
u_{t}(t, x) \text { and } \sum_{j=1}^{n} B_{j}(x, u(t, x)) \partial_{j} u(t, x)+c(x, u(t, x))
$$

are both real analytic functions in $(t, x)$ each having, by construction, the same derivatives at $(0,0)$.

### 39.3 Examples

Corollary 39.13 (Isothermal Coordinates). Suppose that we are given $a$ metric $d s^{2}=E d x^{2}+2 F d x d y+G d y^{2}$ on $\mathbb{R}^{2}$ such that $G / E$ and $F / E$ are real analytic near $(0,0)$. Then there exists a complex function $u$ and $a$ positive function $\rho$ such that $D u(0,0)$ is invertible and $d s^{2}=\rho|d u|^{2}$ where $d u=u_{x} d x+u_{y} d y$.

Proof. Working out $|d u|^{2}$ gives

$$
|d u|^{2}=\left|u_{x}\right|^{2} d x^{2}+2 \operatorname{Re}\left(u_{x} \bar{u}_{y}\right) d x d y+\left|u_{y}\right|^{2} d y^{2}
$$

Writing $u_{y}=\lambda u_{x}$, the previous equation becomes

$$
|d u|^{2}=\left|u_{x}\right|^{2}\left(d x^{2}+2 \operatorname{Re}(\lambda) d x d y+|\lambda|^{2} d y^{2}\right)
$$

Hence we must have

$$
E=\rho\left|u_{x}\right|^{2}, F=\rho\left|u_{x}\right|^{2} \operatorname{Re} \lambda \text { and } G=\rho\left|u_{x}\right|^{2}|\lambda|^{2}
$$

or equivalently

$$
F / E=\operatorname{Re} \lambda \text { and } G / E=|\lambda|^{2}
$$

Writing $\lambda=a+i b$, we find $a=F / E$ and $a^{2}+b^{2}=G / E$ so that

$$
\lambda=\frac{F}{E} \pm i \sqrt{G / E-(F / E)^{2}}=\frac{1}{E}\left(F \pm i \sqrt{G E-F^{2}}\right) .
$$

We make a choice of the sign above, then we are looking for $u(x, y) \in \mathbb{C}$ such that $u_{y}=\lambda u_{x}$. Letting $u=\alpha+i \beta$, the equation $u_{y}=\lambda u_{x}$ may be written as the system of real equations

$$
\begin{aligned}
\alpha_{y} & =\operatorname{Re}\left[(a+i b)\left(\alpha_{x}+i \beta_{x}\right)\right]=a \alpha_{x}-b \beta_{x} \text { and } \\
\beta_{y} & =\operatorname{Im}\left[(a+i b)\left(\alpha_{x}+i \beta_{x}\right)\right]=a \beta_{x}+b \alpha_{x}
\end{aligned}
$$

which is equivalent to

$$
\binom{\alpha}{\beta}_{y}=\left(\begin{array}{cc}
a & -b \\
b & a
\end{array}\right)\binom{\alpha}{\beta}_{x} .
$$

So we may apply the Cauchy Kovalevskaya theorem 39.12 with $t=y$ to find a real analytic solution to this equation with (say) $u(x, 0)=x$, i.e. $\alpha(x, 0)=x$ and $\beta(x, 0)=0$. (We could take $u(x, 0)=f(x)$ for any real analytic function $f$ such that $f^{\prime}(0) \neq 0$.) The only thing that remains to check is that $D u(0,0)$ is invertible. But

$$
\begin{aligned}
D u(0,0) & =\left(\begin{array}{ll}
\operatorname{Re} u_{x} & \operatorname{Re} u_{y} \\
\operatorname{Im} u_{x} & \operatorname{Im} u_{y}
\end{array}\right)=\left(\begin{array}{cc}
\alpha_{x} & \alpha_{y} \\
\beta_{x} & \beta_{y}
\end{array}\right) \\
& =\left(\begin{array}{cc}
\alpha_{x} & a \alpha_{x}-b \beta_{x} \\
\beta_{x} & a \beta_{x}+b \alpha_{x}
\end{array}\right)
\end{aligned}
$$

so that

$$
\operatorname{det}[D u]=b\left(\alpha_{x}^{2}+\beta_{x}^{2}\right)=\operatorname{Im} \lambda\left|u_{x}\right|^{2}
$$

Thus

$$
\operatorname{det}[D u(0,0)]=\operatorname{Im} \lambda(0,0)= \pm\left.\sqrt{G / E-(F / E)^{2}}\right|_{(0,0)} \neq 0
$$

Example 39.14. Consider the linear PDE,

$$
\begin{equation*}
u_{y}=u_{x} \text { with } u(x, 0)=f(x) \tag{39.23}
\end{equation*}
$$

where $f(x)=\sum_{m=0}^{\infty} a_{m} x^{m}$ as real analytic function near $x=0$ with radius of convergence $\rho$. (So for any $r<\rho,\left|a_{m}\right| \leq C r^{-n}$.) Formally the solution to Eq. (39.23) should be given by

$$
u(x, y)=\left.\sum_{n=0}^{\infty} \frac{1}{n!} \partial_{y}^{n} u(x, y)\right|_{y=0} y^{n} .
$$

Now using the PDE (39.23),

$$
\left.\partial_{y}^{n} u(x, y)\right|_{y=0}=\partial_{x}^{n} u(x, 0)=f^{(n)}(x) .
$$

Thus we get

$$
\begin{equation*}
u(x, y)=\sum_{n=0}^{\infty} \frac{1}{n!} f^{(n)}(x) y^{n} \tag{39.24}
\end{equation*}
$$

By the Cauchy estimates,

$$
\left|f^{(n)}(x)\right| \leq \frac{n!\rho}{(\rho-|x|)^{n+1}}
$$

and so

$$
\sum_{n=0}^{\infty} \frac{1}{n!}\left|f^{(n)}(x) y^{n}\right| \leq \rho \sum_{n=0}^{\infty} \frac{|y|^{n}}{(\rho-|x|)^{n+1}}
$$

which is finite provided $|y|<\rho-|x|$, i.e. $|x|+|y|<\rho$. This of course makes sense because we know the solution to Eq. (39.23) is given by

$$
u(x, y)=f(x+y)
$$

Now we can expand Eq. (39.24) out to find

$$
\begin{align*}
u(x, y) & =\sum_{n=0}^{\infty} \frac{1}{n!}\left(\sum_{m \geq n} m(m-1) \ldots(m-n+1) a_{m} x^{m-n}\right) y^{n} \\
& =\sum_{m \geq n \geq 0}\binom{m}{n} a_{m} x^{m-n} y^{n} \tag{39.25}
\end{align*}
$$

Since

$$
\begin{aligned}
\sum_{m \geq n \geq 0}\binom{m}{n}\left|a_{m} x^{m-n} y^{n}\right| & \leq C \sum_{m \geq n \geq 0}\binom{m}{n}\left|r^{-m} x^{m-n} y^{n}\right| \\
& =C \sum_{m \geq 0} r^{-m}(|x|+|y|)^{m}<\infty
\end{aligned}
$$

provided $|x|+|y|<r$. Since $r<\rho$ was arbitrary, it follows that Eq. (39.25) is convergent for $|x|+|y|<\rho$.

Let us redo this example. By the PDE in Eq. (39.23), $\partial_{y}^{m} \partial_{x}^{n} u(x, y)=$ $\partial_{x}^{n+m} u(x, y)$ and hence

$$
\partial_{y}^{m} \partial_{x}^{n} u(0,0)=f^{(m+n)}(0)
$$

Written another way

$$
D^{\alpha} u(0,0)=f^{(|\alpha|)}(0)
$$

and so the power series expansion for $u$ must be given by

$$
\begin{equation*}
u(x, y)=\sum_{\alpha} \frac{f^{(|\alpha|)}(0)}{\alpha!}(x, y)^{\alpha} \tag{39.26}
\end{equation*}
$$

Using $f^{(m)}(0) / m!\leq C r^{-m}$ we learn

$$
\begin{aligned}
\sum_{\alpha}\left|\frac{f^{(|\alpha|)}(0)}{\alpha!}(x, y)^{\alpha}\right| & \left.\leq C \sum_{\alpha} \frac{\left|f^{(|\alpha|)}(0)\right|}{\alpha!}|x|^{\alpha_{1}}|y|^{\alpha_{2}}=C \sum_{m=0}^{\infty} \frac{\left|f^{(|\alpha|)}(0)\right|}{m!} \sum_{|\alpha|=m} \frac{m!}{\alpha!} \right\rvert\, x \\
& \leq C \sum_{m=0}^{\infty} r^{-m}(|x|+|y|)^{m}=C \frac{r}{r-(|x|+|y|)}<\infty
\end{aligned}
$$

if $|x|+|y|<r$. Since $r<\rho$ was arbitrary, it follows that the series in Eq. (39.26) converges for $|x|+|y|<\rho$.

Now it is easy to check directly that Eq. (39.26) solves the PDE. However this is necessary since by construction $D^{\alpha} u_{y}(0,0)=D^{\alpha} u_{x}(0,0)$ for all $\alpha$. This implies, because $u_{y}$ and $u_{x}$ are both real analytic, that $u_{x}=u_{y}$.

## A very short introduction to generalized functions

Let $U$ be an open subset of $\mathbb{R}^{n}$ and

$$
\begin{equation*}
C_{c}^{\infty}(U)=\cup_{K \sqsubset ᄃ U} C^{\infty}(K) \tag{40.1}
\end{equation*}
$$

denote the set of smooth functions on $U$ with compact support in $U$.
Definition 40.1. A sequence $\left\{\phi_{k}\right\}_{k=1}^{\infty} \subset \mathcal{D}(U)$ converges to $\phi \in \mathcal{D}(U)$, iff there is a compact set $K \sqsubset \sqsubset U$ such that $\operatorname{supp}\left(\phi_{k}\right) \subset K$ for all $k$ and $\phi_{k} \rightarrow \phi$ in $C^{\infty}(K)$.

Definition 40.2 (Distributions on $U \subset_{o} \mathbb{R}^{n}$ ). A generalized function $T$ on $U \subset \subset_{o} \mathbb{R}^{n}$ is a continuous linear functional on $\mathcal{D}(U)$, i.e. $T: \mathcal{D}(U) \rightarrow \mathbb{C}$ is linear and $\lim _{n \rightarrow \infty}\left\langle T, \phi_{k}\right\rangle=0$ for all $\left\{\phi_{k}\right\} \subset \mathcal{D}(U)$ such that $\phi_{k} \rightarrow 0$ in $\mathcal{D}(U)$. Here we have written $\langle T, \phi\rangle$ for $T(\phi)$. We denote the space of generalized functions by $\mathcal{D}^{\prime}(U)$.

Example 40.3. Here are a couple of examples of distributions.

1. For $f \in L_{l o c}^{1}(U)$ define $T_{f} \in \mathcal{D}^{\prime}(U)$ by $\left\langle T_{f}, \phi\right\rangle=\int_{U} \phi f d m$ for all $\phi \in \mathcal{D}(U)$. This is called the distribution associated to $f$.
2. More generally let $\mu$ be a complex measure on $U$, then $\langle\mu, \phi\rangle:=\int_{U} \phi d \mu$ is a distribution. For example if $x \in U$, and $\mu=\delta_{x}$ then $\left\langle\delta_{x}, \phi\right\rangle=\phi(x)$ for all $\phi \in \mathcal{D}$.

Lemma 40.4. Let $a_{\alpha} \in C^{\infty}(U)$ and $L=\sum_{|\alpha| \leq m} a_{\alpha} \partial^{\alpha}-a m^{\text {th }}$ order linear differential operator on $\mathcal{D}(U)$. Then for $f \in C^{m}(U)$ and $\phi \in \mathcal{D}(U)$,

$$
\langle L f, \phi\rangle:=\left\langle T_{L f}, \phi\right\rangle=\left\langle T, L^{\dagger} \phi\right\rangle
$$

where $L^{\dagger}$ is the formal adjoint of $L$ defined by

$$
L^{\dagger} \phi=\sum_{|\alpha| \leq m}(-1)^{|\alpha|} \partial^{\alpha}\left[a_{\alpha} \phi\right] .
$$

Proof. This is simply repeated integration by parts. No boundary terms arise since $\phi$ has compact support.

Definition 40.5 (Multiplication by smooth functions). Suppose that $g \in C^{\infty}(U)$ and $T \in \mathcal{D}^{\prime}(U)$ then we define $g T \in \mathcal{D}^{\prime}(U)$ by

$$
\langle g T, \phi\rangle=\langle T, g \phi\rangle \text { for all } \phi \in \mathcal{D}(U)
$$

It is easily checked that $g T$ is continuous.
Definition 40.6 (Differentiation). For $T \in \mathcal{D}^{\prime}(U)$ and $i \in\{1,2, \ldots, n\}$ let $\partial_{i} T \in \mathcal{D}^{\prime}(U)$ be the distribution defined by

$$
\left\langle\partial_{i} T, \phi\right\rangle=-\left\langle T, \partial_{i} \phi\right\rangle \text { for all } \phi \in \mathcal{D}(U)
$$

Again it is easy to check that $\partial_{i} T$ is a distribution.
Definition 40.7. More generally if $L$ is as in Lemma 40.4 and $T \in \mathcal{D}^{\prime}$ we define $L T \in \mathcal{D}^{\prime}$ by

$$
\langle L T, \phi\rangle=\left\langle T, L^{\dagger} \phi\right\rangle .
$$

Example 40.8. Suppose that $f \in L_{l o c}^{1}$ and $g \in C^{\infty}(U)$, then $g T_{f}=T_{g f}$. If further $f \in C^{1}(U)$, then $\partial_{i} T_{f}=T_{\partial_{i} f}$. More generally if $f \in C^{m}(U)$ then, by Lemma 40.4, $L T_{f}=T_{L f}$.

Because of Definition 40.7 we may now talk about distributional or generalized solutions $T$ to PDEs of the form $L T=S$ where $S \in \mathcal{D}^{\prime}$.

Example 40.9. For the moment let us also assume that $U=\mathbb{R} .\left\langle T_{f}, \phi\right\rangle=$ $\int_{U} \phi f d m$. Then we have

1. $\lim _{M \rightarrow \infty} T_{\sin M x}=0$
2. $\lim _{M \rightarrow \infty} T_{M^{-1} \sin M x}=\pi \delta_{0}$ where $\delta_{0}$ is the point measure at 0 .
3. If $f \in L^{1}\left(\mathbb{R}^{n}, d m\right)$ with $\int_{\mathbb{R}^{n}} f d m=1$ and $f_{\epsilon}(x)=\epsilon^{-n} f(x / \epsilon)$, then $\lim _{\epsilon \downarrow 0} T_{f_{\epsilon}}=\delta_{0}$. Indeed,

$$
\begin{aligned}
\lim _{\epsilon \downarrow 0}\left\langle T_{f_{\epsilon}}, \phi\right\rangle & =\lim _{\epsilon \downarrow 0} \int_{\mathbb{R}^{n}} \epsilon^{-n} f(x / \epsilon) \phi(x) d x \\
& =\lim _{\epsilon \downarrow 0} \int_{\mathbb{R}^{n}} f(x) \phi(\epsilon x) d x \stackrel{\text { D.C.T. }}{=} \int_{\mathbb{R}^{n}} f(x) \lim _{\epsilon \downarrow 0} \phi(\epsilon x) d x \\
& =\phi(0) \int_{\mathbb{R}^{n}} f(x) d x=\phi(0)=\left\langle\delta_{0}, \phi\right\rangle .
\end{aligned}
$$

As a concrete example we have

$$
\lim _{\epsilon \downarrow 0} \frac{\epsilon}{\pi\left(x^{2}+\epsilon^{2}\right)}=\delta_{0} \text { on } \mathbb{R},
$$

i.e.

$$
\lim _{\epsilon \downarrow 0} T_{\frac{\epsilon}{\pi\left(x^{2}+\epsilon^{2}\right)}}=\delta_{0} .
$$

$$
\left\langle\partial_{i} \delta_{a}, \phi\right\rangle=-\partial_{i} \phi(a)
$$

and more generally we have

$$
\left\langle L \delta_{a}, \phi\right\rangle=\left(L^{\dagger} \phi\right)(a) .
$$

Lemma 40.11. Suppose $f \in C^{1}([a, b])$ and $g \in P C^{1}([a, b])$, i.e. $g \in$ $C^{1}([a, b] \backslash \Lambda)$ where $\Lambda$ is a finite subset of $(a, b)$ and $g(\alpha+), g(\alpha-)$ exists for $\alpha \in \Lambda$. Then

$$
\begin{align*}
\int_{a}^{b} f^{\prime}(x) g(x) d x & =\left.\left[f^{\prime}(x) g(x)\right]\right|_{a} ^{b}-\int_{a}^{b} f(x) g^{\prime}(x) d x \\
& -\sum_{\alpha \in \Lambda} f(\alpha)(g(\alpha+)-g(\alpha-)) \tag{40.2}
\end{align*}
$$

In particular

$$
\frac{d}{d x} T_{g}=T_{g^{\prime}}+\sum_{\alpha \in A}(g(\alpha+)-g(\alpha-)) \delta_{\alpha}
$$

Proof. Write $\Lambda \cup\{a, b\}$ as $\left\{a=\alpha_{0}<\alpha_{1}<\cdots<\alpha_{n}=b\right\}$, then

$$
\begin{aligned}
\int_{a}^{b} f^{\prime}(x) g(x) d x & =\sum_{k=0}^{n-1} \int_{\alpha_{k}}^{\alpha_{k+1}} f^{\prime}(x) g(x) d x \\
& =\sum_{k=0}^{n-1}\left[\left.[f(x) g(x)]\right|_{\alpha_{k}+} ^{\alpha_{k+1}-}-\int_{\alpha_{k}}^{\alpha_{k+1}} f(x) g^{\prime}(x) d x\right] \\
& =\left.\left[f^{\prime}(x) g(x)\right]\right|_{a} ^{b}-\int_{a}^{b} f(x) g^{\prime}(x) d x-\left.\sum_{k=1}^{n-1}[f(x) g(x)]\right|_{\alpha_{k}-} ^{\alpha_{k}+}
\end{aligned}
$$

which is the same as Eq. (40.2).

## Elliptic Ordinary Differential Operators

Let $\Omega \subset_{o} \mathbb{R}^{n}$ be a bounded connected open region. A function $u \in C^{2}(\Omega)$ is said to satisfy Laplace's equation if

$$
\Delta u=0 \text { in } \Omega .
$$

More generally if $f \in C(\Omega)$ is given we say $u$ solves the Poisson equation if

$$
-\triangle u=f \text { in } \Omega
$$

In order to get a unique solution to either of these equations it is necessary to impose "boundary" conditions on $u$.
Example 41.1. For Dirichlet boundary conditions we impose $u=g$ on $\partial \Omega$ and for Neumann boundary conditions we impose $\frac{\partial u}{\partial \nu}=g$ on $\partial \Omega$, where $g: \partial \Omega \rightarrow \mathbb{R}$ is a given function.
Lemma 41.2. Suppose $f: \Omega \xrightarrow{C^{0}} \mathbb{R}, \partial \Omega$ is $C^{2}$ and $g: \partial \Omega \rightarrow \mathbb{R}$ is continuous. Then if there exists a solution to $-\triangle u=f$ with $u=g$ on $\partial \Omega$ such that $u \in C^{2}\left(\Omega^{o}\right) \cap C^{1}(\Omega)$ then the solution is unique.
Definition 41.3. Given an open set $\Omega \subset \mathbb{R}^{n}$ we say $u \in C^{1}(\bar{\Omega})$ if $u \in$ $C^{1}(\Omega) \cap C(\bar{\Omega})$ and $\nabla u$ extends to a continuous function on $\bar{\Omega}$.

Proof. If $\widetilde{u}$ is another solution then $v=\widetilde{u}-u$ solves $\triangle v=0, v=0$ on $\partial \Omega$. By the divergence theorem,

$$
0=\int_{\Omega} \triangle v \cdot v d m=-\int_{\Omega}|\nabla v|^{2} d m+\int_{\partial \Omega} v \nabla v \cdot n d \sigma=-\int_{\Omega}|\nabla v|^{2} d m
$$

where the boundary terms are zero since $v=0$ on $\partial \Omega$. This identity implies $\int|\nabla u|^{2} d x=0$ which then shows $\nabla v \equiv 0$ and since $\Omega$ is connected we learn $v$ $\Omega$ is constant on $\Omega$. Because $v$ is zero on $\partial \Omega$ we conclude $v \equiv 0$, that is $u=\tilde{u}$.

For the rest of this section we will now restrict to $n=1$. However we will allow for more general operators than $\Delta$ in this case.

### 41.1 Symmetric Elliptic ODE

Let $a \in C^{1}([0,1],(0, \infty))$ and

$$
\begin{equation*}
L f=-\left(a f^{\prime}\right)^{\prime}=-a f^{\prime \prime}-a^{\prime} f^{\prime} \text { for } f \in C^{2}([0,1]) \tag{41.1}
\end{equation*}
$$

In the following theorem we will impose Dirichlet boundary conditions on $L$ by restricting the domain of $L$ to

$$
D(L):=\left\{f \in C^{2}([0,1], \mathbb{R}): f(0)=f(1)=0\right\}
$$

Theorem 41.4. The linear operator $L: D(L) \rightarrow C([0,1], \mathbb{R})$ is invertible and $L^{-1}: C([0,1], \mathbb{R}) \rightarrow D(L) \subset C^{2}([0,1], \mathbb{R})$ is a bounded operator.

## Proof.

1. (Uniqueness) If $f, g \in D(L)$ then by integration by parts

$$
\begin{equation*}
(L f, g):=\int_{0}^{1}(L f)(x) g(x) d x=\int_{0}^{1} a(x) f^{\prime}(x) g^{\prime}(x) d x \tag{41.2}
\end{equation*}
$$

Therefore if $L f=0$ then

$$
0=(L f, f)=\int_{0}^{1} a(x) f^{\prime}(x)^{2} d x
$$

and hence $f^{\prime} \equiv 0$ and since $f(0)=0, f \equiv 0$. This shows $L$ is injective.
2. (Existence) Given $g \in C([0,1], \mathbb{R})$ we are looking for $f \in D(L)$ such that $L f=g$, i.e. $\left(a f^{\prime}\right)^{\prime}=g$. Integrating this equation implies

$$
-a(x) f^{\prime}(x)=-C+\int_{0}^{x} g(y) d y
$$

Therefore

$$
f^{\prime}(z)=\frac{C}{a(z)}-\int 1_{y \leq z} \frac{1}{a(z)} g(y) d y
$$

which upon integration and using $f(0)=0$ gives

$$
f(x)=\int_{0}^{x} \frac{C}{a(z)} d z-\int 1_{y \leq z \leq x} \frac{1}{a(z)} g(y) d z d y .
$$

If we let

$$
\begin{equation*}
\alpha(x):=\int_{0}^{x} \frac{1}{a(z)} d z \tag{41.3}
\end{equation*}
$$

the last equation may be written as

$$
\begin{equation*}
f(x)=C \alpha(x)-\int_{0}^{x}(\alpha(x)-\alpha(y)) g(y) d y \tag{41.4}
\end{equation*}
$$

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It is a simple matter to work backwards to show the function $f$ defined in Eq. (41.4) satisfies $L f=g$ and $f(0)=0$ for any constant $C$. So it only remains to choose $C$ so that

$$
0=f(1)=C \alpha(1)-\int_{0}^{1}(\alpha(1)-\alpha(y)) g(y) d y
$$

Solving for $C$ gives $C=\int_{0}^{1}\left(1-\frac{\alpha(y)}{\alpha(1)}\right) g(y) d y$ and the resulting function $f$ may be written as

$$
\begin{aligned}
f(x) & =\int_{0}^{1}\left[\left(1-\frac{\alpha(y)}{\alpha(1)}\right) \alpha(x)-1_{y \leq x}(\alpha(x)-\alpha(y))\right] g(y) d y \\
& =\int_{0}^{1} G(x, y) g(y) d y
\end{aligned}
$$

where

$$
G(x, y)= \begin{cases}\alpha(x)\left(1-\frac{\alpha(y)}{\alpha(1)}\right) & \text { if } x \leq y  \tag{41.5}\\ \alpha(y)\left(1-\frac{\alpha(x)}{\alpha(1)}\right) & \text { if } y \leq x\end{cases}
$$

For example when $a \equiv 1$,

$$
G(x, y)=\left\{\begin{array}{l}
x(1-y) \text { if } x \leq y \\
y(1-x) \text { if } y \leq x
\end{array}\right.
$$

Definition 41.5. The function $G$ defined in Eq. (41.5) is called the Green's function for the operator $L: D(L) \rightarrow C([0,1], \mathbb{R})$.
Remarks 41.6 The proof of Theorem 41.4 shows

$$
\begin{equation*}
\left(L^{-1} g\right)(x):=\int_{0}^{1} G(x, y) g(y) d y \tag{41.6}
\end{equation*}
$$

where $G$ is defined in Eq. (41.5). The Green's function $G$ has the following properties:

1. Since $L$ is invertible and $G$ is a right inverse, $G$ is also a left inverse, i.e. $G L f=f$ for all $f \in D(L)$.
2. $G$ is continuous.
3. $G$ is symmetric, $G(y, x)=G(x, y)$. (This reflects the symmetry in $L$, $(L f, g)=(f, L g)$ for all $f, g \in D(L)$, which follows from $E q$. (41.2).)
4. $G$ may be written as

$$
G(x, y)=\left\{\begin{array}{l}
u(x) v(y) \text { if } x \leq y \\
u(y) v(x) \text { if } y \leq x
\end{array}\right.
$$

where $u$ and $v$ are $L$-harmonic functions (i.e. and $L u=L v=0$ ) with $u(0)=0$ and $v(1)=0$. In particular $L_{x} G(x, y)=0=L_{y} G(x, y)$ for all $y \neq x$.
5. The first order derivatives of the Green's function have a jump discontinuity on the diagonal. Explicitly,

$$
G_{y}(x, x+)-G_{y}(x, x-)=-\frac{1}{a(x)}
$$

which follows directly from

$$
G_{y}(x, y)=\frac{1}{a(y)}\left\{\begin{array}{cl}
-\frac{\alpha(x)}{\alpha(1)} & \text { if } x<y  \tag{41.7}\\
\left(1-\frac{\alpha(x)}{\alpha(1)}\right) & \text { if } y<x
\end{array}\right.
$$

By symmetry we also have

$$
G_{x}(y+, y)-G_{x}(y-, y)=-\frac{1}{a(y)}
$$

6. By Items 4. and 5. and Lemma 40.11 it follows that

$$
L_{y} G(x, y):=L_{y} T_{G(x, y)}=\frac{d}{d y}\left(a(y) G_{y}(x, y)\right)=\delta(y-x)
$$

and similarly that

$$
L_{x} T_{G(x, y)}=L_{x} G(x, y)=\delta(x-y)
$$

As a consequence of the above remarks we have the following representation theorem for function $f \in C^{2}([0,1])$.
Theorem 41.7 (Representation Theorem). For any $f \in C^{2}([0,1])$,

$$
\begin{equation*}
f(x)=(G L f)(x)-\left.G_{y}(x, y) a(y) f(y)\right|_{y=0} ^{y=1} \tag{41.8}
\end{equation*}
$$

Moreover if we are given $h: \partial[0,1] \rightarrow \mathbb{R}$ and $g \in C([0,1])$, then the unique solution to

$$
L f=g \text { with } f=h \text { on } \partial[0,1]
$$

is

$$
\begin{equation*}
f(x)=(G g)(x)-\left.G_{y}(x, y) a(y) h(y)\right|_{y=0} ^{y=1} \tag{41.9}
\end{equation*}
$$

Proof. By repeated use of Lemma 40.11,

$$
\begin{aligned}
(G L f)(x) & =-\int_{0}^{1} G(x, y) \frac{d}{d y}\left(a(y) f^{\prime}(y)\right) d y \\
& =\int_{0}^{1} G_{y}(x, y) a(y) f^{\prime}(y) d y \\
& =\left.G_{y}(x, y) a(y) f(y)\right|_{y=0} ^{y=1}+\int_{0}^{1} L_{y} G(x, y) f(y) d y \\
& =\left.G_{y}(x, y) a(y) f(y)\right|_{y=0} ^{y=1}+\int_{0}^{1} \delta(x-y) f(y) d y \\
& =\left.G_{y}(x, y) a(y) f(y)\right|_{y=0} ^{y=1}+f(x)
\end{aligned}
$$

which proves Eq. (41.8). There are no boundary terms in the second equality above since $G(x, 0)=G(x, 1)=0$.

Now suppose that $f$ is defined as in Eq. (41.9). Observe from Eq. (41.7) that

$$
\lim _{x \uparrow 1} a(1) G_{y}(x, 1)=-1 \text { and } \lim _{x \downarrow 0} a(0) G_{y}(x, 0)=1
$$

and also notice that $G_{y}(x, 1)$ and $G_{y}(x, 0)$ are $L_{x}$ - harmonic functions. Therefore by these remarks and Eq. (41.6), $f=h$ on $\partial[0,1]$ and

$$
L f(x)=g(x)-\left.L_{x} G_{y}(x, y) a(y) h(y)\right|_{y=0} ^{y=1}=g(x)
$$

as desired.

### 41.2 General Regular 2nd order elliptic ODE

Let $J=[r, s]$ be a closed bounded interval in $\mathbb{R}$.
Definition 41.8. A second order linear operator of the form

$$
\begin{equation*}
L f=-a f^{\prime \prime}+b f^{\prime}+c f \tag{41.10}
\end{equation*}
$$

with $a \in C^{2}(J), b \in C^{1}(J)$ and $c \in C^{2}(J)$ is said to be elliptic if $a>0$, (more generally if a is invertible if we are allowing for vector valued functions).

For this section $L$ will denote an elliptic ordinary differential operator. We will now consider the Dirichlet boundary valued problem for $f \in C^{2}([r, s])$,

$$
\begin{equation*}
L f=-a f^{\prime \prime}+b f^{\prime}+c f=0 \text { with } f=0 \text { on } \partial J . \tag{41.11}
\end{equation*}
$$

Lemma 41.9. Let $u, v \in C^{2}(J)$ be two $L$ - harmonic functions, i.e. $L u=$ $0=L v$ and let

$$
W:=\operatorname{det}\left[\begin{array}{cc}
u & v \\
u^{\prime} & v^{\prime}
\end{array}\right]=u v^{\prime}-v u^{\prime}
$$

be the Wronskian of $u$ and $v$. Then $W$ satisfies

$$
\begin{aligned}
W^{\prime} & =\frac{b}{a} W, \frac{d}{d x} \frac{1}{W}=-\frac{b}{a} \frac{1}{W} \text { and } \\
W(x) & =W(r) e^{\int_{r}^{x} \frac{b}{a}(t) d t}
\end{aligned}
$$

Proof. By direct computation

$$
a W^{\prime}=a\left(u v^{\prime \prime}-v u^{\prime \prime}\right)=u\left(b v^{\prime}+c v\right)-v\left(b u^{\prime}+c u\right)=b W .
$$

Definition 41.10. Let $H^{k}(J)$ denote those $f \in C^{k-1}(J)$ such that $f^{(k-1)}$ is absolutely continuous and $f^{(k)} \in L^{2}(J)$. We also let $H_{0}^{2}(J)=\left\{f \in H^{2}(J):\left.f\right|_{\partial J}=\right.$ We make $H^{k}(J)$ into a Hilbert space using the following inner product

$$
(u, v)_{H^{k}}:=\sum_{j=0}^{k}\left(D^{j} u, D^{j} v\right)_{L^{2}}
$$

Theorem 41.11. As above, let $D(L)=\left\{f \in C^{2}(J): f=0\right.$ on $\left.\partial J\right\}$. If the $\operatorname{Nul}(L) \cap D(L)=\{0\}$, i.e. if the only solution $f \in D(L)$ to $L f=0$ is $f=0$, then $L: D(L) \rightarrow C(J)$ is an invertible. Moreover there exists a continuous function $G$ on $J \times J$ (called the Dirichlet Green's function for $L$ ) such that

$$
\begin{equation*}
\left(L^{-1} g\right)(x)=\int_{J} G(x, y) g(y) d y \text { for all } g \in C(J) \tag{41.12}
\end{equation*}
$$

Moreover if $g \in L^{2}(J)$ then $G g \in H_{0}^{2}(J)$ and $L(G g)=g$ a.e. and more generally if $g \in H^{k}(J)$ then $G g \in H_{0}^{k+2}(J)$

Proof. To prove the surjectivity of $L: D(L) \rightarrow C(J)$, (i.e. existence of solutions $f \in D(L)$ to $L f=g$ with $g \in C(J))$ we are going to construct the Green's function $G$.

1. Formal requirements on the Greens function. Assuming Eq. (41.12) holds and working formally we should have

$$
\begin{equation*}
g(x)=L_{x} \int_{J} G(x, y) g(y) d y=\int_{J} L_{x} G(x, y) g(y) d y \tag{41.13}
\end{equation*}
$$

for all $g \in C(J)$. Hence, again formally, this implies

$$
\begin{equation*}
L_{x} G(x, y)=\delta(y-x) \text { with } G(r, y)=G(s, y)=0 \tag{41.14}
\end{equation*}
$$

This can be made more convincing by as follows. Let $\phi \in \mathcal{D}:=\mathcal{D}(r, s)$, then multiplying

$$
g(x)=L_{x} \int_{J} G(x, y) g(y) d y
$$

by $\phi$, integrating the result and then using integration by parts and Fubini's theorem gives

$$
\begin{aligned}
\int_{J} g(x) \phi(x) d x & =\int_{J} d x \phi(x) L_{x} \int_{J} d y G(x, y) g(y) \\
& =\int_{J} d x L_{x} \phi(x) \int_{J} d y G(x, y) g(y) \\
& =\int_{J} d y g(y) \int_{J} d x L_{x} \phi(x) G(x, y) \text { for all } g \in C(J)
\end{aligned}
$$

From this we conclude

$$
\int_{J} L_{x} \phi(x) G(x, y) d x=\phi(y)
$$

i.e. $L_{x} T_{G(x, y)}=\delta(x-y)$.
2. Constructing $G$. In order to construct a solution to Eq. (41.14), let $u, v$ be two non-zero $L$ - harmonic functions chosen so that $u(r)=0=v(s)$ and $u^{\prime}(r)=1=v^{\prime}(s)$ and let $W$ be the Wronskian of $u$ and $v$. By Lemma 41.9, either $W$ is never zero or is identically zero. If $W=0$, then $\left(u(r), u^{\prime}(r)\right)=\lambda\left(v(r), v^{\prime}(r)\right)$ for some $\lambda \in \mathbb{R}$ and by uniqueness of solutions to ODE it would follow that $u \equiv \lambda v$. In this case $u(r)=0$ and $u(s)=\lambda v(s)=0$, and hence $u \in D(L)$ with $L u=0$. However by assumption, this implies $u=0$ which is impossible since $u^{\prime}(0)=1$. Thus $W$ is never 0 .
By Eq. (41.14) we should require $L_{x} G(x, y)=0$ for $x \neq y$ and $G(r, y)=$ $G(s, y)=0$ which implies that

$$
G(x, y)=\left\{\begin{array}{l}
u(x) \phi(y) \text { if } x<y \\
v(x) \psi(y) \text { if } x>y
\end{array}\right.
$$

for some functions $\phi$ and $\psi$. We now want to choose $\phi$ and $\psi$ so that $G$ is continuous and $L_{x} G(x, y)=\delta(x-y)$. Using

$$
G_{x}(x, y)=\left\{\begin{array}{l}
u^{\prime}(x) \phi(y) \text { if } x<y \\
v^{\prime}(x) \psi(y) \text { if } x>y
\end{array}\right.
$$

Lemma 41.9, we are led to require

$$
\begin{aligned}
& 0=G(y+, y)-G(y-, y)=u(y) \phi(y)-v(y) \psi(y) \\
& 1=-\left.\left[a(x) G_{x}(x, y)\right]\right|_{x=y-} ^{x=y+}=-a(y)\left[v^{\prime}(y) \psi(y)-u^{\prime}(y) \phi(y)\right]
\end{aligned}
$$

Solving these equations for $\phi$ and $\psi$ gives

$$
\binom{\phi}{\psi}=-\frac{1}{a W}\binom{v}{u}
$$

and hence

$$
G(x, y)=-\frac{1}{a(y) W(y)}\left\{\begin{array}{l}
u(x) v(y) \text { if } x \leq y  \tag{41.15}\\
v(x) u(y) \text { if } x \geq y
\end{array}\right.
$$

3. With this $G$, Eq. (41.12) holds. Given $g \in C(J)$, then $f$ in Eq. (41.12) may be written as

$$
\begin{align*}
f(x) & =\int_{J} G(x, y) g(y) d y \\
& =-v(x) \int_{r}^{x} \frac{u(y)}{a(y) W(y)} g(y) d y-u(x) \int_{x}^{s} \frac{v(y)}{a(y) W(y)} g(y) d y \tag{41.16}
\end{align*}
$$

Differentiating this equation twice gives

$$
\begin{equation*}
f^{\prime}(x)=-v^{\prime}(x) \int_{r}^{x} \frac{u(y)}{a(y) W(y)} g(y) d y-u^{\prime}(x) \int_{x}^{s} \frac{v(y)}{a(y) W(y)} g(y) d y \tag{41.17}
\end{equation*}
$$

and

$$
\begin{align*}
f^{\prime \prime}(x) & =-v^{\prime \prime}(x) \int_{r}^{x} \frac{u(y)}{a(y) W(y)} g(y) d y-u^{\prime \prime}(x) \int_{x}^{s} \frac{v(y)}{a(y) W(y)} g(y) d y \\
& -v^{\prime}(x) \frac{u(x)}{a(x) W(x)} g(x)+u^{\prime}(x) \frac{v(x)}{a(x) W(x)} g(x) \tag{41.18}
\end{align*}
$$

Using $L v=0=L u$, the definition of $W$ and the last two equations we find

$$
\begin{aligned}
-a(x) f^{\prime \prime}(x) & =\left[b(x) v^{\prime}(x)+c(x) v(x)\right] \int_{r}^{x} \frac{u(y)}{a(y) W(y)} g(y) d y \\
& +\left[b(x) u^{\prime}(x)+c(x) u(x)\right] \int_{x}^{s} \frac{v(y)}{a(y) W(y)} g(y) d y+g(x) \\
& =-b(x) f^{\prime}(x)-c(x) f(x)+g(x)
\end{aligned}
$$

i.e. $L f=g$.

Hence we have proved $L: D(L) \rightarrow C(J)$ is surjective and $L^{-1}: C(J) \rightarrow$ $D(L)$ is given by Eq. (41.12).

Now suppose $g \in L^{2}(J)$, we will show that $f \in C^{1}(J)$ and Eq. (41.17) is still valid. The difficulty here is that it is clear that $f$ is differentiable almost everywhere and Eq. (41.17) holds for almost every $x$. However this is not good enough, we need Eq. (41.17) to hold for all $x$. To remedy this, choose $g_{n} \in C(J)$ such that $g_{n} \rightarrow g$ in $L^{2}(J)$ and let $f_{n}:=G g_{n}$. Then by what we have just proved,

$$
f_{n}^{\prime}(x)=\int_{J} G_{x}(x, y) g_{n}(y) d y
$$

Now by the Cauchy-Schwarz inequality,

$$
\begin{aligned}
\left|\int_{J} G_{x}(x, y)\left[g(y)-g_{n}(y)\right] d y\right|^{2} & \leq\left\|g-g_{n}\right\|_{L^{2}(J)}^{2} \int_{J}\left|G_{x}(x, y)\right|^{2} d y \\
& \leq C\left\|g-g_{n}\right\|_{L^{2}(J)}^{2}
\end{aligned}
$$

where $C:=\sup _{x \in J} \int_{J}\left|G_{x}(x, y)\right|^{2} d y<\infty$. From this inequality it follows that $f_{n}^{\prime}(x)$ converges uniformly to $\int_{J} G_{x}(x, y) g(y) d y$ as $n \rightarrow \infty$ and hence $f \in C^{1}(J)$ and

$$
f^{\prime}(x)=\int_{J} G_{x}(x, y) g(y) d y \text { for all } x \in J
$$

i.e. Eq. (41.17) is valid for all $x \in J$. It now follows from Eq. (41.17) that $f \in H^{2}(J)$ and Eq. (41.18) holds for almost every $x$. Working as before we may conclude $L f=g$ a.e. Finally if $g \in H^{k}(J)$ for $k \geq 1$, the reader may easily show $f \in H_{0}^{k+2}(J)$ by examining Eqs. (41.17) and (41.18).

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Remark 41.12. When $L$ is given as in Eq. (41.1), $b=-a^{\prime}$ and by Lemma 41.9

$$
W(x)=W(0) e^{-\int_{0}^{x} \frac{a^{\prime}}{a}(t) d t}=W(0) e^{-\ln (a(x) / a(0))}=\frac{W(0) a(0)}{a(x)} .
$$

So in this case

$$
G(x, y)=-\frac{1}{W(0) a(0)}\left\{\begin{array}{l}
u(x) v(y) \text { if } x \leq y \\
v(x) u(y) \text { if } x \geq y
\end{array}\right.
$$

where we may take

$$
u(x)=\alpha(x):=\int_{0}^{x} \frac{1}{a(z)} d z \text { and } v(x)=\left(1-\frac{\alpha(x)}{\alpha(1)}\right) .
$$

Finally for this choice of $u$ and $v$ we have

$$
W(0)=u(0) v^{\prime}(0)-u^{\prime}(0) v(0)=-\frac{1}{a(0)}
$$

giving

$$
G(x, y)=\left\{\begin{array}{l}
u(x) v(y) \text { if } x \leq y \\
v(x) u(y) \text { if } x \geq y
\end{array}\right.
$$

which agrees with Eq. (41.5) above.
Lemma 41.13. Let $L^{*} f:=-(a f)^{\prime \prime}-(b f)^{\prime}+c f$ be the formal adjoint of $L$. Then

$$
\begin{equation*}
(L f, g)=\left(f, L^{*} g\right) \text { for all } f, g \in D(L) \tag{41.19}
\end{equation*}
$$

where $(f, g):=\int_{J} f(x) g(x) d x$. Moreover if $\operatorname{nul}(L)=\{0\}$ then $\operatorname{nul}\left(L^{*}\right)=\{0\}$ and the Greens function for $L^{*}$ is $G^{*}$ defined by $G^{*}(x, y)=G(y, x)$, where $G$ is the Green's function in Eq. (41.15). Consequently $L_{y}^{*} G(x, y)=\delta(x-y)$.

Proof. First observe that $G^{*}$ has been defined so that $\left(G^{*} g, f\right)=(g, G f)$ for all $f \in L^{2}(J)$.Eq. (41.19) follows by two integration by parts after observing the boundary terms are zero because $f=g=0$ on $\partial J$. If $g \in \operatorname{nul}\left(L^{*}\right)$ and $f \in D(L)$, we find

$$
0=\left(L^{*} g, f\right)=(g, L f) \text { for all } f \in D(L)
$$

By Theorem 41.11, if $\operatorname{nul}(L)=\{0\}$ then $L: D(L) \rightarrow C(J)$ is invertible so the above equation implies $\operatorname{nul}\left(L^{*}\right)=\{0\}$. Another application of Theorem 41.11 then shows $L^{*}: D(L) \rightarrow C(J)$ is invertible and has a Green's function which we call $\tilde{G}(x, y)$. We will now complete the proof by showing $\tilde{G}=G^{*}$. To do this observe that

$$
(f, g)=\left(L^{*} \tilde{G} f, g\right)=(\tilde{G} f, L g)=\left(f, \tilde{G}^{*} L g\right) \text { for all } f, g \in D(L)
$$

and this then implies $\tilde{G}^{*} L=I d_{D(L)}=G L$. Cancelling the $L$ from this equation, show $\tilde{G}^{*}=G$ or equivalently that $\tilde{G}=G^{*}$. The remaining assertions of the Lemma follows from this observation.

Here is an alternate proof that $L_{y}^{*} G(x, y)=\delta(x-y)$, also see Using $G L=I_{D(L)}$, we learn for $u \in D(L)$ and $v \in C(J)$ that

$$
(v, u)=(v, G L u)=\left(L^{*} G^{*} v, u\right)
$$

which then implies $L^{*} G^{*} v=v$ for all $v \in C(J)$. This implies

$$
f(x)=\int_{J} G(x, y) L f(y) d y=\left\langle T_{G(x, \cdot)}, L f\right\rangle=\left\langle L^{*} T_{G(x, \cdot)}, f\right\rangle \text { for all } f \in D(L)
$$

from which it follows that $L_{y}^{*} T_{G(x, y)}=\delta(x-y)$.
Definition 41.14. A Green's function for $L$ is a function $G(x, y)$ as defined as in Eq. (41.15) where $u$ and $v$ are any two linearly independent $L$ harmonic functions. ${ }^{1}$

The following theorem in is a generalization of Theorem 41.7.
Theorem 41.15 (Representation Theorem). Suppose and $G$ is a Green's function for $L$ then

1. $L_{x} T_{G(x, y)}=\delta(x-y)$ and $L G=I$ on $L^{2}(J)$. (However $G g$ and $G^{*} g$ may no longer satisfy the given Dirichlet boundary conditions.)
2. $L_{y}^{*} T_{G(x, y)}=\delta(x-y)$. More precisely we have the following representation formula. For any $f \in H^{2}(J)$,

$$
\begin{equation*}
f(x)=(G L f)(x)+\left.\left\{G(x, y) a(y) f^{\prime}(y)-[a(y) G(x, y)]_{y} f(y)\right\}\right|_{y=r} ^{y=s} \tag{41.20}
\end{equation*}
$$

3. Let us now assume $\operatorname{nul}(L)=\{0\}$ and $G$ is the Dirichlet Green's function for L. The Eq. (41.20) specializes to

$$
f(x)=(G L f)(x)-\left.[a(y) G(x, y)]_{y} f(y)\right|_{y=r} ^{y=s}
$$

Moreover if we are given $h: \partial J \rightarrow \mathbb{R}$ and $g \in L^{2}(J)$, then the unique solution $f \in H^{2}(J)$ to

$$
L f=g \text { a.e. with } f=h \text { on } \partial J
$$

is

$$
\begin{equation*}
f(x)=(G g)(x)+H(x) \tag{41.21}
\end{equation*}
$$

where, for $x \in J^{0}$,

$$
\begin{equation*}
H(x):=-\left.[a(y) G(x, y)]_{y} h(y)\right|_{y=r} ^{y=s} \tag{41.22}
\end{equation*}
$$

and $H(r):=H(r+)$ and $H(s):=H(s-)$.

[^1] $v(\alpha)=1$.

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Proof. 1. The first item follows from the proof of Theorem 41.11 with out any modification.
2. Using Lemma 41.9,

$$
\begin{aligned}
L^{*}\left(\frac{u}{a W}\right) & =-\left(\frac{u}{W}\right)^{\prime \prime}-\left(\frac{b u}{a W}\right)^{\prime}+\frac{c u}{a W} \\
& =-\left(\frac{u^{\prime}}{W}-\frac{b}{a} \frac{1}{W} u\right)^{\prime}-\left(\frac{b u}{a W}\right)^{\prime}+\frac{c u}{a W} \\
& =-\left(\frac{u^{\prime}}{W}\right)^{\prime}+\frac{c u}{a W}=-\left(\frac{u^{\prime \prime}}{W}-\frac{b}{a} \frac{1}{W} u\right)+\frac{c u}{a W} \\
& =\frac{1}{a} L u=0 .
\end{aligned}
$$

Similarly $L^{*}\left(\frac{v}{a W}\right)=0$ and therefore $L_{y}^{*} G(x, y)=0$ for $y \neq x$. Since

$$
\begin{align*}
G_{y}(x, y) & =-\left(\frac{d}{d y} \frac{1}{a(y) W(y)}\right)\left\{\begin{array}{l}
u(x) v(y) \text { if } x \leq y \\
v(x) u(y) \text { if } x \geq y
\end{array}\right. \\
& -\frac{1}{a(y) W(y)}\left\{\begin{array}{l}
u(x) v^{\prime}(y) \text { if } x \leq y \\
v(x) u^{\prime}(y) \text { if } x \geq y
\end{array}\right. \tag{41.23}
\end{align*}
$$

we find

$$
\begin{aligned}
G_{y}(x, x+)-G_{y}(x, x-) & =\frac{1}{a(x) W(x)}\left\{v(x) u^{\prime}(x)-u(x) v^{\prime}(x)\right\} \\
& =-\frac{1}{a(x)} .
\end{aligned}
$$

Finally since

$$
L_{y}^{*}=-a \frac{d^{2}}{d y^{2}}+\text { lower order terms }
$$

we may conclude form Lemma 40.11 that $L_{y}^{*} G(x, y)=\delta(x-y)$.Using integration by parts for absolutely continuous functions and Lemma 41.13, for $f \in H^{2}(J)$,

$$
\begin{aligned}
(G L f)(x) & =\int_{J} G(x, y) L f(y) d y \\
& =\int_{J} G(x, y)\left(-a(y) \frac{d^{2}}{d y^{2}}+b(y) \frac{d}{d y}+c(y)\right) f(y) d y \\
& =\int_{J}\left[\begin{array}{c}
\frac{d}{d y}[a(y) G(x, y)] f^{\prime}(y) \\
\left.-\frac{d}{d y}[b(y) G(x, y)] f+c(y)\right) f(y)
\end{array}\right] d y \\
& -\left.G(x, y) a(y) f^{\prime}(y)\right|_{y=r} ^{y=s} \\
& =-\left.G(x, y) a(y) f^{\prime}(y)\right|_{y=r} ^{y=s}+\left.[a(y) G(x, y)]_{y} f(y)\right|_{y=r} ^{y=s} \\
& +\left\langle L_{y}^{*} G(x, y), f(y)\right\rangle \\
& =\left.[a(y) G(x, y)]_{y} f(y)\right|_{y=r} ^{y=s}-\left.G(x, y) a(y) f^{\prime}(y)\right|_{y=r} ^{y=s}+f(x) .
\end{aligned}
$$

This proves Eq. (41.20).
3. Now suppose $G$ is the Dirichlet Green's function for $L$. By Eq. (41.15),

$$
\begin{aligned}
{[-a(y) G(x, y)]_{y} } & =\left(\frac{d}{d y} \frac{1}{W(y)}\right)\left\{\begin{array}{l}
u(x) v(y) \text { if } x \leq y \\
v(x) u(y) \text { if } x \geq y
\end{array}\right. \\
& +\frac{1}{W(y)}\left\{\begin{array}{l}
u(x) v^{\prime}(y) \text { if } x \leq y \\
v(x) u^{\prime}(y) \text { if } x \geq y
\end{array}\right.
\end{aligned}
$$

and hence the function $H$ defined in Eq. (41.22) is more explicitly given by

$$
\begin{equation*}
H(x)=\frac{1}{W(s)}\left(u(x) v^{\prime}(s)\right) h(s)-\frac{1}{W(r)}\left(v(x) u^{\prime}(r)\right) h(r) \tag{41.24}
\end{equation*}
$$

From this equation or the fact that $L_{x} G(x, r)=0=L_{x} G(x, s), H$ is is $L-$ harmonic on $J^{0}$. Moreover, from Eq. (41.24),

$$
\begin{aligned}
H(r) & =-\frac{1}{W(r)}\left(v(r) u^{\prime}(r)\right) h(r) \\
& =\frac{1}{W(r)}\left(u(r) v^{\prime}(r)-v(r) u^{\prime}(r)\right) h(r)=h(r)
\end{aligned}
$$

and

$$
\begin{aligned}
H(s) & =\frac{1}{W(s)}\left(u(s) v^{\prime}(s)\right) h(s) \\
& =\frac{1}{W(s)}\left(u(s) v^{\prime}(s)-v(s) u^{\prime}(s)\right) h(s)=h(s)
\end{aligned}
$$

Therefore if $f$ is defined by Eq. (41.21),

$$
L f=L G g-L H=g \text { a.e. on } J^{0}
$$

because $L G=I$ on $L^{2}(J)$ and

$$
\left.f\right|_{\partial J}=\left.(G g)\right|_{\partial J}+\left.H\right|_{\partial J}=\left.H\right|_{\partial J}=h
$$

since $G g \in H_{0}^{2}(J)$.
Corollary 41.16 (Elliptic Regularity I). Suppose $-\infty \leq r_{0}<s_{0} \leq \infty$, $J_{0}:=\left(r_{0}, s_{0}\right)$ and $L$ is as in Eq. (41.11) with the further assumption that $a, b, c \in C^{\infty}(\mathbb{R})$. If $f \in C^{2}\left(J_{0}\right)$ is a function such that $g:=L f \in C^{k}\left(J_{0}\right)$ for some $k \geq 0$, then $f \in C^{k+2}\left(J_{0}\right)$.

Proof. Let $r<s$ be chosen so that $J:=[r, s]$ is a bounded subinterval of $J_{0}$ and let $G$ be a Green's function as in Definition 41.14. Since $a, b, c$ are smooth, it follows from our general theory of ODE that $G(x, y) \in C^{\infty}(J \times J \backslash \Delta)$ where $\Delta=\{(x, x): x \in J\}$ is the diagonal in $J \times J$. Now by Theorem 41.15, for $x \in J^{0}$,

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$$
f(x)=(G g)(x)+\left.\left\{G(x, y) a(y) f^{\prime}(y)-[a(y) G(x, y)]_{y} f(y)\right\}\right|_{y=r} ^{y=s}
$$

Since

$$
\left.x \rightarrow\left\{G(x, y) a(y) f^{\prime}(y)-[a(y) G(x, y)]_{y} f(y)\right\}\right|_{y=r} ^{y=s} \in C^{\infty}\left(J^{0}\right)
$$

it suffices to show $G g \in C^{k+2}\left(J^{0}\right)$. But this follows by examining the formula for $(G g)^{\prime \prime}$ given on the right side of Eq. (41.18).

In fact we have the following rather striking version of this result.
Theorem 41.17 (Hypoellipticity). Suppose $-\infty \leq r_{0}<s_{0} \leq \infty, J_{0}:=$ $\left(r_{0}, s_{0}\right)$ and $L$ is as in Eq. (41.11) with the further assumption that $a, b, c \in$ $C^{\infty}(\mathbb{R})$. If $u \in \mathcal{D}^{\prime}\left(J_{0}\right)$ is a generalized function such that $v:=L u \in C^{\infty}\left(J_{0}\right)$, then $u \in C^{\infty}\left(J_{0}\right)$.

Proof. As in the proof of Corollary 41.16 let $r<s$ be chosen so that $J:=[r, s]$ is a bounded subinterval of $J_{0}$ and let $G$ be the Green's function constructed above. ${ }^{2}$ Further suppose $\xi \in J^{0}, \theta \in C_{c}^{\infty}\left(J^{0},[0,1]\right)$ such that $\theta=1$ in a neighborhood $U$ of $\xi$ and $\alpha \in C_{c}^{\infty}(V,[0,1])$ such that $\alpha=1$ in a neighborhood $V$ of $\xi$, see Figure 41.1. Finally suppose that $\phi \in C_{c}^{\infty}(V)$, then


Fig. 41.1. Constructing the cutoff functions, $\theta$ and $\alpha$.

$$
\begin{aligned}
\phi & =\theta \phi=\theta L^{*} G^{*} \phi=\theta L^{*}\left(M_{\alpha}+M_{1-\alpha}\right) G^{*} \phi \\
& =L^{*} M_{\alpha} G^{*} \phi+\theta L^{*} M_{1-\alpha} G^{*} \phi
\end{aligned}
$$

and hence

[^2] necessary to know $G^{*}(x, y)=G(y, x)$ where $G$ is a Green's function for $L$.

Now

$$
\left\langle L u, M_{\alpha} G^{*} \phi\right\rangle=\left\langle v, M_{\alpha} G^{*} \phi\right\rangle=\left\langle G M_{\alpha} v, \phi\right\rangle
$$

and writing $u=D^{n} T_{h}$ for some continuous function $h$ (which is always possible locally) we find

$$
\begin{aligned}
& \left\langle u, \theta L^{*} M_{1-\alpha} G^{*} \phi\right\rangle \\
& \quad=(-1)^{n}\left\langle u, D^{n} M_{\theta} L^{*} M_{1-\alpha} G^{*} \phi\right\rangle \\
& \quad=(-1)^{n} \int_{J \times J} h(x) D_{x}^{n}\left[\theta(x) L_{x}^{*}(1-\alpha(x)) G(y, x)\right] \phi(y) d y d x \\
& \quad=\int_{J} \psi(y) \phi(y) d y
\end{aligned}
$$

where

$$
\psi(y):=\int_{J} h(x) D_{x}^{n}\left[\theta(x) L_{x}^{*}(1-\alpha(x)) G(y, x)\right] d x
$$

which is smooth for $y \in V$ because $1-\alpha(x)=0$ on $V$ and so $(1-\alpha(x)) G(y, x)$ is smooth for $(x, y) \in J \times V$. Putting this altogether shows

$$
\langle u, \phi\rangle=\left\langle G M_{\alpha} v+\psi, \phi\right\rangle \text { for all } \phi \in C_{c}^{\infty}(V)
$$

That is to say $u=G M_{\alpha} v+\psi$ on $V$ which proves the theorem since $G M_{\alpha} v+\psi \in$ $C^{\infty}(V)$.
Example 41.18. Let $L=\frac{\partial^{2}}{\partial y^{2}}-\frac{\partial^{2}}{\partial x^{2}}$ be the wave operator on $\mathbb{R}^{2}$ which is not elliptic. Given $f \in C^{2}(\mathbb{R})$ we have already seen that $L f(y-x)=0 \in C^{\infty}\left(\mathbb{R}^{2}\right)$. Clearly since $f$ was arbitrary, it does not follow that $F(x, y):=f(y-x) \in$ $C^{\infty}\left(\mathbb{R}^{2}\right)$. Moreover, if $f$ is merely continuous and $F(x, y):=f(y-x)$, then $L T_{F}=0$ with $F \notin C^{2}\left(\mathbb{R}^{2}\right)$. To check $L T_{F}=0$ we first observe

$$
\begin{aligned}
-\left\langle\left(\partial_{x}+\partial_{y}\right) T_{F}, \phi\right\rangle & =\left\langle T_{F},\left(\partial_{x}+\partial_{y}\right) \phi\right\rangle \\
& =\int_{\mathbb{R}^{2}} f(y-x)\left(\partial_{x}+\partial_{y}\right) \phi(x, y) d x d y \\
& =\int_{\mathbb{R}^{2}} f(y)\left[\phi_{x}(x, y+x)+\phi_{y}(x, y+x)\right] d x d y \\
& =\int_{\mathbb{R}^{2}} f(y) \frac{\partial}{\partial x}[\phi(x, y+x)] d x d y=0 .
\end{aligned}
$$

Therefore $L T_{F}=\left(\partial_{x}-\partial_{y}\right)\left(\partial_{x}+\partial_{y}\right) T_{F}=0$ as well.
Corollary 41.19. Suppose $a, b, c$ are smooth and $u \in \mathcal{D}^{\prime}\left(J^{0}\right)$ is an eigenvector for $L$, i.e. $L u=\lambda u$ for some $\lambda \in \mathbb{C}$. Then $u \in C^{\infty}(J)$.

Proof. Since $L-\lambda$ is an elliptic ordinary differential operator and $(L-\lambda) u=0 \in C^{\infty}\left(J^{0}\right)$, it follows by Theorem 41.17 that $u \in C^{\infty}\left(J^{0}\right)$.

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$$
\mathcal{E}(f, g):=\int_{J}\left[a f^{\prime} g^{\prime}+c f g\right] \rho d m \text { for } f, g \in H^{1}(J) .
$$

Lemma 41.23. For $f, g \in D(L)$,

$$
\begin{equation*}
(L f, g)_{\rho}=\mathcal{E}(f, g)=(f, L g)_{\rho} . \tag{41.25}
\end{equation*}
$$

Moreover

$$
\mathcal{E}(f, f) \geq a_{0}\left\|f^{\prime}\right\|_{2}^{2}+c_{0}\|f\|_{2}^{2} \text { for all } f \in H^{1}(J)
$$

where $c_{0}:=\min _{J} c$ and $a_{0}=\min _{J} a$. If $\lambda_{0} \in \mathbb{R}$ with $\lambda_{0}+c_{0}>0$ then

$$
\begin{equation*}
\|f\|_{H^{1}(J)}^{2} \leq K\left[\mathcal{E}(f, f)+\lambda_{0}\|f\|_{2}^{2}\right] \tag{41.26}
\end{equation*}
$$

where $K=\left[\min \left(a_{0}, c_{0}+\lambda_{0}\right)\right]^{-1}$
Proof. Eq. (41.25) is a simple consequence of integration by parts. By elementary estimates

$$
\mathcal{E}(f, f) \geq a_{0}\left\|f^{\prime}\right\|_{2}^{2}+c_{0}\|f\|_{2}^{2}
$$

and

$$
\mathcal{E}(f, f)+\lambda_{0}\|f\|_{2}^{2} \geq a_{0}\left\|f^{\prime}\right\|_{2}^{2}+\left(c_{0}+\lambda_{0}\right)\|f\|_{2}^{2} \geq \min \left(a_{0}, c_{0}+\lambda_{0}\right)\|f\|_{H^{1}(J)}^{2}
$$

which proves Eq. (41.26).
Corollary 41.24. Suppose $\lambda_{0}+c_{0}>0$ then $\operatorname{Nul}\left(L+\lambda_{0}\right) \cap D(L)=0$ and hence

$$
\left(L+\lambda_{0}\right): H_{0}^{2}(J) \rightarrow L^{2}(J)
$$

is invertible and the resolvent $\left(L+\lambda_{0}\right)^{-1}$ has a continuous integral kernel $G(x, y)$, i.e.

$$
\left(L+\lambda_{0}\right)^{-1} u(x)=\int_{J} G(x, y) u(y) d y
$$

Moreover if we define $D\left(L^{k}\right)$ inductively by

$$
D\left(L^{k}\right):=\left\{u \in D\left(L^{k-1}\right): L^{k-1} u \in D(L)\right\}
$$

we have $D\left(L^{k}\right)=H_{0}^{2 k}(J)$.
Proof. By Lemma 41.23, for all $u \in D(L)$,

$$
\|u\|_{H^{1}(J)}^{2} \leq K\left((L u, u)+\lambda_{0}\|u\|_{2}^{2}\right)=K\left(\left(\left(L+\lambda_{0}\right) u, u\right)\right)
$$

so that if $\left(L+\lambda_{0}\right) u=0$, then $\|u\|_{H^{1}(J)}^{2}=0$ and hence $u=0$. The remaining assertions except for $D\left(L^{k}\right)=H_{0}^{k}(J)$ now follow directly from Theorem 41.11

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applied with $L$ replaced by $L+\lambda_{0}$. Finally if $u \in D(L)$ then $\left(L+\lambda_{0}\right) u=$ $L u+\lambda_{0} u \in L^{2}(J)$ and therefore

$$
u=\left(L+\lambda_{0}\right)^{-1}\left(L u+\lambda_{0} u\right) \in H_{0}^{2}(J) .
$$

Now suppose we have shown, $D\left(L^{k}\right)=H_{0}^{2 k}(J)$ and $u \in D\left(L^{k+1}\right)$, then

$$
\left(L+\lambda_{0}\right) u=L u+\lambda_{0} u \in D\left(L^{k}\right)+D\left(L^{k+1}\right) \subset D\left(L^{k}\right)=H_{0}^{2 k}(J)
$$

and so by Theorem 41.11, $u \in\left(L+\lambda_{0}\right)^{-1} H_{0}^{2 k}(J) \subset H_{0}^{2 k+2}(J)$.
Corollary 41.25. There exists an orthonormal basis $\left\{\phi_{n}\right\}_{n=0}^{\infty}$ for $L^{2}(J, \rho d m)$ of eigenfunctions of $L$ with eigenvalues $\lambda_{n} \in \mathbb{R}$ such that $-c_{0} \leq \lambda_{0}<\lambda_{1}<$ $\lambda_{2}<\ldots$..

Proof. Let $\lambda_{0}>-c_{0}$ and let $G:=\left(L+\lambda_{0}\right)^{-1}: L^{2}(J) \rightarrow H_{0}^{2}(J)=D(L) \subset$ $L^{2}(J)$. From the theory of compact operators to be developed later, $G$ is a compact symmetric positive definite operator on $L^{2}(J)$ and hence there exists an orthonormal basis $\left\{\phi_{n}\right\}_{n=0}^{\infty}$ for $L^{2}(J, \rho d m)$ of eigenfunctions of $G$ with eigenvalues $\mu_{n}>0$ such that $\mu_{0} \geq \mu_{1} \geq \mu_{2} \geq \ldots \rightarrow 0 .{ }^{3}$ Since

$$
\mu_{n} \phi_{n}=G \phi_{n}=\left(L+\lambda_{0}\right)^{-1} \phi_{n}
$$

it follows that $\mu_{n}\left(L+\lambda_{0}\right) \phi_{n}=\phi_{n}$ for all $n$ and therefore $L \phi_{n}=\lambda_{n} \phi_{n}$ with $\lambda_{n}=\left(\mu_{n}^{-} 1-\lambda_{0}\right) \uparrow \infty$. Finally since $L$ is a second order ordinary differential equation there can be at most one linearly independent eigenvector for a given eigenvalue $\lambda_{n}$ and hence $\lambda_{n}<\lambda_{n}+1$ for all $n$.

Example 41.26. Let $J=[0, \pi], \rho=1$ and $L=-D^{2}$ on $H_{0}^{2}(J)$. Then $L \phi=\lambda \phi$ implies $\phi^{\prime \prime}+\lambda \phi=0$. Since $L$ is positive, we need only consider the case where $\lambda \geq 0$ in which case $\phi(x)=a \cos (\sqrt{\lambda} x)+b \sin (\sqrt{\lambda} x)$. The boundary conditions for $f$ imply $a=0$ and $0=\sin (\sqrt{\lambda} \pi)$, i.e. $\sqrt{\lambda} \in \mathbb{N}_{+}$. Therefore in this example

$$
\phi_{k}(x)=\sqrt{\frac{2}{\pi}} \sin (k x) \text { with } \lambda_{k}=k^{2}
$$

The collection of functions $\left\{\phi_{k}\right\}_{k=1}^{\infty}$ is an orthonormal basis for $L^{2}(J)$.
Theorem 41.27. Let $J=[r, s]$ and $\rho, a \in C^{2}(J,(0, \infty)), c \in C^{2}(J)$ and $L$ be defined by

$$
L f=-\frac{1}{\rho} D\left(\rho a f^{\prime}\right)+c f .
$$

$\overline{{ }^{3} \text { In fact } G}$ is "Hilbert Schmidt" which then implies

$$
\sum_{n=0}^{\infty} \mu_{n}^{2}<\infty
$$

and for $\lambda \in \mathbb{R}$ let

$$
E^{\lambda}:=\left\{\phi \in H_{0}^{2}(J): L \phi=\alpha \phi \text { for some } \alpha<\lambda\right\} .
$$

Then there are constants $d_{1}, d_{2}>0$ such that

$$
\begin{equation*}
\operatorname{dim}\left(E^{\lambda}\right) \leq d_{1} \lambda+d_{2} . \tag{41.27}
\end{equation*}
$$

Proof. For $\lambda \in \mathbb{R}$ let $E_{\lambda}:=\left\{\phi \in H_{0}^{2}(J): L \phi=\lambda \phi\right\}$. By Corollary 41.24, $E_{\lambda}=\{0\}$ if $\lambda<c_{0}$ and since $(L f, g)_{\rho}=(f, L g)_{\rho}$ for all $f, g \in H_{0}^{2}(J)$ it follows that $E_{\lambda} \perp E_{\beta}$ for all $\lambda \neq \beta$. Indeed, if $f \in E_{\lambda}$ and $g \in E_{\beta}$, then

$$
(\beta-\lambda)(f, g)_{\rho}=(f, L g)_{\rho}-(L f, g)_{\rho}=0
$$

Thus it follow that any finite dimensional subspace $W \subset E^{\lambda}$ has an orthonormal basis (relative to $(\cdot, \cdot)_{\rho}$ - inner product) of eigenvectors $\left\{\phi_{k}\right\}_{k=1}^{n} \subset E^{\lambda}$ of $L$, say $L \phi_{k}=\lambda_{k} \phi_{k}$. Let $u=\sum_{k=1}^{n} u_{k} \phi_{k}$ where $u_{k} \in \mathbb{R}$. By Proposition 41.21 and Lemma 41.23,

$$
\|u\|_{u}^{2} \leq C\|u\|_{H^{1}(J)}^{2} \leq C\left(\left(L+\lambda_{0}\right) u, u\right)_{\rho}=C\left(\sum_{k=1}^{n} u_{k}\left(\lambda_{k}+\lambda_{0}\right) \phi_{k}, u\right)_{\rho}
$$

(where $C$ is a constant varying from place to place but independent of $u$ ) and hence for any $x \in J$,

$$
\left|\sum_{k=1}^{n} u_{k} \phi_{k}(x)\right|^{2} \leq\|u\|_{u}^{2} \leq C\left(\lambda+\lambda_{0}\right) \sum_{k=1}^{n}\left|u_{k}\right|^{2} .
$$

Now choose $u_{k}=\phi_{k}(x)$ in this equation to find

$$
\left.\left.\left|\sum_{k=1}^{n}\right| \phi_{k}(x)\right|^{2}\right|^{2} \leq C\left(\lambda+\lambda_{0}\right) \sum_{k=1}^{n}\left|\phi_{k}(x)\right|^{2}
$$

or equivalently that

$$
\sum_{k=1}^{n}\left|\phi_{k}(x)\right|^{2} \leq C\left(\lambda+\lambda_{0}\right) .
$$

Multiplying this equation by $\rho$ and then integrating shows

$$
\operatorname{dim}(W)=n=\sum_{k=1}^{n}\left(\phi_{k}, \phi_{k}\right)_{\rho} \leq C\left(\lambda+\lambda_{0}\right) \int_{J} \rho d m=C^{\prime}\left(\lambda+\lambda_{0}\right)
$$

Since $W \subset E^{\lambda}$ is arbitrary, it follows that

$$
\operatorname{dim}\left(E^{\lambda}\right) \leq C^{\prime}\left(\lambda+\lambda_{0}\right)
$$

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Remarks 41.28 Notice that for all $\lambda \in \mathbb{R}, \operatorname{dim}\left(E_{\lambda}\right) \leq 1$ because if $u, v \in E_{\lambda}$ then by uniqueness of solutions to $O D E, u=\left[u^{\prime}(r) / v^{\prime}(r)\right] v$. Let $\left\{\phi_{k}\right\}_{k=1}^{\infty} \subset$ $H_{0}^{2}(J) \cap C^{\infty}(J)$ be the eigenvectors of $L$ ordered so that the corresponding eigenvalues are increasing. With this ordering we have $k=\operatorname{dim}\left(E^{\lambda_{k}}\right) \leq d_{1} \lambda_{k}+$ $d_{2}$ and therefore,

$$
\begin{equation*}
\lambda_{k} \geq d_{1}^{-1}\left(k-d_{2}\right) \tag{41.28}
\end{equation*}
$$

The estimates in Eqs. (41.27) and (41.28) are not particularly good as Example 41.26 illustrates.

### 41.4 Associated Heat and Wave Equations

Lemma 41.29. $L$ is a closed operator, i.e. if $s_{n} \in D(L)$ and $s_{n} \rightarrow s$ and $L s_{n} \rightarrow g$ in $L^{2}$, then $s \in D(L)$ and $L s=g$. In particular if $f_{k} \in D(L)$ and $\sum_{k=1}^{\infty} f_{k}$ and $\sum_{k=1}^{\infty} L f_{k}$ exists in $L^{2}$, then $\sum_{k=1}^{\infty} f_{k} \in D(L)$ and

$$
L \sum_{k=1}^{\infty} f_{k}=\sum_{k=1}^{\infty} L f_{k}
$$

Proof. Let $\lambda_{0}+c_{0}>0$ and $G=\left(L+\lambda_{0}\right)^{-1}$. Then by assumption $\left(L+\lambda_{0}\right) s_{n} \rightarrow g+\lambda_{0} s$ and so

$$
s \leftarrow s_{n}=G\left(L+\lambda_{0}\right) s_{n} \rightarrow G\left(g+\lambda_{0} s\right) \text { as } n \rightarrow \infty
$$

showing $s=G g \in D\left(L+\lambda_{0}\right)=D(L)$ and

$$
\left(L+\lambda_{0}\right) s=\left(L+\lambda_{0}\right) G\left(g+\lambda_{0} s\right)=g+\lambda_{0} s
$$

and hence $L s=g$ as desired. The assertions about the sums follow by applying the sequence results to $s_{n}=\sum_{k=1}^{n} f_{k}$.
Theorem 41.30. Given $f \in L^{2}$, let

$$
\begin{equation*}
u(t)=e^{-t L} f=\sum_{n=0}^{\infty}\left(f, \phi_{n}\right) e^{-t \lambda_{n}} \phi_{n} \tag{41.29}
\end{equation*}
$$

Then for $t>0, u(t, x)$ is smooth in $(t, x)$ and solves the heat equation

$$
\begin{align*}
u_{t}(t, x) & =-L u(t, x), u(t, x)=0 \text { for } x \in \partial J  \tag{41.30}\\
\text { and } f & =L^{2}-\lim _{t \downarrow 0} u(t) \tag{41.31}
\end{align*}
$$

Moreover, $u(t, x)=\int_{J} p_{t}(x, y) f(y) \rho(y) d y$ where

$$
\begin{equation*}
p_{t}(x, y):=\sum_{n=0}^{\infty} e^{-t \lambda_{n}} \phi_{n}(x) \phi_{n}(y) \tag{41.32}
\end{equation*}
$$

is a smooth function in $t>0$ and $x, y \in J$. The function $p_{t}$ is called the Diurichlet Heat Kernel for $L$.

Proof. (Sketch.) For any $t>0$ and $k \in \mathbb{N}, \sup _{n}\left(e^{-t \lambda_{n}} \lambda_{n}^{k}\right)<\infty$ and so by Lemma 41.29, for $t>0, u(t) \in D\left(L^{k}\right)=H_{0}^{2 k}(J)^{4}$ (Corollary 41.24) and

$$
L^{k} u(t)=\sum_{n=0}^{\infty}\left(f, \phi_{n}\right) e^{-t \lambda_{n}} \lambda_{n}^{k} \phi_{n}
$$

Also we have $L^{k} u^{(m)}(t)$ exists in $L^{2}$ for all $k, m \in \mathbb{N}$ and

$$
L^{k} u^{(m)}(t)=(-1)^{m} \sum_{n=0}^{\infty}\left(f, \phi_{n}\right) e^{-t \lambda_{n}} \lambda_{n}^{k+m} \phi_{n}
$$

By Sobolev inequalities and elliptic estimates such as Proposition 41.21 and Lemma 41.23, one concludes that $u \in C^{\infty}\left((0, \infty), H_{0}^{k}(J)\right)$ for all $k$ and then that $u \in C^{\infty}((0, \infty) \times J, \mathbb{R})$. Eq. (41.30) is now relatively easy to prove and Eq. (41.31) follows from the following computation

$$
\|f-u(t)\|_{2}^{2}=\sum_{n=1}^{\infty}\left|\left(f, \phi_{n}\right)\right|^{2}\left|1-e^{-t \lambda_{n}}\right|^{2}
$$

which goes to 0 as $t \downarrow 0$ by the D.C.T. for sums.
Finally from Eq. (41.29)

$$
\begin{aligned}
u(t, x) & =\sum_{n=0}^{\infty} \int_{J} f(y) \phi(y) \rho(y) d y e^{-t \lambda_{n}} \phi_{n}(x) \\
& =\int_{J} \sum_{n=0}^{\infty} e^{-t \lambda_{n}} \phi_{n}(x) \phi(y) f(y) \rho(y) d y
\end{aligned}
$$

where the interchange of the sum and the integral is permissible since

$$
\begin{aligned}
& \int_{J} \sum_{n=0}^{\infty} e^{-t \lambda_{n}}\left|\phi_{n}(x) \phi(y) f(y)\right| \rho(y) d y \\
& \leq C \int_{J} \sum_{n=0}^{\infty} e^{-t \lambda_{n}}\left(\lambda_{0}+\lambda_{n}\right)^{2}|f(y)| \rho(y) d y<\infty
\end{aligned}
$$

since $\sum_{n=0}^{\infty} e^{-t \lambda_{n}}\left(\lambda_{0}+\lambda_{n}\right)^{2}<\infty$ because $\lambda_{n}$ grows linearly in $n$. Moreover one similarly shows

$$
\left(\frac{\partial}{\partial t}\right)^{j} \partial_{x}^{2 k-1} \partial_{y}^{2 l-1} p_{t}(x, y)=\sum_{n=0}^{\infty}\left(-\lambda_{n}\right)^{j} e^{-t \lambda_{n}} \partial_{x}^{2 k-1} \phi_{n}(x) \partial_{y}^{2 l-1} \phi(y)
$$

where the above operations are permissible since

[^3]$$
\left\|\phi_{n}^{(2 k-1)}\right\|_{u} \leq C\left\|\phi_{n}\right\|_{H_{0}^{2 k}(J)} \leq C\left\|\left(L+\lambda_{0}\right)^{k} \phi_{n}\right\|_{2}=C\left(\lambda_{n}+\lambda_{0}\right)^{k}
$$
and therefore
$\sum_{n=0}^{\infty}\left|\left(-\lambda_{n}\right)^{j} e^{-t \lambda_{n}} \partial_{x}^{2 k-1} \phi_{n}(x) \partial_{y}^{2 l-1} \phi(y)\right| \leq C \sum_{n=0}^{\infty}\left|\lambda_{n}\right|^{j}\left(\lambda_{n}+\lambda_{0}\right)^{k+l} e^{-t \lambda_{n}}<\infty$.
Again we use $\lambda_{n}$ grows linearly with $n$. From this one may conclude that $p_{t}(x, y)$ is smooth for $t>0$ and $x, y \in J$. (We will do this in more detail when we work out the higher dimensional analogue.)
Remark 41.31 (Wave Equation). Suppose $f \in D\left(L^{k}\right)$, then
$$
\left|\left(f, \phi_{n}\right)\right|=\left|\frac{1}{\lambda_{n}^{k}}\left(f, L^{k} \phi_{n}\right)\right|=\left|\frac{1}{\lambda_{n}^{k}}\left(L^{k} f, \phi_{n}\right)\right| \leq \frac{1}{\left|\lambda_{n}^{k}\right|}\left\|L^{k} f\right\|_{2}
$$
and therefore
$$
\cos (t \sqrt{L}) f:=\sum_{n=0}^{\infty} \cos \left(t \sqrt{\lambda_{n}}\right)\left(f, \phi_{n}\right) \phi_{n}
$$
will be convergent in $L^{2}$ but moreover
\[

$$
\begin{aligned}
L^{k} \cos (t \sqrt{L}) f & =\sum_{n=0}^{\infty} \cos \left(t \sqrt{\lambda_{n}}\right)\left(f, \phi_{n}\right) \lambda_{n}^{k} \phi_{n} \\
& =\sum_{n=0}^{\infty} \cos \left(t \sqrt{\lambda_{n}}\right)\left(L^{k} f, \phi_{n}\right) \phi_{n}
\end{aligned}
$$
\]

will also be convergent. Therefore if we let

$$
u(t)=\cos (t \sqrt{L}) f+\frac{\sin (t \sqrt{L})}{\sqrt{L}} g
$$

where $f, g \in D\left(L^{k}\right)$ for all $k$. Then we will get a solution to the wave equation

$$
u_{t t}(t, x)+L u(t, x)=0 \text { with } u(0)=f \text { and } \dot{u}(0)=g .
$$

More on all of this later.

### 41.5 Extensions to Other Boundary Conditions

In this section, we will assume $\rho \in C^{2}(J,(0, \infty))$,

$$
\begin{equation*}
L u=-\rho^{-1}\left(\rho a u^{\prime}\right)^{\prime}+b u^{\prime}+c u \tag{41.33}
\end{equation*}
$$

is an elliptic ODE on $L^{2}(J)$ with smooth coefficients and

$$
\begin{equation*}
(u, v)=(u, v)_{\rho}=\int_{J} u(x) v(x) \rho(x) d x . \tag{41.34}
\end{equation*}
$$

Theorem 41.32. For $v \in H^{2}(J)$ let

$$
\begin{equation*}
L^{*} v=-\rho^{-1}\left(\rho a v^{\prime}\right)^{\prime}-b v^{\prime}+\left[c-\rho^{-1}(\rho b)^{\prime}\right] v . \tag{41.35}
\end{equation*}
$$

Then for $u, v \in H^{2}(J)$,

$$
\begin{equation*}
(L u, v)=\left(u, L^{*} v\right)+\left.\mathcal{B}(u, v)\right|_{\partial J} \tag{41.36}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathcal{B}(u, v)=\rho a\left\{\left(u^{\prime}, u\right) \cdot\left(-v, v^{\prime}+\frac{b}{a} v\right)\right\} . \tag{41.37}
\end{equation*}
$$

Proof. This is an exercise in integration by parts,

$$
\begin{aligned}
(L u, v) & =\int_{J}\left(-\left(\rho a u^{\prime}\right)^{\prime}+\rho b u^{\prime}+\rho c u\right) v d m \\
& =\int_{J}\left(\rho a u^{\prime} v^{\prime}-(\rho b v)^{\prime} u+\rho c u\right) d m+\left.\left[\rho b u v-\rho a u^{\prime} v\right]\right|_{\partial J} \\
& =\int_{J}\left(-u\left(\rho a v^{\prime}\right)^{\prime}-(\rho b v)^{\prime} u+\rho c v u\right) d m \\
& +\left.\left[\rho b u v+\rho a u v^{\prime}-\rho a u^{\prime} v\right]\right|_{\partial J} \\
& =\int_{J}\left(-u \rho^{-1}\left(\rho a v^{\prime}\right)^{\prime}-\rho^{-1}(\rho b v)^{\prime} u+c v u\right) \rho d m \\
& +\left.\left[\rho a\left(\frac{b}{a} u v+u v^{\prime}-v u^{\prime}\right)\right]\right|_{\partial J} \\
& =\left(u, L^{*} v\right)+\left.\left[\rho a\left(u^{\prime}, u\right) \cdot\left(-v, v^{\prime}+\frac{b}{a} v\right)\right]\right|_{\partial J} .
\end{aligned}
$$

Notation 41.33 Given $(\alpha, \beta): \partial J \rightarrow \mathbb{R}^{2} \backslash\{0\}$ and $u, v \in H^{2}(J)$ let

$$
B u=\alpha u^{\prime}+\beta u=(\alpha, \beta) \cdot\left(u^{\prime}, u\right) \text { on } \partial J
$$

and

$$
B^{*} v=\alpha v^{\prime}+\left(\beta+\frac{b}{a} \alpha\right) v=\alpha v^{\prime}+\tilde{\beta} v \text { on } \partial J
$$

where $\tilde{\beta}:=\left(\beta+\frac{b}{a} \alpha\right)$.
Remarks 41.34 The function $(\alpha, \tilde{\beta}): \partial J \rightarrow \mathbb{R}^{2}$ also takes values in $\mathbb{R}^{2} \backslash\{0\}$ because $(\alpha, \tilde{\beta})=0$ iff $(\alpha, \beta)=0$. Furthermore if $\alpha=0$ then $\tilde{\beta}=\beta$.

Proposition 41.35. Let $B$ and $B^{*}$ be as defined in Notation 41.33 and define

$$
\begin{aligned}
D(L) & =\left\{u \in H^{2}(J): B u=0 \text { on } \partial J\right\}, \\
D\left(L^{*}\right) & =\left\{u \in H^{2}(J): B^{*} u=0 \text { on } \partial J\right\},
\end{aligned}
$$

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Then $v \in H^{2}(J)$ satisfies

$$
\begin{equation*}
(L u, v)=\left(u, L^{*} v\right) \text { for all } u \in D(L) \tag{41.38}
\end{equation*}
$$

iff $v \in D\left(L^{*}\right)$. (This result will be substantially improved on in Theorem 41.41 below.)

Proof. We have to check that $\mathcal{B}(u, v)$ appearing in Eq. (41.36) is 0 . (Actually we must check that $\left.\mathcal{B}(u, v)\right|_{\partial J}=0$ which we might arrange by using something like "periodic boundary conditions." I am not considering this type of condition at the moment. Since $u$ may be chosen to be zero near $r$ or $s$ we must require $\mathcal{B}(u, v)=0$ on $\partial J$.) Now $\mathcal{B}(u, v)=0$ iff

$$
\begin{equation*}
\left(u^{\prime}, u\right) \cdot\left(-v, v^{\prime}+\frac{b}{a} v\right)=0 \tag{41.39}
\end{equation*}
$$

which happens iff $\left(u^{\prime}, u\right)$ is parallel to $\left(v^{\prime}+\frac{b}{a} v, v\right)$. The boundary condition $B u=0$ may be rewritten as saying $\left(u^{\prime}, u\right) \cdot(\alpha, \beta)=0$ or equivalently that $\left(u^{\prime}, u\right)$ is parallel to $(-\beta, \alpha)$ on $\partial J$. Therefore the condition in Eq. (41.39) is equivalent to $(-\beta, \alpha)$ is parallel to $\left(v^{\prime}+\frac{b}{a} v, v\right)$ or equivalently that

$$
0=(\alpha, \beta) \cdot\left(v^{\prime}+\frac{b}{a} v, v\right)=B^{*} v
$$

Corollary 41.36. The formulas for $L$ and $L^{*}$ agree iff $b=0$ in which case

$$
L u=-\rho^{-1} D\left(a \rho u^{\prime}\right)+c u,
$$

$B=B^{*}, D(L)=D\left(L^{*}\right)$ and

$$
\begin{equation*}
(L u, v)=(u, L v) \text { for all } u, v \in D(L) . \tag{41.40}
\end{equation*}
$$

(In fact L is a "self-adjoint operator," as we will see later by showing $\left(L+\lambda_{0}\right)^{-1}$ exists for $\lambda_{0}$ sufficiently large. Eq. (41.40) then may be used to deduce $\left(L+\lambda_{0}\right)^{-1}$ is a bounded self-adjoint operator with a symmetric Green's functions $G$.)

### 41.5.1 Dirichlet Forms Associated to $(L, D(L))$

For the rest of this section let $a, b_{1}, b_{2}, c_{0}, \rho \in C^{2}(J)$, with $a>0$ and $\rho>0$ on $J$ and for $u, v \in H^{1}(J)$, let

$$
\begin{equation*}
\mathcal{E}(u, v):=\int_{J}\left(a u^{\prime} v^{\prime}+b_{1} u v^{\prime}+b_{2} u^{\prime} v+c_{0} u v\right) \rho d m \text { and } \tag{41.41}
\end{equation*}
$$

$$
\|u\|_{H^{1}(J)}:=\left(\left\|u^{\prime}\right\|^{2}+\|u\|^{2}\right)^{1 / 2}
$$

where $\|u\|^{2}=(u, u)_{\rho}$ as defined in Eq. (41.34).

Lemma 41.37 (A Coercive inequality for $\mathcal{E}$ ). There is a constant $K<\infty$ such that

$$
\begin{equation*}
|\mathcal{E}(u, v)| \leq K\|u\|_{H^{1}(J)}\|v\|_{H^{1}(J)} \text { for } u, v \in H^{1}(J) \tag{41.42}
\end{equation*}
$$

Let $a_{0}=\min _{J} a, \bar{c}=\min _{J} c_{0}$ and $B:=\max _{J}\left|b_{1}+b_{2}\right|$, then for $u \in H^{1}(J)$,

$$
\begin{equation*}
\mathcal{E}(u, u) \geq \frac{a_{0}}{2}\left\|u^{\prime}\right\|^{2}+\left(\bar{c}-\frac{B^{2}}{2 a_{0}}\right)\|u\|^{2} \tag{41.43}
\end{equation*}
$$

Proof. Let $A=\max _{J} a, B_{i}=\max _{J}\left|b_{i}\right|$ and $C_{0}:=\max _{J}\left|c_{0}\right|$, then

$$
\begin{aligned}
|\mathcal{E}(u, v)| & \leq \int_{J}\left(a\left|u^{\prime}\right|\left|v^{\prime}\right|+\left|b_{1}\right||u|\left|v^{\prime}\right|+\left|b_{2}\right|\left|u^{\prime}\right||v|+\left|c_{0}\right||u||v|\right) \rho d m \\
& \leq A\left\|u^{\prime}\right\|\left\|v^{\prime}\right\|+B_{1}\|u\|\left\|v^{\prime}\right\|+B_{2}\left\|u^{\prime}\right\|\|v\|+C_{0}\|u\|\|v\| \\
& \leq K\left(\left\|u^{\prime}\right\|^{2}+\|u\|^{2}\right)^{1 / 2}\left(\left\|v^{\prime}\right\|^{2}+\|v\|^{2}\right)^{1 / 2}
\end{aligned}
$$

Let $a_{0}=\min _{J} a, \bar{c}=\min _{J} c$ and $B:=\max _{J}\left|b_{1}+b_{2}\right|$, then for any $\delta>0$,

$$
\begin{aligned}
\mathcal{E}(u, u) & =\int_{J}\left(a\left|u^{\prime}\right|^{2}+\left(b_{1}+b_{2}\right) u u^{\prime}+c_{0}|u|^{2}\right) \rho d m \\
& \geq a_{0}\left\|u^{\prime}\right\|^{2}+\bar{c}\|u\|^{2}-B \int_{J}|u|\left|u^{\prime}\right| \rho d m \\
& \geq a_{0}\left\|u^{\prime}\right\|^{2}+\bar{c}\|u\|^{2}-\frac{B}{2}\left(\delta\left\|u^{\prime}\right\|^{2}+\delta^{-1}\|u\|^{2}\right) \\
& =\left(a_{0}-\frac{B \delta}{2}\right)\left\|u^{\prime}\right\|^{2}+\left(\bar{c}-\frac{B}{2} \delta^{-1}\right)\|u\|^{2}
\end{aligned}
$$

Taking $\delta=a_{0} / B$ in this equation proves Eq. (41.43).
Theorem 41.38. Let

$$
\begin{align*}
b & =\left(b_{2}-b_{1}\right), c:=c_{0}-\rho^{-1}\left(\rho b_{1}\right)^{\prime},  \tag{41.44}\\
L u & =-\rho^{-1}\left(a \rho u^{\prime}\right)^{\prime}+b u^{\prime}+c u \text { and } \\
B u & =\left.\left(\rho a u^{\prime}+\rho b_{1} u\right)\right|_{\partial J} .
\end{align*}
$$

Then for $u \in H^{2}(J)$ and $v \in H^{1}(J)$

$$
\mathcal{E}(u, v)=(L u, v)+[(B u) v]_{\partial J}
$$

and for $u \in H^{1}(J)$ and $v \in H^{2}(J)$,

$$
\mathcal{E}(u, v)=\left(u, L^{*} v\right)+\left[\left(B^{*} v\right) u\right]_{\partial J} .
$$

Here (as in Eq. (41.35)
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$$
L^{*} v=-\rho^{-1}\left(a \rho u^{\prime}\right)^{\prime}-\rho^{-1}[\rho b u]^{\prime}+c u
$$

and (as in Notation 41.33)

$$
B^{*} v=\rho a v^{\prime}+\left(\rho b_{1}+\frac{b}{a} \rho a\right) v=\rho a v^{\prime}+\rho b_{2} v
$$

Proof. Let $u \in H^{2}(J)$ and $v \in H^{1}(J)$ and integrating Eq. (41.41) by parts to find

$$
\begin{align*}
\mathcal{E}(u, v) & =\int_{J}\left(-\rho^{-1}\left(a \rho u^{\prime}\right)^{\prime} v-\rho^{-1}\left(\rho b_{1} u\right)^{\prime} v+b_{2} u^{\prime} v+c_{0} u v\right) \rho d m \\
& +\left[\rho a u^{\prime} v+\rho b_{1} u v\right]_{\partial J} \\
& =(L u, v)+[B u \cdot v]_{\partial J} \tag{41.45}
\end{align*}
$$

where

$$
\begin{aligned}
L u & =-\rho^{-1}\left(a \rho u^{\prime}\right)^{\prime}-\rho^{-1}\left(\rho b_{1} u\right)^{\prime}+b_{2} u^{\prime}+c_{0} u \\
& =-\rho^{-1}\left(a \rho u^{\prime}\right)^{\prime}+\left(b_{2}-b_{1}\right) u^{\prime}+\left[c_{0}-\rho^{-1}\left(\rho b_{1}\right)^{\prime}\right] u \\
& =-\rho^{-1}\left(a \rho u^{\prime}\right)^{\prime}+b u^{\prime}+c u
\end{aligned}
$$

and

$$
B u=\rho a u^{\prime}+\rho b_{1} u
$$

Similarly

$$
\begin{aligned}
\mathcal{E}(u, v) & =\int_{J}\left(-u \rho^{-1}\left(a \rho v^{\prime}\right)^{\prime}+b_{1} u v^{\prime}-u \rho^{-1}\left(\rho b_{2} v\right)^{\prime}+c_{0} u v\right) \rho d m \\
& +\left[\left(\rho a u v^{\prime}+\rho b_{2} u v\right)\right]_{\partial J} \\
& =\left(u, L^{\dagger} v\right)+\left[B^{\dagger} v \cdot u\right]_{\partial J}
\end{aligned}
$$

where

$$
\begin{aligned}
L^{\dagger} v & =-\rho^{-1}\left(a \rho v^{\prime}\right)^{\prime}+b_{1} v^{\prime}-\rho^{-1}\left(\rho b_{2} v\right)^{\prime}+c_{0} v \\
& =-\rho^{-1}\left(a \rho v^{\prime}\right)^{\prime}+\left(b_{1}-b_{2}\right) v^{\prime}+\left[c_{0}-\rho^{-1}\left(\rho b_{2}\right)^{\prime}\right] v \\
& =-\rho^{-1}\left(a \rho v^{\prime}\right)^{\prime}-b v^{\prime}+\left[c+\rho^{-1}\left(\rho\left(b_{1}-b_{2}\right)\right)^{-1}\right] v \\
& =-\rho^{-1}\left(a \rho v^{\prime}\right)^{\prime}-b v^{\prime}+\left[c-\rho^{-1}(\rho b)^{\prime}\right] v=L^{*} v
\end{aligned}
$$

and

$$
B^{\dagger} v=\left(\rho a v^{\prime}+\rho b_{2} v\right)=B^{*} v
$$

Remark 41.39. As a consequence of Theorem 41.38, the mapping

$$
\left(a, b_{1}, b_{2}, c_{0}\right) \rightarrow\left[\begin{array}{c}
(u, v) \rightarrow \\
\mathcal{E}(u, v)=\int_{J}\left(a u^{\prime} v^{\prime}+b_{1} u v^{\prime}+b_{2} u^{\prime} v+c_{0} u v\right) \rho d m
\end{array}\right]
$$

is highly non-injective. In fact $\mathcal{E}$ depends only on $a, b=b_{2}-b_{1}$ and $c:=$ $c_{0}-\rho^{-1}\left(\rho b_{1}\right)^{\prime}$ on $J$ and $b_{1}$ on $\partial J$.

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Corollary 41.40. As above let $(\alpha, \beta): \partial J \rightarrow \mathbb{R}^{2} \backslash\{0\}$ and let

$$
\begin{aligned}
D(L) & =\left\{u \in H^{2}(J): B u=\alpha u^{\prime}+\beta u=0 \text { on } \partial J\right\} \text { and } \\
L u & =-\rho^{-1}\left(a \rho u^{\prime}\right)^{\prime}+b u^{\prime}+c u .
\end{aligned}
$$

Given $\lambda_{0}>0$ sufficiently large, $\left(L+\lambda_{0}\right): D(L) \rightarrow L^{2}(J)$ and $\left(L^{*}+\lambda_{0}\right):$ $D\left(L^{*}\right) \rightarrow L^{2}(J)$ are invertible and there is a continuous Green's function $G(x, y)$ such that

$$
\left(L+\lambda_{0}\right)^{-1} f(x)=\int_{J} G(x, y) f(y) d y
$$

Proof. Let us normalize $\alpha$ so that $\alpha=a$ whenever $\alpha \neq 0$. The boundary term in Eq. (41.45) will be zero whenever

$$
a u^{\prime}+b_{1} u=0 \text { when } v \neq 0 \text { on } \partial J
$$

This suggests that we define a subspace $\chi$ of $H^{1}(J)$ by

$$
\chi:=\left\{u \in H^{1}(J): u=0 \text { on } \partial J \text { where } \alpha=0 \text { on } \partial J\right\} .
$$

Hence $\chi$ is either $H_{0}^{1}(J), H^{1}(J),\left\{u \in H^{1}(J): u(r)=0\right\}$ or $\left\{u \in H^{1}(J): u(s)=0\right.$ Now choose a function $b_{1} \in C^{2}(J)$ such that $b_{1}=\beta$ on $\partial J$, then set $b_{2}:=b+b_{1}$ and $c_{0}=c+\rho^{-1}\left(\rho b_{1}\right)^{\prime}$, then

$$
D(L)=\chi \cap\left\{u \in H^{2}(J): B u=a u^{\prime}+b_{1} u=0 \text { on } \partial J\right\}
$$

and

$$
(L u, v)=\mathcal{E}(u, v) \text { for all } u \in D(L) \text { and } v \in \chi
$$

Using this observation, it follows from Eq. (41.43) of Lemma 41.37, for $\lambda_{0}$ sufficiently large and any $u \in D(L)$, that

$$
\begin{aligned}
\left(\left(L+\lambda_{0}\right) u, u\right) & =\mathcal{E}(u, u)+\lambda_{0}(u, u) \\
& \geq \frac{a_{0}}{2}\left\|u^{\prime}\right\|^{2}+\left(\bar{c}-\frac{B^{2}}{2 a_{0}}+\lambda_{0}\right)\|u\|^{2} \geq \frac{a_{0}}{2}\|u\|_{H^{1}(J)}^{2}
\end{aligned}
$$

As usual this equation shows $\operatorname{Nul}\left(L+\lambda_{0}\right)=\{0\}$. Similarly on shows

$$
\left(u, L^{*} v\right)=\mathcal{E}(u, v) \text { for all } v \in D\left(L^{*}\right) \text { and } u \in \chi
$$

and working as above we conclude that $\operatorname{Nul}\left(L^{*}+\lambda_{0}\right)=\{0\}$. The remaining assertions are now proved as in the proof of Corollary 41.24.

With this result in hand we may now improve on Proposition 41.35.
Theorem 41.41. Let $L, B, D(L), L^{*}, B^{*}$ and $D\left(L^{*}\right)$ be as in Proposition 41.35 and $v \in L^{2}(J)$. Then there exists $g \in L^{2}(J)$ such that

$$
\begin{equation*}
(L u, v)=(u, g) \text { for all } u \in D(L) \tag{41.46}
\end{equation*}
$$

iff $v \in D\left(L^{*}\right)$ and in which case $g=L^{*} v$.

Proof. Choose $\lambda_{0}>0$ so that $L_{0}:=\left(L+\lambda_{0} I\right): D(L) \rightarrow L^{2}(J)$ is invertible. Then Eq. (41.46) is equivalent to

$$
\begin{equation*}
\left(L_{0} u, v\right)=\left(u, g+\lambda_{0} v\right) \text { for all } u \in D(L) . \tag{41.47}
\end{equation*}
$$

Taking $u=L_{0}^{-1} w$ with $w \in L^{2}(J)$ in this equation implies

$$
(w, v)=\left(L_{0}^{-1} w, g+\lambda_{0} v\right)=\left(w,\left(L_{0}^{-1}\right)^{*}\left(g+\lambda_{0} v\right)\right) \text { for all } w \in L^{2}(J)
$$

which shows

$$
\begin{equation*}
v=\left(L_{0}^{-1}\right)^{*}\left(g+\lambda_{0} v\right) . \tag{41.48}
\end{equation*}
$$

Since $\left(L_{0} u, v\right)=\left(u, L_{0}^{*} v\right)$ for all $u \in D(L)$ and $v \in D\left(L^{*}\right)$, by replacing $u$ by $L_{0}^{-1} u$ and $v$ by $\left(L_{0}^{*}\right)^{-1} v$ in this equation we learn

$$
\left(u,\left(L_{0}^{*}\right)^{-1} v\right)=\left(L_{0}^{-1} u, v\right) \text { for all } u, v \in L^{2}(J)
$$

From this equation it follows that $\left(L_{0}^{*}\right)^{-1}=\left(L_{0}^{-1}\right)^{*}$ and hence from Eq. (41.48) it follows that $v \in D\left(L_{0}^{*}\right)=D\left(L^{*}\right)$

Constant Coefficient Equations

## Convolutions, Test Functions and Partitions of Unity

### 42.1 Convolution and Young's Inequalities

Letting $\delta_{x}$ denote the "delta-function" at $x$, we wish to define a product $(*)$ on functions on $\mathbb{R}^{n}$ such that $\delta_{x} * \delta_{y}=\delta_{x+y}$. Now formally any function $f$ on $\mathbb{R}^{n}$ is of the form

$$
f=\int_{\mathbb{R}^{n}} f(x) \delta_{x} d x
$$

so we should have

$$
\begin{aligned}
f * g & =\int_{\mathbb{R}^{n} \times \mathbb{R}^{n}} f(x) g(y) \delta_{x} * \delta_{y} d x d y=\int_{\mathbb{R}^{n} \times \mathbb{R}^{n}} f(x) g(y) \delta_{x+y} d x d y \\
& =\int_{\mathbb{R}^{n} \times \mathbb{R}^{n}} f(x-y) g(y) \delta_{x} d x d y \\
& =\int_{\mathbb{R}^{n}}\left[\int_{\mathbb{R}^{n}} f(x-y) g(y) d y\right] \delta_{x} d x
\end{aligned}
$$

which suggests we make the following definition.
Definition 42.1. Let $f, g: \mathbb{R}^{n} \rightarrow \mathbb{C}$ be measurable functions. We define

$$
f * g(x)=\int_{\mathbb{R}^{n}} f(x-y) g(y) d y
$$

whenever the integral is defined, i.e. either $f(x-\cdot) g(\cdot) \in L^{1}\left(\mathbb{R}^{n}, m\right)$ or $f(x-$ $\cdot) g(\cdot) \geq 0$. Notice that the condition that $f(x-\cdot) g(\cdot) \in L^{1}\left(\mathbb{R}^{n}, m\right)$ is equivalent to writing $|f| *|g|(x)<\infty$.

Notation 42.2 Given a multi-index $\alpha \in \mathbb{Z}_{+}^{n}$, let $|\alpha|=\alpha_{1}+\cdots+\alpha_{n}$,

$$
x^{\alpha}:=\prod_{j=1}^{n} x_{j}^{\alpha_{j}}, \text { and } \partial_{x}^{\alpha}=\left(\frac{\partial}{\partial x}\right)^{\alpha}:=\prod_{j=1}^{n}\left(\frac{\partial}{\partial x_{j}}\right)^{\alpha_{j}} .
$$

Remark 42.3 (The Significance of Convolution). Suppose that $L=\sum_{|\alpha| \leq k} a_{\alpha} \partial^{\alpha}$ is a constant coefficient differential operator and suppose that we can solve (uniquely) the equation $L u=g$ in the form

$$
u(x)=K g(x):=\int_{\mathbb{R}^{n}} k(x, y) g(y) d y
$$

where $k(x, y)$ is an "integral kernel." (This is a natural sort of assumption since, in view of the fundamental theorem of calculus, integration is the inverse operation to differentiation.) Since $\tau_{z} L=L \tau_{z}$ for all $z \in \mathbb{R}^{n}$, (this is another way to characterize constant coefficient differential operators) and $L^{-1}=K$ we should have $\tau_{z} K=K \tau_{z}$. Writing out this equation then says

$$
\begin{aligned}
\int_{\mathbb{R}^{n}} k(x-z, y) g(y) d y & =(K g)(x-z)=\tau_{z} K g(x)=\left(K \tau_{z} g\right)(x) \\
& =\int_{\mathbb{R}^{n}} k(x, y) g(y-z) d y=\int_{\mathbb{R}^{n}} k(x, y+z) g(y) d y
\end{aligned}
$$

Since $g$ is arbitrary we conclude that $k(x-z, y)=k(x, y+z)$. Taking $y=0$ then gives

$$
k(x, z)=k(x-z, 0)=: \rho(x-z) .
$$

We thus find that $K g=\rho * g$. Hence we expect the convolution operation to appear naturally when solving constant coefficient partial differential equations. More about this point later.

The following proposition is an easy consequence of Minkowski's inequality for integrals.

Proposition 42.4. Suppose $q \in[1, \infty], f \in L^{1}$ and $g \in L^{q}$, then $f * g(x)$ exists for almost every $x, f * g \in L^{q}$ and

$$
\|f * g\|_{p} \leq\|f\|_{1}\|g\|_{p} .
$$

For $z \in \mathbb{R}^{n}$ and $f: \mathbb{R}^{n} \rightarrow \mathbb{C}$, let $\tau_{z} f: \mathbb{R}^{n} \rightarrow \mathbb{C}$ be defined by $\tau_{z} f(x)=$ $f(x-z)$.

Proposition 42.5. Suppose that $p \in[1, \infty)$, then $\tau_{z}: L^{p} \rightarrow L^{p}$ is an isometric isomorphism and for $f \in L^{p}, z \in \mathbb{R}^{n} \rightarrow \tau_{z} f \in L^{p}$ is continuous.

Proof. The assertion that $\tau_{z}: L^{p} \rightarrow L^{p}$ is an isometric isomorphism follows from translation invariance of Lebesgue measure and the fact that $\tau_{-z} \circ \tau_{z}=i d$. For the continuity assertion, observe that

$$
\left\|\tau_{z} f-\tau_{y} f\right\|_{p}=\left\|\tau_{-y}\left(\tau_{z} f-\tau_{y} f\right)\right\|_{p}=\left\|\tau_{z-y} f-f\right\|_{p}
$$

from which it follows that it is enough to show $\tau_{z} f \rightarrow f$ in $L^{p}$ as $z \rightarrow 0 \in \mathbb{R}^{n}$.

When $f \in C_{c}\left(\mathbb{R}^{n}\right), \tau_{z} f \rightarrow f$ uniformly and since the $K:=\cup_{|z| \leq 1} \operatorname{supp}\left(\tau_{z} f\right)$ is compact, it follows by the dominated convergence theorem that $\tau_{z} f \rightarrow f$ in $L^{p}$ as $z \rightarrow 0 \in \mathbb{R}^{n}$. For general $g \in L^{p}$ and $f \in C_{c}\left(\mathbb{R}^{n}\right)$,

$$
\begin{aligned}
\left\|\tau_{z} g-g\right\|_{p} & \leq\left\|\tau_{z} g-\tau_{z} f\right\|_{p}+\left\|\tau_{z} f-f\right\|_{p}+\|f-g\|_{p} \\
& =\left\|\tau_{z} f-f\right\|_{p}+2\|f-g\|_{p}
\end{aligned}
$$

and thus

$$
\lim \sup _{z \rightarrow 0}\left\|\tau_{z} g-g\right\|_{p} \leq \lim \sup _{z \rightarrow 0}\left\|\tau_{z} f-f\right\|_{p}+2\|f-g\|_{p}=2\|f-g\|_{p}
$$

Because $C_{c}\left(\mathbb{R}^{n}\right)$ is dense in $L^{p}$, the term $\|f-g\|_{p}$ may be made as small as we please.

Definition 42.6. Suppose that $(X, \tau)$ is a topological space and $\mu$ is a measure on $\mathcal{B}_{X}=\sigma(\tau)$. For a measurable function $f: X \rightarrow \mathbb{C}$ we define the essential support of $f$ by

$$
\begin{equation*}
\operatorname{supp}_{\mu}(f)=\{x \in U: \mu(\{y \in V: f(y) \neq 0\})>0 \forall \tau \ni V \ni x\} \tag{42.1}
\end{equation*}
$$

Lemma 42.7. Suppose $(X, \tau)$ is second countable and $f: X \rightarrow \mathbb{C}$ is a measurable function and $\mu$ is a measure on $\mathcal{B}_{X}$. Then $X:=U \backslash \operatorname{supp}_{\mu}(f)$ may be described as the largest open set $W$ such that $f 1_{W}(x)=0$ for $\mu$ - a.e. $x$. Equivalently put, $C:=\operatorname{supp}_{\mu}(f)$ is the smallest closed subset of $X$ such that $f=f 1_{C}$ a.e.

Proof. To verify that the two descriptions of $\operatorname{supp}_{\mu}(f)$ are equivalent, suppose $\operatorname{supp}_{\mu}(f)$ is defined as in Eq. (42.1) and $W:=X \backslash \operatorname{supp}_{\mu}(f)$. Then

$$
\begin{aligned}
W & =\left\{\begin{array}{c}
x \in X: \mu(\{y \in V: f(y) \neq 0\})=0 \\
\text { for some neighborhood } V \text { of } x
\end{array}\right\} \\
& =\cup\left\{V \subset_{o} X: \mu\left(f 1_{V} \neq 0\right)=0\right\} \\
& =\cup\left\{V \subset_{o} X: f 1_{V}=0 \text { for } \mu \text { - a.e. }\right\} .
\end{aligned}
$$

So to finish the argument it suffices to show $\mu\left(f 1_{W} \neq 0\right)=0$. To to this let $\mathcal{U}$ be a countable base for $\tau$ and set

$$
\mathcal{U}_{f}:=\left\{V \in \mathcal{U}: f 1_{V}=0 \text { a.e. }\right\} .
$$

Then it is easily seen that $W=\cup \mathcal{\mathcal { U } _ { f }}$ and since $\mathcal{U}_{f}$ is countable $\mu\left(f 1_{W} \neq 0\right) \leq$ $\sum_{V \in \mathcal{U}_{f}} \mu\left(f 1_{V} \neq 0\right)=0$.
Lemma 42.8. Suppose $f, g, h: \mathbb{R}^{n} \rightarrow \mathbb{C}$ are measurable functions and assume that $x$ is a point in $\mathbb{R}^{n}$ such that $|f| *|g|(x)<\infty$ and $|f| *(|g| *|h|)(x)<\infty$, then

$$
\text { 1. } f * g(x)=g * f(x)
$$

2. $f *(g * h)(x)=(f * g) * h(x)$
3. If $z \in \mathbb{R}^{n}$ and $\tau_{z}(|f| *|g|)(x)=|f| *|g|(x-z)<\infty$, then

$$
\tau_{z}(f * g)(x)=\tau_{z} f * g(x)=f * \tau_{z} g(x)
$$

4. If $x \notin \operatorname{supp}_{m}(f)+\operatorname{supp}_{m}(g)$ then $f * g(x)=0$ and in particular, $\operatorname{supp}_{m}(f *$ $g) \subset \operatorname{supp}_{m}(f)+\operatorname{supp}_{m}(g)$ where in defining $\operatorname{supp}_{m}(f * g)$ we will use the convention that " $f * g(x) \neq 0$ " when $|f| *|g|(x)=\infty$.

Proof. For item 1.,
$|f| *|g|(x)=\int_{\mathbb{R}^{n}}|f|(x-y)|g|(y) d y=\int_{\mathbb{R}^{n}}|f|(y)|g|(y-x) d y=|g| *|f|(x)$
where in the second equality we made use of the fact that Lebesgue measure invariant under the transformation $y \rightarrow x-y$. Similar computations prove all of the remaining assertions of the first three items of the lemma.

Item 4. Since $f * g(x)=\tilde{f} * \tilde{g}(x)$ if $f=\tilde{f}$ and $g=\tilde{g}$ a.e. we may, by replacing $f$ by $f 1_{\text {supp }_{m}(f)}$ and $g$ by $g 1_{\text {supp }_{m}(g)}$ if necessary, assume that $\{f \neq 0\} \subset$ $\operatorname{supp}_{m}(f)$ and $\{g \neq 0\} \subset \operatorname{supp}_{m}(g)$. So if $x \notin\left(\operatorname{supp}_{m}(f)+\operatorname{supp}_{m}(g)\right)$ then $x \notin(\{f \neq 0\}+\{g \neq 0\})$ and for all $y \in \mathbb{R}^{n}$, either $x-y \notin\{f \neq 0\}$ or $y \notin$ $\{g \neq 0\}$. That is to say either $x-y \in\{f=0\}$ or $y \in\{g=0\}$ and hence $f(x-y) g(y)=0$ for all $y$ and therefore $f * g(x)=0$. This shows that $f * g=0$ on $\mathbb{R}^{n} \backslash\left(\overline{\operatorname{supp}_{m}(f)+\operatorname{supp}_{m}(g)}\right)$ and therefore

$$
\mathbb{R}^{n} \backslash\left(\overline{\operatorname{supp}_{m}(f)+\operatorname{supp}_{m}(g)}\right) \subset \mathbb{R}^{n} \backslash \operatorname{supp}_{m}(f * g)
$$

i.e. $\operatorname{supp}_{m}(f * g) \subset \operatorname{supp}_{m}(f)+\operatorname{supp}_{m}(g)$.

Remark 42.9. Let $A, B$ be closed sets of $\mathbb{R}^{n}$, it is not necessarily true that $A+B$ is still closed. For example, take

$$
A=\{(x, y): x>0 \text { and } y \geq 1 / x\} \text { and } B=\{(x, y): x<0 \text { and } y \geq 1 /|x|\}
$$

then every point of $A+B$ has a positive $y$-component and hence is not zero. On the other hand, for $x>0$ we have $(x, 1 / x)+(-x, 1 / x)=(0,2 / x) \in A+B$ for all $x$ and hence $0 \in \overline{A+B}$ showing $A+B$ is not closed. Nevertheless if one of the sets $A$ or $B$ is compact, then $A+B$ is closed again. Indeed, if $A$ is compact and $x_{n}=a_{n}+b_{n} \in A+B$ and $x_{n} \rightarrow x \in \mathbb{R}^{n}$, then by passing to a subsequence if necessary we may assume $\lim _{n \rightarrow \infty} a_{n}=a \in A$ exists. In this case

$$
\lim _{n \rightarrow \infty} b_{n}=\lim _{n \rightarrow \infty}\left(x_{n}-a_{n}\right)=x-a \in B
$$

exists as well, showing $x=a+b \in A+B$.
Proposition 42.10. Suppose that $p, q \in[1, \infty]$ and $p$ and $q$ are conjugate exponents, $f \in L^{p}$ and $g \in L^{q}$, then $f * g \in B C\left(\mathbb{R}^{n}\right),\|f * g\|_{u} \leq\|f\|_{p}\|g\|_{q}$ and if $p, q \in(1, \infty)$ then $f * g \in C_{0}\left(\mathbb{R}^{n}\right)$.

Proof. The existence of $f * g(x)$ and the estimate $|f * g|(x) \leq\|f\|_{p}\|g\|_{q}$ for all $x \in \mathbb{R}^{n}$ is a simple consequence of Holders inequality and the translation invariance of Lebesgue measure. In particular this shows $\|f * g\|_{u} \leq\|f\|_{p}\|g\|_{q}$. By relabeling $p$ and $q$ if necessary we may assume that $p \in[1, \infty)$. Since

$$
\begin{aligned}
\left\|\tau_{z}(f * g)-f * g\right\|_{u} & =\left\|\tau_{z} f * g-f * g\right\|_{u} \\
& \leq\left\|\tau_{z} f-f\right\|_{p}\|g\|_{q} \rightarrow 0 \text { as } z \rightarrow 0
\end{aligned}
$$

it follows that $f * g$ is uniformly continuous. Finally if $p, q \in(1, \infty)$, we learn from Lemma 42.8 and what we have just proved that $f_{m} * g_{m} \in C_{c}\left(\mathbb{R}^{n}\right)$ where $f_{m}=f 1_{|f| \leq m}$ and $g_{m}=g 1_{|g| \leq m}$. Moreover,

$$
\begin{aligned}
\left\|f * g-f_{m} * g_{m}\right\|_{u} & \leq\left\|f * g-f_{m} * g\right\|_{u}+\left\|f_{m} * g-f_{m} * g_{m}\right\|_{u} \\
& \leq\left\|f-f_{m}\right\|_{p}\|g\|_{q}+\left\|f_{m}\right\|_{p}\left\|g-g_{m}\right\|_{q} \\
& \leq\left\|f-f_{m}\right\|_{p}\|g\|_{q}+\|f\|_{p}\left\|g-g_{m}\right\|_{q} \rightarrow 0 \text { as } m \rightarrow \infty
\end{aligned}
$$

showing $f * g \in C_{0}\left(\mathbb{R}^{n}\right)$.
Theorem 42.11 (Young's Inequality). Let $p, q, r \in[1, \infty]$ satisfy

$$
\begin{equation*}
\frac{1}{p}+\frac{1}{q}=1+\frac{1}{r} \tag{42.2}
\end{equation*}
$$

If $f \in L^{p}$ and $g \in L^{q}$ then $|f| *|g|(x)<\infty$ for $m$ - a.e. $x$ and

$$
\begin{equation*}
\|f * g\|_{r} \leq\|f\|_{p}\|g\|_{q} . \tag{42.3}
\end{equation*}
$$

In particular $L^{1}$ is closed under convolution. (The space $\left(L^{1}, *\right)$ is an example of a "Banach algebra" without unit.)

Remark 42.12. Before going to the formal proof, let us first understand Eq. (42.2) by the following scaling argument. For $\lambda>0$, let $f_{\lambda}(x):=f(\lambda x)$, then after a few simple change of variables we find

$$
\left\|f_{\lambda}\right\|_{p}=\lambda^{-1 / p}\|f\| \text { and }(f * g)_{\lambda}=\lambda f_{\lambda} * g_{\lambda} .
$$

Therefore if Eq. (42.3) holds for some $p, q, r \in[1, \infty]$, we would also have

$$
\begin{aligned}
\|f * g\|_{r} & =\lambda^{1 / r}\left\|(f * g)_{\lambda}\right\|_{r} \leq \lambda^{1 / r} \lambda\left\|f_{\lambda}\right\|_{p}\left\|g_{\lambda}\right\|_{q} \\
& =\lambda^{(1+1 / r-1 / p-1 / q)}\|f\|_{p}\|g\|_{q}
\end{aligned}
$$

for all $\lambda>0$. This is only possible if Eq. (42.2) holds.
Proof. Let $\alpha, \beta \in[0,1]$ and $p_{1}, p_{2} \in[0, \infty]$ satisfy $p_{1}^{-1}+p_{2}^{-1}+r^{-1}=1$. Then by Hölder's inequality,

$$
\begin{aligned}
|f * g(x)| & =\left|\int f(x-y) g(y) d y\right| \\
& \leq \int|f(x-y)|^{(1-\alpha)}|g(y)|^{(1-\beta)}|f(x-y)|^{\alpha}|g(y)|^{\beta} d y \\
& \leq\left(\int|f(x-y)|^{(1-\alpha) r}|g(y)|^{(1-\beta) r} d y\right)^{1 / r}\left(\int|f(x-y)|^{\alpha p_{1}} d y\right)^{1 / p_{1}} \times \\
& \left(\int|g(y)|^{\beta p_{2}} d y\right)^{1 / p_{2}} \\
& =\left(\int|f(x-y)|^{(1-\alpha) r}|g(y)|^{(1-\beta) r} d y\right)^{1 / r}\|f\|_{\alpha p_{1}}^{\alpha}\|g\|_{\beta p_{2}}^{\beta}
\end{aligned}
$$

Taking the $r^{\text {th }}$ power of this equation and integrating on $x$ gives

$$
\begin{align*}
\|f * g\|_{r}^{r} & \leq \int\left(\int|f(x-y)|^{(1-\alpha) r}|g(y)|^{(1-\beta) r} d y\right) d x \cdot\|f\|_{\alpha p_{1}}^{\alpha}\|g\|_{\beta p_{2}}^{\beta} \\
& =\|f\|_{(1-\alpha) r}^{(1-\alpha) r}\|g\|_{(1-\beta) r}^{(1-\beta) r}\|f\|_{\alpha p_{1}}^{\alpha r}\|g\|_{\beta p_{2}}^{\beta r} . \tag{42.4}
\end{align*}
$$

Let us now suppose, $(1-\alpha) r=\alpha p_{1}$ and $(1-\beta) r=\beta p_{2}$, in which case Eq. (42.4) becomes,

$$
\|f * g\|_{r}^{r} \leq\|f\|_{\alpha p_{1}}^{r}\|g\|_{\beta p_{2}}^{r}
$$

which is Eq. (42.3) with

$$
\begin{equation*}
p:=(1-\alpha) r=\alpha p_{1} \text { and } q:=(1-\beta) r=\beta p_{2} . \tag{42.5}
\end{equation*}
$$

So to finish the proof, it suffices to show $p$ and $q$ are arbitrary indices in $[1, \infty]$ satisfying $p^{-1}+q^{-1}=1+r^{-1}$.

If $\alpha, \beta, p_{1}, p_{2}$ satisfy the relations above, then

$$
\alpha=\frac{r}{r+p_{1}} \text { and } \beta=\frac{r}{r+p_{2}}
$$

and

$$
\frac{1}{p}+\frac{1}{q}=\frac{1}{p_{1}} \frac{r+p_{1}}{r}+\frac{1}{p_{2}} \frac{r+p_{2}}{r}=\frac{1}{p_{1}}+\frac{1}{p_{2}}+\frac{2}{r}=1+\frac{1}{r} .
$$

Conversely, if $p, q, r$ satisfy Eq. (42.2), then let $\alpha$ and $\beta$ satisfy $p=(1-\alpha) r$ and $q=(1-\beta) r$, i.e.

$$
\alpha:=\frac{r-p}{r}=1-\frac{p}{r} \leq 1 \text { and } \beta=\frac{r-q}{r}=1-\frac{q}{r} \leq 1 .
$$

From Eq. (42.2), $\alpha=p\left(1-\frac{1}{q}\right) \geq 0$ and $\beta=q\left(1-\frac{1}{p}\right) \geq 0$, so that $\alpha, \beta \in[0,1]$. We then define $p_{1}:=p / \alpha$ and $p_{2}:=q / \beta$, then

$$
\frac{1}{p_{1}}+\frac{1}{p_{2}}+\frac{1}{r}=\beta \frac{1}{q}+\alpha \frac{1}{p}+\frac{1}{r}=\frac{1}{q}-\frac{1}{r}+\frac{1}{p}-\frac{1}{r}+\frac{1}{r}=1
$$

as desired.

Theorem 42.13 (Approximate $\delta$ - functions). Let $p \in[1, \infty], \phi \in$ $L^{1}\left(\mathbb{R}^{n}\right), a:=\int_{\mathbb{R}^{n}} f(x) d x$, and for $t>0$ let $\phi_{t}(x)=t^{-n} \phi(x / t)$. Then

1. If $f \in L^{p}$ with $p<\infty$ then $\phi_{t} * f \rightarrow a f$ in $L^{p}$ as $t \downarrow 0$.
2. If $f \in B C\left(\mathbb{R}^{n}\right)$ and $f$ is uniformly continuous then $\left\|\phi_{t} * f-f\right\|_{\infty} \rightarrow 0$ as $t \downarrow 0$.
3. If $f \in L^{\infty}$ and $f$ is continuous on $U \subset_{o} \mathbb{R}^{n}$ then $\phi_{t} * f \rightarrow a f$ uniformly on compact subsets of $U$ as $t \downarrow 0$.

See Theorem 8.15 if Folland for a statement about almost everywhere convergence.

Proof. Making the change of variables $y=t z$ implies

$$
\phi_{t} * f(x)=\int_{\mathbb{R}^{n}} f(x-y) \phi_{t}(y) d y=\int_{\mathbb{R}^{n}} f(x-t z) \phi(z) d z
$$

so that

$$
\begin{align*}
\phi_{t} * f(x)-a f(x) & =\int_{\mathbb{R}^{n}}[f(x-t z)-f(x)] \phi(z) d z \\
& =\int_{\mathbb{R}^{n}}\left[\tau_{t z} f(x)-f(x)\right] \phi(z) d z \tag{42.6}
\end{align*}
$$

Hence by Minkowski's inequality for integrals, Proposition 42.5 and the dominated convergence theorem,

$$
\left\|\phi_{t} * f-a f\right\|_{p} \leq \int_{\mathbb{R}^{n}}\left\|\tau_{t z} f-f\right\|_{p}|\phi(z)| d z \rightarrow 0 \text { as } t \downarrow 0
$$

Item 2. is proved similarly. Indeed, form Eq. (42.6)

$$
\left\|\phi_{t} * f-a f\right\|_{\infty} \leq \int_{\mathbb{R}^{n}}\left\|\tau_{t z} f-f\right\|_{\infty}|\phi(z)| d z
$$

which again tends to zero by the dominated convergence theorem because $\lim _{t \downarrow 0}\left\|\tau_{t z} f-f\right\|_{\infty}=0$ uniformly in $z$ by the uniform continuity of $f$.

Item 3. Let $B_{R}=B(0, R)$ be a large ball in $\mathbb{R}^{n}$ and $K \sqsubset ᄃ U$, then

$$
\begin{aligned}
\sup _{x \in K}\left|\phi_{t} * f(x)-a f(x)\right| \leq & \left|\int_{B_{R}}[f(x-t z)-f(x)] \phi(z) d z\right| \\
& +\left|\int_{B_{R}^{c}}[f(x-t z)-f(x)] \phi(z) d z\right| \\
\leq & \int_{B_{R}}|\phi(z)| d z \cdot \sup _{x \in K, z \in B_{R}}|f(x-t z)-f(x)| \\
& +2\|f\|_{\infty} \int_{B_{R}^{c}}|\phi(z)| d z \\
\leq & \|\phi\|_{1} \cdot \sup _{x \in K, z \in B_{R}}|f(x-t z)-f(x)| \\
& +2\|f\|_{\infty} \int_{|z|>R}|\phi(z)| d z
\end{aligned}
$$

so that using the uniform continuity of $f$ on compact subsets of $U$,
$\lim \sup _{t \downarrow 0} \sup _{x \in K}\left|\phi_{t} * f(x)-a f(x)\right| \leq 2\|f\|_{\infty} \int_{|z|>R}|\phi(z)| d z \rightarrow 0$ as $R \rightarrow \infty$.

Remark 42.14 (Another Proof of part of Theorem 42.13). By definition of the convolution and Hölder's or Jensen's inequality we have

$$
\begin{aligned}
\int_{\mathbb{R}^{n}}\left|v * \phi_{t}(x)\right|^{p} d x & \leq \int_{\mathbb{R}^{n}}\left(\int_{\mathbb{R}^{n}}(v(x-y))\left|\phi_{t}(y)\right| d y\right)^{p} d x \\
& \leq \int_{\mathbb{R}^{n} \times \mathbb{R}^{n}}|v(x-y)|^{p} \phi_{t}(y) d y d x=\|v\|_{L^{p}}^{p}
\end{aligned}
$$

Therefore $\left\|v * \phi_{t}\right\|_{L^{p}} \leq\|v\|_{L^{p}}$ which implies $v * \phi_{t} \in L^{p}$. If $\phi_{t} \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$, by differentiating under the integral (see Theorem 42.18 below) it is easily seen that $v * \phi_{t} \in C^{\infty}$. Finally for $u \in C_{c}\left(\mathbb{R}^{n}\right)$,

$$
\begin{aligned}
\left\|v-v * \phi_{t}\right\|_{L^{p}} & \leq\|v-u\|_{L^{p}}+\left\|u-u * \phi_{t}\right\|_{L^{p}}+\left\|u * \phi_{t}-v * \phi_{t}\right\|_{L^{p}} \\
& \leq\left\|u-u * \phi_{t}\right\|_{L^{p}}+2\|v-u\|_{L^{p}}
\end{aligned}
$$

and hence

$$
\limsup _{t \downarrow 0}\left\|v-v * \phi_{t}\right\|_{L^{p}} \leq 2\|v-u\|_{L^{p}}
$$

which may be made arbitrarily small since $C_{c}\left(\mathbb{R}^{n}\right)$ is dense in $L^{p}\left(\mathbb{R}^{n}, m\right)$.
Exercise 42.15. Let

$$
f(t)=\left\{\begin{array}{cc}
e^{-1 / t} & \text { if } t>0 \\
0 & \text { if } t \leq 0
\end{array}\right.
$$

Show $f \in C^{\infty}(\mathbb{R},[0,1])$.

Lemma 42.16. There exists $\phi \in C_{c}^{\infty}\left(\mathbb{R}^{n},[0, \infty)\right)$ such that $\phi(0)>0$, $\operatorname{supp}(\phi) \subset \bar{B}(0,1)$ and $\int_{\mathbb{R}^{n}} \phi(x) d x=1$.

Proof. Define $h(t)=f(1-t) f(t+1)$ where $f$ is as in Exercise 42.15. Then $h \in C_{c}^{\infty}(\mathbb{R},[0,1]), \operatorname{supp}(h) \subset[-1,1]$ and $h(0)=e^{-2}>0$. Define $c=$ $\int_{\mathbb{R}^{n}} h\left(|x|^{2}\right) d x$. Then $\phi(x)=c^{-1} h\left(|x|^{2}\right)$ is the desired function.

Definition 42.17. Let $X \subset \mathbb{R}^{n}$ be an open set. A Radon measure on $\mathcal{B}_{X}$ is a measure $\mu$ which is finite on compact subsets of $X$. For a Radon measure $\mu$, we let $L_{l o c}^{1}(\mu)$ consists of those measurable functions $f: X \rightarrow \mathbb{C}$ such that $\int_{K}|f| d \mu<\infty$ for all compact subsets $K \subset X$.

Theorem 42.18 (Differentiation under integral sign). Let $\Omega \subset \mathbb{R}^{n}$ and $f: \mathbb{R}^{m} \times \Omega \rightarrow \mathbb{R}$ be given. Assume:

1. $x \rightarrow f(x, y)$ is differentiable for all $y \in \Omega$.
2. $\left|\frac{\partial f}{\partial x^{i}}(x, y)\right| \leq g(y)$ for some $g$ such that $\int_{\Omega}|g(y)| d y<\infty$.
3. $\int|f(x, y)| d y<\infty$.

Then $\frac{\partial}{\partial x^{i}} \int_{\Omega} f(x, y) d y=\int_{\Omega} \frac{\partial f}{\partial x^{i}}(x, y) d y$ and moreover if $x \rightarrow \frac{\partial f}{\partial x^{i}}(x, y)$ is continuous then so is $x \rightarrow \int \frac{\partial f}{\partial x^{i}}(x, y) d y$.

The reader asked to use Theorem 42.18 to verify the following proposition.
Proposition 42.19. Suppose that $f \in L_{\text {loc }}^{1}\left(\mathbb{R}^{n}, m\right)$ and $\phi \in C_{c}^{1}\left(\mathbb{R}^{n}\right)$, then $f * \phi \in C^{1}\left(\mathbb{R}^{n}\right)$ and $\partial_{i}(f * \phi)=f * \partial_{i} \phi$. Moreover if $\phi \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ then $f * \phi \in C^{\infty}\left(\mathbb{R}^{n}\right)$.

Corollary 42.20 ( $C^{\infty}$ - Uryhson's Lemma). Given $K \sqsubset \sqsubset U \subset \subset_{o} \mathbb{R}^{n}$, there exists $f \in C_{c}^{\infty}\left(\mathbb{R}^{n},[0,1]\right)$ such that $\operatorname{supp}(f) \subset U$ and $f=1$ on $K$.

Proof. Let $\phi$ be as in Lemma 42.16, $\phi_{t}(x)=t^{-n} \phi(x / t)$ be as in Theorem $42.13, d$ be the standard metric on $\mathbb{R}^{n}$ and $\epsilon=d\left(K, U^{c}\right)$. Since $K$ is compact and $U^{c}$ is closed, $\epsilon>0$. Let $V_{\delta}=\left\{x \in \mathbb{R}^{n}: d(x, K)<\delta\right\}$ and $f=\phi_{\epsilon / 3} * 1_{V_{\epsilon / 3}}$, then

$$
\operatorname{supp}(f) \subset \overline{\operatorname{supp}\left(\phi_{\epsilon / 3}\right)+V_{\epsilon / 3}} \subset \bar{V}_{2 \epsilon / 3} \subset U
$$

Since $\bar{V}_{2 \epsilon / 3}$ is closed and bounded, $f \in C_{c}^{\infty}(U)$ and for $x \in K$,

$$
f(x)=\int_{\mathbb{R}^{n}} 1_{d(y, K)<\epsilon / 3} \cdot \phi_{\epsilon / 3}(x-y) d y=\int_{\mathbb{R}^{n}} \phi_{\epsilon / 3}(x-y) d y=1 .
$$

The proof will be finished after the reader (easily) verifies $0 \leq f \leq 1$.
Here is an application of this corollary whose proof is left to the reader.

Lemma 42.21 (Integration by Parts). Suppose $f$ and $g$ are measurable functions on $\mathbb{R}^{n}$ such that $t \rightarrow f\left(x_{1}, \ldots, x_{i-1}, t, x_{i+1}, \ldots, x_{n}\right)$ and $t \rightarrow$ $g\left(x_{1}, \ldots, x_{i-1}, t, x_{i+1}, \ldots, x_{n}\right)$ are continuously differentiable functions on $\mathbb{R}$ for each fixed $x=\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{n}$. Moreover assume $f \cdot g, \frac{\partial f}{\partial x_{i}} \cdot g$ and $f \cdot \frac{\partial g}{\partial x_{i}}$ are in $L^{1}\left(\mathbb{R}^{n}, m\right)$. Then

$$
\int_{\mathbb{R}^{n}} \frac{\partial f}{\partial x_{i}} \cdot g d m=-\int_{\mathbb{R}^{n}} f \cdot \frac{\partial g}{\partial x_{i}} d m
$$

Exercise 42.22 (Integration by Parts). Suppose that $(x, y) \in \mathbb{R} \times \mathbb{R}^{n-1} \rightarrow$ $f(x, y) \in \mathbb{C}$ and $(x, y) \in \mathbb{R} \times \mathbb{R}^{n-1} \rightarrow g(x, y) \in \mathbb{C}$ are measurable functions such that for each fixed $y \in \mathbb{R}^{n-1}, x \rightarrow f(x, y)$ and $x \rightarrow g(x, y)$ are continuously differentiable. Also assume $f \cdot g, \partial_{x} f \cdot g$ and $f \cdot \partial_{x} g$ are integrable relative to Lebesgue measure on $\mathbb{R} \times \mathbb{R}^{n-1}$, where $\partial_{x} f(x, y):=\left.\frac{d}{d t} f(x+t, y)\right|_{t=0}$. Show

$$
\begin{equation*}
\int_{\mathbb{R} \times \mathbb{R}^{n-1}} \partial_{x} f(x, y) \cdot g(x, y) d x d y=-\int_{\mathbb{R} \times \mathbb{R}^{n-1}} f(x, y) \cdot \partial_{x} g(x, y) d x d y \tag{42.7}
\end{equation*}
$$

(Note: this result and Fubini's theorem proves Lemma 42.21.)
Hints: Let $\psi \in C_{c}^{\infty}(\mathbb{R})$ be a function which is 1 in a neighborhood of $0 \in \mathbb{R}$ and set $\psi_{\epsilon}(x)=\psi(\epsilon x)$. First verify Eq. (42.7) with $f(x, y)$ replaced by $\psi_{\epsilon}(x) f(x, y)$ by doing the $x$ - integral first. Then use the dominated convergence theorem to prove Eq. (42.7) by passing to the limit, $\epsilon \downarrow 0$.
Solution 42.23 (42.22). By assumption, $\partial_{x}\left[\psi_{\epsilon}(x) f(x, y)\right] \cdot g(x, y)$ and $\psi_{\epsilon}(x) f(x, y) \partial_{x} g(x, y)$ are in $L^{1}\left(\mathbb{R}^{n}\right)$, so we may use Fubini's theorem and follow the hint to learn

$$
\begin{align*}
& \int_{\mathbb{R}^{n-1}} d y \int_{\mathbb{R}} \partial_{x}\left[\psi_{\epsilon}(x) f(x, y)\right] \cdot g(x, y) d x \\
& =-\int_{\mathbb{R}^{n-1}} d y \int_{\mathbb{R}}\left[\psi_{\epsilon}(x) f(x, y)\right] \cdot \partial_{x} g(x, y) d x \tag{42.8}
\end{align*}
$$

wherein we have done and integration by parts. (There are no boundary terms because $\psi_{\epsilon}$ is compactly supported.) Now

$$
\begin{aligned}
\partial_{x}\left[\psi_{\epsilon}(x) f(x, y)\right] & =\partial_{x} \psi_{\epsilon}(x) \cdot f(x, y)+\psi_{\epsilon}(x) \partial_{x} f(x, y) \\
& =\epsilon \psi^{\prime}(\epsilon x) f(x, y)+\psi_{\epsilon}(x) \partial_{x} f(x, y)
\end{aligned}
$$

and by the dominated convergence theorem and the given assumptions we have, as $\epsilon \downarrow 0$, that

$$
\begin{aligned}
\left|\int_{\mathbb{R}^{n}} \epsilon \psi^{\prime}(\epsilon x) f(x, y) \cdot g(x, y) d x d y\right| & \leq C \epsilon \int_{\mathbb{R}^{n}}|f(x, y) \cdot g(x, y)| d x d y \rightarrow 0 \\
\int_{\mathbb{R}^{n}} \psi_{\epsilon}(x) \partial_{x} f(x, y) \cdot g(x, y) d x d y & \rightarrow \int_{\mathbb{R}^{n}} \partial_{x} f(x, y) \cdot g(x, y) d x d y \text { and } \\
\int_{\mathbb{R}^{n}} \psi_{\epsilon}(x) f(x, y) \cdot \partial_{x} g(x, y) d x d y & \rightarrow \int_{\mathbb{R}^{n}} f(x, y) \cdot \partial_{x} g(x, y) d x d y
\end{aligned}
$$

where $C=\sup _{x \in \mathbb{R}}\left|\psi^{\prime}(x)\right|$. Combining the last three equations with Eq. (42.8) shows

$$
\int_{\mathbb{R} \times \mathbb{R}^{n-1}} \partial_{x} f(x, y) \cdot g(x, y) d x d y=-\int_{\mathbb{R} \times \mathbb{R}^{n-1}} f(x, y) \cdot \partial_{x} g(x, y) d x d y
$$

as desired.
With this result we may give another proof of the Riemann Lebesgue Lemma.

Lemma 42.24. For $f \in L^{1}\left(\mathbb{R}^{n}, m\right)$ let

$$
\hat{f}(\xi):=(2 \pi)^{-n / 2} \int_{\mathbb{R}^{n}} f(x) e^{-i \xi \cdot x} d m(x)
$$

be the Fourier transform of $f$. Then $\hat{f} \in C_{0}\left(\mathbb{R}^{n}\right)$ and $\|\hat{f}\|_{u} \leq(2 \pi)^{-n / 2}\|f\|_{1}$. (The choice of the normalization factor, $(2 \pi)^{-n / 2}$, in $\hat{f}$ is for later convenience.)

Proof. The fact that $\hat{f}$ is continuous is a simple application of the dominated convergence theorem. Moreover,

$$
|\hat{f}(\xi)| \leq \int|f(x)| d m(x) \leq(2 \pi)^{-n / 2}\|f\|_{1}
$$

so it only remains to see that $\hat{f}(\xi) \rightarrow 0$ as $|\xi| \rightarrow \infty$.
First suppose that $f \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ and let $\Delta=\sum_{j=1}^{n} \frac{\partial^{2}}{\partial x_{j}^{2}}$ be the Laplacian on $\mathbb{R}^{n}$. Notice that $\frac{\partial}{\partial x_{j}} e^{-i \xi \cdot x}=-i \xi_{j} e^{-i \xi \cdot x}$ and $\Delta e^{-i \xi \cdot x}=-|\xi|^{2} e^{-i \xi \cdot x}$. Using Lemma 42.21 repeatedly,

$$
\begin{aligned}
\int \Delta^{k} f(x) e^{-i \xi \cdot x} d m(x) & =\int f(x) \Delta_{x}^{k} e^{-i \xi \cdot x} d m(x)=-|\xi|^{2 k} \int f(x) e^{-i \xi \cdot x} d m(x) \\
& =-(2 \pi)^{n / 2}|\xi|^{2 k} \hat{f}(\xi)
\end{aligned}
$$

for any $k \in \mathbb{N}$. Hence $(2 \pi)^{n / 2}|\hat{f}(\xi)| \leq|\xi|^{-2 k}\left\|\Delta^{k} f\right\|_{1} \rightarrow 0$ as $|\xi| \rightarrow \infty$ and $\hat{f} \in C_{0}\left(\mathbb{R}^{n}\right)$. Suppose that $f \in L^{1}(m)$ and $f_{k} \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ is a sequence such that $\lim _{k \rightarrow \infty}\left\|f-f_{k}\right\|_{1}=0$, then $\lim _{k \rightarrow \infty}\left\|\hat{f}-\hat{f}_{k}\right\|_{u}=0$ and hence Hence $\hat{f} \in C_{0}\left(\mathbb{R}^{n}\right)$ because $C_{0}\left(\mathbb{R}^{n}\right)$ is complete.
Corollary 42.25. Let $X \subset \mathbb{R}^{n}$ be an open set and $\mu$ be a Radon measure on $\mathcal{B}_{X}$.

1. Then $C_{c}^{\infty}(X)$ is dense in $L^{p}(\mu)$ for all $1 \leq p<\infty$.
2. If $h \in L_{\text {loc }}^{1}(\mu)$ satisfies

$$
\begin{equation*}
\int_{X} f h d \mu=0 \text { for all } f \in C_{c}^{\infty}(X) \tag{42.9}
\end{equation*}
$$

then $h(x)=0$ for $\mu-$ a.e. $x$.
Proof. Let $f \in C_{c}(X), \phi$ be as in Lemma 42.16, $\phi_{t}$ be as in Theorem 42.13 and set $\psi_{t}:=\phi_{t} *\left(f 1_{X}\right)$. Then by Proposition $42.19 \psi_{t} \in C^{\infty}(X)$ and by Lemma 42.8 there exists a compact set $K \subset X$ such that $\operatorname{supp}\left(\psi_{t}\right) \subset K$ for all $t$ sufficiently small. By Theorem $42.13, \psi_{t} \rightarrow f$ uniformly on $X$ as $t \downarrow 0$

1. The dominated convergence theorem (with dominating function being $\|f\|_{\infty} 1_{K}$ ), shows $\psi_{t} \rightarrow f$ in $L^{p}(\mu)$ as $t \downarrow 0$. This proves Item 1 . because of the measure theoretic fact that $C_{c}(X)$ is dense in $L^{p}(\mu)$.
2. Keeping the same notation as above, the dominated convergence theorem (with dominating function being $\|f\|_{\infty}|h| 1_{K}$ ) implies

$$
0=\lim _{t \downarrow 0} \int_{X} \psi_{t} h d \mu=\int_{X} \lim _{t \downarrow 0} \psi_{t} h d \mu=\int_{X} f h d \mu .
$$

Since this is true for all $f \in C_{c}(X)$, it follows by measure theoretic arguments that $h=0$ a.e.

### 42.2 Smooth Partitions of Unity

Theorem 42.26. Let $V_{1}, \ldots, V_{k} \subset_{0} \mathbb{R}^{n}$ and $\phi \in C_{c}^{\infty}\left(\cup_{i=1}^{k} V_{i}\right)$. Then there exists $\phi_{j} \in C_{c}^{\infty}\left(V_{j}\right)$ such that $\phi=\sum_{i}^{k} \phi_{j}$. If $\phi \geq 0$ one can choose $\phi_{j} \geq 0$.

Proof. The proof will be split into two steps.

1. There exists $K_{j} \sqsubset \sqsubset V_{j}$ such that $\operatorname{supp} \phi \subset \cup K_{j}$. Indeed, for all $x \in \operatorname{supp} \phi$ there exists an open neighborhood $N_{x}$ of $x$ such that $\bar{N}_{x} \subset V_{j}$ for some $j$ and $\bar{N}_{x}$ is compact. Now $\left\{N_{x}\right\}_{x \in \operatorname{supp} \phi}$ covers $K:=\operatorname{supp} \phi$ and hence there exists a finite set $\Lambda \subset \subset K$ such that $K \subset \cup_{x \in \Lambda} N_{x}$.. Let $K_{j}:=$ $\cup\left\{\bar{N}_{x}: x \in \Lambda\right.$ and $\left.\bar{N}_{x} \subset V_{j}\right\}$. Then each $K_{j}$ is compact, $K_{j} \subset V_{j}$ and $\operatorname{supp} \phi=K \subset \bigcup_{j=1}^{k} K_{j}$.
2. By Corollary 42.20 there exists $\psi_{j} \in C_{c}^{\infty}\left(V_{j},[0,1]\right)$ such that $\psi_{j}:=1$ in the neighborhood of $K_{j}$. Now define

$$
\begin{aligned}
\phi_{1} & =\phi \psi_{1} \\
\phi_{2} & =\left(\phi-\phi_{1}\right) \psi_{2}=\phi\left(1-\psi_{1}\right) \psi_{2} \\
\phi_{3} & =\left(\phi-\phi_{1}-\phi_{2}\right) \psi_{3}=\phi\left\{\left(1-\psi_{1}\right)-\left(1-\psi_{1}\right) \psi_{2}\right\} \psi_{3} \\
& =\phi\left(1-\psi_{1}\right)\left(1-\psi_{2}\right) \psi_{3} \\
& \vdots \\
\phi_{k} & =\left(\phi-\phi_{1}-\phi_{2}-\cdots-\phi_{k-1}\right) \psi_{k} \\
& =\phi\left(1-\psi_{1}\right)\left(1-\psi_{2}\right) \ldots\left(1-\psi_{k-1}\right) \psi_{k}
\end{aligned}
$$

By the above computations one finds that (a) $\phi_{i} \geq 0$ if $\phi \geq 0$ and (b)

$$
\phi-\phi_{1}-\phi_{2}-\cdots-\phi_{k}=\phi\left(1-\psi_{1}\right)\left(1-\psi_{2}\right) \ldots\left(1-\psi_{k}\right)=0 .
$$

since either $\phi(x)=0$ or $x \notin \operatorname{supp} \phi=K$ and $1-\psi_{i}(x)=0$ for some $i$.

Corollary 42.27. Let $V_{1}, \ldots, V_{k} \subset_{0} \mathbb{R}^{n}$ and $K$ be a compact subset of $\cup_{i=1}^{k} V_{i}$. Then there exists $\phi_{i} \in C_{c}^{\infty}\left(V_{i},[0,1]\right)$ such $\sum_{i=1}^{k} \phi_{i} \leq 1$ with $\sum_{i=1}^{k} \phi_{i}=1$ on a neighborhood of $K$.

Proof. By Corollary 42.20 there exists $\phi \in C_{c}^{\infty}\left(\cup_{i=1}^{k} V_{i},[0,1]\right)$ such that $\phi=1$ on a neighborhood of $K$. Now let $\left\{\phi_{i}\right\}_{i=1}^{k}$ be the functions constructed in Theorem 42.26.

In particular $\Delta F(|x|)=0$ implies $\frac{d}{d r}\left(r^{n-1} F^{\prime}(r)\right)=0$ and hence $F^{\prime}(r)=$ $\tilde{A} r^{1-n}$. That is to say

$$
F(r)=\left\{\begin{array}{c}
A r^{2-n}+B \text { if } n \neq 2 \\
A \ln r+B \text { if } n=2
\end{array}\right.
$$

Proof. Since $\left(\partial_{v} f\right)(x)=F^{\prime}(|x|) \partial_{v}|x|=F^{\prime}(|x|) \hat{x} \cdot v$ where $\hat{x}=\frac{x}{|x|}$, $\nabla f(x)=F^{\prime}(|x|) \hat{x}$. Hence for $g \in C_{c}^{1}\left(\mathbb{R}^{n}\right)$,

$$
\begin{aligned}
\int_{\mathbb{R}^{n}} \triangle f(x) g(x) d x & =-\int_{\mathbb{R}^{n}} \nabla f(x) \cdot \nabla g(x) d x \\
& =-\int_{\mathbb{R}^{n}} F^{\prime}(r) \hat{x} \cdot \nabla g(r \hat{x}) d x \\
& =-\int_{S^{n-1} \times[0, \infty)} F^{\prime}(r) \frac{d}{d r} g(r \omega) r^{n-1} d r d \sigma(\omega) \\
& =\int_{S^{n-1} \times[0, \infty)} \frac{d}{d r}\left(r^{n-1} F^{\prime}(r)\right) g(r \omega) d r d \sigma(\omega) \\
& =\int_{S^{n-1} \times[0, \infty)} \frac{1}{r^{n-1}} \frac{d}{d r}\left(r^{n-1} F^{\prime}(r)\right) g(r \omega) r^{n-1} d r d \sigma(\omega) \\
& =\left.\int_{\mathbb{R}^{n}} \frac{1}{r^{n-1}} \frac{d}{d r}\left(r^{n-1} F^{\prime}(r)\right)\right|_{r=|x|} g(x) d x
\end{aligned}
$$

Since this is valid for all $g \in C_{c}^{1}\left(\mathbb{R}^{n}\right)$, Eq. (43.3) is valid. Alternatively, we may simply compute directly as follows:

$$
\begin{aligned}
\triangle f(x) & =\nabla \cdot\left[F^{\prime}(|x|) \hat{x}\right]=\nabla F^{\prime}(|x|) \cdot \hat{x}+F^{\prime}(|x|) \nabla \cdot \hat{x} \\
& =F^{\prime \prime}(|x|) \hat{x} \cdot \hat{x}+F^{\prime}(|x|) \nabla \cdot \frac{x}{|x|} \\
& =F^{\prime \prime}(|x|)+F^{\prime}(|x|)\left\{\frac{n}{|x|}-\frac{x}{|x|^{2}} \cdot \hat{x}\right\} \\
& =F^{\prime \prime}(|x|)+\frac{(n-1)}{|x|} F^{\prime}(|x|) .
\end{aligned}
$$

Notation 43.3 For $t>0$, let

$$
\alpha(t):=\alpha_{n}(t):=c_{n} \begin{cases}\frac{1}{t^{n-2}} & \text { if } n \neq 2  \tag{43.4}\\ \ln t & \text { if } n=2,\end{cases}
$$

where $c_{n}=\left\{\begin{array}{cl}\frac{1}{(n-2) \sigma\left(S^{n-1}\right)} & \text { if } n \neq 2 \\ -\frac{1}{2 \pi} & \text { if } n=2 .\end{array}\right.$ Also let

$$
\phi(y)=\phi_{n}(y):=\alpha(|y|)=c_{n} \begin{cases}\frac{1}{\mid y^{n-2}} & \text { if } n \neq 2  \tag{43.5}\\ \ln |y| & \text { if } n=2 .\end{cases}
$$

An important feature of $\alpha$ is that

$$
\alpha^{\prime}(t)=c_{n}\left\{\begin{array}{cc}
-(n-2) \frac{1}{t^{n-1}} & \text { if } n \neq 2  \tag{43.6}\\
\frac{1}{t} & \text { if } n=2
\end{array}=-\frac{1}{\sigma\left(S^{n-1}\right)} \frac{1}{t^{n-1}}\right.
$$

for all $n$. This then implies, for all $n$, that

$$
\begin{equation*}
\nabla \phi(x)=\nabla[\alpha(|x|)]=\alpha^{\prime}(|x|) \hat{x}=-\frac{1}{\sigma\left(S^{n-1}\right)} \frac{1}{|x|^{n-1}} \hat{x}=-\frac{1}{\sigma\left(S^{n-1}\right)} \frac{1}{|x|^{n}} x . \tag{43.7}
\end{equation*}
$$

One more piece of notation will be useful in the sequel.
Notation 43.4 (Averaging operator) Suppose $\mu$ is a finite measure on some space $\Omega$, we will define

$$
f_{\Omega} f d \mu:=\frac{1}{\mu(\Omega)} \int_{\Omega} f d \mu
$$

For example if $\Omega$ is a compact manifold with $C^{2}$ - boundary in $\mathbb{R}^{n}$ then

$$
f_{\Omega} f(x) d x=\frac{1}{m(\Omega)} \int_{\Omega} f(x) d x=\frac{1}{\operatorname{Vol}(\Omega)} \int_{\Omega} f(x) d x
$$

and

$$
f_{\partial \Omega} f d \sigma=\frac{1}{\sigma(\partial \Omega)} \int_{\partial \Omega} f(x) d x=\frac{1}{\operatorname{Area}(\partial \Omega)} \int_{\partial \Omega} f(x) d x
$$

Theorem 43.5. Let $\Omega$ be a compact manifold with $C^{2}$ - boundary, $u \in$ $C^{2}\left(\Omega^{o}\right) \cap C^{1}(\Omega)$ with $\int_{\Omega}|\Delta u(y)| d y<\infty$. Then for $x \in \Omega$

$$
\begin{equation*}
u(x)=\int_{\partial \Omega}\left(\phi(x-y) \frac{\partial u}{\partial n}(y)-u(y) \frac{\partial \phi(x-y)}{\partial n_{y}}\right) d \sigma-\int_{\Omega} \phi(x-y) \triangle u(y) d y \tag{43.8}
\end{equation*}
$$

Proof. Let $\psi(y):=\phi(x-y)$ and $\epsilon>0$ be small so that $B_{x}(\epsilon) \subset \Omega$ and let $\Omega_{\epsilon}:=\Omega \backslash B_{x}(\epsilon)$, see Figure 43.1 below.

Let us begin by observing

$$
\begin{aligned}
\int_{|x-y| \leq \epsilon} \psi(y) d y & =\int_{|y| \leq \epsilon} \frac{1}{|y|^{n-2}} d y=\sigma\left(S^{n-1}\right) \int_{0}^{\epsilon} \frac{1}{r^{n-2}} r^{n-1} d r \\
& =\sigma\left(S^{n-1}\right) \int_{0}^{\epsilon} r d r=\sigma\left(S^{n-1}\right) \frac{\epsilon^{2}}{2}
\end{aligned}
$$

when $n \neq 2$ and for $n=2$ that


Fig. 43.1. Removing the region where $\psi$ is singular from $\Omega$.

$$
\begin{aligned}
\int_{|x-y| \leq \epsilon} \psi(y) d y & =\int_{|y| \leq \epsilon} \ln |y| d y=\sigma\left(S^{1}\right) \int_{0}^{\epsilon} r \ln r d r \\
& =2 \pi\left[\frac{1}{2} r^{2} \ln r-\frac{1}{4} r^{2}\right]_{0}^{\epsilon}=\pi \epsilon^{2}[\ln \epsilon-1 / 2]
\end{aligned}
$$

This shows $\psi \in L_{l o c}^{1}(\Omega)$ and hence that $\psi \Delta u \in L^{1}(\Omega)$ and by dominated convergence theorem,

$$
\int_{\Omega} \psi(y) \Delta u(y) d y=\lim _{\epsilon \downarrow 0} \int_{\Omega_{\epsilon}} \psi(y) \Delta u(y) d y .
$$

Using Green's identity (Eq. (43.2) and Proposition 43.2) and $\Delta \psi=0$ on $\Omega_{\epsilon}$, we find

$$
\begin{align*}
\int_{\Omega_{\epsilon}} \triangle u(y) \psi(y) d y= & \int_{\Omega_{\epsilon}} \Delta \psi(y) u(y) d y+\int_{\partial \Omega_{\epsilon}}\left(\psi \frac{\partial u}{\partial n}-\frac{\partial \psi}{\partial n} u\right) d \sigma \\
= & \int_{\partial \Omega}\left(\psi \frac{\partial u}{\partial n}-\frac{\partial \psi}{\partial n} u\right) d \sigma \\
& +\int_{\partial \Omega_{\epsilon} \backslash \partial \Omega}\left(\psi \frac{\partial u}{\partial n}-\frac{\partial \psi}{\partial n} u\right) d \sigma . \tag{43.9}
\end{align*}
$$

Working on the last term in Eq. (43.9) we have, for $n \neq 2$,

43 Poisson and Laplace's Equation

$$
\begin{aligned}
\int_{\partial B(x, \epsilon)} \psi(y) \frac{\partial u}{\partial \mathbf{n}}(y) & =\int_{|y|=\epsilon} \psi(x+\omega) \frac{\partial u}{\partial \mathbf{n}}(x+\omega) d \sigma(\omega) \\
& =\int_{|\omega|=1} \psi(x+\epsilon \omega) \frac{\partial u}{\partial \mathbf{n}}(x+\epsilon \omega) \epsilon^{n-1} d \sigma(\omega) \\
& =\int_{|\omega|=1} \frac{1}{\epsilon^{n-2}} \frac{\partial u}{\partial \mathbf{n}}(x+\epsilon \omega) \epsilon^{n-1} d \sigma(\omega) \\
& =\epsilon \int_{|\omega|} \frac{\partial u}{\partial \mathbf{n}}(x+\epsilon \omega) d \sigma(\omega) \rightarrow 0 \text { as } \epsilon \downarrow 0 .
\end{aligned}
$$

Similarly when $n=2$,

$$
\int_{\partial B(x, \epsilon)} \psi(y) \frac{\partial u}{\partial \mathbf{n}}(y)=\epsilon \ln \epsilon \int_{|\omega|=1} \frac{\partial u}{\partial \mathbf{n}}(x+\epsilon \omega) d \sigma(\omega) \rightarrow 0 \text { as } \epsilon \downarrow 0 \text {. }
$$

Using Eq. (43.7) and $n(y)=-(\widehat{y-x})$ as in Figure 43.2 we find


Fig. 43.2. The outward normal to $\Omega_{\epsilon}$ is the inward normal to $B(x, \epsilon)$.

$$
\begin{align*}
\frac{\partial \psi}{\partial \mathbf{n}}(y) & =\nabla_{y} \phi(y-x) \cdot n(y)=-\frac{1}{\sigma\left(S^{n-1}\right)} \frac{1}{|y-x|^{n}}(y-x) \cdot(-(\widehat{y-x})) \\
& =\frac{1}{\sigma\left(S^{n-1}\right)} \frac{1}{\epsilon^{n-1}} \tag{43.10}
\end{align*}
$$

and therefore

$$
\begin{aligned}
-\int_{\partial \Omega_{\epsilon} \backslash \partial \Omega} u \frac{\partial \psi}{\partial \mathbf{n}} d \sigma(y) & =-\frac{1}{\sigma\left(S^{n-1}\right)} \frac{1}{\epsilon^{n-1}} \int_{\partial B(x, \epsilon)} u(y) d \sigma(y) \\
& =-\int_{|\omega|=1} u(x+\epsilon \omega) d \sigma(\omega) \rightarrow-u(x) \text { as } \epsilon \downarrow 0
\end{aligned}
$$

by the dominated convergence theorem. So we may pass to the limit in Eq. (43.9) to find

$$
\int_{\Omega} \psi(y) \Delta u(y) d y=\int_{\partial \Omega}\left(\psi(y) \frac{\partial u}{\partial \mathbf{n}}-u \frac{\partial \psi}{\partial \mathbf{n}}\right) d \sigma(y)-u(x)
$$

which is equivalent to Eq. (43.8).
The following Corollary gives an easy but useful extension of Theorem 43.5 It will be us

Corollary 43.6. Keeping the same notation as in Theorem 43.5. Further assume that $h \in C^{2}\left(\Omega^{o}\right) \cap C^{1}(\Omega)$ and $\Delta h=0$ and set $G(y):=\phi(x-y)+h(y)$. Then we still have the representation formula

$$
\begin{equation*}
u(x)=\int_{\partial \Omega}\left(G(y) \frac{\partial u}{\partial \mathbf{n}}(y)-u(y) \frac{\partial G(y)}{\partial \mathbf{n}}\right) d \sigma-\int_{\Omega} G(y) \triangle u(y) d y . \tag{43.11}
\end{equation*}
$$

Proof. By Green's identity (Proposition 22.30) with $v=h$,

$$
\int_{\Omega} \triangle u h d m=\int_{\Omega}(\triangle u h-\triangle h u) d m=\int_{\partial \Omega}\left(h \frac{\partial u}{\partial n}-\frac{\partial h}{\partial n} u\right) d \sigma,
$$

i.e.

$$
\begin{equation*}
0=-\int_{\Omega} \triangle u h d m+\int_{\partial \Omega}\left(h \frac{\partial u}{\partial n}-\frac{\partial h}{\partial n} u\right) d \sigma \tag{43.12}
\end{equation*}
$$

Eq. (43.11) now follows by adding Eqs. (43.8) and (43.12).
Corollary 43.7. For all $u \in C_{c}^{2}\left(\mathbb{R}^{n}\right)$,

$$
\begin{equation*}
-\int_{\mathbb{R}^{n}} \triangle u(y) \phi(y) d y=u(0) \tag{43.13}
\end{equation*}
$$

Proof. Let $\Omega=B(0, R)$ where $R$ is chosen so large that $\operatorname{supp}(g) \subset \Omega$, then by Theorem 43.5,

$$
\begin{aligned}
u(0) & =\int_{\partial \Omega}\left(\phi(y) \frac{\partial u}{\partial v}(y)-u(y) \frac{\partial \phi(y)}{\partial v_{y}}\right) d \sigma-\int_{\Omega} \phi(y) \triangle u(y) d y \\
& =-\int_{\Omega} \phi(y) \triangle u(y) d y
\end{aligned}
$$

Remark 43.8. We summarize (43.13) by saying $-\triangle \phi=\delta$.

Formally we expect for reasonable functions $\rho$ that

$$
\Delta(\phi * \rho)=\Delta \phi * \rho=-\delta * \rho=-\rho
$$

Theorem 43.9. Suppose $\Omega \subset_{o} \mathbb{R}^{n}, \rho \in C^{2}(\Omega) \cap L^{1}(\Omega)$ and

$$
u(x):=\int_{\Omega} \phi(x-y) \rho(y) d y=\left(\phi * 1_{\Omega} \rho\right)(x)
$$

then

$$
-\triangle u=\rho \text { on } \Omega
$$

Proof. First assume that $\rho \in C_{c}^{2}(\Omega)$ in which case we may set $\rho:=1_{\Omega} \rho \in$ $C_{c}^{2}\left(\mathbb{R}^{n}\right)$. Therefore

$$
u(x)=\int_{\mathbb{R}^{n}} \rho(y) \frac{1}{|x-y|^{n-2}} d y=\int_{\mathbb{R}^{n}} \rho(x-y) \frac{1}{|y|^{n-2}} d y
$$

and so we may differentiate under the integral to find

$$
\triangle u(x)=\int_{\mathbb{R}^{n}} \triangle_{x} \rho(x-y) \frac{1}{|y|^{n-2}} d y=-\rho(x)
$$

where the last equality follows from Corollary 43.7.
For $\rho \in C^{2}(\Omega) \cap L^{1}(\Omega)$ and $x_{0} \in \Omega$, choose $\alpha \in C_{c}^{\infty}(\Omega,[0,1])$ such that $\alpha=$ 1 in a neighborhood of $x_{0}$ and let $\beta:=1-\alpha$. Then $u=(\phi * \alpha \rho)+\left(\phi * \beta 1_{\Omega} \rho\right)$ and so

$$
\begin{equation*}
\Delta u=\Delta(\phi * \alpha \rho)+\Delta\left(\phi * \beta 1_{\Omega} \rho\right) . \tag{43.14}
\end{equation*}
$$

By what we have just proved

$$
\begin{equation*}
\Delta(\phi * \alpha \rho)(x)=-(\alpha \rho)(x)=-\rho(x) \text { for } x \text { near } x_{0} \tag{43.15}
\end{equation*}
$$

Since $\beta=0$ near $x_{0}$ and

$$
\left(\phi * \beta 1_{\Omega} \rho\right)(x)=\int_{\Omega} \phi(x-y) \beta(y) \rho(y) d y
$$

we may differentiate past the integral to learn

$$
\begin{equation*}
\Delta\left(\phi * \beta 1_{\Omega} \rho\right)(x)=\int_{\Omega} \Delta_{x} \phi(x-y) \beta(y) \rho(y) d y=0 \tag{43.16}
\end{equation*}
$$

for $x$ near $x_{0}$. and this completes the proof. The combination of Eqs. (43.14 -43.16 ) completes the proof.

### 43.1 Harmonic and Subharmonic Functions

Definition 43.10 (HarmonicFunctions). Let $\Omega \subset_{o} \mathbb{R}^{n}$. A function $u \in$ $C^{2}(\Omega)$ is said to be harmonic (subharmonic) on $\Omega$ if $\Delta u=0(\Delta u \geq 0)$ on $\Omega$.

Because of the Cauchy Riemann equations, the real and imaginary parts of holomorphic functions are harmonic. For example $z^{2}=\left(x^{2}-y^{2}\right)+2 i x y$ implies $\left(x^{2}-y^{2}\right)$ and $x y$ are harmonic functions on the plane. Similarly,

$$
\begin{aligned}
e^{z} & =e^{x} \cos y+i e^{x} \sin y \text { and } \\
\ln (z) & =\ln r+i \theta
\end{aligned}
$$

implies

$$
e^{x} \cos y, e^{x} \sin y, \ln r, \text { and } \theta(x, y)
$$

are harmonic functions on their domains of definition.
Remark 43.11. If we can choose $h$ in Corollary 43.6 so that $G=0$ on $\partial \Omega$, then Eq. (43.11) gives

$$
\begin{equation*}
u(x)=-\int_{\Omega} G(y) \triangle u(y) d y-\int_{\partial \Omega} u \frac{\partial G(y)}{\partial v} d \sigma \tag{43.17}
\end{equation*}
$$

which shows how to recover $u(x)$ from $\Delta u$ on $\Omega$ and $u$ on $\partial \Omega$. The next theorem is a consequence of this remark.

Theorem 43.12 (Mean Value Property). If $\triangle u=0$ on $\Omega$ and $\overline{B(x, r)} \subset$ $\Omega$ then

$$
\begin{equation*}
u(x)=\frac{1}{\sigma(\partial B(x, r))} \int_{\partial B(x, r)} u(y) d \sigma(y)=: \int_{\partial B(x, r)} u d \sigma \tag{43.18}
\end{equation*}
$$

More generally if $\Delta u \geq 0$ on $\Omega$, then

$$
\begin{equation*}
u(x) \leq \int_{\partial B(x, r)} u d \sigma \tag{43.19}
\end{equation*}
$$

Proof. For $y \in B(x, r)$,

$$
G(y)=\phi(x-y)-\alpha(r)=\alpha(|x-y|)-\alpha(r)
$$

where $\alpha$ is defined as in Eq. (43.4). Then $G(y)=0$ for $y \in \partial B(x, r)$ and $G(y)>0$ for all $y \in B(x, r)$ because $\alpha$ is decreasing as is seen from Eq. (43.6). From Eq. (43.10) (using now that $n$ is the outward normal to $B(x, r)$ ),

$$
\frac{\partial G}{\partial \mathbf{n}}(x+r \omega)=-\frac{1}{\sigma\left(S^{n-1}\right) r^{n-1}} \text { for }|\omega|=1
$$

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and so according to Eq. (43.17) we have

$$
\begin{align*}
u(x) & =\frac{1}{r^{n-1} \sigma\left(S^{n-1}\right)} \int_{\partial B(x, r)} u d \sigma-\int_{B(x, r)} G(y) \triangle u(y) d y \\
& =\int_{\partial B(x, r)} u d \sigma-\int_{B(x, r)} G(y) \triangle u(y) d y \tag{43.20}
\end{align*}
$$

This completes the proof since $G(y)>0$ for all $y \in B(x, r)$.
Remark 43.13 (Mean value theorem). Assuming $\overline{B(x, R)} \subset \Omega$ and multiplying Eq. (43.18) (Eq. (43.19)) by

$$
\sigma(\partial B(x, r))=\sigma\left(S^{n-1}\right) r^{n-1}
$$

and then integrating on $0 \leq r \leq R$, implies

$$
\begin{aligned}
u(x) m(B(x, R)) & =(\text { or } \leq) \int_{0}^{R} d r \int_{\partial B(x, r)} u(y) d \sigma(y) \\
& =\int_{0}^{R} d r r^{n-1} \int_{S^{n-1}} u(x+r \omega) d \sigma(\omega)=\int_{B(x, R)} u d m
\end{aligned}
$$

Therefore if $\Delta u=0$ or $\Delta u \geq 0$ then

$$
\begin{equation*}
u(x)=\int_{B(x, R)} u d m \text { or } u(x) \leq \int_{B(x, R)} u d m \text { respectively } \tag{43.21}
\end{equation*}
$$

for all $\overline{B(x, R)} \subset \Omega$.
Proposition 43.14 (Converse of the mean value property). If $u \in$ $C(\Omega)$ (or more generally measurable and locally bounded) and

$$
\begin{equation*}
u(x)=f_{\partial B(x, r)} u(y) d \sigma(y) \tag{43.22}
\end{equation*}
$$

for all $x \in \Omega$ and $r>0$ such that $\overline{B(x, r)} \subset \Omega$, then $u \in C^{\infty}(\Omega)$ and $\triangle u=0$. Similarly, if $u \in C^{2}(\Omega)$ and $x \in \Omega$ and

$$
\begin{equation*}
u(x) \leq \int_{\partial B(x, r)} u(y) d \sigma(y) \tag{43.23}
\end{equation*}
$$

for all $r$ sufficiently small, then $\Delta u(x) \geq 0$.

Proof. First assume $u \in C(\Omega)$ and Eq. (43.22) hold which implies

$$
\begin{equation*}
u(x)=f_{S} u(x+r \omega) d \sigma(\omega) \tag{43.24}
\end{equation*}
$$

for all $x \in \Omega$ and $r$ sufficiently small, where $S=S^{n-1}$ denotes the unit sphere in $\mathbb{R}^{n}$. Let $\eta \in C_{c}^{\infty}\left(\mathbb{R}^{n},[0, \infty)\right)$ such that $\eta(0)>0$ and

$$
1=\int_{\mathbb{R}^{n}} \eta\left(|x|^{2}\right) d x=\sigma(S) \int_{0}^{\infty} \eta\left(r^{2}\right) r^{n-1} d r
$$

and for $\epsilon>0$ let $\eta_{\epsilon}(x)=\epsilon^{-n} \eta\left(\frac{|x|^{2}}{\epsilon^{2}}\right) \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ and $u_{\epsilon}(x)=\eta_{\epsilon} * u(x)$. Then for any $x_{0} \in \Omega$ and $\epsilon>0$ sufficiently small, $u_{\epsilon}$ is a well defined smooth function near $x_{0}$. Moreover for $x$ near $x_{0}$ we have

$$
\begin{aligned}
u_{\epsilon}(x) & =\int_{\mathbb{R}^{n}} \eta_{\epsilon}(x-y) u(y) d y=\int_{0}^{\infty} d r r^{n-1} \int_{|\omega|=1} \eta_{\epsilon}(r \omega) u(x+r \omega) d \sigma(\omega) \\
& =\int_{0}^{\infty} d r r^{n-1} \int_{|\omega|=1} \epsilon^{-n} \eta\left(\frac{r^{2}}{\epsilon^{2}}\right) u(x+r \omega) d \sigma(\omega) \\
& =u(x) \sigma(S) \int_{0}^{\infty} d r r^{n-1} \epsilon^{-n} \eta\left(\frac{r^{2}}{\epsilon^{2}}\right)=u(x)
\end{aligned}
$$

which shows $u$ is smooth near $x_{0}$.
Now suppose that $u \in C^{2}$, and $u$ satisfies Eq. (43.23), $x \in \Omega$ and $|r|<\epsilon$ with $\epsilon$ sufficiently small so that

$$
f(r):=\int_{\partial B(x, r)} u d \sigma=\int_{S^{n-1}} u(x+r \omega) d \sigma(\omega)
$$

is well defined. Clearly $f \in C^{2}(-\epsilon, \epsilon), f$ is an even function of $r$ so $f^{\prime}(0)=0$, $f(0)=u(x)$ and $f(r) \geq f(0)$. From these conditions it follows that $f^{\prime \prime}(0) \geq 0$ for otherwise we would find from Taylor's theorem that $f(r)<f(0)$ for $0<$ $|r|<\epsilon$. On the other hand

$$
\begin{align*}
0 & \leq f^{\prime \prime}(0)=f_{S^{n-1}}\left(\partial_{\omega}^{2} u\right)(x) d \sigma(\omega) \\
& =f_{S^{n-1}}\left(\partial_{i} \partial_{j} u\right)(x) \omega_{i} \omega_{j} d \sigma(\omega) \\
& =\left(\partial_{i} \partial_{j} u\right)(x) \delta_{i j} f_{S^{n-1}} \omega_{i}^{2} d \sigma(\omega)=\frac{1}{n} \Delta u(x) . \tag{43.25}
\end{align*}
$$

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wherein we have used the symmetry of $d \sigma$ on $S^{n-1}$ to conclude

$$
\int_{S^{n-1}} \omega_{i} \omega_{j} d \sigma(\omega)=0 \text { if } i \neq j
$$

and

$$
\begin{aligned}
\int_{S^{n-1}} \omega_{i}^{2} d \sigma(\omega) & =\frac{1}{n} \sum_{j=1}^{n} \int_{S^{n-1}} \omega_{j}^{2} d \sigma(\omega) \\
& =\frac{1}{n} \int_{S^{n-1}}|\omega|^{2} d \sigma(\omega)=\frac{1}{n} \forall i .
\end{aligned}
$$

Alternatively, by the divergence theorem,

$$
\begin{aligned}
\int_{S^{n-1}} \omega_{i} \omega_{j} d \sigma(\omega) & =\int_{S^{n-1}} \omega_{i} e_{j} \cdot n(\omega) d \sigma(\omega) \\
& =\frac{1}{\sigma\left(S^{n-1}\right)} \int_{B(0,1)} \nabla \cdot\left(x_{i} e_{j}\right) d m \\
& =\frac{1}{\sigma\left(S^{n-1}\right)} m(B(0,1)) \delta_{i j}=\frac{1}{n} \delta_{i j}
\end{aligned}
$$

This completes the proof since if $u$ satisfies (43.22) then $f$ is constant and it follows from Eq. (43.25) that $\Delta u(x)=0$

Second proof of the last statement. Now that we know $u$ is $C^{2}$ we have by Eq. (43.20) that

$$
\int_{B(x, r)} G(y) \triangle u(y) d y=\int_{\partial B(x, r)} u d \sigma-u(x) \geq 0
$$

and since with $\alpha$ as in Eq. (43.4),

$$
\begin{aligned}
\int_{B(x, r)} G(y) \triangle u(y) d y & =\int_{B(0, r)} G(x+y) \triangle u(x+y) d y \\
& =\int_{0}^{r} \rho^{n-1} d \rho \int_{S^{n}} d \omega G(x+\rho \omega) \triangle u(x+\rho \omega) \\
& =\int_{0}^{r} \rho^{n-1} d \rho(\alpha(\rho)-\alpha(r)) \int_{S^{n}} d \omega \triangle u(x+\rho \omega) \\
& \cong \Delta u(x) \sigma\left(S^{n-1}\right) \int_{0}^{r} \rho^{n-1} d \rho(\alpha(\rho)-\alpha(r)) \\
& =\Delta u(x) \sigma\left(S^{n-1}\right) c_{n}\left\{\frac{r^{2}}{2}-\frac{r^{n}}{n r^{n-2}}\right\} \\
& =b_{n} r^{2} \Delta u(x)
\end{aligned}
$$

where $b_{n}$ is a positive constant. From this it follows that $\Delta u(x) \geq 0$.
Third proof of the last statement. If $u \in C^{2}(\Omega)$ satisfies expand $u(x+r \omega)$ in a Taylor series

$$
u(x+r \omega)=u(x)+r \nabla u(x) \cdot \omega+\frac{r^{2}}{2} \partial_{\omega}^{2} u(x)+o\left(r^{3}\right)
$$

and integrate on $\omega$ to find

$$
\begin{aligned}
f_{\partial B(x, r)} u d \sigma & =\int_{S^{n-1}} u(x+r \omega) d \sigma(\omega) \\
& =f_{S^{n-1}}\left[u(x)+r \nabla u(x) \cdot \omega+r^{2} \frac{1}{2} \partial_{\omega}^{2} u(x)+\ldots\right] d \sigma(\omega) \\
& =u(x)+\frac{1}{2} r^{2} \Delta u(x)+o\left(r^{2}\right) .
\end{aligned}
$$

Thus if $u$ satisfies Eq. (43.22) Eq. (43.23) we conclude

$$
\begin{aligned}
& u(x)=u(x)+\frac{1}{2} r^{2} \Delta u(x)+o\left(r^{2}\right) \text { or } \\
& u(x) \leq u(x)+\frac{1}{2} r^{2} \Delta u(x)+o\left(r^{2}\right)
\end{aligned}
$$

from which we conclude $\Delta u(x)=0$ or $\Delta u(x) \geq 0$ respectively.
Fourth proof of the statement: If $u$ satisfies Eq. (43.22) then $\Delta u=0$ Since we already know $u$ is smooth, it is permissible to differentiate Eq. (43.24) in $r$ to learn,

$$
\begin{aligned}
0 & =\int_{S^{n-1}} \nabla u(x+r \omega) \cdot \omega d \sigma(\omega)=\int_{S^{n-1}} \frac{\partial u}{\partial \mathbf{n}}(x+r \omega) d \sigma(\omega) \\
& =\frac{1}{\sigma\left(S^{n-1}\right) r^{n-1}} \int_{\partial B(x, r)} \nabla u \cdot \mathbf{n} d \sigma=\frac{1}{\sigma\left(S^{n-1}\right) r^{n-1}} \int_{B(x, r)} \Delta u d m .
\end{aligned}
$$

Dividing this equation by $r$ and letting $r \downarrow 0$ shows $\Delta u(x)=0$.
Corollary 43.15 (Smoothness of Harmonic Functions). If $u \in C^{2}(\Omega)$ and $\Delta u=0$ then $u \in C^{\infty}(\Omega)$. (Soon we will show $u$ is real analytic, see Theorem 43.16 of Corollary 43.34 below.)
Theorem 43.16 (Bounds on Harmonic functions). Suppose $u$ is a Harmonic function on $\Omega \subset \mathbb{R}^{n}, x_{0} \in \Omega, \alpha$ is a multi-index with $k:=|\alpha|$ and $0<r<\operatorname{dist}\left(x_{o}, \partial \Omega\right)$. Then

$$
\begin{equation*}
\left|D^{\alpha} u\left(x_{0}\right)\right| \leq \frac{C_{k}}{r^{n+k}}\|u\|_{L^{1}\left(B\left(x_{0}, r\right)\right)} \leq \frac{C_{k}}{\operatorname{dist}\left(x_{o}, \partial \Omega\right)^{n+k}}\|u\|_{L^{1}(\Omega)} \tag{43.26}
\end{equation*}
$$

where $C_{k}=\frac{\left(2^{n+1} n k\right)^{k}}{\alpha(n)}$. In particular one shows that $u$ is real analytic in $\Omega$.

Proof. Let $\eta_{\epsilon}(x)$ be constructed as in the proof of Proposition 43.14 so that $u(x)=u * \eta_{\epsilon}(x)$. Therefore, $D^{\alpha} u(x)=u_{*} D^{\alpha} \eta_{\epsilon}(x)$ and hence

$$
\left|D^{\alpha} u\left(x_{0}\right)\right| \leq\|u\|_{L^{1}\left(B\left(x_{0}, \epsilon\right)\right)}\left\|D^{\alpha} \eta_{\epsilon}\right\|_{L^{\infty}} .
$$

Now

$$
D^{\alpha} \eta_{\epsilon}(x)=\epsilon^{-n} \frac{1}{\epsilon^{|\alpha|}}\left(D^{\alpha} \eta\right)\left(\frac{x}{\epsilon}\right)
$$

so that

$$
\left|D^{\alpha} \eta_{\epsilon}(x)\right|=\epsilon^{-n} \frac{1}{\epsilon^{|\alpha|}}\left|\left(D^{\alpha} \eta\right)\left(\frac{x}{\epsilon}\right)\right| \leq C_{\alpha} \frac{1}{\epsilon^{|\alpha|+n}} \cong C_{\alpha} \frac{1}{r^{|\alpha|+n}}
$$

where the last identity is gotten by taking $\epsilon$ comparable to $r$. Putting this all together then implies that

$$
\left|D^{\alpha} u\left(x_{0}\right)\right| \leq \frac{1}{r^{n+|\alpha|}}\left\|D^{\alpha} \eta\right\|_{L^{\infty}}\|u\|_{L^{1}\left(B\left(x_{0}, r\right)\right)}
$$

which is an inequality of the form in Eq. (43.26). To get the desired constant we will have to work harder. This is done in Theorem 7. on p. 29 of the book. The idea is to use $D^{\alpha} u$ is harmonic for all $\alpha$ and therefore,

$$
\begin{aligned}
D^{\alpha} u\left(x_{0}\right) & =f_{B\left(x_{0}, \rho\right)} D^{\alpha} u d m=f_{B\left(x_{0}, \rho\right)} \partial_{i} D^{\beta} u d m \\
& =\frac{n}{\sigma\left(S^{n-1}\right) \rho^{n}} \int_{B\left(x_{0}, \rho\right)} \partial_{i} D^{\beta} u d m \\
& =\frac{n}{\sigma\left(S^{n-1}\right) \rho^{n}} \int_{\partial B\left(x_{0}, \rho\right)} D^{\beta} u n_{i} d \sigma
\end{aligned}
$$

so that

$$
\left|D^{\alpha} u\left(x_{0}\right)\right| \leq \frac{n}{\rho}\left\|D^{\beta} u\right\|_{L^{\infty}\left(B\left(x_{0}, \rho\right)\right)}
$$

and for $\alpha=0$ and $x \in B\left(x_{0}, r / 2\right)$ we have

$$
|u(x)| \leq \int_{B(x, r / 2)}|u| d m \leq \frac{1}{|B(0,1)|}\left(\frac{2}{r}\right)^{n}\|u\|_{L^{1}\left(B\left(x_{0}, r\right)\right.}
$$

Using this and similar inequalities along with a tricky induction argument one gets the desired constants. The details are in Theorem 7. p. 29 and Theorem 10 p. 31 of the book. (See also Corollary 43.34 below for another proof of analyticity of $u$.)

Corollary 43.17 (Liouville's Theorem). Suppose $u \in C^{2}\left(\mathbb{R}^{n}\right), \Delta u=0$ on $\mathbb{R}^{n}$ and $|u(x)| \leq C\left(1+|x|^{N}\right)$ for all $x \in \mathbb{R}^{n}$. Then $u$ is a polynomial of degree at most $N$.

Proof. We have seen there are constants $C_{|\alpha|}<\infty$ such that

$$
\begin{aligned}
\left|D^{\alpha} u\left(x_{0}\right)\right| & \leq C_{|\alpha|}\|u\|_{L^{1}\left(B\left(x_{0}, r\right)\right)} \frac{1}{r^{n+|\alpha|}} \\
& \leq \widetilde{C}_{|\alpha|} r^{n}\|u\|_{L^{\infty}\left(B\left(x_{0}, r\right)\right)} \cdot \frac{1}{r^{n+|\alpha|}} \\
& \cong C \frac{\left(1+r^{N}\right)}{r^{|\alpha|}} \rightarrow 0 \text { as } r \rightarrow \infty
\end{aligned}
$$

when if $|\alpha|>N$. Therefore $D^{\alpha} u:=0$ for all $|\alpha|>N$ and the the result follows by Taylor's Theorem with remainder,

$$
u(x)=\sum_{|\alpha| \leq N} \frac{D^{\alpha} u\left(x_{0}\right)\left(x-x_{0}\right)^{\alpha}}{\alpha!}
$$

Corollary 43.18 (Compactness of Harmonic Functions). Suppose $\Omega \subset_{o}$ $\mathbb{R}^{n}$ and $u_{n} \in C^{2}(\Omega)$ is a sequence of harmonic functions such that for each compact set $K \subset \Omega$,

$$
C_{K}:=\sup \left\{\int_{K}\left|u_{n}\right| d m: n \in \mathbb{N}\right\}<\infty
$$

Then there is a subsequence $\left\{v_{n}\right\} \subset\left\{u_{n}\right\}$ which converges, along with all of its derivatives, uniformly on compact subsets of $\Omega$ to a harmonic function $u$.

Proof. An application of Theorem 43.16 shows that for each compact set $K \subset \Omega, \sup _{n}\left|\nabla u_{n}\right|_{L^{\infty}(K)}<\infty$ and hence by the locally compact form of the Arzela-Ascolli theorem, there is a subsequence $\left\{v_{n}\right\} \subset\left\{u_{n}\right\}$ which converges uniformly on compact subsets of $\Omega$ to a continuous function $u \in C(\Omega)$. Passing to the limit in the mean value theorem for harmonic functions along with the converse to the mean value theorem, Proposition 43.14, shows $u$ is harmonic on $\Omega$. Since $v_{m} \rightarrow u$ uniformly on compacts it follows for any $K \sqsubset \sqsubset \Omega$ that $\int_{K}\left|u-v_{n}\right| d m \rightarrow 0$. Another application of Theorem 43.16 then shows $D^{\alpha} v_{n} \rightarrow D^{\alpha} u$ uniformly on compacts.

In light of Proposition 43.14, we will extend the notion of subharmonicity as follows.

Definition 43.19 (Subharmonic Functions). A function $u \in C(\Omega)$ is said to be subharmonic if for all $x \in \Omega$ and all $r>0$ sufficiently small,

$$
u(x) \leq \int_{\partial B(x, r)} u d \sigma
$$

The reason for the name subharmonic should become apparent from Corollary 43.26 below.

Remark 43.20. Suppose that $u, v \in C(\Omega)$ are subharmonic functions then so is $u+v$. Indeed,

$$
u(x)+v(x) \leq f_{\partial B(x, r)} u d \sigma+f_{\partial B(x, r)} v d \sigma=f_{\partial B(x, r)}(u+v) d \sigma .
$$

Theorem 43.21 (Harnack's Inequality). Let $V$ be a precompact open and connected subset of $\Omega$. Then there exists $C=C(V, \Omega)$ such that

$$
\begin{equation*}
\sup _{V} u \leq C \inf _{V} u \tag{43.27}
\end{equation*}
$$

for all non-negative harmonic functions, $u$, on $\Omega$.
Proof. Let $r=\frac{1}{4} \operatorname{dist}\left(V, \Omega^{c}\right)$ and $x \in V$ (as in Figure 43.3) and $|y-x| \leq r$,


Fig. 43.3. A pre-compact region $V \subset \Omega$.
then by the mean value equality in Eq. (43.21) of Remark 43.13,

$$
\begin{aligned}
u(x) & =\int_{B(x, 2 r)} u(z) d z=\frac{1}{m(B(0,1))(2 r)^{n}} \int_{B(x, 2 r)} u(z) d z \\
& \geq \frac{1}{m(B(0,1))(2 r)^{n}} \int_{B(y, r)} u(z) d z=\frac{1}{2^{n}} \int_{B(y, r)} u(z) d z=\frac{1}{2^{n}} u(y),
\end{aligned}
$$

see Figure 43.4. Therefore

$$
\begin{equation*}
u(x) \geq \frac{1}{2^{n}} u(y) \text { for all } x, y \in V \text { with }|x-y| \leq r \tag{43.28}
\end{equation*}
$$

Since $\bar{V}$ is compact there exists a finite cover $\mathcal{S}:=\left\{W_{i}\right\}_{i=1}^{M}$ of $\bar{V}$ consisting of balls with of radius $r$ with centers $x_{i} \in \bar{V}$.


Fig. 43.4. Nested balls.

Claim: For all $x, y \in V$, there exists a chain $\left\{B_{i}\right\}_{i=1}^{k} \subset \mathcal{S}$ of distinct balls such that $x \in B_{1}, y \in B_{k}$ and $B_{i} \cap B_{i+1} \neq \phi$ for all $i=1, \ldots, k-1$.

Indeed, by connectedness of $V$ there exists $\gamma \in C([0,1], V)$ such that $\gamma(0)=x$ while $\gamma(1)=y$. For sake of contradiction, suppose

$$
T:=\sup \left\{t \in[0,1]: \exists \text { a chain as above } \ni \gamma(t) \in B_{k}\right\}<1
$$

Since there are only finitely many possible chains (at most $\sum_{k=1}^{M} \frac{M!}{(M-k)!}$ ) there must be a chain $\left\{B_{i}\right\}_{i=1}^{k} \subset \mathcal{S}$ such that $\gamma(T) \in \bar{B}_{k}$. Let $B_{k+1} \in \mathcal{S}$ such that $B_{k+1} \ni \gamma(T)$. If $B_{k+1}=B_{j}$ with $j \leq k$, then $\left\{B_{i}\right\}_{i=1}^{j} \subset \mathcal{S}$ is a chain such that $\gamma(T) \in B_{j}$. Otherwise, since $\gamma(T) \in B_{k+1} \cap \bar{B}_{k}$, it follows that $B_{k+1} \cap B_{k} \neq \emptyset$ and $\left\{B_{i}\right\}_{i=1}^{k+1} \subset \mathcal{S}$ is a chain such that $\gamma(T) \in B_{k+1}$. In either case we will have violated the definition of $T$ and hence we must conclude $T=1$. This proves the claim, since again using the fact that there are only a finite number of possible chains, there must be at least one chain for which $\gamma(1) \in B_{k}$.

To complete the proof, for any $x, y \in V$ use a chain as in the above claim to find a sequence of points $\left\{x_{j}\right\}_{j=1}^{N} \subset V$ with $N \leq 2 M, x_{1}=x, x_{N}=y$ and $\left|x_{i+1}-x_{i}\right|<r$ for all $i$. Then by repeated use of Eq. (43.28) we may conclude

$$
u(y) \leq\left(2^{n}\right)^{2 M} u(x)
$$

Since $x, y \in V$ are arbitrary, this equation implies Eq. (43.27) with $C:=2^{4 M}$.

Remark 43.22. It is not sufficient to assume $u$ is sub-harmonic in Theorem 43.21. For example if $M>0$, then $u(x)=M x^{2}+1$ is sub-harmonic on $\mathbb{R}$, $\inf _{(-1,1)} u=1$ while $\sup _{(-1,1)} u=M+1$. Since $M$ is arbitrary, this would force $C=\infty$.

Also it is important that $\Omega$ is connected. Indeed, if $\Omega=\Omega_{1} \cup \Omega_{2}$ with $\Omega_{1}$ and $\Omega_{2}$ being disjoint open sets, then let $u \equiv 1$ on $\Omega_{1}$ and $u \equiv M$ on $\Omega_{2}$ for any $M>0$. This function is harmonic on $\Omega$ and hence $C \geq M$ for all $M>0$, i.e. $C=\infty$.

Theorem 43.23 (Strong Maximum Principle). Let $\Omega \subset \mathbb{R}^{n}$ be connected and open and $u \in C(\Omega)$ be a subharmonic function (see Definition 43.19). If $M=\sup u(x)$ is attained in $\Omega$ then $u:=M$. (Notice that $u \in C^{2}(\Omega)$ and $x \in \Omega$
$\Delta u=0$, then $u$ is harmonic and hence in particular sub-harmonic.)
Proof. Suppose there exists $x \in \Omega$ such that $M=u(x)$. If $\epsilon>0$ is chosen so that $\overline{B(x, \epsilon)} \subset \Omega$ as in Figure 43.1 and $u(y)<M$ for some $y \in \partial B(x, \epsilon)$, then by the mean value inequality,

$$
M=u(x) \leq \int_{\partial B(x, \epsilon)} u(y) d \sigma(y)<M
$$

which is nonsense. Therefore $u:=M$ on $\partial B(x, \epsilon)$ and since $\epsilon \in(0, \operatorname{dist}(x, \partial \Omega))$ we concluded that $u:=M$ on $B(x, R)$ provided $\overline{B(x, R)} \subset \Omega$. Therefore $\{x \in \Omega: u(x)=M\}$ is both open and relatively closed in $\Omega$ and hence $\{x \in \Omega: u(x)=M\}=\Omega$ because $\Omega$ is connected.

Corollary 43.24. If $\Omega$ is bounded open set $u \in C(\bar{\Omega})$ is subharmonic, then

$$
M:=\max _{x \in \bar{\Omega}} u(x)=\max _{x \in \operatorname{bd}(\Omega)} u(x) .
$$

Again this corollary applies to $u \in C(\bar{\Omega}) \cap C^{2}(\Omega)$ such that $\Delta u=0$.
Proof. By Theorem 43.23, if $x \in \Omega$ is an interior maximum of $u$, then $u=M$ on the connected component $\Omega_{x}$ of $\Omega$ which contains $x$. By continuity, $u$ is constant on $\bar{\Omega}_{x}$ and in particular $u$ takes on the value $M$ on $\operatorname{bd}(\Omega)$.

Corollary 43.25. Given $g \in C(\operatorname{bd}(\Omega)), f \in C(\Omega)$ there exists at most one function $u \in C^{2}(\Omega) \cap C(\bar{\Omega})$ such that $\triangle u=f$ on $\Omega$ and $u=g$ on $\operatorname{bd}(\Omega)$.

Proof. If $v \in C^{2}(\Omega) \cap C(\bar{\Omega})$ is another such function then $w:=u-v \in$ $C^{2}(\Omega) \cap C(\bar{\Omega})$ satisfies $\Delta w=0$ in $\Omega$ and $w=0$ on $\operatorname{bd}(\Omega)$. Therefore applying Corollary 43.24 to $w$ and $-w$ implies

$$
\max _{x \in \bar{\Omega}} w(x)=\max _{x \in \operatorname{bd}(\Omega)} w(x)=0 \text { and } \min _{x \in \bar{\Omega}} w(x)=\min _{x \in \operatorname{bd}(\Omega)} w(x)=0
$$

- 

Corollary 43.26. Suppose $g \in C(\operatorname{bd}(\Omega))$ and $u \in C^{2}(\Omega) \cap C(\bar{\Omega})$ such that $\Delta u=0$ on $\Omega$. Then $w \leq u$ for any subharmonic function $w \in C(\bar{\Omega})$ such that $w \leq g$ on $\operatorname{bd}(\Omega)$.

Proof. The function $-u$ is subharmonic and so is $v=w-u$ by Remark 43.20. Since $v=w-g \leq 0$ on $\operatorname{bd}(\Omega)$, it follows by Corollary 43.24 that $v \leq 0$ on $\Omega$, i.e. $w \leq g$ on $\Omega$.

### 43.2 Green's Functions

Notation 43.27 Unless otherwise stated, for the rest of this section assume $\Omega \subset \mathbb{R}^{n}$ is a compact manifold with $C^{2}$ - boundary.

For $x \in \Omega$, suppose there exists $h \in C^{2}\left(\Omega^{o}\right) \cap C^{1}(\Omega)$ which solves

$$
\begin{equation*}
\Delta h_{x}=0 \text { on } \Omega \text { with } h_{x}(y)=\phi(x-y) \text { for } y \in \partial \Omega . \tag{43.29}
\end{equation*}
$$

Hence if we define

$$
\begin{equation*}
G(x, y)=\phi_{x}(y)-h_{x}(y) \tag{43.30}
\end{equation*}
$$

then by the representation formula (Eq. (43.11) also see Remark 43.11) implies

$$
\begin{equation*}
u(x)=-\int_{\Omega} G(x, y) \triangle u(y) d y-\int_{\partial \Omega} \frac{\partial G}{\partial \mathbf{n}_{y}}(x, y) u(y) d \sigma(y) \tag{43.31}
\end{equation*}
$$

for all $u \in C^{2}\left(\Omega^{o}\right) \cap C^{1}(\Omega)$.
Throughout the rest of this subsection we will make the following assumption.
Assumption 2 (Solvability of Dirichlet Problem) We assume that for each $g \in C(\partial \Omega)$ there exists $h=h_{g} \in C^{2}\left(\Omega^{o}\right) \cap C^{1}(\Omega)$ such that

$$
\Delta h=0 \text { on } \Omega \text { with } h=g \text { on } \partial \Omega .
$$

In this case we define $G(x, y)$ as in $E q$. (43.30). We will (almost) verify that this assumption holds in Section 43.5 below. The full verification will come later when we study Hilbert space methods.
Theorem 43.28. Let $G(x, y)$ be given as in $E q$ (43.30). Then

1. $G(x, y)$ is smooth on $\left(\Omega^{o} \times \Omega^{o}\right) \backslash \triangle$ where $\triangle=\left\{(x, x): x \in \Omega^{o}\right\}$.
2. $G(x, y)=G(y, x)$ for all $x, y \in \Omega$. In particular the function $h(x, y):=$
$h_{x}(y)$ is symmetric in $x, y$ and $x \in \Omega^{o} \rightarrow h_{x} \in C(\Omega)$ is a smooth mapping.
3. If $\Omega$ is connected, then $G(x, y)>0$ for all $(x, y) \in\left(\Omega^{o} \times \Omega^{o}\right) \backslash \triangle$.

Proof. Let $\epsilon>0$ be small and $\Omega_{\epsilon}:=\Omega \backslash(B(x, \epsilon) \cup B(z, \epsilon))$ as in Figure 43.5, then by Green's theorem and the fact that $\Delta_{y} G(x, y)=0$ if $y \neq x$,

$$
\begin{aligned}
0 & =\int_{\Omega_{\epsilon}} \triangle_{y} G(x, y) G(z, y) d y \\
& =\int_{\partial \Omega_{\epsilon}}\left(\frac{\partial}{\partial \mathbf{n}} G(x, y) G(z, y)-G(x, y) \frac{\partial G}{\partial \mathbf{n}}(z, y)\right) d \sigma \\
& +\int_{\Omega_{\epsilon}} G(x, y) \triangle_{y} G(z, y) d y \\
& =\int_{\partial \Omega_{\epsilon}}\left(\frac{\partial}{\partial \mathbf{n}} G(x, y) G(z, y)-G(x, y) \frac{\partial G}{\partial \mathbf{n}}(z, y)\right) d \sigma
\end{aligned}
$$

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$$
\begin{aligned}
h(x, y) & =\int_{\Omega} h(\widetilde{x}, y) \eta(x-\widetilde{x}) d \widetilde{x} \\
& =\int_{\Omega \times \Omega} h(\widetilde{x}, \widetilde{y}) \eta(y-\widetilde{y}) \eta(x-\widetilde{x}) d \widetilde{x} d \widetilde{y}
\end{aligned}
$$

from which it follows that in fact $h$ is smooth on $\Omega \times \Omega$.
It only remains to show $x \rightarrow h_{x} \in C(\Omega)$ is smooth as well. Fix $x \in \Omega$ and for $v \in \mathbb{R}^{n}$, let $H_{v} \in C^{2}\left(\Omega^{o}\right) \cap C^{1}(\Omega)$ denote the solution to

$$
\Delta H_{v}=0 \text { on } \Omega \text { with } H_{v}(y)=v \cdot \nabla \phi(x-y) \text { for } y \in \partial \Omega
$$

Notice that $v \rightarrow H_{v}$ is linear and by the maximum principle,

$$
\begin{aligned}
\left\|h_{x+v}-h_{x}-H_{v}\right\|_{L^{\infty}(\Omega)} & \leq\left\|h_{x+v}-h_{x}-H_{v}\right\|_{L^{\infty}(\partial \Omega)} \\
& =\|\phi(x+v-\cdot)-\phi(x-\cdot)-v \cdot \nabla \phi(x-\cdot)\|_{L^{\infty}(\partial \Omega)}
\end{aligned}
$$

Now,

$$
\begin{aligned}
\phi(x+v-y) & -\phi(x-y)-v \cdot \nabla \phi(x-y) \\
& =\int_{0}^{1}[\nabla \phi(x+t v-y)-\nabla \phi(x-y)] \cdot v d t
\end{aligned}
$$

so that, by the dominated convergence theorem,

$$
\begin{aligned}
\| \phi(x+v-\cdot) & -\phi(x-\cdot)-v \cdot \nabla \phi(x-\cdot) \|_{L^{\infty}(\partial \Omega)} \\
& \leq|v| \int_{0}^{1}\|\nabla \phi(x+t v-\cdot)-\nabla \phi(x-\cdot)\|_{L^{\infty}(\partial \Omega)} d t=o(|v|) .
\end{aligned}
$$

This proves $x \rightarrow h_{x}$ is differentiable and that $\partial_{v} h_{x}=H_{v}$. Similarly one shows that $x \rightarrow h_{x}$ has higher derivatives as well.

For the last item, let $x \in \Omega^{\circ}$ and choose $\epsilon>0$ sufficiently small so that $\overline{B(x, \epsilon)} \subset \Omega^{o} \backslash\{y\}$ and $G(x, z)>0$ for all $z \in \overline{B(x, \epsilon)}$. Then the function $u(y):=G(x, y)$ is Harmonic on $\Omega^{0} \backslash \overline{B(x, \epsilon)}, u \in C(\Omega \backslash B(x, \epsilon)), u=0$ on $\partial \Omega$ and $u>0$ on $\partial \bar{B}(x, \epsilon)$. Hence by the maximum principle, $0 \leq u$ on $\Omega \backslash B(x, \epsilon)$ and since $u$ is not constant we must also have $u>0$ on $\Omega^{0} \backslash \overline{B(x, \epsilon)}$. Since $\epsilon>0$ was any sufficiently small number, it follows $G(x, y)>0$ for all $y \in \Omega^{o} \backslash\{x\}$.

Corollary 43.29. Keeping the above hypothesis and assuming $\rho \in C^{2}\left(\Omega^{o}\right) \cap$ $L^{1}(\Omega)$ and $g \in C(\partial \Omega)$, then there is (a necessarily unique) solution $u \in$ $C^{2}\left(\Omega^{o}\right) \cap C(\Omega)$ to

$$
\begin{equation*}
\Delta u=-\rho \text { with } u=g \text { on } \partial \Omega \tag{43.33}
\end{equation*}
$$

which is given by Eq. (43.31).
Proof. According to the remarks just before Eq. (43.31), if a solution to Eq. (43.33) exists it must be given by

$$
\begin{equation*}
u(x)=\int_{\Omega} G(x, y) \rho(y) d y-\int_{\partial \Omega} \frac{\partial G}{\partial \mathbf{n}_{y}}(x, y) g(y) d \sigma(y) \tag{43.34}
\end{equation*}
$$

From Assumption 2, there exists a solution $v \in C^{2}(\Omega) \cap C^{1}(\bar{\Omega})$ such that $\Delta v=0$ and $v=g$ on $\partial \Omega$. So replacing $u$ by $u-v$ if necessary, it suffices to prove there is a solution $u \in C^{2}\left(\Omega^{\circ}\right) \cap C(\Omega)$ such that Eq. (43.33) holds with $g \equiv 0$. To produce this solution, let

$$
u(x):=\int_{\Omega} G(x, y) \rho(y) d y=\int_{\Omega} \phi(x-y) \rho(y) d y-H(x)
$$

where

$$
H(x):=\int_{\Omega} h(x, y) \rho(y) d y
$$

Using the result in Theorem 43.28, one easily shows $H \in C^{\infty}\left(\Omega^{o}\right) \cap C^{1}(\Omega)$ and $\Delta H=0$. By Theorem 43.9,

$$
\Delta_{x} \int_{\Omega} \phi(x-y) \rho(y) d y=-\rho(x) \text { for } x \in \Omega
$$

and therefore $u \in C^{2}(\Omega)$ and $\Delta u=-\rho$.
Remark 43.30. Because of the maximum principle, for any $x \in \Omega$ the map $g \in C(\partial \Omega) \rightarrow h_{g}(x) \in C(\bar{\Omega})$ is a positive linear functional. So by the Riesz representation theorem, there exists a unique positive probability measure $\sigma_{x}$ on $\partial \Omega$ such that

$$
h_{g}(x)=\int_{\partial \Omega} g(y) d \sigma_{x}(y) \text { for all } g \in C(\partial \Omega)
$$

Evidently this measure is given by

$$
d \sigma_{x}(y)=-\frac{\partial G}{\partial \mathbf{n}_{y}}(x, y) d \sigma(y)
$$

and in particular $-\frac{\partial G}{\partial \mathbf{n}_{y}}(x, y) \geq 0$ for all $x \in \Omega$ and $y \in \partial \Omega$. It is in fact easy to see that $-\frac{\partial G}{\partial \mathbf{n}_{y}}(x, y)>0$ for all $x \in \Omega$ and $y \in \partial \Omega$.

### 43.3 Explicit Green's Functions and Poisson Kernels

In this section we will use the method of images to construct explicit formula for the Green's functions and Poisson Kernels for the half plane ${ }^{1}, \overline{\mathbb{H}}^{n}=\{x \in$

[^4]$\left.\mathbb{R}^{n}: x_{n} \geq 0\right\}$ and Balls $\overline{B(0, a)}$. For $x=\left(x^{\prime}, z\right) \in \mathbb{R}^{n-1} \times(0, \infty)=\mathbb{H}^{n}$ let $R x:=\left(x^{\prime},-z\right)$. It is simple to verify $|x-y|=|R x-y|$ for all $x \in \mathbb{H}^{n}$ and $y \in \partial \overline{\mathbb{H}}^{n}$. Form this and the properties of $\phi$, one concluded, for $x \in \mathbb{H}^{n}$, that $h_{x}(y):=\phi(y-R x)$ is Harmonic in $y \in \mathbb{H}^{n}$ and $h_{x}(y)=\phi(x-y)$ for all $y \in \partial \bar{H}^{n}$. These remarks give rise to the following theorem.
Theorem 43.31. For $x, y \in \overline{\mathbb{H}}^{n}$, let
$$
G(x, y):=\phi(y-x)-\phi(y-R x)=\phi(y-x)-\phi(R y-x)
$$

Then $G$ is the Greens function for $\Delta$ on $\mathbb{H}^{n}$ and

$$
K(x, y):=-\frac{\partial G}{\partial \mathbf{n}}(x, y)=\frac{2 x_{n}}{\sigma\left(S^{n-1}\right)} \frac{1}{|x-y|^{n}} \text { for } x \in \mathbb{H}^{n} \text { and } y \in \partial \overline{\mathbb{H}}^{n}
$$

is the Poisson kernel for $\mathbb{H}^{n}$. Furthermore if $\rho \in C^{2}\left(\mathbb{H}^{n}\right) \cap L^{1}\left(\mathbb{H}^{n}\right)$ and $f \in$ $B C\left(\partial \overline{\mathbb{H}}^{n}\right)$, then

$$
u(x)=\int_{\mathbb{H}^{n}} G(x, y) \rho(y) d y+\int_{\partial \overline{\mathbb{H}}^{n}} K(x, y) f(y) d \sigma(y)
$$

solves the equation

$$
\Delta u=-\rho \text { on } \mathbb{H}^{n} \text { with } u=f \text { on } \partial \overline{\mathbb{H}}^{n} .
$$

Proof. First notice that

$$
G(y, x)=\phi(x-y)-\phi(x-R y)=\phi(x-y)-\phi(R x-R R y)=G(x, y)
$$

since $\phi$ is a function of $|\cdot|$. Therefore, if

$$
u(x)=\int_{\mathbb{H}^{n}} G(x, y) \rho(y) d y=\int_{\mathbb{H}^{n}} \phi(x-y) \rho(y) d y-\int_{\mathbb{H}^{n}} \phi(x-R y) \rho(y) d y
$$

we have from Theorem 43.9 that

$$
\Delta u(x)=-\rho(x)-\int_{\mathbb{H}^{n}} \Delta_{x} \phi(x-R y) \rho(y) d y=-\rho(x) .
$$

Since $G(x, y)=0$ for $x \in \partial \overline{\mathbb{H}}^{n}$ and so $u(x)=0$ for $x \in \partial \bar{H}^{n}$. It is left to the reader to show $u$ is continuous on $\overline{\mathbb{H}^{n}}$.

For $x \in \mathbb{H}^{n}$ and $y \in \partial \overline{\mathbb{H}}^{n}$, we find form Eq. (43.7),

$$
\begin{aligned}
K(x, y) & :=-\frac{\partial G}{\partial \mathbf{n}_{y}}(x, y)=\frac{\partial}{\partial y_{n}} G(x, y) \\
& =\frac{\partial}{\partial y_{n}}[\phi(y-x)-\phi(y-R x)] \\
& =-\frac{1}{\sigma\left(S^{n-1}\right)} \frac{1}{|y-x|^{n}}(y-x) \cdot e_{n} \\
& +\frac{1}{\sigma\left(S^{n-1}\right)} \frac{1}{|y-R x|^{n}}(y-R x) \cdot e_{n} \\
& =\frac{1}{\sigma\left(S^{n-1}\right)} \frac{2 x_{n}}{|y-x|^{n}} .
\end{aligned}
$$

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$$
\lim _{M \uparrow \infty} \int_{\partial \overline{\mathbb{H}}^{n}} K(x, y) u(y / M) d \sigma(y)=\int_{\partial \overline{\mathbb{H}}^{n}} K(x, y) d \sigma(y)
$$

and therefore passing the limit in Eq. (43.35) gives

$$
1=\int_{\partial \overline{\mathbb{H}}^{n}} K(x, y) d \sigma(y)-\lim _{M \uparrow \infty}\left[M^{-2} \int_{\mathbb{H}^{n}} G(x, y) \triangle u(y / M) d y\right]
$$

This latter limit is zero, since

$$
\begin{aligned}
& M^{-2} \int_{\mathbb{H}^{n}} G(x, y) \triangle u(y / M) d y \\
& \quad=c_{n} M^{-2} \int_{\mathbb{H}^{n}}\left[\frac{1}{|x-y|^{n-2}}-\frac{1}{|R x-y|^{n-2}}\right](\triangle u)(y / M) d y \\
& \quad=c_{n} M^{-2} M^{n} \int_{\mathbb{H}^{n}}\left[\frac{1}{|x-M y|^{n-2}}-\frac{1}{|R x-M y|^{n-2}}\right] \triangle u(y) d y \\
& \quad=c_{n} \int_{\mathbb{H}^{n}}\left[\frac{1}{|x / M-y|^{n-2}}-\frac{1}{|R x / M-y|^{n-2}}\right] \triangle u(y) d y .
\end{aligned}
$$

This latter expression tends to zero and $M \rightarrow \infty$ by the dominated convergence and this proves the claim. (Alternatively, for $y$ large,

$$
\begin{aligned}
\frac{1}{|x-y|^{n-2}} & -\frac{1}{|R x-y|^{n-2}} \\
& =\frac{1}{|y|^{n-2}}\left[\frac{1}{\left|\frac{x}{|y|}-\hat{y}\right|^{n-2}}-\frac{1}{\left|R \frac{x}{|y|}-\hat{y}\right|^{n-2}}\right] \\
& =\frac{1}{|y|^{n-2}}\left[\left(1+2 \frac{x}{|y|} \cdot \hat{y}+\ldots\right)-\left(1+2 \frac{R x}{|y|} \cdot \hat{y}+\ldots\right)\right] \\
& =O\left(\frac{1}{|y|^{n-1}}\right)
\end{aligned}
$$

and therefore

$$
\begin{aligned}
& M^{-2} \int_{\mathbb{H}^{n}}\left[\frac{1}{|x-y|^{n-2}}-\frac{1}{|R x-y|^{n-2}}\right](\Delta u)(y / M) d y \\
& \quad=O\left(M^{-2} \frac{1}{M^{n-1}} M^{n}\right)=O(1 / M) \rightarrow 0
\end{aligned}
$$

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$$
0=r^{2} \beta^{4}-\left(|x|^{2}+r^{2}\right) \beta^{2}+|x|^{2}
$$

as $M \rightarrow \infty$.
Since $G(x, y)$ is harmonic in $x$, it follows that $K(x, y)=-\frac{\partial}{\partial \mathbf{n}_{y}} G(x, y)$ is still Harmonic in $x$. and therefore

$$
u(x):=\int_{\partial \mathbb{H}^{n}} K(x, y) f(y) d \sigma(y)=\frac{2}{\sigma\left(S^{n-1}\right)} \int_{\partial \mathbb{H}^{n}} \frac{x_{n}}{|x-y|^{n}} f(y) d \sigma(y)
$$

is harmonic as well. Since

$$
\begin{align*}
u(x) & =\frac{2}{\sigma\left(S^{n-1}\right)} \int_{\partial \overline{\mathbb{H}}^{n}} \frac{x_{n}}{\left(\left|x^{\prime}-y\right|^{2}+x_{n}^{2}\right)^{n / 2}} f(y) d y  \tag{43.36}\\
& =\frac{2}{\sigma\left(S^{n-1}\right)} \frac{1}{x_{n}^{n-1}} \int_{\partial \overline{\mathbb{H}}^{n}} \frac{1}{\left(\left|\frac{x^{\prime}-y}{x_{n}}\right|^{2}+1\right)^{n / 2}} f(y) d y
\end{align*}
$$

it follows from Theorem 42.13 that $u\left(\left(x^{\prime}, x_{n}\right)\right) \rightarrow f\left(x^{\prime}\right)$ as $x_{n} \downarrow 0$ uniformly for $x^{\prime}$ in compact subsets of $\partial \overline{\mathbb{H}}^{n}$.

### 43.4 Green's function for Ball

Let $r>0$ be fixed, we will construct the Green's function for the ball $B(0, r)$. The idea for a given $x \in B(0, r)$, we should find a mirror location, say $\rho \hat{x}$ and a charge $q$ so that

$$
\phi(x-y)=q \phi(\rho \hat{x}-y) \text { for all }|y|=r .
$$

Assuming for the moment that $n \geq 3$ and writing $q=\beta^{(2-n)}$, this leads to the equations

$$
|x-y|^{2}=|\beta \rho \hat{x}-\beta y|^{2}=\beta^{2}|\rho \hat{x}-y|^{2}
$$

or equivalently squaring out both sides and using $|y|=r$,

$$
|x|^{2}-2 x \cdot y+r^{2}=\beta^{2}\left(\rho^{2}-2 \rho \hat{x} \cdot y+r^{2}\right)
$$

Choosing $y \perp x$ and $y=r \hat{x}$ leads to the conditions

$$
\begin{aligned}
|x|^{2}+r^{2} & =\beta^{2}\left(\rho^{2}+r^{2}\right) \text { and } \\
|x|^{2}-2 r|x|+r^{2} & =\beta^{2}\left(\rho^{2}-2 \rho r+r^{2}\right) .
\end{aligned}
$$

Subtracting these two equations implies $-2 r|x|=-2 \rho \beta^{2} r$ or equivalently that $\rho=|x| / \beta^{2}$. Putting this into the first equation above then implies

$$
|x|^{2}+r^{2}=\frac{|x|^{2}}{\beta^{2}}+\beta^{2} r^{2}
$$

or equivalently that

By the quadratic formula, this implies

$$
\begin{aligned}
\beta^{2} & =\frac{\left(|x|^{2}+r^{2}\right) \pm \sqrt{\left(|x|^{2}+r^{2}\right)^{2}-4 r^{2}|x|^{2}}}{2 r^{2}} \\
& =\frac{\left(|x|^{2}+r^{2}\right) \pm \sqrt{\left(|x|^{2}-r^{2}\right)^{2}}}{2 r^{2}}=\frac{\left(|x|^{2}+r^{2}\right) \pm\left(r^{2}-|x|^{2}\right)}{2 r^{2}} \\
& =1 \text { or } \frac{|x|^{2}}{r^{2}}
\end{aligned}
$$

Clearly the charge $\beta=1$ will not work so we must take $\beta=|x| / r$ in which case, $\rho=r^{2} /|x|$ and hence

$$
q \phi(\rho \hat{x}-y)=(|x| / r)^{(2-n)} \phi\left(r^{2} \frac{\hat{x}}{|x|}-y\right)=\phi\left(r \hat{x}-\frac{|x|}{r} y\right) .
$$

Let us now verify that our guess has worked. Let us begin by noting the following identities for $x, y \in \mathbb{R}^{n}$,

$$
\begin{equation*}
\left|r \hat{x}-r^{-1}\right| x|y|^{2}=\left(r^{2}-2 x \cdot y+r^{-2}|x|^{2}|y|^{2}\right) \tag{43.37}
\end{equation*}
$$

and in particular when $|y|=r$ this implies

$$
|\hat{x} r-|x| \hat{y}|^{2}=\left(r^{2}-2 x \cdot y+|x|^{2}\right)=|x-y|^{2}
$$

so that

$$
\begin{equation*}
|x-y|=|\hat{x} r-|x| \hat{y}|=\left|\frac{x}{|x|} r-|x| \frac{y}{r}\right| . \tag{43.38}
\end{equation*}
$$

Now the function

$$
h_{x}(y)=\phi\left(r \hat{x}-\frac{|x|}{r} y\right)=\phi\left(\frac{|x|}{r}\left(y-r^{2} \frac{x}{|x|^{2}}\right)\right)
$$

is harmonic in $y$ and by Eq. (43.38),

$$
h_{x}(y)=\phi\left(\hat{x} r-|x| \frac{y}{r}\right)=\phi(\hat{x} r-|x| \hat{y})=\phi(x-y) \text { when }|y|=r \text {. }
$$

Hence we should define the Green's function for the ball to be given by

$$
\begin{aligned}
G(x, y) & =\phi(x-y)-h_{x}(y)=\phi(x-y)-\phi\left(\hat{x} r-|x| \frac{y}{r}\right) \\
& =\phi(x-y)-\phi\left(\hat{x} r-r^{-1}|x||y| \hat{y}\right) \\
& =\phi(x-y)-\phi\left(\frac{|x|}{r}\left(y-r^{2} \frac{\hat{x}}{|x|}\right)\right) .
\end{aligned}
$$

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From Eq. (43.37), it follows that $h_{x}(y)=h_{y}(x)$ and therefore $G(x, y)$ is again symmetric under the interchange of $x$ and $y$.

For $y \in \partial \overline{B(0, r)}$, using Eq. (43.7) we find

$$
\begin{aligned}
-K(x, y) & =\frac{\partial G}{\partial \mathbf{n}_{y}}(x, y)=\partial_{\hat{y}} G(x, y)=\nabla_{y} G(x, y) \cdot \hat{y} \\
& =\nabla_{y}\left[\phi(x-y)-\phi\left(\frac{|x|}{r}\left(y-r^{2} \frac{\hat{x}}{|x|}\right)\right)\right] \cdot \hat{y} \\
& =-\frac{1}{\sigma\left(S^{n-1}\right)}\left[\frac{1}{|y-x|^{n}}(y-x)-\frac{|x|}{r} \frac{1}{\left(\frac{|x|}{r}\right)^{n}\left|y-r^{2} \frac{\hat{x}}{|x|}\right|^{n}} \frac{|x|}{r}\left(y-r^{2} \frac{\hat{x}}{|x|},\right.\right. \\
& =-\frac{1}{\sigma\left(S^{n-1}\right)}\left[\frac{1}{|y-x|^{n}}(y-x)-\frac{|x|}{r} \frac{1}{\left|\frac{|x|}{r} y-r \hat{x}\right|^{n}}\left(\frac{|x|}{r} y-r \hat{x}\right)\right] \cdot \hat{y} \\
& =-\frac{1}{\sigma\left(S^{n-1}\right)}\left[\frac{1}{|y-x|^{n}}(y-x)-\frac{|x|}{r} \frac{1}{|x-y|^{n}}(|x| \hat{y}-r \hat{x})\right] \cdot \hat{y} \\
& =-\frac{1}{\sigma\left(S^{n-1}\right)|y-x|^{n}}\left[(y-x)-\left(\frac{|x|^{2}}{r} \hat{y}-x\right)\right] \cdot \hat{y} \\
& =-\frac{1}{\sigma\left(S^{n-1}\right)|y-x|^{n}}\left[y-\frac{|x|^{2}}{r} \hat{y}\right] \cdot \hat{y} \\
& =-\frac{1}{\sigma\left(S^{n-1}\right) r|y-x|^{n}}\left[r^{2}-|x|^{2}\right] .
\end{aligned}
$$

These computations lead to the following theorem.
Theorem 43.32. For $x, y \in B(0, r)$, let

$$
G(x, y):=\phi(x-y)-\phi\left(\hat{x} r-|x| \frac{y}{r}\right)
$$

and if $y \in \partial \overline{B(0, r)}$, let

$$
K(x, y):=-\frac{\partial G}{\partial \mathbf{n}}(x, y)=\frac{r^{2}-|x|^{2}}{\sigma\left(S^{n-1}\right) r}|x-y|^{-n}
$$

Then $\rho \in C^{2}(B(0, r)) \cap L^{1}(B(0, r))$ and $f \in C(\partial \overline{B(0, r)})$, then

$$
\begin{equation*}
u(x)=\int_{B(0, r)} G(x, y) \rho(y) d y+\int_{\partial \overline{B(0, r)}} K(x, y) f(y) d \sigma(y) \tag{43.39}
\end{equation*}
$$

solves the equation

$$
\Delta u=-\rho \text { on } B(0, r) \text { with } u=f \text { on } \partial \overline{B(0, r)} .
$$

Remark 43.33. Letting $\rho=|x|$, we may write $K(x, y)$ as

$$
\begin{equation*}
K(x, y)=\frac{r^{2}-\rho^{2}}{\sigma\left(S^{n-1}\right) r} \frac{1}{\left(\rho^{2}+r^{2}-2 \rho r \hat{x} \cdot \hat{y}\right)^{n / 2}} \tag{43.40}
\end{equation*}
$$

In particular when $r=1, n=2, x=\rho e^{i \theta}$ and $y=r e^{i \alpha}$, this gives

$$
K(x, y)=\frac{1-\rho^{2}}{2 \pi\left(\rho^{2}+1-2 \rho \cos (\theta-\alpha)\right)}
$$

which agrees with the Poisson kernel $P_{\rho}(\theta-\alpha)$ of Eq. (22.40) which we derived earlier by Fourier series methods.

Proof. The proof is essentially the same as Theorem 43.31 but a bit easier. For these reasons, we will only prove here the assertion

$$
\begin{equation*}
\lim _{x \rightarrow x_{0}} \int_{\partial \overline{B(0, r)}} K(x, y) f(y) d \sigma(y)=f\left(x_{0}\right) \text { for all } x_{0} \in \partial B(0, r) \tag{43.41}
\end{equation*}
$$

From Theorem 43.5 with $u=1$ it follows again that

$$
\int_{\partial \overline{B(0, r)}} K(x, y) d \sigma(y)=1
$$

For any $\alpha \in(0,1)$ let

$$
\epsilon(\alpha):=\sup \left\{|f(r y)-f(r x)|: y, x \in S^{n-1} \text { with } \hat{y} \cdot \hat{x} \leq \alpha\right\} .
$$

Then by uniform continuity of $f$ on $\partial B(0, r)$, it follows that $\epsilon(\alpha) \rightarrow 0$ as $\alpha \uparrow 1$ and hence

$$
\begin{aligned}
\mid \int_{\partial \overline{B(0, r)}} & K(x, y) f(y) d \sigma(y)-f(r \hat{x}) \mid \\
& \leq \int_{\partial \overline{B(0, r)}} K(x, y)|f(y)-f(r \hat{x})| d \sigma(y) \\
& =\int_{\partial \overline{B(0, r)}} 1_{\hat{y} \cdot \hat{x}>\alpha} K(x, y)|f(y)-f(r \hat{x})| d \sigma(y) \\
& +\int_{\partial \overline{B(0, r)}} 1_{\hat{y} \cdot \hat{x} \leq \alpha} K(x, y)|f(y)-f(r \hat{x})| d \sigma(y) \\
& \leq 2\|f\|_{u} \int_{\partial \overline{B(0, r)}} 1_{\hat{y} \cdot \hat{x}>\alpha} K(x, y) d \sigma(y)+\epsilon(\alpha) \\
& \leq C(\alpha)\|f\|_{u}\left(r^{2}-|x|^{2}\right)+\epsilon(\alpha)
\end{aligned}
$$

where $C(\alpha)$ is some constant only depending on $\alpha$, see Eq. (43.40). Therefore,

$$
\sup _{|x|=\rho}\left|\int_{\partial \frac{B(0, r)}{}} K(x, y) f(y) d \sigma(y)-f(r \hat{x})\right| \leq C(\alpha)\|f\|_{u}\left(r^{2}-\rho^{2}\right)+\epsilon(\alpha)
$$

and hence

$$
\lim _{\rho \uparrow r} \sup \sup _{|x|=\rho}\left|\int_{\partial \overline{B(0, r)}} K(x, y) f(y) d \sigma(y)-f(r \hat{x})\right| \leq \epsilon(\alpha) \rightarrow 0 \text { as } \alpha \downarrow 0
$$

which implies Eq. (43.41).
Corollary 43.34. Suppose that $u$ is a harmonic function on $\Omega$, then $u$ is real analytic on $\Omega$.

Proof. The condition of being real analytic is local and invariant under translations as is the notion of being harmonic. Hence we may assume $0 \in$ $\overline{B(0, r)} \subset \Omega$ for some $r>0$, in which case we have, for $|x|<r$ and $f=$ $\left.u\right|_{\partial \overline{B(0, r)}}$, that

$$
\begin{align*}
u(x) & =\int_{\partial \overline{B(0, r)}} K(x, y) f(y) d \sigma(y)=\frac{r^{2}-|x|^{2}}{\sigma\left(S^{n-1}\right) r} \int_{\partial \overline{B(0, r)}}|x-y|^{-n} f(y) d \sigma(y) \\
& =\frac{r^{2}-|x|^{2}}{\sigma\left(S^{n-1}\right) r} \int_{\partial \overline{B(0, r)}}|x-\hat{y} r|^{-n} f(y) d \sigma(y) \tag{43.42}
\end{align*}
$$

Now

$$
\begin{aligned}
|x-y|^{-n} & =|x-\hat{y} r|^{-n}=r^{-n}\left|\hat{y}-r^{-1} x\right|^{-n} \\
& =r^{-n}\left(1-2 r^{-1} \hat{y} \cdot x+\frac{|x|^{2}}{r^{2}}\right)^{-n / 2} \\
& =r^{-n}(1-\alpha(x, y))^{-n / 2}
\end{aligned}
$$

where

$$
\alpha(x, y):=2 r^{-1} \hat{y} \cdot x-\frac{|x|^{2}}{r^{2}} .
$$

Since

$$
|\alpha(x, y)| \leq 2 r^{-1}|x|+\frac{|x|^{2}}{r^{2}} \leq 2 \alpha_{0}+\alpha_{0}^{2}<1
$$

if $|x| \leq \alpha_{0} r$ and $\alpha_{0}<\sqrt{2}-1$, we find that $|x-y|^{-n}$ has a convergent power series expansion,

$$
|x-y|^{-n}=r^{-n} \sum_{m=0}^{\infty} a_{m} \alpha(x, y)^{m} \text { for }|x| \leq \alpha_{0} r
$$

Plugging this into Eq. (43.42) shows $u(x)$ has a convergent power series expansion in $x$ for $|x| \leq(\sqrt{2}-1) r$.

### 43.5 Perron's Method for solving the Dirichlet Problem

For this section let $\Omega \subset_{o} \mathbb{R}^{n}$ be a bounded open set and $f \in C(\operatorname{bd}(\Omega), \mathbb{R})$ be a given function. We are going to investigate the solvability of the Dirichlet problem:

$$
\begin{equation*}
\Delta u=0 \text { on } \Omega \text { with } u=f \text { on } \operatorname{bd}(\Omega) . \tag{43.43}
\end{equation*}
$$

Let $\mathcal{S}(\Omega)$ denote those $w \in C(\bar{\Omega})$ such that $w$ is subharmonic on $\Omega$ and let $\mathcal{S}_{f}(\Omega)$ denote those $w \in \mathcal{S}(\Omega)$ such that $w \leq f$ on $\operatorname{bd}(\Omega)$. As we have seen in Corollary 43.26, if there is a solution to $u \in C^{2}(\Omega) \cap C(\bar{\Omega})$, then $w \leq u$ for all $w \in \mathcal{S}_{f}(\Omega)$. This suggests we try to define

$$
\begin{equation*}
u(x):=u_{f}(x):=\sup \left\{w(x): w \in \mathcal{S}_{f}(\Omega)\right\} \text { for all } x \in \bar{\Omega} . \tag{43.44}
\end{equation*}
$$

Notation 43.35 Given $w \in \mathcal{S}(\Omega), \xi \in \Omega$ and $r>0$ such that $\overline{B(\xi, r)} \subset \Omega$, let (see Figure 43.6)

$$
w_{\xi, r}(y)=\left\{\begin{array}{l}
w(y) \text { for } y \in \Omega \backslash B(\xi, r) \\
h(y) \text { for } \quad y \in \overline{B(\xi, r)}
\end{array}\right.
$$

where $h \in C(\overline{B(\xi, r)})$ is the unique solution to

$$
\Delta h=0 \text { on } B(\xi, r) \text { with } h=w \text { on } \partial B(\xi, r) .
$$

The existence of $h$ is guaranteed by Theorem 43.32.



Fig. 43.6. The construction of $w_{\xi, r}$ in the one-dimensional case.

Proposition 43.36. Let $w \in \mathcal{S}(\Omega)$ and $w_{\xi, r}$ be as above. Then

1. $w \leq w_{\xi, r}$.
2. $w_{\xi, r} \in \mathcal{S}(\Omega)$, i.e. $w_{\xi, r}$ is subharmonic.
3. For any $\xi \in \Omega$ and $r>0$ such that $\overline{B(\xi, r)} \subset \Omega$, the mean value inequality is valid,

$$
w(\xi) \leq \int_{\partial B(\xi, r)} w d \sigma
$$

Proof. 1. Since $w=w_{\xi, r}$ on $\Omega \backslash B(\xi, r)$, it suffices to show $w \leq h$ on $B(\xi, r)$. But this follows from Corollary 43.26.
2. Since $w_{\xi, r}$ is harmonic on $B(\xi, r)$ and subharmonic on $\Omega \backslash \overline{B(\xi, r)}$, we need only show

$$
w_{\xi, r}(y) \leq \int_{\partial B(y, \rho)} w_{\xi, r} d \sigma
$$

for all $y \in \partial B(\xi, r)$ and $\rho$ sufficiently small. This is easily checked, since $w$ is subharmonic,

$$
w_{\xi, r}(y)=w(y) \leq \int_{\partial B(y, \rho)} w d \sigma \leq \int_{\partial B(y, \rho)} w_{\xi, r} d \sigma
$$

wherein the last equality we made use of Item 1.
3. By item 1. and the mean value property for the harmonic function, $w_{\xi, r}$, we have

$$
w(\xi) \leq w_{\xi, r}(\xi)=\int_{\partial B(\xi, r)} w_{\xi, r} d \sigma=\int_{\partial B(\xi, r)} w d \sigma
$$

Theorem 43.37. The function $u=u_{f}$ defined in Eq. (43.44) is harmonic on $\Omega$ and $u \leq f$ on $\operatorname{bd}(\Omega)$.

Proof. Let us begin with a couple of observations. In what follows

$$
m:=\min \{f(x): x \in \operatorname{bd}(\Omega)\} \text { and } M:=\min \{f(x): x \in \operatorname{bd}(\Omega)\}
$$

1. The function $u=u_{f} \geq m$ on $\Omega$ since $m \in \mathcal{S}_{f}(\Omega)$.
2. By the maximum principle $w \leq M$ on $\Omega$ for all $w \in \mathcal{S}_{f}(\Omega)$ and therefore $u_{f} \leq M$ on $\Omega$.
3. If $w_{1}, \ldots, w_{m} \in \mathcal{S}_{f}(\Omega)$, then $w=\max \left\{w_{1}, \ldots, w_{m}\right\} \in \mathcal{S}_{f}(\Omega)$. Indeed for $\xi \in \Omega$ and $r$ small,

$$
f_{\partial B(\xi, r)} w d \sigma \geq \int_{\partial B(\xi, r)} w_{i} d \sigma \geq w_{i}(\xi)
$$

for all $i$.
4. Now suppose $\xi \in \Omega$ and $R>0$ be chosen so that $\overline{B(\xi, R)} \subset \Omega$ and $D \subset B(\xi, R)$ is a countable set. Then there is a harmonic function $w_{D}$ on $B(\xi, R)$ such that $w_{D}=u_{f}$ on $D$.
To prove this last item let $D:=\left\{y_{k}\right\}_{k=1}^{\infty}$ and choose $\left\{w_{k}^{m}\right\} \subset \mathcal{S}_{f}(\Omega)$ such that $w_{k}^{m}\left(y_{k}\right) \rightarrow u\left(y_{k}\right)$ as $m \rightarrow \infty$ for each $k$. By replacing $w_{k}^{m}$ by $\max \left\{w_{k}^{1}, \ldots, w_{k}^{m}\right\}$ if necessary we may assume for each $k$ that $w_{k}^{m}$ is increasing in $m$ for each $k$. Letting

$$
W_{m}:=\max \left\{w_{1}^{m}, \ldots, w_{m}^{m}\right\}
$$

we find an increasing sequence $\left\{W_{m}\right\} \subset \mathcal{S}_{f}(\Omega)$ such that $W_{m}(y) \uparrow u_{f}(y)$ for all $y \in D$. Finally define a sequence $\left\{w_{m}\right\} \subset \mathcal{S}_{f}(\Omega)$ by $w_{m}:=\left(W_{m}\right)_{\xi, 2 R}$. By the maximum principle, $w_{m}$ is still increasing and since $W_{m} \leq w_{m}$ and we still have $w_{m}(y) \uparrow u_{f}(y)$ for all $y \in D$. We now define $w_{D}:=\left.\lim _{m \rightarrow \infty} w_{m}\right|_{B(\xi, R)}$ which exists because $w_{m}$ is increasing and $w$. We have $w_{D}=u_{f}$ on $D$ and because $\left\{w_{m}\right\}$ is a bounded and convergent sequence of harmonic functions on $B(\xi, R)$, it follows from Corollary 43.18 that $w$ is harmonic on $B(\xi, R)$. This completes the proof of item 4.

We now use item 4. to prove $u_{f}$ is continuous at $\xi \in \Omega$. To do this let $\left\{y_{k}\right\}_{k=1}^{\infty} \subset B(\xi, R)$ be any sequence such that $y_{k} \rightarrow \xi$ as $k \rightarrow \infty$ and let $D=\{\xi\} \cup\left\{y_{k}\right\}_{k=1}^{\infty} \subset B(\xi, R)$. Since $w_{D}$ is harmonic and hence continuous,

$$
\lim _{k \rightarrow \infty} u_{f}\left(y_{k}\right)=\lim _{k \rightarrow \infty} w_{D}\left(y_{k}\right)=w_{D}(\xi)=u_{f}(\xi)
$$

showing $u_{f}$ is continuous.
To show $u_{f}$ is harmonic on $B(\xi, R)$, let $D$ be a countable dense subset of $B(\xi, R)$. Then the continuity of $u_{f}$ and the fact that $u_{f}=w_{D}$ on $D$, it follows that $u_{f}=w_{D}$ on $B(\xi, R)$. In particular $u_{f}$ is harmonic on $B(\xi, R)$. Since $\xi$ is arbitrary, we have shown $u_{f}$ is harmonic.

To complete our program, we want to show that $u_{f}$ extends to a function in $C(\bar{\Omega})$ and that $u_{f}=f$ on $\operatorname{bd}(\Omega)$. For this we will need some assumption on $\operatorname{bd}(\Omega)$.

Definition 43.38. A function $Q \in C(\bar{\Omega})$ is a barrier function for $\eta \in$ $\operatorname{bd}(\Omega)$ if $Q$ is subharmonic on $\Omega, Q(\eta)=0$ and $Q(x)<0$ for all $x \in \operatorname{bd}(\Omega) \backslash$ $\{\eta\}$.
Example 43.39. Suppose that $\eta \in \operatorname{bd}(\Omega)$ and there exists $\xi \in \mathbb{R}^{n}$ such that $(x-\eta) \cdot \xi<0$ for all $x \in \operatorname{bd}(\Omega) \backslash\{\eta\}$ (see Figure 43.7 below), then the function $Q(x):=(x-\eta) \cdot \xi$ is a barrier function of $\eta$.

Example 43.40. Suppose that $\eta \in \operatorname{bd}(\Omega)$ and there exists a ball $\overline{B(\xi, r)} \cap \bar{\Omega}=$ $\{\eta\}$ (see Figure 43.8), then $Q(x):=\alpha(r)-\alpha(|x-\xi|)$ is a barrier function for $\eta$, where $\alpha$ is defined in Eq. 43.4.

Theorem 43.41. Suppose $f \in C(\operatorname{bd}(\Omega))$ and $u=u_{f}$ is the harmonic function defined by $E q$. (43.44) and there exists a barrier function $Q$ for $\eta \in \operatorname{bd}(\Omega)$. Then $\lim _{x \rightarrow \eta} u_{f}(x)=f(\eta)$. In particular if every point $\eta \in \operatorname{bd}(\Omega)$ admits a barrier function, then there is a unique solution $u \in C(\bar{\Omega}) \cap C^{2}(\Omega)$ to $\Delta u=0$ with $u=f$ on $\operatorname{bd}(\Omega)$.

Proof. Given $\epsilon>0$ and $K>0$, let $w(x):=f(\eta)-\epsilon-K Q(x)$ for all $x \in \bar{\Omega}$. For any $\epsilon>0$ we may choose (using continuity of $f$ and compactness

## 43 Poisson and Laplace's Equation



Fig. 43.7. Constructing a barrier function at point where $\eta$ where $\partial \Omega$ lies in a half plane.


Fig. 43.8. Another $\eta$ for which there exists a barrier function.
of $\operatorname{bd}(\Omega)) K$ sufficiently large so that $w \leq f$ on $\operatorname{bd}(\Omega)$, i.e. $w \in \mathcal{S}_{f}(\Omega)$. Therefore $w \leq u_{f}$ and hence

$$
f(\eta)-\epsilon=w(\eta)=\lim _{x \rightarrow \eta} w(x) \leq \liminf _{x \rightarrow \eta} u_{f}(x) .
$$

Since $\epsilon>0$ is arbitrary, this shows

$$
\begin{equation*}
\liminf _{x \rightarrow \eta} u_{f}(x) \geq f(\eta) . \tag{43.45}
\end{equation*}
$$

We now consider the function

$$
\begin{align*}
-u_{-f}(x) & =-\sup \left\{w(x): w \in \mathcal{S}_{-f}(\Omega)\right\}=\inf \left\{-w(x): w \in \mathcal{S}_{-f}(\Omega)\right\} \\
& =\inf \left\{W(x):-W \in \mathcal{S}_{-f}(\Omega)\right\} \tag{43.46}
\end{align*}
$$

If $w \in \mathcal{S}_{f}(\Omega)$ and $-W \in \mathcal{S}_{-f}(\Omega)$, then $w-W$ is sub harmonic and $w-W \leq$ $f-f=0$ on $\operatorname{bd}(\Omega)$, therefore by the maximum principle it follows that $w \leq W$ on $\bar{\Omega}$. Coupling this fact with Eq. (43.46) shows $-u_{-f}(x) \geq w(x)$ for all $w \in \mathcal{S}_{f}(\Omega)$ and then taking the supremum over $w$ shows $u_{f}(x) \leq-u_{-f}(x)$. Therefore using Eq. (43.45) with $f$ replaced by $-f$ shows

$$
\limsup _{x \rightarrow \eta} u_{f}(x) \leq \limsup _{x \rightarrow \eta}\left(-u_{-f}(x)\right)=-\liminf _{x \rightarrow \eta}\left(u_{-f}(x)\right) \leq-(-f(\eta))=f(\eta)
$$

which combined with Eq. (43.45) shows

$$
\lim _{x \rightarrow \eta} u_{f}(x)=f(\eta)
$$

Exercise 43.42. Suppose that $R$ is an $n \times n$ orthogonal matrix ( $R^{\operatorname{tr}} R=I=$ $\left.R R^{\operatorname{tr}}\right)$ viewed as a linear transformation on $\mathbb{R}^{n}$. Show for $f \in C^{2}\left(\mathbb{R}^{n}\right)$ that $\Delta(f \circ R)=\Delta f \circ R$, i.e. $\Delta$ is invariant under rotations.

Exercise 43.43. Show that every point $\eta \in \operatorname{bd}(\Omega)$ has a barrier function when $\operatorname{bd}(\Omega)$ is $C^{2}$. Hint: By making a change of coordinated involving rotations and translations change of coordinates, it suffices to assume $\eta=0 \in \operatorname{bd}(\Omega)$ and that $B(0, r) \cap \operatorname{bd}(\Omega)$ is the graph of a $C^{2}$ - function $g: B(0, r) \cap \mathbb{R}^{n-1} \rightarrow \mathbb{R}^{n}$ such that $g(0)=0$ and $\nabla g(0)=0$. Show for $\delta>0$ sufficiently small that

$$
d_{\delta}(x):=\left|\delta e_{n}-x\right|^{2} \text { for } x \in \operatorname{bd}(\Omega)
$$

has a unique global minimum at $x=0$. Use this fact and Example 43.40 to complete the proof.

Remark 43.44. To make Barrier functions for cones $C$, let $D:=C \cap S^{n-1}$ and let $u(\omega)$ for $\omega \in D$ denote the Dirichlet eigenfunction on $D$ for the spherical Laplacian with smallest eigenvalue $\lambda>0$, i.e. $-\Delta_{S^{n-1}} u=\lambda u$. This function is positive on $D^{o}$ and vanishes on the boundary. If $C$ has sperical symmetry, the function should be describable explicitly. At any rate, we can now considier the function $U(r \omega)=r^{\alpha} u$, then

$$
\begin{aligned}
\Delta U & =\frac{1}{r^{n-1}} \partial_{r}\left(r^{n-1} \partial_{r}\left[r^{\alpha} u\right]\right)+\frac{r^{\alpha}}{r^{2}} \Delta_{S^{n-1}} u \\
& =\alpha(\alpha+n-2) r^{\alpha-2} u-\lambda r^{\alpha-2} u
\end{aligned}
$$

which will be zero if $\lambda=\alpha(\alpha+n-2)$, i.e.

$$
\alpha^{2}+(n-2) \alpha-\lambda=0
$$

where

$$
\alpha=\frac{-(n-2) \pm \sqrt{(n-2)^{2}+4 \lambda}}{2}
$$

Hence taking

$$
\alpha=\frac{\sqrt{(n-2)^{2}+4 \lambda}-(n-2)}{2}
$$

the function $U(r \omega)=r^{\alpha} u(\omega)$ is harmonic in $C$ and 0 on $\partial C$. (Probably should be doing these considerations for the exterior of $C$.)

### 43.6 Solving the Dirichlet Problem by Integral Equations

Another method for solving the Dirichlet problem to reduce it to a question of solvability of a certain integral equation in $\operatorname{bd}(\Omega)$. For a nice sketch of how this goes the reader is referred to Reed and Simon [11], included below. For a more detailed account the reader may consult Sobolev [16] or Guenther and Lee [6].

The following text is taken from Reed and Simon Volume 1.

Uses Results for the BANACH SPACEC Example (Dirichlet problem) $\sqrt{ }$ The main impetus for the study of compact operators arose from the use of integral equations in attempting to solve the classical boundary value problems of mathematical physics. We briefly describe this method. Let $D$ be an open bounded region in $\mathbb{R}^{3}$ with a smooth boundary surface $\partial D$. The Dirichlet problem for Laplace's equation is: given a continuous function $f$ on $\partial D$, find a function $u$, twice differentiable in $D$ and continuous on $\bar{D}$, which satisfies

$$
\begin{aligned}
\Delta u(x) & =0 & & x \in D \\
u(x) & =f(x) & & x \in \partial D
\end{aligned}
$$

Let $K(x, y)=\left(x-y, n_{y}\right) / 2 \pi|x-y|^{3}$ where $n_{y}$ is the outer normal to $\partial D$ at the point $y \in \partial D$. Then, as a function of $x, K(x, y)$ satisfies $\Delta_{x} K(x, y)=0$ in the interior which suggests that we try to write $u$ as a superposition

$$
\begin{equation*}
u(x)=\int_{\partial D} K(x, y) \varphi(y) d S(y) \tag{VI.6a}
\end{equation*}
$$

where $\varphi(y)$ is some continuous function on $\partial D$ and $d S$ is the usual surface measure. Indeed, for $x \in D$, the integral makes perfectly good sense and
$\Delta u(x)=0$ in $D$. Furthermore, if $x_{0}$ is any point in $\partial D$ and $x \rightarrow x_{0}$ from inside $D$, it can be proven that

$$
\begin{equation*}
u(x) \rightarrow-\varphi\left(x_{0}\right)+\int_{\partial D} K\left(x_{0}, y\right) \varphi(y) d S(y) \tag{VI.6b}
\end{equation*}
$$

If $x \rightarrow x_{0}$ from outside $D$, the minus is replaced by a plus. Also,

$$
\int_{\partial D} K\left(x_{0}, y\right) \varphi(y) d S(y)
$$

ists and is a continuous function on $\partial D$ if $\varphi$ is a continuous function on $\partial D$. The proof depends on the fact that the boundary of $D$ is smooth which implies that for $x, y \in \partial D,\left(x-y, n_{y}\right) \approx c|x-y|^{2}$ as $x \rightarrow y$.
Since we wish $u(x)=f(x)$ on $\partial D$, the whole question reduces to whether we can find $\varphi$ so that

$$
f(x)=-\varphi(x)+\int_{\partial D} K(x, y) \varphi(y) d S(y), \quad x \in \partial D
$$

Let $T: C(\partial D) \rightarrow C(\partial D)$ be defined by

$$
T \varphi=\int_{\partial D} K(x, y) \varphi(y) d S(y)
$$

Not only is $T$ bounded but (as we will shortly see) $T$ is also compact. Thus, by the Fredholm alternative, either $\lambda=1$ is in the point spectrum of $T$ in which case there is a $\psi \in C(\partial D)$ such that $(I-T) \psi=0$, or $-f=(I-T) \varphi$ has a unique solution for each $f \in C(\partial D)$. If $u$ is defined by (VI.6a) with $\psi$ replacing $\varphi$, then $u \equiv 0$ in $D$ by the maximum principle. Further, $\partial u / \partial n$ is continuous across $\partial D$ and therefore equals zero on $\partial D$. By an integration by parts this implies that $u \equiv 0$ outside $\partial D$. Therefore, by (VI.6b), $2 \psi(x) \equiv 0$ on $\partial D$, so the first alternative does not hold.
The idea of the compactness proof is the following. Let

$$
K_{\delta}(x, z)=\frac{\left(x-z, n_{z}\right)}{|x-z|^{3}+\delta}
$$

If $\delta>0$, the kernel $K_{\delta}$ is continuous, so, by the discussion at the beginning of this section, the corresponding integral operators $T_{b}$, are compact. To prove that $T$ is compact, we need only show that $\left\|T-T_{\delta}\right\| \rightarrow 0$ as $\delta \rightarrow 0$. By the estimate

$$
\left|\left(T_{\delta} f\right)(x)-(T f)(x)\right| \leq\|f\|_{\infty} \int_{\partial D}\left|K(x, z)-K_{d}(x, z)\right| d S(z)
$$

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we must only show that the integral converges to zero uniformly in $x$ as $\delta \rightarrow 0$. To prove this, divide the integration region into the set where $|x-z| \geq \varepsilon$ and its complement. For fixed $\varepsilon$, the kernels converge uniformly on the first region. By using the fact that $K$ is integrable, the contribution from the second region can be made arbitrarily small for $\varepsilon$ sufficiently small.

44 Introduction to the Spectral Theorem

## Introduction to the Spectral Theorem

The following spectral theorem is a minor variant of the usual spectral theorem for matrices. This reformulation has the virtue of carrying over to general (unbounded) self adjoint operators on infinite dimensional Hilbert spaces.

Theorem 44.1. Suppose $A$ is an $n \times n$ complex self adjoint matrix, i.e. $A^{*}=$ $A$ or equivalently $A_{j i}=\bar{A}_{i j}$ and let $\mu$ be counting measure on $\{1,2, \ldots, n\}$. Then there exists a unitary map $U: \mathbb{C}^{n} \rightarrow L^{2}(\{1,2, \ldots, n\}, d \mu)$ and a real function $\lambda:\{1,2, \ldots, n\} \rightarrow \mathbb{R}$ such that $U A \xi=\lambda \cdot U \xi$ for all $\xi \in \mathbb{C}^{n}$. We summarize this equation by writing $U A U^{-1}=M_{\lambda}$ where

$$
M_{\lambda}: L^{2}(\{1,2, \ldots, n\}, d \mu) \rightarrow L^{2}(\{1,2, \ldots, n\}, d \mu)
$$

is the linear operator, $g \in L^{2}(\{1,2, \ldots, n\}, d \mu) \rightarrow \lambda \cdot g \in L^{2}(\{1,2, \ldots, n\}, d \mu)$.
Proof. By the usual form of the spectral theorem for self-adjoint matrices, there exists an orthonormal basis $\left\{e_{i}\right\}_{i=1}^{n}$ of eigenvectors of $A$, say $A e_{i}=\lambda_{i} e_{i}$ with $\lambda_{i} \in \mathbb{R}$. Define $U: \mathbb{C}^{n} \rightarrow L^{2}(\{1,2, \ldots, n\}, d \mu)$ to be the unique (unitary) map determined by $U e_{i}=\delta_{i}$ where

$$
\delta_{i}(j)=\left\{\begin{array}{l}
1 \text { if } i=j \\
0 \text { if } i \neq j
\end{array}\right.
$$

and let $\lambda:\{1,2, \ldots, n\} \rightarrow \mathbb{R}$ be defined by $\lambda(i):=\lambda_{i}$.
Definition 44.2. Let $A: H \rightarrow H$ be a possibly unbounded operator on $H$. We let

$$
D\left(A^{*}\right)=\{y \in H: \quad \exists z \in H \ni(A x, y)=(x, z) \forall x \in D(A)\}
$$

and for $y \in D\left(A^{*}\right)$ set $A^{*} y=z$.
Definition 44.3. An operator $A$ on $H$ is symmetric if $A \subset A^{*}$ and is selfadjoint iff $A=A^{*}$.

The reader should check that $A: H \rightarrow H$ is symmetric iff $(A x, y)=(x, A y)$ for all $x, y \in D(A)$.

Proposition 44.4. Let $(X, \mu)$ be $\sigma$ - finite measure space, $H=L^{2}(X, d \mu)$ and $f: X \rightarrow \mathbb{C}$ be a measurable function. Set $A g=f g=M_{f} g$ for all

$$
g \in D\left(M_{f}\right)=\{g \in H: f g \in H\} .
$$

Then $D\left(M_{f}\right)$ is a dense subspace of $H$ and $M_{f}^{*}=M_{\bar{f}}$.
Proof. For any $g \in H=L^{2}(X, d \mu)$ and $m \in \mathbb{N}$, let $g_{m}:=g 1_{|f| \leq m}$. Since $\left|f g_{m}\right| \leq m|g|$ it follows that $f g_{m} \in H$ and hence $g_{m} \in D\left(M_{f}\right)$. By the dominated convergence theorem, it follows that $g_{m} \rightarrow g$ in $H$ as $m \rightarrow \infty$, hence $D\left(M_{f}\right)$ is dense in $H$.

Suppose $h \in \mathcal{D}\left(M_{f}^{*}\right)$ then there exists $k \in L^{2}$ such that $\left(M_{f} g, h\right)=(g, k)$ for all $g \in D\left(M_{f}\right)$, i.e.

$$
\int_{X} f g \bar{h} d \mu=\int_{X} g \bar{k} d \mu \text { for all } g \in D\left(M_{f}\right)
$$

or equivalently

$$
\begin{equation*}
\int_{X} g(\overline{\bar{f} h-k}) d \mu=0 \text { for all } g \in D\left(M_{f}\right) \tag{44.1}
\end{equation*}
$$

Choose $X_{n} \subset X$ such that $X_{n} \uparrow X$ and $\mu\left(X_{n}\right)<\infty$ for all $n$. It is easily checked that

$$
g_{n}:=1_{X_{n}} \frac{\bar{f} h-k}{|\bar{f} h-k|} 1_{|f| \leq n}
$$

is in $D\left(M_{f}\right)$ and putting this function into Eq. (44.1) shows

$$
\int_{X}|\bar{f} h-k| 1_{|f| \leq n} d \mu=0 \text { for all } n .
$$

Using the monotone convergence theorem, we may let $n \rightarrow \infty$ in this equation to find $\int_{X}|\bar{f} h-k| d \mu=0$ and hence that $\bar{f} h=k \in L^{2}$. This shows $h \in$ $D\left(M_{\bar{f}}\right)$ and $M_{f}^{*} h=\bar{f} h$.
Theorem 44.5 (Spectral Theorem). Suppose $A^{*}=A$ then there exists $(X, \mu)$ a $\sigma$ - finite measure space, $f: X \rightarrow \mathbb{R}$ measurable, and $U: H \rightarrow$ $L^{2}(x, \mu)$ unitary such that $U A U^{-1}=M_{f}$. Note this is a statement about domains as well, i.e. $U D\left(M_{f}\right)=D(A)$.

I would like to give some examples of computing $A^{*}$ and Theorem 44.5 as well. We will consider here the case of constant coefficient differential operators on $L^{2}\left(\mathbb{R}^{n}\right)$. First we need the following definition.

Definition 44.6. Let $a_{\alpha} \in C^{\infty}(U), L=\sum_{|\alpha| \leq m} a_{\alpha} \partial^{\alpha}-a m^{\text {th }}$ order linear differential operator on $\mathcal{D}(U)$ and

$$
L^{\dagger} \phi=\sum_{|\alpha| \leq m}(-1)^{|\alpha|} \partial^{\alpha}\left[a_{\alpha} \phi\right]
$$

denote the formal adjoint of $L$ as in Lemma 40.4 above. For $f \in L^{p}(U)$ we say $L f \in L^{p}(U)$ or $L_{l o c}^{p}(U)$ if the generalized function $L f$ may be represented by an element of $L^{p}(U)$ or $L_{\text {loc }}^{p}(U)$ respectively, i.e. $L f=g \in L_{\text {loc }}^{p}(U)$ iff

$$
\begin{equation*}
\int_{U} f \cdot L^{\dagger} \phi d m=\int_{U} g \phi d m \text { for all } \phi \in C_{c}^{\infty}(U) \tag{44.2}
\end{equation*}
$$

In terms of the complex inner product,

$$
(f, g):=\int_{U} f(x) \bar{g}(x) d m(x)
$$

Eq. (44.2) is equivalent to

$$
\left(f \cdot L^{\circledast} \phi\right)=(g, \phi) \text { for all } \phi \in C_{c}^{\infty}(U)
$$

where

$$
L^{\circledast} \phi:=\sum_{|\alpha| \leq m}(-1)^{|\alpha|} \partial^{\alpha}\left[\bar{a}_{\alpha} \phi\right] .
$$

Notice that $L^{\circledast}$ satisfies $L^{\circledast} \bar{\phi}=\overline{L^{\dagger} \phi}$. (We do not write $L^{*}$ here since $L^{\circledast}$ is to be considered an operator on the space on $\mathcal{D}^{\prime}(U)$.)
Remark 44.7. Recall that if $f, h \in L^{2}\left(\mathbb{R}^{n}\right)$, then the following are equivalent

1. $\hat{f}=h$.
2. $(h, g)=\left(f, \mathcal{F}^{-1} g\right)$ for all $g \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$.
3. $(h, g)=\left(f, \mathcal{F}^{-1} g\right)$ for all $g \in \mathcal{S}\left(\mathbb{R}^{n}\right)$.
4. $(h, g)=\left(f, \mathcal{F}^{-1} g\right)$ for all $g \in L^{2}\left(\mathbb{R}^{n}\right)$.

Indeed if $\hat{f}=h$ and $g \in L^{2}\left(\mathbb{R}^{n}\right)$, the unitarity of $\mathcal{F}$ implies

$$
(h, g)=(\hat{f}, g)=(\mathcal{F} f, g)=\left(f, \mathcal{F}^{-1} g\right)
$$

Hence $1 \Longrightarrow 4$ and it is clear that $4 \Longrightarrow 3 \Longrightarrow 2$. If 2 holds, then again since $\mathcal{F}$ is unitary we have

$$
(h, g)=\left(f, \mathcal{F}^{-1} g\right)=(\hat{f}, g) \text { for all } g \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)
$$

which implies $h=\hat{f}$ a.e., i.e. $h=\hat{f}$ in $L^{2}\left(\mathbb{R}^{n}\right)$.
Proposition 44.8. Let $p(x)=\sum_{|\alpha| \leq m} a_{\alpha} x^{\alpha}$ be a polynomial on $\mathbb{C}^{n}$,

$$
\begin{equation*}
L:=p(\partial):=\sum_{|\alpha| \leq m} a_{\alpha} \partial^{\alpha} \tag{44.3}
\end{equation*}
$$

and $f \in L^{2}\left(\mathbb{R}^{n}\right)$. Then $L f \in L^{2}\left(\mathbb{R}^{n}\right)$ iff $p(i \xi) \hat{f}(\xi) \in L^{2}\left(\mathbb{R}^{n}\right)$ and in which case

Put more concisely, letting

$$
D(B)=\left\{f \in L^{2}\left(\mathbb{R}^{n}\right): L f \in L^{2}\left(\mathbb{R}^{n}\right)\right\}
$$

with $B f=L f$ for all $f \in D(B)$, we have

$$
\mathcal{F} B \mathcal{F}^{-1}=M_{p(i \xi)}
$$

Proof. As above, let

$$
\begin{equation*}
L^{\dagger}:=\sum_{|\alpha| \leq m} a_{\alpha}(-\partial)^{\alpha} \text { and } L^{\circledast}:=\sum_{|\alpha| \leq m} \bar{a}_{\alpha}(-\partial)^{\alpha} \tag{44.5}
\end{equation*}
$$

For $\phi \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$,

$$
\begin{aligned}
L^{\circledast} \phi^{\vee}(x) & =L^{\circledast} \int \phi(\xi) e^{i x \cdot \xi} d \lambda(\xi)=\sum_{|\alpha| \leq m} \bar{a}_{\alpha}\left(-\partial_{x}\right)^{\alpha} \int \phi(\xi) e^{i x \cdot \xi} d \lambda(\xi) \\
& =\int \overline{p(i \xi)} \phi(\xi) e^{i x \cdot \xi} d \lambda(\xi)=\mathcal{F}^{-1}[\overline{p(i \xi)} \phi(\xi)](x)
\end{aligned}
$$

So if $f \in L^{2}\left(\mathbb{R}^{n}\right)$ such that $L f \in L^{2}\left(\mathbb{R}^{n}\right)$. Then by Remark 44.7,

$$
\begin{aligned}
(\widehat{L f}, \phi) & =\left(L f, \phi^{\vee}\right)=\left\langle f, L^{\circledast} \phi^{\vee}\right\rangle=\left\langle f(x), \mathcal{F}^{-1}[\overline{p(i \xi)} \phi(\xi)](x)\right\rangle \\
& =\langle\hat{f}(\xi),[\overline{p(i \xi)} \phi(\xi)]\rangle=\langle p(i \xi) \hat{f}(\xi), \phi(\xi)\rangle \text { for all } \phi \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)
\end{aligned}
$$

from which it follows that Eq. (44.4) holds and that $p(i \xi) \hat{f}(\xi) \in L^{2}\left(\mathbb{R}^{n}\right)$.
Conversely, if $f \in L^{2}\left(\mathbb{R}^{n}\right)$ is such that $p(i \xi) \hat{f}(\xi) \in L^{2}\left(\mathbb{R}^{n}\right)$ then for $\phi \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$,

$$
\begin{equation*}
\left(f, L^{\circledast} \phi\right)=\left(\hat{f}, \mathcal{F} L^{\circledast} \phi\right) . \tag{44.6}
\end{equation*}
$$

Since

$$
\begin{aligned}
\mathcal{F}\left(L^{\circledast} \phi\right)(\xi) & =\int L^{\circledast} \phi(x) e^{-i x \cdot \xi} d \lambda(x)=\int \phi(x) \overline{L_{x}} e^{-i x \cdot \xi} d \lambda(x) \\
& =\int \phi(x) \overline{a_{\alpha}} \partial_{x}^{\alpha} e^{-i x \cdot \xi} d \lambda(x)=\int \phi(x) \overline{a_{\alpha}}(-i \xi)^{\alpha} e^{-i x \cdot \xi} d \lambda(x) \\
& =\overline{n(i \xi)} \hat{\phi}(\xi) .
\end{aligned}
$$

Eq. (44.6) becomes

$$
\begin{aligned}
\left(f, L^{\circledast} \phi\right) & =(\hat{f}(\xi), \overline{p(i \xi)} \hat{\phi}(\xi))=(p(i \xi) \hat{f}(\xi), \hat{\phi}(\xi)) \\
& =\left(\mathcal{F}^{-1}[p(i \xi) \hat{f}(\xi)](x), \phi(x)\right) .
\end{aligned}
$$

This shows $L f=\mathcal{F}^{-1}[p(i \xi) \hat{f}(\xi)] \in L^{2}\left(\mathbb{R}^{n}\right)$.

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Lemma 44.9. Suppose $p(x)=\sum_{|\alpha| \leq m} a_{\alpha} x^{\alpha}$ is a polynomial on $\mathbb{R}^{n}$ and $L=$ $p(\partial)$ is the constant coefficient differential operator $B=\sum_{|\alpha| \leq m} a_{\alpha} \partial^{\alpha}$ with $D(B):=\mathcal{S}\left(\mathbb{R}^{n}\right) \subset L^{2}\left(\mathbb{R}^{n}\right)$. Then

$$
\mathcal{F} B \mathcal{F}^{-1}=\left.M_{p(i \xi)}\right|_{\mathcal{S}\left(\mathbb{R}^{n}\right)}
$$

Proof. This is result of the fact that $\mathcal{F}\left(\mathcal{S}\left(\mathbb{R}^{n}\right)\right)=\mathcal{S}\left(\mathbb{R}^{n}\right)$ and for $f \in$ $\mathcal{S}\left(\mathbb{R}^{n}\right)$ we have

$$
f(x)=\int_{\mathbb{R}^{n}} \hat{f}(\xi) e^{i \xi \cdot x} d \lambda(\xi)
$$

so that

$$
B f(x)=\int_{\mathbb{R}^{n}} \hat{f}(\xi) L_{x} e^{i \xi \cdot x} d \lambda(\xi)=\int_{\mathbb{R}^{n}} \hat{f}(\xi) p(i \xi) e^{i \xi \cdot x} d \lambda(\xi)
$$

so that

$$
(B f) \hat{}(\xi)=p(i \xi) \hat{f}(\xi) \text { for all } f \in \mathcal{S}\left(\mathbb{R}^{n}\right)
$$

Lemma 44.10. Suppose $g: \mathbb{R}^{n} \rightarrow \mathbb{C}$ is a measurable function such that $|g(x)| \leq C\left(1+|x|^{M}\right)$ for some constants $C$ and $M$. Let $A$ be the unbounded operator on $L^{2}\left(\mathbb{R}^{n}\right)$ defined by $D(A)=\mathcal{S}\left(\mathbb{R}^{n}\right)$ and for $f \in \mathcal{S}\left(\mathbb{R}^{n}\right)$, $A f=g f$. Then $A^{*}=M_{\bar{g}}$.

Proof. If $h \in D\left(M_{\bar{g}}\right)$ and $f \in D(A)$, we have

$$
(A f, h)=\int_{\mathbb{R}^{n}} g f \bar{h} d m=\int_{\mathbb{R}^{n}} f \overline{g h} d m=\left(f, M_{\bar{g}} h\right)
$$

which shows $M_{\bar{g}} \subset A^{*}$, i.e. $h \in D\left(A^{*}\right)$ and $A^{*} h=M_{\bar{g}} h$. Now suppose $h \in$ $D\left(A^{*}\right)$ and $A^{*} h=k$, i.e.

$$
\int_{\mathbb{R}^{n}} g f \bar{h} d m=(A f, h)=(f, k)=\int_{\mathbb{R}^{n}} f \bar{k} d m \text { for all } f \in \mathcal{S}\left(\mathbb{R}^{n}\right)
$$

or equivalently that

$$
\int_{\mathbb{R}^{n}}(g \bar{h}-\bar{k}) f d m=0 \text { for all } f \in \mathcal{S}\left(\mathbb{R}^{n}\right)
$$

Since the last equality (even just for $f \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ ) implies $g \bar{h}-\bar{k}=0$ a.e. we may conclude that $h \in D\left(M_{\bar{g}}\right)$ and $k=M_{\bar{g}} h$, i.e. $A^{*} \subset M_{\bar{g}}$.

Theorem 44.11. Suppose $p(x)=\sum_{|\alpha| \leq m} a_{\alpha} x^{\alpha}$ is a polynomial on $\mathbb{R}^{n}$ and $A=p(\partial)$ is the constant coefficient differential operator with $D(A):=$ $C_{c}^{\infty}\left(\mathbb{R}^{n}\right) \subset L^{2}\left(\mathbb{R}^{n}\right)$ such that $A=L=p(\partial)$ on $D(A)$, see Eq. (44.3). Then $A^{*}$ is the operator described by

$$
\begin{aligned}
D\left(A^{*}\right) & =\left\{f \in L^{2}\left(\mathbb{R}^{n}\right): L^{\dagger} f \in L^{2}\left(\mathbb{R}^{n}\right)\right\} \\
& =\left\{f \in L^{2}\left(\mathbb{R}^{n}\right): p(i \xi) \hat{f}(\xi) \in L^{2}\left(\mathbb{R}^{n}\right)\right\}
\end{aligned}
$$

and $A^{*} f=L^{\dagger} f$ for $f \in D\left(A^{*}\right)$ where $L^{\dagger}$ is defined in Eq. (44.5) above. Moreover we have $\mathcal{F} A^{*} \mathcal{F}^{-1}=M_{\overline{p(i \xi)}}$.

Proof. Let $D(B)=\mathcal{S}\left(\mathbb{R}^{n}\right)$ and $B:=L$ on $D(B)$ so that $A \subset B$. We are first going to show $A^{*}=B^{*}$. As is easily verified, in general if $A \subset B$ then $B^{*} \subset A^{*}$. So we need only show $A^{*} \subset B^{*}$. Now by definition, if $g \in D\left(A^{*}\right)$ with $k=A^{*} g$, then

$$
(A f, g)=(f, k) \text { for all } f \in D(A):=C_{c}^{\infty}\left(\mathbb{R}^{n}\right)
$$

Suppose that $f \in \mathcal{S}\left(\mathbb{R}^{n}\right)$ and $\phi \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ such that $\phi=1$ in a neighborhood of 0 . Then $f_{n}(x):=\phi(x / n) f(x)$ is in $\mathcal{S}\left(\mathbb{R}^{n}\right)$ and hence

$$
\begin{equation*}
\left(f_{n}, k\right)=\left(L f_{n}, g\right) \tag{44.7}
\end{equation*}
$$

An exercise in the product rule and the dominated convergence theorem shows $f_{n} \rightarrow f$ and $L f_{n} \rightarrow L f$ in $L^{2}\left(\mathbb{R}^{n}\right)$ as $n \rightarrow \infty$. Therefore we may pass to the limit in Eq. (44.7) to learn

$$
(f, k)=(B f, g) \text { for all } f \in \mathcal{S}\left(\mathbb{R}^{n}\right)
$$

which shows $g \in D\left(B^{*}\right)$ and $B^{*} g=k$.
By Lemma 44.10 , we may conclude that $A^{*}=B^{*}=M_{\overline{p(i \xi)}}$ and by Proposition 44.8 we then conclude that

$$
\begin{aligned}
D\left(A^{*}\right) & =\left\{f \in L^{2}\left(\mathbb{R}^{n}\right): p(i \xi) \hat{f}(\xi) \in L^{2}\left(\mathbb{R}^{n}\right)\right\} \\
& =\left\{f \in L^{2}\left(\mathbb{R}^{n}\right): L^{\dagger} f \in L^{2}\left(\mathbb{R}^{n}\right)\right\}
\end{aligned}
$$

and for $f \in D\left(A^{*}\right)$ we have $A^{*} f=L^{\dagger} f$.
Example 44.12. If we take $L=\Delta$ with $D(L):=C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$, then

$$
L^{*}=\bar{\Delta}=\mathcal{F} M_{-|\xi|^{2}} \mathcal{F}^{-1}
$$

where $D(\bar{\Delta})=\left\{f \in L^{2}\left(\mathbb{R}^{n}\right): \Delta f \in L^{2}\left(\mathbb{R}^{n}\right)\right\}$ and $\bar{\Delta} f=\Delta f$.
Theorem 44.13. Suppose $A=A^{*}$ and $A \leq 0$. Then for all $u_{0} \in D(A)$ there exists a unique solution $u \in C^{1}([0, \infty))$ such that $u(t) \in D(A)$ for all $t$ and

$$
\begin{equation*}
\dot{u}(t)=A u(t) \text { with } u(0)=u_{0} . \tag{44.8}
\end{equation*}
$$

Writing $u(t)=e^{t A} u_{0}$, the map $u_{0} \rightarrow e^{t A} u_{0}$ is a linear contraction semi-group, i.e.

$$
\begin{equation*}
\left\|e^{t A} u_{0}\right\| \leq\left\|u_{0}\right\| \text { for all } t \geq 0 \tag{44.9}
\end{equation*}
$$

So $e^{t A}$ extends uniquely to $H$ by continuity. This extension satisfies:

1. Strong Continuity: the map $t \in[0, \infty) \rightarrow e^{t A} u_{0}$ is continuous for all $u_{0} \in H$.
2. Smoothing property: $t>0$

$$
e^{t A} u_{0} \in \bigcap_{n=0}^{\infty} D\left(A^{n}\right)=: C^{\infty}(A)
$$

and

$$
\begin{equation*}
\left\|A^{k} e^{t A}\right\| \leq\left(\frac{k}{t}\right)^{k} e^{-k} \text { for all } k \in \mathbb{N} \tag{44.10}
\end{equation*}
$$

Proof. Uniqueness. Suppose $u$ solves Eq. (44.8), then

$$
\frac{d}{d t}(u(t), u(t))=2 \operatorname{Re}(\dot{u}, u)=2 \operatorname{Re}(A u, u) \leq 0
$$

Hence $\|u(t)\|$ is decreasing so that $\|u(t)\| \leq\left\|u_{0}\right\|$.This implies the uniqueness assertion in the theorem and the norm estimate in Eq. (44.9).

Existence: By the spectral theorem we may assume $A=M_{f}$ acting on $L^{2}(X, \mu)$ for some $\sigma$ - finite measure space $(X, \mu)$ and some measurable function $f: X \rightarrow(-\infty, 0]$. We wish to show $u(t)=e^{t f} u_{0} \in L^{2}$ solves

$$
\dot{u}(t)=f u(t) \text { with } u(0)=u_{0} \in D\left(M_{f}\right) \subset L^{2} .
$$

Let $t>0$ and $|\Delta|<t$. Then by the mean value inequality
$\left|\frac{e^{(t+\Delta) f}-e^{t f}}{\triangle} u_{0}\right|=\max \left\{\left|f e^{(t+\tilde{\Delta}) f} u_{0}\right|: \tilde{\Delta}\right.$ between 0 and $\left.\Delta\right\} \leq\left|f u_{0}\right| \in L^{2}$.
This estimated along with the fact that

$$
\frac{u(t+\Delta)-u(t)}{\triangle}=\frac{e^{(t+\Delta) f}-e^{t f}}{\triangle} u_{0} \xrightarrow{\text { point wise }} f e^{t f} u_{0} \text { as } \Delta \rightarrow 0
$$

enables us to use the dominated convergence theorem to conclude

$$
\dot{u}(t)=L^{2}-\lim _{\Delta \rightarrow 0} \frac{u(t+\triangle)-u(t)}{\triangle}=e^{t f} f u_{0}=f u(t)
$$

as desired. i.e. $\dot{u}(t)=f u(t)$.
The extension of $e^{t A}$ to $H$ is given by $M_{e^{t f}}$. For $g \in L^{2},\left|e^{t f} g\right| \leq|g| \in L^{2}$ and $e^{t f} g \rightarrow e^{\tau g} f$ pointwise as $t \rightarrow \tau$, so the Dominated convergence theorem shows $t \in[0, \infty) \rightarrow e^{t A} g \in H$ is continuous. For the last two assertions, let $t>0$ and $f(x)=x^{k} e^{t x}$. Then $(\ln f)^{\prime}(x)=\frac{k}{x}+t$ which is zero when $x=-k / t$ and therefore

$$
\max _{x \leq 0}\left|x^{k} e^{t x}\right|=|f(-k / t)|=\left(\frac{k}{t}\right)^{k} e^{-k}
$$

Hence

$$
\left\|A^{k} e^{t A}\right\|_{o p} \leq \max _{x \leq 0}\left|x^{k} e^{t x}\right| \leq\left(\frac{k}{t}\right)^{k} e^{-k}<\infty
$$

Theorem 44.14. Take $A=\mathcal{F} M_{-|\xi|^{2}} \mathcal{F}^{-1}$ so $\left.A\right|_{\mathcal{S}}=\Delta$ then

$$
C^{\infty}(A):=\bigcap_{n=1}^{\infty} D\left(A^{n}\right) \subset C^{\infty}\left(\mathbb{R}^{d}\right)
$$

i.e. for all $f \in C^{\infty}(A)$ there exists a version $\tilde{f}$ of $f$ such that $\tilde{f} \in C^{\infty}\left(\mathbb{R}^{d}\right)$.

Proof. By assumption $|\xi|^{2 n} \hat{f}(\xi) \in L^{2}$ for all $n$. Therefore $\hat{f}(\xi)=\frac{g_{n}(\xi)}{1+\mid \xi \xi^{2 n}}$ for some $g_{n} \in L^{2}$ for all $n$. Therefore for $n$ chosen so that $2 n>m+d$, we have

$$
\int_{\mathbb{R}^{d}}|\xi|^{m}|\hat{f}(\xi)| d \xi \leq\left\|g_{n}\right\|_{L^{2}}\left\|\frac{|\xi|^{m}}{1+|\xi|^{2 n}}\right\|_{2}<\infty
$$

which shows $|\xi|^{m}|\hat{f}(\xi)| \in L^{1}$ for all $m=0,1,2, \ldots$ We may now differentiate the inversion formula, $f(x)=\int \hat{f}(\xi) e^{i x \cdot \xi} d \xi$ to find

$$
D^{\alpha} f(x)=\int(i \xi)^{\alpha} \hat{f}(\xi) e^{i x \cdot \xi} d \xi \text { for any } \alpha
$$

and thus conclude $f \in C^{\infty}$.
Exercise 44.15. Some Exercises: Section $2.54,5,6,8,9,11,12,17$.

### 44.1 Du Hammel's principle again

Lemma 44.16. Suppose $A$ is an operator on $H$ such that $A^{*}$ is densely defined then $A^{*}$ is closed.

Proof. If $f_{n} \in D\left(A^{*}\right) \rightarrow f \in H$ and $A^{*} f_{n} \rightarrow g$ then for all $h \in D(A)$

$$
(g, h)=\lim _{n \rightarrow \infty}\left(A^{*} f_{n}, h\right)
$$

while

$$
\lim _{n \rightarrow \infty}\left(A^{*} f_{n}, h\right)=\lim _{n \rightarrow \infty}\left(f_{n}, A h\right)=(f, A h)
$$

i.e. $(A h, f)=(h, g)$ for all $h \in D(A)$. Thus $f \in D\left(A^{*}\right)$ and $A^{*} f=g$.

Corollary 44.17. If $A^{*}=A$ then $A$ is closed.
Corollary 44.18. Suppose $A$ is closed and $u(t) \in D(A)$ is a path such that $u(t)$ and $A u(t)$ are continuous in $t$. Then

$$
A \int_{0}^{T} u(\tau) d \tau=\int_{0}^{T} A u(\tau) d \tau
$$

44 Introduction to the Spectral Theorem

Proof. Let $\pi_{n}$ be a sequence of partitions of $[0, T]$ such that $\operatorname{mesh}\left(\pi_{n}\right) \rightarrow 0$ as $n \rightarrow \infty$ and set

$$
f_{n}=\sum_{\pi_{n}} u\left(\tau_{i}\right)\left(\tau_{i+1}-\tau_{i}\right) \in D(A) .
$$

Then $f_{n} \rightarrow \int_{0}^{T} u(\tau) d \tau$ and

$$
A f_{n}=\sum_{\pi_{n}} A u\left(\tau_{i}\right)\left(\tau_{i+1}-\tau_{i}\right) \rightarrow \int_{0}^{T} A u(\tau) d \tau
$$

Therefore $\int_{0}^{T} u(\tau) d \tau \in D(A)$ and $A \int_{0}^{T} u(\tau) d \tau=\int_{0}^{T} A u(\tau) d \tau$.
Lemma 44.19. Suppose $A=A^{*}, A \leq 0$, and $h:[0, \infty] \rightarrow H$ is continuous. Then

$$
\begin{gathered}
(s, t) \in[0, \infty) \times[0, \infty) \rightarrow e^{s A} h(t) \\
(s, t) \in(0, \infty) \times[0, \infty) \rightarrow A^{k} e^{s A} h(t)
\end{gathered}
$$

are continuous maps into $H$.
Proof. Let $k \geq 0$, then if $s \geq \sigma$,

$$
\begin{aligned}
\| A^{k} & \left(e^{s A} h(t)-e^{\sigma A} h(\tau)\right) \| \\
& =\left\|A^{k} e^{\sigma A}\left(e^{(s-\sigma) A} h(t)-h(\tau)\right)\right\| \\
& \leq\left\|A^{k} e^{\sigma A}\right\|\left\|e^{(s-\sigma) A}[h(t)-h(\tau)]+e^{(s-\sigma) A} h(\tau)-h(\tau)\right\| \\
& \leq\left(\frac{k}{\sigma}\right)^{k} e^{-k} \cdot\left[\|h(t)-h(\tau)\|+\left\|e^{(s-\sigma) A} h(\tau)-h(\tau)\right\|\right]
\end{aligned}
$$

So

$$
\lim _{s \downarrow \sigma \text { and } t \rightarrow \tau}\left\|A^{k}\left(e^{s A} h(t)-e^{\sigma A} h(\tau)\right)\right\|=0
$$

and we may take $\sigma=0$ if $k=0$. Similarly, if $s \leq \sigma$,

$$
\begin{aligned}
& \| A^{k}\left(e^{s A} h(t)-e^{\sigma A} h(\tau)\right) \| \\
&=\left\|A^{k} e^{s A}\left(h(t)-e^{(\sigma-s) A} h(\tau)\right)\right\| \\
& \quad \leq\left\|A^{k} e^{s A}\right\|\left[\|h(t)-h(\tau)\|+\left\|h(\tau)-e^{(\sigma-s) A} h(\tau)\right\|\right] \\
& \leq\left(\frac{k}{s}\right)^{k} e^{-k}\left[\|h(t)-h(\tau)\|+\left\|h(\tau)-e^{(\sigma-s) A} h(\tau)\right\|\right]
\end{aligned}
$$

and the latter expression tends to zero as $s \uparrow \sigma$ and $t \rightarrow \tau$.
Lemma 44.20. Let $h \in C([0, \infty), H), D:=\left\{(s, t) \in \mathbb{R}^{2}: s>t \geq 0\right\}$ and $F(s, t):=\int_{0}^{t} e^{(s-\tau) A} h(\tau) d \tau$ for $(s, t) \in D$. Then

1. $F \in C^{1}(D, H)\left(\right.$ in fact $\left.F \in C^{\infty}(D, H)\right)$,

$$
\begin{equation*}
\frac{\partial}{\partial t} F(s, t)=e^{(s-t) A} h(t) \tag{44.11}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{\partial F(s, t)}{\partial s}=\int_{0}^{t} A e^{(s-\tau) A} h(\tau) d \tau \tag{44.12}
\end{equation*}
$$

2. Given $\epsilon>0$ let

$$
u_{\epsilon}(t):=F(t+\epsilon, t)=\int_{0}^{t} e^{(t+\epsilon-\tau) A} h(\tau) d \tau
$$

Then $u_{\epsilon} \in C^{1}((-\epsilon, \infty), H), u_{\epsilon}(t) \in D(A)$ for all $t>-\epsilon$ and

$$
\begin{equation*}
\dot{u}_{\epsilon}(t)=e^{\epsilon A} h(t)+A u_{\epsilon}(t) \tag{44.13}
\end{equation*}
$$

Proof. We claim the function

$$
(s, t) \in D \rightarrow F(s, t):=\int_{0}^{t} e^{(s-\tau) A} h(\tau) d \tau
$$

is continuous. Indeed if $\left(s^{\prime}, t^{\prime}\right) \in D$ and $(s, t) \in D$ is sufficiently close to $\left(s^{\prime}, t^{\prime}\right)$ so that $s>t^{\prime}$, we have

$$
\begin{aligned}
F(s, t)-F\left(s^{\prime}, t^{\prime}\right) & =\int_{0}^{t} e^{(s-\tau) A} h(\tau) d \tau-\int_{0}^{t^{\prime}} e^{\left(s^{\prime}-\tau\right) A} h(\tau) d \tau \\
& =\int_{0}^{t} e^{(s-\tau) A} h(\tau) d \tau-\int_{0}^{t^{\prime}} e^{(s-\tau) A} h(\tau) d \tau \\
& +\int_{0}^{t^{\prime}}\left[e^{(s-\tau) A}-e^{\left(s^{\prime}-\tau\right) A}\right] h(\tau) d \tau
\end{aligned}
$$

so that

$$
\begin{align*}
\left\|F(s, t)-F\left(s^{\prime}, t^{\prime}\right)\right\| \leq & \left|\int_{t^{\prime}}^{t}\left\|e^{(s-\tau) A} h(\tau)\right\| d \tau\right| \\
& +\int_{0}^{t^{\prime}}\left\|\left[e^{(s-\tau) A}-e^{\left(s^{\prime}-\tau\right) A}\right] h(\tau)\right\| d \tau \\
\leq & \left|\int_{t^{\prime}}^{t}\|h(\tau)\| d \tau\right| \\
& +\int_{0}^{t^{\prime}}\left\|\left[e^{(s-\tau) A}-e^{\left(s^{\prime}-\tau\right) A}\right] h(\tau)\right\| d \tau \tag{44.14}
\end{align*}
$$

By the dominated convergence theorem,

$$
\lim _{(s, t) \rightarrow\left(s^{\prime}, t^{\prime}\right)}\left|\int_{t^{\prime}}^{t}\|h(\tau)\| d \tau\right|=0
$$

$$
\lim _{(s, t) \rightarrow\left(s^{\prime}, t^{\prime}\right)} \int_{0}^{t^{\prime}}\left\|\left[e^{(s-\tau) A}-e^{\left(s^{\prime}-\tau\right) A}\right] h(\tau)\right\| d \tau=0
$$

which along with Eq. (44.14) shows $F$ is continuous.
By the fundamental theorem of calculus,

$$
\frac{\partial}{\partial t} F(s, t)=e^{(s-t) A} h(t)
$$

and as we have seen this expression is continuous on $D$. Moreover, since

$$
\frac{\partial}{\partial s} e^{(s-\tau) A} h(\tau)=A e^{(s-\tau) A} h(\tau)
$$

is continuous and bounded for on $s>t>\tau$, we may differentiate under the integral to find

$$
\frac{\partial F(s, t)}{\partial s}=\int_{0}^{t} A e^{(s-\tau) A} h(\tau) d \tau \text { for } s>t
$$

A similar argument (making use of Eq. (44.10) with $k=1$ ) shows $\frac{\partial F(s, t)}{\partial s}$ is continuous for $(s, t) \in D$.

By the chain rule, $u_{\epsilon}(t):=F(t+\epsilon, t)$ is $C^{1}$ for $t>-\epsilon$ and

$$
\begin{aligned}
\dot{u}_{\epsilon}(t) & =\frac{\partial F(t+\epsilon, t)}{\partial s}+\frac{\partial F(t+\epsilon, t)}{\partial t} \\
& =e^{\epsilon A} h(t)+\int_{0}^{t} A e^{(s-\tau+\epsilon) A} h(\tau) d \tau=e^{\epsilon A} h(t)+u_{\epsilon}(t) .
\end{aligned}
$$

Theorem 44.21. Suppose $A=A^{*}, A \leq 0, u_{0} \in H$ and $h:[0, \infty) \rightarrow H$ is continuous. Assume further that $h(t) \in D(A)$ for all $t \in[0, \infty)$ and $t \rightarrow A h(t)$ is continuous, then

$$
\begin{equation*}
u(t):=e^{t A} u_{0}+\int_{0}^{t} e^{(t-\tau) A} h(\tau) d \tau \tag{44.15}
\end{equation*}
$$

is the unique function $u \in C^{1}((0, \infty), H) \cap C([0, \infty), H)$ such that $u(t) \in D(A)$ for all $t>0$ satisfying the differential equation

$$
\dot{u}(t)=A u(t)+h(t) \text { for } t>0 \text { and } u(0+)=u_{0} .
$$

Proof. Uniqueness: If $v(t)$ is another such solution then $w(t):=u(t)-$ $v(t)$ satisfies,

$$
\dot{w}(t)=A w(t) \text { with } w(0+)=0
$$

which we have already seen implies $w=0$.

Existence: By linearity and Theorem 44.13 we may assume with out loss of generality that $u_{0}=0$ in which case

$$
u(t)=\int_{0}^{t} e^{(t-\tau) A} h(\tau) d \tau
$$

By Lemma 44.19, we know $\tau \in[0, t] \rightarrow e^{(t-\tau) A} h(\tau) \in H$ is continuous, so the integral in Eq. (44.15) is well defined. Similarly by Lemma 44.19,

$$
\tau \in[0, t] \rightarrow e^{(t-\tau) A} A h(\tau)=A e^{(t-\tau) A} h(\tau) \in H
$$

and so by Corollary 44.18, $u(t) \in D(A)$ for all $t \geq 0$ and

$$
A u(t)=\int_{0}^{t} A e^{(t-\tau) A} h(\tau) d \tau=\int_{0}^{t} e^{(t-\tau) A} A h(\tau) d \tau
$$

Let

$$
u_{\epsilon}(t)=\int_{0}^{t} e^{(t+\epsilon-\tau) A} h(\tau) d \tau
$$

be defined as in Lemma 44.20. Then using the dominated convergence theorem,

$$
\begin{aligned}
\sup _{t \leq T}\left\|u_{\epsilon}(t)-u(t)\right\| & \leq \sup _{t \leq T} \int_{0}^{t}\left\|\left(e^{(t+\epsilon-\tau) A}-e^{(t-\tau) A}\right) h(\tau)\right\| d \tau \\
& \leq \int_{0}^{T}\left\|\left(e^{\epsilon A}-I\right) h(\tau)\right\| d \tau \rightarrow 0 \text { as } \epsilon \downarrow 0
\end{aligned}
$$

$$
\sup _{t \leq T}\left\|A u_{\epsilon}(t)-A u(t)\right\| \leq \int_{0}^{T}\left\|\left(e^{\epsilon A}-I\right) A h(\tau)\right\| d \tau \rightarrow 0 \text { as } \epsilon \downarrow 0
$$

and

$$
\left\|\int_{0}^{t} e^{\epsilon A} h(\tau) d \tau-\int_{0}^{t} h(\tau) d \tau\right\| \leq \int_{0}^{t}\left\|\left(e^{\epsilon A}-I\right) h(\tau)\right\| d \tau \rightarrow 0 \text { as } \epsilon \downarrow 0 .
$$

Integrating Eq. (44.13) shows

$$
\begin{equation*}
u_{\epsilon}(t)=\int_{0}^{t} e^{\epsilon A} h(\tau) d \tau+\int_{0}^{t} A u_{\epsilon}(\tau) d \tau \tag{44.16}
\end{equation*}
$$

and then passing to the limit as $\epsilon \downarrow 0$ in this equations shows

$$
u(t)=\int_{0}^{t} h(\tau) d \tau+\int_{0}^{t} A u(\tau) d \tau
$$

This shows $u$ is differentiable and $\dot{u}(t)=h(t)+A u(t)$ for all $t>0$.

Theorem 44.22. Let $\alpha>0, h:[0, \infty) \rightarrow H$ be a locally $\alpha$ - Holder continuous function, $A=A^{*}, A \leq 0$ and $u_{0} \in H$. The function

$$
u(t):=e^{t A} u_{0}+\int_{0}^{t} e^{(t-\tau) A} h(\tau) d \tau
$$

is the unique function $u \in C^{1}((0, \infty), H) \cap C([0, \infty), H)$ such that $u(t) \in D(A)$ for all $t>0$ satisfying the differential equation

$$
\dot{u}(t)=A u(t)+h(t) \text { for } t>0 \text { and } u(0+)=u_{0} .
$$

(For more details see Pazy [9, §5.7].)
Proof. The proof of uniqueness is the same as in Theorem 44.21 and for existence we may assume $u_{0}=0$.

With out loss of generality we may assume $u_{0}=0$ so that

$$
u(t)=\int_{0}^{t} e^{(t-\tau) A} h(\tau) d \tau
$$

By Lemma 44.19, we know $\tau \in[0, t] \rightarrow e^{(t-\tau) A} h(\tau) \in H$ is continuous, so the integral defining $u$ is well defined. For $\epsilon>0$, let

$$
u_{\epsilon}(t):=\int_{0}^{t} e^{(t+\epsilon-\tau) A} h(\tau) d \tau=\int_{0}^{t} e^{(t-\tau) A} e^{\epsilon A} h(\tau) d \tau
$$

Notice that $v(\tau):=e^{\epsilon A} h(\tau) \in C^{\infty}(A)$ for all $\tau$ and moreover since $A e^{\epsilon A}$ is a bounded operator, it follows that $\tau \rightarrow A v(\tau)$ is continuous. So by Lemma 44.19, it follows that $\tau \in[0, t] \rightarrow A e^{(t-\tau) A} v(\tau) \in H$ is continuous as well. Hence we know $u_{\epsilon}(t) \in D(A)$ and

$$
A u_{\epsilon}(t)=\int_{0}^{t} A e^{(t-\tau) A} e^{\epsilon A} h(\tau) d \tau
$$

Now

$$
\begin{gathered}
A u_{\epsilon}(t)=\int_{0}^{t} A e^{(t+\epsilon-\tau) A} h(t) d \tau+\int_{0}^{t} A e^{(t+\epsilon-\tau) A}[h(\tau)-h(t)] d \tau \\
\int_{0}^{t} A e^{(t+\epsilon-\tau) A} h(t) d \tau=-\left.e^{(t+\epsilon-\tau) A} h(t)\right|_{\tau=0} ^{\tau=t}=e^{(t+\epsilon) A} h(t)-e^{\epsilon A} h(t)
\end{gathered}
$$

and

$$
\begin{aligned}
\left\|A e^{(t+\epsilon-\tau) A}[h(\tau)-h(t)]\right\| & \leq e^{-1} \frac{1}{(t+\epsilon-\tau)}\|h(\tau)-h(t)\| \\
& \leq C e^{-1} \frac{1}{(t+\epsilon-\tau)}|t-\tau|^{\alpha} \leq C e^{-1}|t-\tau|^{\alpha-1}
\end{aligned}
$$

These results along with the dominated convergence theorem shows $\lim _{\epsilon \downarrow 0} A u_{\epsilon}(t)$ exists and is given by

$$
\begin{aligned}
\lim _{\epsilon \downarrow 0} A u_{\epsilon}(t) & =\lim _{\epsilon \downarrow 0}\left[e^{(t+\epsilon) A} h(t)-e^{\epsilon A} h(t)\right] \\
& +\lim _{\epsilon \downarrow 0} \int_{0}^{t} A e^{(t+\epsilon-\tau) A}[h(\tau)-h(t)] d \tau \\
& =e^{t A} h(t)-h(t)+\int_{0}^{t} A e^{(t-\tau) A}[h(\tau)-h(t)] d \tau
\end{aligned}
$$

Because $A$ is a closed operator, it follows that $u(t) \in D(A)$ and

$$
A u(t)=e^{t A} h(t)-h(t)+\int_{0}^{t} A e^{(t-\tau) A}[h(\tau)-h(t)] d \tau
$$

Claim: $t \rightarrow A u(t)$ is continuous. To prove this it suffices to show

$$
v(t):=A \int_{0}^{t} e^{(t-\tau) A}(h(\tau)-h(t)) d \tau
$$

is continuous and for this we have

$$
\begin{aligned}
v(t+\triangle)-v(t) & =\int_{0}^{t+\Delta} A e^{(t+\Delta-\tau) A}(h(\tau)-h(t+\triangle)) d \tau \\
& -\int_{0}^{t} A e^{(t-\tau) A}(h(\tau)-h(t)) d \tau \\
& =I+I I
\end{aligned}
$$

where

$$
\begin{aligned}
I & =\int_{t}^{t+\Delta} A e^{(t+\Delta-\tau) A}(h(\tau)-h(t+\triangle)) d \tau \text { and } \\
I I & =\int_{0}^{t}\left[A e^{(t+\triangle-\tau) A}(h(\tau)-h(t+\triangle))-A e^{(t-\tau) A}(h(\tau)-h(t))\right] d \tau \\
& =\int_{0}^{t}\left[A e^{(t+\triangle-\tau) A}(h(\tau)-h(t))-A e^{(t-\tau) A}(h(\tau)-h(t))\right] d \tau \\
& +\int_{0}^{t}\left[A e^{(t+\triangle-\tau) A}(h(t)-h(t+\triangle))\right] d \tau \\
& =I I_{1}+I I_{2}
\end{aligned}
$$

and

$$
\begin{aligned}
& I I_{1}=\int_{0}^{t} A\left[e^{(t+\Delta-\tau) A}-e^{(t-\tau) A}\right](h(\tau)-h(t)) d \tau \text { and } \\
& I I_{2}=\left[e^{(t+\Delta) A}-e^{\Delta A}\right](h(t)-h(t+\triangle)) .
\end{aligned}
$$

We estimate $I$ as

$$
\begin{aligned}
\|I\| & \leq\left|\int_{t}^{t+\Delta}\left\|A e^{(t+\triangle-\tau) A}(h(\tau)-h(t+\triangle))\right\| d \tau\right| \\
& \left.\leq C\left|\int_{t}^{t+\Delta} \frac{1}{t+\triangle-\tau}\right| t+\triangle-\left.\tau\right|^{\alpha} d \tau \right\rvert\, \\
& =C \int_{0}^{|\Delta|} x^{\alpha-1} d x=C \alpha^{-1}|\Delta|^{\alpha} \rightarrow 0 \text { as } \Delta \rightarrow 0
\end{aligned}
$$

It is easily seen that $\left\|I I_{2}\right\| \leq 2 C|\Delta|^{\alpha} \rightarrow 0$ as $\Delta \rightarrow 0$ and

$$
\left\|A\left[e^{(t+\triangle-\tau) A}-e^{(t-\tau) A}\right](h(\tau)-h(t))\right\| \leq C|t-\tau|^{\alpha-1}
$$

which is integrable, so by the dominated convergence theorem,

$$
\left\|I I_{1}\right\| \leq \int_{0}^{t}\left\|A\left[e^{(t+\Delta-\tau) A}-e^{(t-\tau) A}\right](h(\tau)-h(t))\right\| d \tau \rightarrow 0 \text { as } \Delta \rightarrow 0
$$

This completes the proof of the claim.
Moreover,

$$
\begin{aligned}
A u_{\epsilon}(t)-A u(t) & =e^{(t+\epsilon) A} h(t)-e^{t A} h(t)+h(t)-e^{\epsilon A} h(t) \\
& +\int_{0}^{t} A\left(e^{(t+\epsilon-\tau) A}-e^{(t-\tau) A}\right)[h(\tau)-h(t)] d \tau
\end{aligned}
$$

so that

$$
\begin{aligned}
\left\|A u_{\epsilon}(t)-A u(t)\right\| \leq & 2\left\|h(t)-e^{\epsilon A} h(t)\right\| \\
& +\int_{0}^{t}\left\|A e^{(t-\tau) A}\left(e^{\epsilon A}-I\right)[h(\tau)-h(t)]\right\| d \tau \\
\leq & 2\left\|h(t)-e^{\epsilon A} h(t)\right\| \\
& +e^{-1} \int_{0}^{t} \frac{1}{|t-\tau|}\left\|\left(e^{\epsilon A}-I\right)[h(\tau)-h(t)]\right\| d \tau
\end{aligned}
$$

from which it follows $A u_{\epsilon}(t) \rightarrow A u(t)$ boundedly. We may now pass to the limit in Eq. (44.16) to find

$$
\begin{aligned}
u(t)=\lim _{\epsilon \downarrow 0} u_{\epsilon}(t) & =\lim _{\epsilon \downarrow 0}\left[\int_{0}^{t} e^{\epsilon A} h(\tau) d \tau+\int_{0}^{t} A u_{\epsilon}(\tau) d \tau\right] \\
& =\int_{0}^{t} h(\tau) d \tau+\int_{0}^{t} A u(\tau) d \tau
\end{aligned}
$$

from which it follows that $u \in C^{1}((0, \infty), H)$ and $\dot{u}(t)=h(t)+A u(t)$.

## Heat Equation

The heat equation for a function $u: \mathbb{R}_{+} \times \mathbb{R}^{n} \rightarrow \mathbb{C}$ is the partial differential equation

$$
\begin{equation*}
\left(\partial_{t}-\frac{1}{2} \Delta\right) u=0 \text { with } u(0, x)=f(x) \tag{45.1}
\end{equation*}
$$

where $f$ is a given function on $\mathbb{R}^{n}$. By Fourier transforming Eq. (45.1) in the $x$ - variables only, one finds that (45.1) implies that

$$
\begin{equation*}
\left(\partial_{t}+\frac{1}{2}|\xi|^{2}\right) \hat{u}(t, \xi)=0 \text { with } \hat{u}(0, \xi)=\hat{f}(\xi) . \tag{45.2}
\end{equation*}
$$

and hence that $\hat{u}(t, \xi)=e^{-t|\xi|^{2} / 2} \hat{f}(\xi)$. Inverting the Fourier transform then shows that

$$
u(t, x)=\mathcal{F}^{-1}\left(e^{-t|\xi|^{2} / 2} \hat{f}(\xi)\right)(x)=\left(\mathcal{F}^{-1}\left(e^{-t|\xi|^{2} / 2}\right) \star f\right)(x)=: e^{t \Delta / 2} f(x)
$$

From Example 32.4,

$$
\mathcal{F}^{-1}\left(e^{-t|\xi|^{2} / 2}\right)(x)=p_{t}(x)=t^{-n / 2} e^{-\frac{1}{2 t}|x|^{2}}
$$

and therefore,

$$
u(t, x)=\int_{\mathbb{R}^{n}} p_{t}(x-y) f(y) \mathbf{d} y
$$

This suggests the following theorem.

## Theorem 45.1. Let

$$
\begin{equation*}
p_{t}(x-y):=(2 \pi t)^{-n / 2} e^{-|x-y|^{2} / 2 t} \tag{45.3}
\end{equation*}
$$

be the heat kernel on $\mathbb{R}^{n}$. Then

$$
\begin{equation*}
\left(\partial_{t}-\frac{1}{2} \Delta_{x}\right) p_{t}(x-y)=0 \text { and } \lim _{t \downarrow 0} p_{t}(x-y)=\delta_{x}(y), \tag{45.4}
\end{equation*}
$$

where $\delta_{x}$ is the $\delta$-function at $x$ in $\mathbb{R}^{n}$. More precisely, if $f$ is a continuous bounded function on $\mathbb{R}^{n}$, then

$$
u(t, x)=\int_{\mathbb{R}^{n}} p_{t}(x-y) f(y) d y
$$

is a solution to $E q$. (45.1) where $u(0, x):=\lim _{t \downarrow 0} u(t, x)$.
Proof. Direct computations show that $\left(\partial_{t}-\frac{1}{2} \Delta_{x}\right) p_{t}(x-y)=0$ and an application of Theorem 11.21 shows $\lim _{t \downarrow 0} p_{t}(x-y)=\delta_{x}(y)$ or equivalently that $\lim _{t \downarrow 0} \int_{\mathbb{R}^{n}} p_{t}(x-y) f(y) d y=f(x)$ uniformly on compact subsets of $\mathbb{R}^{n}$. This shows that $\lim _{t \downarrow 0} u(t, x)=f(x)$ uniformly on compact subsets of $\mathbb{R}^{n}$.

Proposition 45.2 (Properties of $e^{t \Delta / 2}$ ).

1. For $f \in L^{2}\left(\mathbb{R}^{n}, d x\right)$, the function

$$
\left(e^{t \Delta / 2} f\right)(x)=\left(P_{t} f\right)(x)=\int_{\mathbb{R}^{n}} f(y) \frac{e^{-\frac{1}{2 t}|x-y|^{2}}}{(2 \pi t)^{n / 2}} d y
$$

is smooth in $(t, x)$ for $t>0$ and $x \in \mathbb{R}^{n}$ and is in fact real analytic. 2. $e^{t \Delta / 2}$ acts as a contraction on $L^{p}\left(\mathbb{R}^{n}, d x\right)$ for all $p \in[0, \infty]$ and $t>0$. Indeed,
3. Moreover, $p_{t} * f \rightarrow f$ in $L^{p}$ as $t \rightarrow 0$.

Proof. Item 1. is fairly easy to check and is left the reader. One just notices that $p_{t}(x-y)$ analytically continues to $\operatorname{Re} t>0$ and $x \in \mathbb{C}^{n}$ and then shows that it is permissible to differentiate under the integral.

Item 2.

$$
\left|\left(p_{t} * f\right)(x)\right| \leq \int_{\mathbb{R}^{n}}|f(y)| p_{t}(x-y) d y
$$

and hence with the aid of Jensen's inequality we have,

$$
\left\|p_{t} * f\right\|_{L^{p}}^{p} \leq \int_{\mathbb{R}^{n}} \int_{\mathbb{R}^{n}}|f(y)|^{p} p_{t}(x-y) d y d x=\|f\|_{L^{p}}^{p}
$$

So $P_{t}$ is a contraction $\forall t>0$.
Item 3. It suffices to show, because of the contractive properties of $p_{t}{ }^{*}$, that $p_{t} * f \rightarrow f$ as $t \downarrow 0$ for $f \in C_{c}\left(\mathbb{R}^{n}\right)$. Notice that if $f$ has support in the ball of radius $R$ centered at zero, then

$$
\begin{aligned}
\left|\left(p_{t} * f\right)(x)\right| & \leq \int_{\mathbb{R}^{n}}|f(y)| P_{t}(x-y) d y \leq\|f\|_{\infty} \int_{|y| \leq R} P_{t}(x-y) d y \\
& =\|f\|_{\infty} C R^{n} e^{-\frac{1}{2 t}(|x|-R)^{2}}
\end{aligned}
$$

and hence

$$
\left\|p_{t} * f-f\right\|_{L^{p}}^{p}=\int_{|y| \leq R}\left|p_{t} * f-f\right|^{p} d y+\|f\|_{\infty} C R^{n} e^{-\frac{1}{2 t}(|x|-R)^{2}}
$$

Therefore $p_{t} * f \rightarrow f$ in $L^{p}$ as $t \downarrow 0 \quad \forall f \in C_{c}\left(\mathbb{R}^{n}\right)$.

Theorem 45.3 (Forced Heat Equation). Suppose $g \in C_{b}\left(\mathbb{R}^{d}\right)$ and $f \in$ $C_{b}^{1,2}\left([0, \infty) \times \mathbb{R}^{d}\right)$ then

$$
u(t, x):=p_{t} * g(x)+\int_{0}^{t} p_{t-\tau} * f(\tau, x) d \tau
$$

solves

$$
\frac{\partial u}{\partial t}=\frac{1}{2} \triangle u+f \text { with } u(0, \cdot)=g
$$

Proof. Because of Theorem 45.1, we may with out loss of generality assume $g=0$ in which case

$$
u(t, x)=\int_{0}^{t} p_{t} * f(t-\tau, x) d \tau
$$

Therefore

$$
\begin{aligned}
\frac{\partial u}{\partial t}(t, x) & =p_{t} * f(0, x)+\int_{0}^{t} p_{\tau} * \frac{\partial}{\partial t} f(t-\tau, x) d \tau \\
& =p_{t} * f_{0}(x)-\int_{0}^{t} p_{\tau} * \frac{\partial}{\partial \tau} f(t-\tau, x) d \tau
\end{aligned}
$$

and

$$
\frac{\triangle}{2} u(t, x)=\int_{0}^{t} p_{t} * \frac{\triangle}{2} f(t-\tau, x) d \tau
$$

Hence we find, using integration by parts and approximate $\delta$ - function arguments, that

$$
\begin{aligned}
\left(\frac{\partial}{\partial t}-\frac{\triangle}{2}\right) u(t, x)= & p_{t} * f_{0}(x)+\int_{0}^{t} p_{\tau} *\left(-\frac{\partial}{\partial \tau}-\frac{1}{2} \triangle\right) f(t-\tau, x) d \tau \\
= & p_{t} * f_{0}(x) \\
& +\lim _{\epsilon \downarrow 0} \int_{\epsilon}^{t} p_{\tau} *\left(-\frac{\partial}{\partial \tau}-\frac{1}{2} \triangle\right) f(t-\tau, x) d \tau \\
= & p_{t} * f_{0}(x)-\left.\lim _{\epsilon \downharpoonright 0} p_{\tau} * f(t-\tau, x)\right|_{\epsilon} ^{t} \\
+ & \lim _{\epsilon \downarrow 0} \int_{\epsilon}^{t}\left(\frac{\partial}{\partial \tau}-\frac{1}{2} \triangle\right) p_{\tau} * f(t-\tau, x) d \tau \\
= & p_{t} * f_{0}(x)-p_{t} * f_{0}(x)+\lim _{\epsilon \downarrow 0} p_{\epsilon} * f(t-\epsilon, x) \\
= & f(t, x)
\end{aligned}
$$

### 45.1 Extensions of Theorem 45.1

Proposition 45.4. Suppose $f: \mathbb{R}^{n} \rightarrow \mathbb{R}$ is a measurable function and there exists constants $c, C<\infty$ such that

$$
|f(x)| \leq C e^{\frac{c}{2}|x|^{2}}
$$

Then $u(t, x):=p_{t} * f(x)$ is smooth for $(t, x) \in\left(0, c^{-1}\right) \times \mathbb{R}^{n}$ and for all $k \in \mathbb{N}$ and all multi-indices $\alpha$,

$$
\begin{equation*}
D^{\alpha}\left(\frac{\partial}{\partial t}\right)^{k} u(t, x)=\left(D^{\alpha}\left(\frac{\partial}{\partial t}\right)^{k} p_{t}\right) * f(x) \tag{45.5}
\end{equation*}
$$

In particular $u$ satisfies the heat equation $u_{t}=\Delta u / 2$ on $\left(0, c^{-1}\right) \times \mathbb{R}^{n}$.
Proof. The reader may check that

$$
D^{\alpha}\left(\frac{\partial}{\partial t}\right)^{k} p_{t}(x)=q\left(t^{-1}, x\right) p_{t}(x)
$$

where $q$ is a polynomial in its variables. Let $x_{0} \in \mathbb{R}^{n}$ and $\epsilon>0$ be small, then for $x \in B\left(x_{0}, \epsilon\right)$ and any $\beta>0$,

$$
\begin{aligned}
|x-y|^{2} & =|x|^{2}-2|x||y|+|y|^{2} \geq|y|^{2}+|x|^{2}-\left(\beta^{-2}|x|^{2}+\beta^{2}|y|^{2}\right) \\
& \geq\left(1-\beta^{2}\right)|y|^{2}-\left(\beta^{-2}-1\right)\left(\left|x_{0}\right|^{2}+\epsilon\right)
\end{aligned}
$$

Hence

$$
\begin{aligned}
g(y) & :=\sup \left\{\left|D^{\alpha}\left(\frac{\partial}{\partial t}\right)^{k} p_{t}(x-y) f(y)\right|: \epsilon \leq t \leq c-\epsilon \& x \in B\left(x_{0}, \epsilon\right)\right\} \\
& \leq \sup \left\{\left|q\left(t^{-1}, x-y\right) \frac{e^{-\frac{1}{2 t}|x-y|^{2}}}{(2 \pi t)^{n / 2}} C e^{\frac{c}{2}|y|^{2}}\right|: \epsilon \leq t \leq c-\epsilon \& x \in B\left(x_{0}, \epsilon\right)\right\} \\
& \leq C\left(\beta, x_{0}, \epsilon\right) \sup \left\{\left|q\left(t^{-1}, x-y\right) \frac{e^{\left[-\frac{1}{2 t}\left(1-\beta^{2}\right)+\frac{c}{2}\right]|y|^{2}}}{(2 \pi t)^{n / 2}}\right|: \begin{array}{c}
\epsilon \leq t \leq c-\epsilon \text { and } \\
x \in B\left(x_{0}, \epsilon\right)
\end{array}\right\}
\end{aligned}
$$

By choosing $\beta$ close to 0 , the reader should check using the above expression that for any $0<\delta<(1 / t-c) / 2$ there is a $\tilde{C}<\infty$ such that $g(y) \leq \tilde{C} e^{-\delta|y|^{2}}$. In particular $g \in L^{1}\left(\mathbb{R}^{n}\right)$. Hence one is justified in differentiating past the integrals in $p_{t} * f$ and this proves Eq. (45.5).

Lemma 45.5. There exists a polynomial $q_{n}(x)$ such that for any $\beta>0$ and $\delta>0$,

$$
\int_{\mathbb{R}^{n}} 1_{|y| \geq \delta} e^{-\beta|y|^{2}} d y \leq \delta^{n} q_{n}\left(\frac{1}{\beta \delta^{2}}\right) e^{-\beta \delta^{2}}
$$

Proof. Making the change of variables $y \rightarrow \delta y$ and then passing to polar coordinates shows
$\int_{\mathbb{R}^{n}} 1_{|y| \geq \delta} e^{-\beta|y|^{2}} d y=\delta^{n} \int_{\mathbb{R}^{n}} 1_{|y| \geq 1} e^{-\beta \delta^{2}|y|^{2}} d y=\sigma\left(S^{n-1}\right) \delta^{n} \int_{1}^{\infty} e^{-\beta \delta^{2} r^{2}} r^{n-1} d r$.
Letting $\lambda=\beta \delta^{2}$ and $\phi_{n}(\lambda):=\int_{r=1}^{\infty} e^{-\lambda r^{2}} r^{n} d r$, integration by parts shows

$$
\begin{aligned}
\phi_{n}(\lambda) & =\int_{r=1}^{\infty} r^{n-1} d\left(\frac{e^{-\lambda r^{2}}}{-2 \lambda}\right)=\frac{1}{2 \lambda} e^{-\lambda}+\frac{1}{2} \int_{r=1}^{\infty}(n-1) r^{(n-2)} \frac{e^{-\lambda r^{2}}}{\lambda} d r \\
& =\frac{1}{2 \lambda} e^{-\lambda}+\frac{n-1}{2 \lambda} \phi_{n-2}(\lambda)
\end{aligned}
$$

Iterating this equation implies

$$
\phi_{n}(\lambda)=\frac{1}{2 \lambda} e^{-\lambda}+\frac{n-1}{2 \lambda}\left(\frac{1}{2 \lambda} e^{-\lambda}+\frac{n-3}{2 \lambda} \phi_{n-4}(\lambda)\right)
$$

and continuing in this way shows

$$
\phi_{n}(\lambda)=e^{-\lambda} r_{n}\left(\lambda^{-1}\right)+\frac{(n-1)!!}{2^{\delta} \lambda^{\delta}} \phi_{i}(\lambda)
$$

where $\delta$ is the integer part of $n / 2, i=0$ if $n$ is even and $i=1$ if $n$ is odd and $r_{n}$ is a polynomial. Since

$$
\phi_{0}(\lambda)=\int_{r=1}^{\infty} e^{-\lambda r^{2}} d r \leq \phi_{1}(\lambda)=\int_{r=1}^{\infty} r e^{-\lambda r^{2}} d r=\frac{e^{-\lambda}}{2 \lambda}
$$

it follows that

$$
\phi_{n}(\lambda) \leq e^{-\lambda} q_{n}\left(\lambda^{-1}\right)
$$

for some polynomial $q_{n}$.
Proposition 45.6. Suppose $f \in C\left(\mathbb{R}^{n}, \mathbb{R}\right)$ such that $|f(x)| \leq C e^{\frac{c}{2}|x|^{2}}$ then $p_{t} * f \rightarrow f$ uniformly on compact subsets as $t \downarrow 0$. In particular in view of Proposition 45.4, $u(t, x):=p_{t} * f(x)$ is a solution to the heat equation with $u(0, x)=f(x)$.

Proof. Let $M>0$ be fixed and assume $|x| \leq M$ throughout. By uniform continuity of $f$ on compact set, given $\epsilon>0$ there exists $\delta=\delta(t)>0$ such that $|f(x)-f(y)| \leq \epsilon$ if $|x-y| \leq \delta$ and $|x| \leq M$. Therefore, choosing $a>c / 2$ sufficiently small,
$\left|p_{t} * f(x)-f(x)\right|=\left|\int p_{t}(y)[f(x-y)-f(x)] d y\right|$

$$
\leq \int p_{t}(y)|f(x-y)-f(x)| d y
$$

$$
\leq \epsilon \int_{|y| \leq \delta} p_{t}(y) d y+\frac{C}{(2 \pi t)^{n / 2}} \int_{|y| \geq \delta}\left[e^{\frac{c}{2}|x-y|^{2}}+e^{\frac{c}{2}|x|^{2}}\right] e^{-\frac{1}{2 t}|y|^{2}} d y
$$

$$
\leq \epsilon+\tilde{C}(2 \pi t)^{-n / 2} \int_{|y| \geq \delta} e^{-\left(\frac{1}{2 t}-a\right)|y|^{2}} d y .
$$

So by Lemma 45.5, it follows that

$$
\left|p_{t} * f(x)-f(x)\right| \leq \epsilon+\tilde{C}(2 \pi t)^{-n / 2} \delta^{n} q_{n}\left(\frac{1}{\beta\left(\frac{1}{2 t}-a\right)^{2}}\right) e^{-\left(\frac{1}{2 t}-a\right) \delta^{2}}
$$

and therefore

$$
\limsup _{t \downarrow 0} \sup _{|x| \leq M}\left|p_{t} * f(x)-f(x)\right| \leq \epsilon \rightarrow 0 \text { as } \epsilon \downarrow 0 .
$$

Lemma 45.7. If $q(x)$ is a polynomial on $\mathbb{R}^{n}$, then

$$
\int_{\mathbb{R}^{n}} p_{t}(x-y) q(y) d y=\sum_{n=0}^{\infty} \frac{t^{n}}{n!} \frac{\Delta^{n}}{2^{n}} q(x) .
$$

## Proof. Since

$f(t, x):=\int_{\mathbb{R}^{n}} p_{t}(x-y) q(y) d y=\int_{\mathbb{R}^{n}} p_{t}(y) \sum a_{\alpha \beta} x^{\alpha} y^{\beta} d y=\sum C_{\alpha}(t) x^{\alpha}$,
$f(t, x)$ is a polynomial in $x$ of degree no larger than that of $q$. Moreover $f(t, x)$ solves the heat equation and $f(t, x) \rightarrow q(x)$ as $t \downarrow 0$. Since $g(t, x):=$ $\sum_{n=0}^{\infty} \frac{t^{n}}{n!} \frac{\Delta^{n}}{2^{n}} q(x)$ has the same properties of $f$ and $\Delta$ is a bounded operator when acting on polynomials of a fixed degree we conclude $f(t, x)=g(t, x)$.
Example 45.8. Suppose $q(x)=x_{1} x_{2}+x_{3}^{4}$, then

$$
\begin{aligned}
e^{t \Delta / 2} q(x) & =x_{1} x_{2}+x_{3}^{4}+\frac{t}{2} \Delta\left(x_{1} x_{2}+x_{3}^{4}\right)+\frac{t^{2}}{2!\cdot 4} \Delta^{2}\left(x_{1} x_{2}+x_{3}^{4}\right) \\
& =x_{1} x_{2}+x_{3}^{4}+\frac{t}{2} 12 x_{3}^{2}+\frac{t^{2}}{2!\cdot 4} 4! \\
& =x_{1} x_{2}+x_{3}^{4}+6 t x_{3}^{2}+3 t^{2} .
\end{aligned}
$$

Proposition 45.9. Suppose $f \in C^{\infty}\left(\mathbb{R}^{n}\right)$ and there exists a constant $C<\infty$ such that

$$
\sum_{|\alpha|=2 N+2}\left|D^{\alpha} f(x)\right| \leq C e^{C|x|^{2}},
$$

then

$$
\left(p_{t} * f\right)(x)=" e^{t \Delta / 2} f(x) "=\sum_{k=0}^{N} \frac{t^{k}}{k!} \Delta^{k} f(x)+O\left(t^{N+1}\right) \text { as } t \downarrow 0
$$

Proof. Fix $x \in \mathbb{R}^{n}$ and let

$$
f_{N}(y):=\sum_{|\alpha| \leq 2 N+1} \frac{1}{\alpha!} D^{\alpha} f(x) y^{\alpha} .
$$

Then by Taylor's theorem with remainder

$$
\begin{aligned}
\left|f(x+y)-f_{N}(y)\right| & \leq C|y|^{2 N+2} \sup _{t \in[0,1]} e^{C|x+t y|^{2}} \\
& \leq C|y|^{2 N+2} e^{2 C\left[|x|^{2}+|y|^{2}\right]} \leq \tilde{C}|y|^{2 N+2} e^{2 C|y|^{2}}
\end{aligned}
$$

and thus

$$
\begin{aligned}
\mid \int_{\mathbb{R}^{n}} & p_{t}(y) f(x+y) d y-\int_{\mathbb{R}^{n}} p_{t}(y) f_{N}(y) d y \mid \\
& \leq \tilde{C} \int_{\mathbb{R}^{n}} p_{t}(y)|y|^{2 N+2} e^{2 C|y|^{2}} d y \\
& =\tilde{C} t^{N+1} \int_{\mathbb{R}^{n}} p_{1}(y)|y|^{2 N+2} e^{2 t^{2} C|y|^{2}} d y=O\left(t^{N+1}\right) .
\end{aligned}
$$

Since $f(x+y)$ and $f_{N}(y)$ agree to order $2 N+1$ for $y$ near zero, it follows that

$$
\begin{aligned}
\int_{\mathbb{R}^{n}} p_{t}(y) f_{N}(y) d y & =\sum_{k=0}^{N} \frac{t^{k}}{k!} \Delta^{k} f_{N}(0)=\left.\sum_{k=0}^{N} \frac{t^{k}}{k!} \Delta_{y}^{k} f(x+y)\right|_{y=0} \\
& =\sum_{k=0}^{N} \frac{t^{k}}{k!} \Delta^{k} f(x)
\end{aligned}
$$

which completes the proof.

### 45.2 Representation Theorem and Regularity

In this section, suppose that $\Omega$ is a bounded domain such that $\bar{\Omega}$ is a $C^{2}-$ submanifold with $C^{2}$ boundary and for $T>0$ let $\Omega_{T}:=(0, T) \times \Omega$, and

$$
\Gamma_{T}:=([0, T] \times \partial \Omega) \cup(\{0\} \times \Omega) \subset \operatorname{bd}\left(\Omega_{T}\right)=([0, T] \times \partial \Omega) \cup(\{0, T\} \times \Omega)
$$

as in Figure 45.1 below.
Theorem 45.10 (Representation Theorem). Suppose $u \in C^{2,1}\left(\bar{\Omega}_{T}\right)$ $\left(\bar{\Omega}_{T}=\bar{\Omega}_{T}=[0, T] \times \bar{\Omega}\right)$ solves $u_{t}=\frac{1}{2} \Delta u+f$ on $\bar{\Omega}_{T}$. Then

$$
\begin{aligned}
u(T, x) & =\int_{\Omega} p_{T}(x-y) u(0, y) d y+\int_{[0, T] \times \Omega} p_{T-t}(x-y) f(t, y) d y d t \\
& +\frac{1}{2} \int_{[0, T] \times \partial \Omega}\left[\frac{\partial p_{T-t}}{\partial n_{y}}(x-y) u(t, y)-p_{T-t}(x-y) \frac{\partial u}{\partial n}(y)\right] d \sigma(y) d t
\end{aligned}
$$



Fig. 45.1. A cylindrical region $\Omega_{T}$ and the parabolic boundary $\Gamma_{T}$.

Proof. For $v \in C^{2,1}([0, T] \times \bar{\Omega})$, integration by parts shows

$$
\begin{aligned}
\int_{\Omega_{T}} f v d y d t= & \int_{\Omega_{T}} v\left(u_{t}-\frac{1}{2} \triangle v\right) d y d t \\
= & \int_{\Omega_{T}}\left(-v_{t}+\frac{1}{2} \nabla v \cdot \nabla u\right) d y d t+\left.\int_{\Omega} v u\right|_{t=0} ^{t=T} d y \\
& +\frac{1}{2} \int_{[0, T] \times \partial \Omega} v \frac{\partial v}{\partial n} d t d \sigma \\
= & \int_{\Omega_{T}}\left(-v_{t}-\frac{1}{2} \triangle v\right) u d y d t+\left.\int_{\Omega} v u\right|_{0} ^{T} d y \\
& +\frac{1}{2} \int_{[0, T] \times \partial \Omega}\left(\frac{\partial u}{\partial n} u-v \frac{\partial u}{\partial n}\right) d \sigma d t .
\end{aligned}
$$

Given $\epsilon>0$, taking $v(t, y):=p_{T+\epsilon-t}(x-y)$ (note that $v_{t}+\frac{1}{2} \triangle v=0$ and $v \in C^{2,1}([0, T] \times \Omega)$ ) implies

$$
\begin{aligned}
& \int_{[0, T] \times \Omega} f(t, y) p_{T+\epsilon-t}(x-y) d y d t=0+\int_{\Omega} p_{\epsilon}(x-y) u(t, y) d y \\
& \quad-\int_{\Omega} p_{T+\epsilon}(x-y) u(t, y) d y \\
& \quad+\frac{1}{2} \int_{[0, T] \times \partial \Omega}\left[\begin{array}{c}
\frac{\partial p_{T+\epsilon-t}(x-y)}{\partial n_{y}} u(t, y) \\
-p_{T+\epsilon-t}(x-y) \frac{\partial u}{\partial n}(y)
\end{array}\right] d \sigma(y) d t
\end{aligned}
$$

Let $\epsilon \downarrow 0$ above to complete the proof.
Corollary 45.11. Suppose $f:=0$ so $u_{t}(t, x)=\frac{1}{2} \triangle u(t, x)$. Then $u \in$ $C^{\infty}((0, T) \times \Omega)$.

45 Heat Equation

Proof. Extend $p_{t}(x)$ for $t \leq 0$ by setting $p_{t}(x):=0$ if $t \leq 0$. It is not to hard to check that this extension is $C^{\infty}$ on $\mathbb{R} \times \mathbb{R}^{n} \backslash\{0\}$. Using this notation we may write Eq. (45.6) as

$$
\begin{aligned}
& u(t, x)=\int_{\Omega} p_{t}(x-y) u(0, y) d y \\
& +\frac{1}{2} \int_{[0, \infty) \times \partial \Omega}\left[\frac{\partial p_{t-\tau}}{\partial n_{y}}(x-y) u(t, y)-p_{T-t}(x-y) \frac{\partial u}{\partial n}(y)\right] d \sigma(y) d \tau
\end{aligned}
$$

The result follows since now it permissible to differentiate under the integral to show $u \in C^{\infty}((0, T) \times \Omega)$.

Remark 45.12. Since $x \rightarrow p_{t}(x)$ is analytic one may show that $x \rightarrow u(t, x)$ is analytic for all $x \in \Omega$.

### 45.3 Weak Max Principles

Notation 45.13 Let $a_{i j}, b_{j} \in C\left(\bar{\Omega}_{T}\right)$ satisfy $a_{i j}=a_{j i}$ and for $u \in C^{2}(\Omega)$ let

$$
\begin{equation*}
L u(t, x)=\sum_{i, j=1}^{n} a_{i j}(t, x) u_{x_{i} x_{j}}(x)+\sum_{i=1}^{n} b_{i}(t, x) u_{x_{i}}(x) . \tag{45.7}
\end{equation*}
$$

We say $L$ is elliptic if there exists $\theta>0$ such that

$$
\sum a_{i j}(t, x) \xi_{i} \xi_{j} \geq \theta|\xi|^{2} \text { for all } \xi \in \mathbb{R}^{n} \text { and }(t, x) \in \bar{\Omega}_{T}
$$

Assumption 3 In this section we assume $L$ is elliptic. As an example $L=$ $\frac{1}{2} \Delta$ is elliptic.

Lemma 45.14. Let $L$ be an elliptic operator as above and suppose $u \in C^{2}(\Omega)$ and $x_{0} \in \Omega$ is a point where $u(x)$ has a local maximum. Then $L u\left(t, x_{0}\right) \leq 0$ for all $t \in[0, T]$.

Proof. Fix $t \in[0, T]$ and set $B_{i j}=u_{x_{i} x_{j}}\left(x_{0}\right), A_{i j}:=a_{i j}\left(t, x_{0}\right)$ and let $\left\{e_{i}\right\}_{i=1}^{n}$ be an orthonormal basis for $\mathbb{R}^{n}$ such that $A e_{i}=\lambda_{i} e_{i}$. Notice that $\lambda_{i} \geq \theta>0$ for all $i$. By the first derivative test, $u_{x_{i}}\left(x_{0}\right)=0$ for all $i$ and hence

$$
\begin{aligned}
L u\left(t, x_{0}\right) & =\sum A_{i j} B_{i j}=\sum A_{j i} B_{i j}=\operatorname{tr}(A B) \\
& =\sum e_{i} \cdot A B e_{i}=\sum A e_{i} \cdot B e_{i}=\sum_{i} \lambda_{i} e_{i} \cdot B e_{i} \\
& =\sum_{i} \lambda_{i} \partial_{e_{i}}^{2} u\left(t, x_{0}\right) \leq 0 .
\end{aligned}
$$

The last inequality if a consequence of the second derivative test which asserts $\partial_{v}^{2} u\left(t, x_{0}\right) \leq 0$ for all $v \in \mathbb{R}^{n}$.

Theorem 45.15 (Elliptic weak maximum principle). Let $\Omega$ be a bounded domain and $L$ be an elliptic operator as in Eq. (45.7). We now assume that $a_{i j}$ and $b_{j}$ are functions of $x$ alone. For each $u \in C(\bar{\Omega}) \cap C^{2}(\Omega)$ such that $L u \geq 0$ on $\Omega$ (i.e. $u$ is $L$-subharmonic) we have

$$
\begin{equation*}
\max _{\bar{\Omega}} u \leq \max _{\operatorname{bd}(\Omega)} u \tag{45.8}
\end{equation*}
$$

Proof. Let us first assume $L u>0$ on $\Omega$. If $u$ and had an interior local maximum at $x_{0} \in \Omega$ then by Lemma 45.14, $L u\left(x_{0}\right) \leq 0$ which contradicts the assumption that $L u\left(x_{0}\right)>0$. So if $L u>0$ on $\Omega$ we conclude that Eq. (45.8) holds.

Now suppose that $L u \geq 0$ on $\Omega$. Let $\phi(x):=e^{\lambda x_{1}}$ with $\lambda>0$, then

$$
L \phi(x)=\left(\lambda^{2} a_{11}(x)+b_{1}(x) \lambda\right) e^{\lambda x_{1}} \geq \lambda\left(\lambda \theta+b_{1}(x)\right) e^{\lambda x_{1}}
$$

By continuity of $b(x)$ we may choose $\lambda$ sufficiently large so that $\lambda \theta+b_{1}(x)>0$ on $\bar{\Omega}$ in which case $L \phi>0$ on $\Omega$. The results in the first paragraph may now be applied to $u_{\epsilon}(x):=u(x)+\epsilon \phi(x)$ (for any $\epsilon>0$ ) to learn

$$
u(x)+\epsilon \phi(x)=u_{\epsilon}(x) \leq \max _{\operatorname{bd}(\Omega)} u_{\epsilon} \leq \max _{\mathrm{bd}(\Omega)} u+\epsilon \max _{\mathrm{bd}(\Omega)} \phi \text { for all } x \in \bar{\Omega}
$$

Letting $\epsilon \downarrow 0$ in this expression then implies

$$
u(x) \leq \max _{\operatorname{bd}(\Omega)} u \text { for all } x \in \bar{\Omega}
$$

which is equivalent to Eq. (45.8).
Theorem 45.16 (Parabolic weak maximum principle). Assume $u \in$ $C^{1,2}\left(\bar{\Omega}_{T} \backslash \Gamma_{T}\right) \cap C\left(\bar{\Omega}_{T}\right)$.

1. If $u_{t}-L u \leq 0$ in $\Omega_{T}$ then

$$
\begin{equation*}
\max _{\bar{\Omega}_{T}} u=\max _{\Gamma_{T}} u \tag{45.9}
\end{equation*}
$$

2. If $u_{t}-L u \geq 0$ in $\Omega_{T}$ then $\min _{\bar{\Omega}_{T}} u=\min _{\Gamma_{T}} u$.

Proof. Item 1. follows from Item 2. by replacing $u \rightarrow-u$, so it suffices to prove item 1. We begin by assuming $u_{t}-L u<0$ on $\bar{\Omega}_{T}$ and suppose for the sake of contradiction that there exists a point $\left(t_{0}, x_{0}\right) \in \bar{\Omega}_{T} \backslash \Gamma_{T}$ such that $u\left(t_{0}, x_{0}\right)=\max _{\bar{\Omega}_{T}} u$.

1. If $\left(t_{0}, x_{0}\right) \in \Omega_{T}$ (i.e. $\left.0<t_{0}<T\right)$ then by the first derivative test $\frac{\partial u}{\partial t}\left(t_{0}, x_{0}\right)=0$ and by Lemma 45.14 Lu( $\left.t_{0}, x_{0}\right) \leq 0$. Therefore,

$$
\left(u_{t}-L u\right)\left(t_{0}, x_{0}\right)=-L u\left(t_{0}, x_{0}\right) \geq 0
$$

which contradicts the assumption that $u_{t}-L u<0$ in $\Omega_{T}$.
2. If $\left(t_{0}, x_{0}\right) \in \bar{\Omega}_{T} \backslash \Gamma_{T}$ with $t_{0}=T$, then by the first derivative test $\frac{\partial u}{\partial t}\left(T, x_{0}\right) \geq 0$ and by Lemma $45.14 L u\left(t_{0}, x_{0}\right) \leq 0$. So again

$$
\left(u_{t}-L u\right)\left(t_{0}, x_{0}\right) \geq 0
$$

which contradicts the assumption that $u_{t}-L u<0$ in $\Omega_{T}$.
Thus we have proved Eq. (45.9) holds if $u_{t}-L u<0$ on $\bar{\Omega}_{T}$. Finally if $u_{t}-L u \leq 0$ on $\Omega_{T}$ and $\epsilon>0$, the function $u^{\epsilon}(t, x):=u(t, x)-\epsilon t$ satisfies $u_{t}^{\epsilon}-L u^{\epsilon} \leq-\epsilon<0$. Therefore by what we have just proved

$$
u(t, x)-\epsilon t \leq \max _{\bar{\Omega}_{T}} u^{\epsilon}=\max _{\Gamma_{T}} u^{\epsilon} \leq \max _{\Gamma_{T}} u \text { for all }(t, x) \in \bar{\Omega}_{T} .
$$

Letting $\epsilon \downarrow 0$ in the last equation shows that Eq. (45.9) holds.
Corollary 45.17. There is at most one solution $u \in C^{1,2}\left(\bar{\Omega}_{T} \backslash \Gamma_{T}\right) \cap C\left(\bar{\Omega}_{T}\right)$ to the partial differential equation

$$
\frac{\partial u}{\partial t}=L u \text { with } u=f \text { on } \Gamma_{T} .
$$

Proof. If there were another solution $v$, then $w:=u-v$ would solve $\frac{\partial w}{\partial t}=L w$ with $w=0$ on $\Gamma_{T}$. So by the maximum principle in Theorem 45.16, $w=0$ on $\bar{\Omega}_{T}$.

We now restrict back to $L=\frac{1}{2} \triangle$ and we wish to see what can be said when $\Omega=\mathbb{R}^{n}$ - an unbounded set.

Theorem 45.18. Suppose $u \in C\left([0, T] \times \mathbb{R}^{n}\right) \cap C^{2,1}\left((0, T) \times \mathbb{R}^{n}\right)$,

$$
u_{t}-\frac{1}{2} \triangle u \leq 0 \text { on }[0, T] \times \mathbb{R}^{n}
$$

and there exists constants $A, a<\infty$ such that

$$
u(t, x) \leq A e^{a|x|^{2}} \text { for }(t, x) \in(0, T) \times \mathbb{R}^{n}
$$

Then

$$
\sup _{(t, x) \in[0, T] \times \mathbb{R}^{n}} u(t, x) \leq K:=\sup _{x \in \mathbb{R}^{n}} u(0, x)
$$

Proof. Recall that

$$
p_{t}(x)=\left(\frac{1}{t}\right)^{n / 2} e^{-\frac{1}{2 t}|x|^{2}}=\left(\frac{1}{t}\right)^{n / 2} e^{-\frac{1}{2 t} x \cdot x}
$$

solves the heat equation

$$
\begin{equation*}
\partial_{t} p_{t}(x)=\frac{1}{2} \triangle p_{t}(x) \tag{45.10}
\end{equation*}
$$

Since both sides of Eq. (45.10) are analytic as functions in $x$, so ${ }^{1}$

$$
\frac{\partial p_{t}}{\partial t}(i x)=\frac{1}{2}\left(\triangle p_{t}\right)(i x)=-\frac{1}{2} \triangle_{x} p_{t}(i x)
$$

and therefore for all $\tau>0$ and $t<\tau$

$$
\frac{\partial p_{\tau-t}}{\partial t}(i x)=-\dot{p}_{\tau-t}(i x)=\frac{1}{2} \triangle_{x} p_{\tau-t}(i x) .
$$

That is to say the function

$$
\rho(t, x):=p_{\tau-t}(i x)=\left(\frac{1}{\tau-t}\right)^{n / 2} e^{\frac{1}{2(\tau-t)}|x|^{2}} \text { for } 0 \leq t<\tau
$$

solves the heat equation. (This can be checked directly as well.)
Let $\epsilon, \tau>0$ (to be chosen later) and set

$$
v(t, x)=u(t, x)-\epsilon \rho(t, x) \text { for } 0 \leq t \leq \tau / 2 .
$$

Since $\rho(t, x)$ is increasing in $t$,

$$
v(t, x) \leq A e^{a|x|^{2}}-\epsilon\left(\frac{1}{\tau}\right)^{n / 2} e^{\frac{1}{2 \tau}|x|^{2}} \text { for } 0 \leq t \leq \tau / 2
$$

Hence if we require $\frac{1}{2 \tau}>a$ or $\tau<\frac{1}{2 a}$ it will follows that

$$
\lim _{|x| \rightarrow \infty}\left[\sup _{0 \leq t \leq \tau / 2} v(t, x)\right]=-\infty
$$

Therefore we may choose $M$ sufficiently large so that

$$
v(t, x) \leq K:=\sup _{z} u(0, z) \text { for all }|x| \geq M \text { and } 0 \leq t \leq \tau / 2 .
$$

Since

$$
\left(\partial_{t}-\frac{\triangle}{2}\right) v=\left(\partial_{t}-\frac{\triangle}{2}\right) u \leq 0
$$

we may apply the maximum principle with $\Omega=B(0, M)$ and $T=\tau / 2$ to conclude for $(t, x) \in \Omega_{T}$ that

$$
u(t, x)-\epsilon \rho(t, x)=v(t, x) \leq \sup _{z \in \Omega} v(0, z) \leq K \text { if } 0 \leq t \leq \tau / 2
$$

[^5]$$
\frac{\partial}{\partial t} p_{-t}(x)=-\dot{p}_{t}(x)=-\frac{1}{2} \triangle p_{-t} .
$$

We may now let $\epsilon \downarrow 0$ in this equation to conclude that

$$
\begin{equation*}
u(t, x) \leq K \text { if } 0 \leq t \leq \tau / 2 \tag{45.11}
\end{equation*}
$$

By applying Eq. (45.11) to $u(t+\tau / 2, x)$ we may also conclude

$$
u(t, x) \leq K \text { if } 0 \leq t \leq \tau
$$

Repeating this argument then enables us to show $u(t, x) \leq K$ for all $0 \leq t \leq$ $T$.

## Corollary 45.19. The heat equation

$$
u_{t}-\frac{1}{2} \triangle u=0 \text { on }[0, T] \times \mathbb{R}^{n} \text { with } u(0, \cdot)=f(\cdot) \in C\left(\mathbb{R}^{n}\right)
$$

has at most one solution in the class of functions $u \in C\left([0, T] \times \mathbb{R}^{n}\right) \cap$ $C^{2,1}\left((0, T) \times \mathbb{R}^{n}\right)$ which satisfy

$$
u(t, x) \leq A e^{a|x|^{2}} \text { for }(t, x) \in(0, T) \times \mathbb{R}^{n}
$$

for some constants $A$ and $a$.
Theorem 45.20 (Max Principle a la Hamilton). Suppose $u \in C^{1,2}([0, T] \times \mathbb{R}$

1. $u(t, x) \leq A e^{a|x|^{2}}$ for some $A$, $a($ for all $t \leq T)$
2. $u(0, x) \leq 0$ for all $x$
3. $\frac{\partial u}{\partial t} \leq \Delta u$ i.e. $\left(\partial_{t}-\triangle\right) u \leq 0$.

Then $u(t, x) \leq 0$ for all $(t, x) \in[0, T] \times \mathbb{R}^{d}$.
Proof. Special Case. Assume $\frac{\partial u}{\partial t}<\Delta u$ on $[0, T] \times \mathbb{R}^{d}, u(0, x)<0$ for all $x \in \mathbb{R}^{d}$ and there exists $M>0$ such that $u(t, x)<0$ if $|x| \geq M$ and $t \in[0, T]$.For the sake of contradiction suppose there is some point $(t, x) \in$ $[0, T] \times \mathbb{R}^{d}$ such that $u(t, x)>0$.

By the intermediate value theorem there exists $\tau \in[0, t]$ such that $u(\tau, x)=$ 0 . In particular the set $\{u=0\}$ is a non-empty closed compact subset of $(0, T] \times B(0, M)$. Let

$$
\pi:(0, T] \times B(0, M) \rightarrow(0, T]
$$

be projection onto the first factor, since $\{u \neq 0\}$ is a compact subset of $(0, T] \times$ $B(0, M)$ if follows that

$$
t_{0}:=\min \{t \in \pi(\{u=0\})\}>0 .
$$

Choose a point $x_{0} \in B(0, M)$ such that $\left(t_{0}, x_{0}\right) \in\{u=0\}$, i.e. $u\left(t_{0}, x_{0}\right)=0$, see Figure 45.2 below. Since $u(t, x)<0$ for all $0 \leq t<t_{0}$ and $x \in \mathbb{R}^{d}$,


Fig. 45.2. Finding a point $\left(t_{0}, x_{0}\right)$ such that $t_{0}$ is as small as possible and $u\left(t_{0}, x_{0}\right)=$ 0.
$u\left(t_{0}, x\right) \leq 0$ for all $x \in \mathbb{R}^{d}$ with $u\left(t_{0}, x_{0}\right)=0$. This information along with the first and second derivative tests allows us to conclude

$$
\nabla u\left(t_{0}, x_{0}\right)=0, \Delta u\left(t_{0}, x_{0}\right) \leq 0 \text { and } \frac{\partial u}{\partial t}\left(t_{0}, x_{0}\right) \geq 0
$$

This then implies that

$$
0 \leq \frac{\partial u}{\partial t}\left(t_{0}, x_{0}\right)<\triangle u\left(t_{0}, x_{0}\right) \leq 0
$$

which is absurd. Hence we conclude that $u \leq 0$ on $[0, T] \times \mathbb{R}^{d}$.
General Case: Let $p_{t}(x)=\frac{1}{t^{d / 2}} e^{-\frac{1}{4 t}|x|^{2}}$ be the fundamental solution to the heat equation

$$
\partial_{t} p_{t}=\triangle p_{t}
$$

Let $\tau>0$ to be determined later. As in the proof of Theorem 45.18, the function

$$
\rho(t, x):=p_{\tau-t}(i x)=\left(\frac{1}{\tau-t}\right)^{d / 2} e^{\frac{1}{4(\tau-t)}|x|^{2}} \text { for } 0 \leq t<\tau
$$

is still a solution to the heat equation. Given $\epsilon>0$, define, for $t \leq \tau / 2$,

$$
u_{\epsilon}(t, x)=u(t, x)-\epsilon-\epsilon t-\epsilon \rho(t, x)
$$

Then

$$
\begin{aligned}
\left(\partial_{t}-\triangle\right) u_{\epsilon} & =\left(\partial_{t}-\triangle\right) u-\epsilon \leq-\epsilon<0 \\
u_{\epsilon}(0, x) & =u(0, x)-\epsilon \leq 0-\epsilon \leq-\epsilon<0
\end{aligned}
$$

$$
u_{\epsilon}(t, x) \leq A e^{a|x|^{2}}-\epsilon-\epsilon \frac{1}{\tau^{d / 2}}{ }^{\frac{1}{4 \tau}|x|^{2}} .
$$

Hence if we choose $\tau$ such that $\frac{1}{4 \tau}>a$, we will have $u_{\epsilon}(t, x)<0$ for $|x|$ sufficiently large. Hence by the special case already proved, $u_{\epsilon}(t, x) \leq 0$ for all $0 \leq t \leq \frac{\tau}{2}$ and $\epsilon>0$. Letting $\epsilon \downarrow 0$ implies that $u(t, x) \leq 0$ for all $0 \leq t \leq \tau / 2$. As in the proof of Theorem 45.18 we may step our way up by applying the previous argument to $u(t+\tau / 2, x)$ and then to $u(t+\tau, x)$, etc. to learn $u(t, x) \leq 0$ for all $0 \leq t \leq T$.

### 45.4 Non-Uniqueness of solutions to the Heat Equation

Theorem 45.21 (See Fritz John §7). For any $\alpha>1$, let

$$
g(t):= \begin{cases}e^{-t^{-\alpha}} & t>0  \tag{45.12}\\ 0 & t \leq 0\end{cases}
$$

and define

$$
u(t, x)=\sum_{k=0}^{\infty} \frac{g^{(k)}(t) x^{2 k}}{(2 k)!}
$$

Then $u \in C^{\infty}\left(\mathbb{R}^{2}\right)$ and

$$
\begin{equation*}
u_{t}=u_{x x} \text { and } u(0, x):=0 . \tag{45.13}
\end{equation*}
$$

In particular, the heat equation does not have unique solutions.
Proof. We are going to look for a solution to Eq. (45.13) of the form

$$
u(t, x)=\sum_{n=0}^{\infty} g_{n}(t) x^{n}
$$

in which case we have (formally) that

$$
\begin{aligned}
0 & =u_{t}-u_{x x}=\sum_{n=0}^{\infty}\left(\dot{g}_{n}(t) x^{n}-g_{n}(t) n(n-1) x^{n-2}\right) \\
& =\sum_{n=0}^{\infty}\left[\dot{g}_{n}(t)-(n+2)(n+1) g_{n+2}(t)\right] x^{n}
\end{aligned}
$$

This implies

$$
\begin{equation*}
g_{n+2}=\frac{\dot{g}_{n}}{(n+2)(n+1)} . \tag{45.14}
\end{equation*}
$$

To simplify the final answer, we will now assume $u_{x}(0, x)=0$, i.e. $g_{1} \equiv 0$ in which case Eq. (45.14) implies $g_{n} \equiv 0$ for all $n$ odd. We also have with $g:=g_{0}$,

$$
g_{2}=\frac{\dot{g}_{0}}{2 \cdot 1}=\frac{\dot{g}}{2!}, g_{4}=\frac{\dot{g}_{2} 0}{4 \cdot 3}=\frac{\ddot{g}}{4!}, g_{6}=\frac{g^{(3)}}{6!} \ldots g_{2 k}=\frac{g^{(k)}}{(2 k)!}
$$

and hence

$$
\begin{equation*}
u(t, x)=\sum_{k=0}^{\infty} \frac{g^{(k)}(t) x^{2 k}}{(2 k)!} \tag{45.15}
\end{equation*}
$$

The function $u(t, x)$ will solve $u_{t}=u_{x x}$ for $(t, x) \in \mathbb{R}^{2}$ with $u(0, x)=0$ provided the convergence in the sum is adequate to justify the above computations.

Now let $g(t)$ be given by Eq. (45.12) and extend $g$ to $\mathbb{C} \backslash(-\infty, 0]$ via $g(z)=$ $e^{-z^{-\alpha}}$ where

$$
z^{-\alpha}=e^{-\alpha \log (z)}=e^{-\alpha(\ln r+i \theta)} \text { for } z=r e^{i \theta} \text { with }-\pi<\theta<\pi
$$

In order to estimate $g^{(k)}(t)$ we will use of the Cauchy estimates on the contour $|z-t|=\gamma t$ where $\gamma$ is going to be chosen sufficiently close to 0 . Now

$$
\operatorname{Re}\left(z^{-\alpha}\right)=e^{-\alpha \ln r} \cos (\alpha \theta)=|z|^{-\alpha} \cos (\alpha \theta)
$$

and hence

$$
|g(z)|=e^{-\operatorname{Re}\left(z^{-\alpha}\right)}=e^{-|z|^{-\alpha} \cos (\alpha \theta)}
$$

From Figure 45.3, we see


Fig. 45.3. Here is a picture of the maximum argument $\theta_{m}$ that a point $z$ on $\partial B(t, \gamma t)$ may attain. Notice that $\sin \theta_{m}=\gamma t / t=\gamma$ is independent of $t$ and $\theta_{m} \rightarrow 0$ as $\gamma \rightarrow 0$.

$$
\beta(\gamma):=\min \left\{\cos (\alpha \theta):-\pi<\theta<\pi \text { and }\left|r e^{i \theta}-t\right|=\gamma t\right\}
$$

is independent of $t$ and $\beta(\gamma) \rightarrow 1$ as $\gamma \rightarrow 0$. Therefore for $|z-t|=\gamma t$ we have

$$
|g(z)| \leq e^{-|z|^{-\alpha} \beta(\gamma)} \leq e^{-([\gamma+1] t)^{-\alpha} \beta(\gamma)}=e^{-\frac{\beta(\gamma)}{1+\gamma} t^{-\alpha}} \leq e^{-\frac{1}{2} t^{-\alpha}}
$$

provided $\gamma$ is chosen so small that $\frac{\beta(\gamma)}{1+\gamma} \geq \frac{1}{2}$.

$$
\frac{d^{2}}{d x^{2}} \chi_{k}^{L}=-\left(\frac{k}{L}\right)^{2} \chi_{k}^{L}
$$

Therefore the solution to the heat equation on $S_{L}$,

$$
u_{t}=\frac{1}{2} u_{x x} \text { with } u(0, \cdot)=f \in H_{L}
$$

is given by

$$
\begin{aligned}
u(t, x) & =\sum_{k \in \mathbb{Z}}\left(f, \chi_{k}^{L}\right) e^{-\frac{1}{2}\left(\frac{k}{L}\right)^{2} t} e^{i k x / L} \\
& =\sum_{k \in \mathbb{Z}}\left(\frac{1}{2 \pi L} \int_{-\pi L}^{\pi L} f(y) e^{-i k y / L} d y\right) e^{-\frac{1}{2}\left(\frac{k}{L}\right)^{2} t} e^{i k x / L} \\
& =\int_{-\pi L}^{\pi L} p_{t}^{L}(x-y) f(y) d y
\end{aligned}
$$

where

$$
p_{t}^{L}(x)=\frac{1}{2 \pi L} \sum_{k \in \mathbb{Z}} e^{-\frac{1}{2}\left(\frac{k}{L}\right)^{2} t} e^{i k x / L}
$$

If $f$ is $L$ periodic then it is $n L$ - periodic for all $n \in \mathbb{N}$, so we also would learn

$$
u(t, x)=\int_{-\pi n L}^{\pi n L} p_{t}^{n L}(x-y) f(y) d y \text { for all } n \in \mathbb{N}
$$

this suggest that we might pass to the limit as $n \rightarrow \infty$ in this equation to learn

$$
u(t, x)=\int_{\mathbb{R}} p_{t}(x-y) f(y) d y
$$

where

$$
\begin{aligned}
p_{t}(x) & :=\lim _{n \rightarrow \infty} p_{t}^{n L}(x)=\lim _{L \rightarrow \infty} \frac{1}{2 \pi L} \sum_{k \in \mathbb{Z}} e^{-\frac{1}{2}\left(\frac{k}{L}\right)^{2} t} e^{i\left(\frac{k}{L}\right) x} \\
& =\frac{1}{2 \pi} \int_{\mathbb{R}} e^{-\frac{1}{2} \xi^{2} t} e^{i \xi x} d \xi=\frac{1}{\sqrt{2 \pi t}} e^{-\frac{x^{2}}{2 t}}
\end{aligned}
$$

From this we conclude

$$
u(t, x)=\int_{\mathbb{R}} p_{t}(x-y) f(y) d y=\int_{-\pi L}^{\pi L} \sum_{n \in \mathbb{Z}} p_{t}(x-y+2 \pi n L) f(y) d y
$$

and we arrive at the identity

$$
\sum_{n \in \mathbb{Z}} \frac{1}{\sqrt{2 \pi t}} e^{-\frac{(x+2 \pi n L)^{2}}{2 t}}=\sum_{n \in \mathbb{Z}} p_{t}(x+2 \pi n L)=\frac{1}{2 \pi L} \sum_{k \in \mathbb{Z}} e^{-\frac{1}{2}\left(\frac{k}{L}\right)^{2} t} e^{i k x / L}
$$

which is a special case of Poisson's summation formula.

## Abstract Wave Equation

In the next section we consider

$$
\begin{equation*}
u_{t t}-\triangle u=0 \text { with } u(x, 0)=f(x) \text { and } u_{t}(x, 0)=g(x) \text { for } x \in \mathbb{R}^{n} . \tag{46.1}
\end{equation*}
$$

Before working with this explicit equation we will work out an abstract Hilbert space theory first.

Theorem 46.1 (Existence). Suppose $A: H \rightarrow H$ is a self-adjoint nonpositive operator, i.e. $A^{*}=A$ and $A \leq 0$ and $f \in D(A)$ and $g \in g \in$ $D(\sqrt{-A})$ are given. Then

$$
\begin{equation*}
u(t)=\cos (t \sqrt{-A}) f+\frac{\sin (t \sqrt{-A})}{\sqrt{-A}} g \tag{46.2}
\end{equation*}
$$

satisfies:

1. $\dot{u}(t)=\cos (t \sqrt{-A}) \sqrt{-A} f+\sin (t \sqrt{-A}) g$ exists and is continuous.
2. $\ddot{u}(t)$ exists and is continuous

$$
\begin{equation*}
\ddot{u}(t)=A u \text { with } u(0)=f \text { and } \dot{u}(0)=g \text {. } \tag{46.3}
\end{equation*}
$$

3. $\frac{d}{d t} \sqrt{-A} u(t)=-\cos (t \sqrt{-A}) A f+\sin (t \sqrt{-A}) \sqrt{-A} g$ exists and is continuous.

Eq. (46.3) is Newton's equation of motion for an infinite dimensional harmonic oscillation. Given any solution $u$ to Eq. (46.3) it is natural to define its energy by

$$
E(t, u):=\frac{1}{2}\left[\|\dot{u}(t)\|^{2}+\|\omega u(t)\|^{2}\right]=\text { K.E. }+ \text { P.E. }
$$

where $\omega:=\sqrt{-A}$. Notice that Eq. (46.3) becomes $\ddot{u}+\omega^{2} u=0$ with this definition of $\omega$.

Lemma 46.2 (Conservation of Energy). Suppose $u$ is a solution to Eq. (46.3) such that $\frac{d}{d t} \sqrt{-A} u(t)$ exists and is continuous. Then $\dot{E}(t)=0$.

## Proof.

$$
\dot{E}(t)=\operatorname{Re}(\dot{u}, \ddot{u})+\operatorname{Re}(\omega u, \omega \dot{u})=\operatorname{Re}\left(\dot{u},-\omega^{2} u\right)-\operatorname{Re}\left(\omega^{2} u, \dot{u}\right)=0 .
$$

Theorem 46.3 (Uniqueness of Solutions). The only function $u \in$ $C^{2}(\mathbb{R}, H)$ satisfying 1) $u(t) \in D(A)$ for all $t$ and 2$)$

$$
\ddot{u}=A u \text { with } u(0)=0=\dot{u}(0)
$$

is the $u(t) \equiv 0$ for all $t$.
Proof. Let $\chi_{M}(x)=1_{|x| \leq M}$ and define $P_{M}=\chi_{M}(A)$ so that $P_{M}$ is orthogonal projection onto the spectral subspace of $H$ where $-M \leq A \leq 0$. Then for all $f \in D(A)$ we have $P_{M} A f=A P_{M} f$ and for all $f \in H$ we have $P_{M} f \in D\left((-A)^{\alpha}\right)$ for any $\alpha \geq 0$. Let $u_{M}(t):=P_{M} u(t)$, then $u_{M} \in C^{2}(\mathbb{R}, H)$, $u_{M}(t) \in D\left((-A)^{\alpha}\right)$ for all $t$ and $\alpha, t \rightarrow \sqrt{-A} u_{M}(t)$ is continuous and

$$
\ddot{u}_{M}=\frac{d^{2}}{d t^{2}}\left(P_{M} u\right)=P_{M} \ddot{u}=P_{M} A u=A P_{M} u=A u_{M}
$$

with $u_{M}(0)=0=\dot{u}_{M}(0)$. By Lemma 46.2,

$$
\frac{1}{2}\left[\left\|\dot{u}_{M}(t)\right\|^{2}+\left\|\omega u_{M}(t)\right\|^{2}\right]=\frac{1}{2}\left[\left\|\dot{u}_{M}(0)\right\|^{2}+\left\|\omega u_{M}(0)\right\|^{2}\right]=0
$$

for all $t$. In particular this implies $\dot{u}_{M}(t)=0$ and hence $P_{M} u(t)=u_{M}(t) \equiv 0$ Letting $M \rightarrow \infty$ then shows $u(t) \equiv 0$.

Corollary 46.4. Any solution to $\ddot{u}=A u$ with $u(0) \in D(A)$ and $\dot{u}(0) \in$ $D(\sqrt{-A})$ must satisfy $t \rightarrow \sqrt{-A} u(t)$ is $C^{1}$.

### 46.1 Corresponding first order O.D.E.

Let $v(t)=\dot{u}(t)$, and

$$
x(t)=\binom{u(t)}{v(t)}=\binom{u}{\dot{u}}
$$

then

$$
\begin{aligned}
\dot{x} & =\binom{\dot{u}}{\ddot{u}}=\binom{v}{A u}=\left(\begin{array}{cc}
0 & I \\
A & 0
\end{array}\right) x=B x \text { with } \\
x(0) & =\binom{f}{g},
\end{aligned}
$$

where

$$
B:=\left(\begin{array}{ll}
0 & I \\
A & 0
\end{array}\right)=\left(\begin{array}{rr}
0 & I \\
-\omega^{2} & 0
\end{array}\right)
$$

Note formally that

$$
\begin{align*}
e^{t B}\binom{f}{g} & =\binom{u(t)}{\dot{u}(t)}=\binom{\cos \omega t f+\frac{\sin t \omega}{\omega} g}{-\omega \sin \omega t f+\cos \omega t g} \\
& =\binom{\cos \omega t}{-\omega \sin \omega t \cos \omega t}\binom{f}{g} \tag{46.4}
\end{align*}
$$

and this suggests that

$$
e^{t B}=\left(\begin{array}{ll}
\cos \omega t & \frac{\sin t \omega}{\omega} \\
-\omega \sin \omega t & \cos \omega t
\end{array}\right)
$$

which is formally correct since

$$
\begin{aligned}
\frac{d}{d t} e^{t B} & =\left(\begin{array}{ll}
-\omega \sin \omega t & \cos \omega t \\
-\omega^{2} \cos \omega t & -\omega \sin \omega t
\end{array}\right) \\
& =\left(\begin{array}{rr}
0 & I \\
-\omega^{2} & 0
\end{array}\right)\left(\begin{array}{ll}
\cos \omega t & \frac{\sin t \omega}{\omega} \\
-\omega \sin \omega t & \cos \omega t
\end{array}\right)=B e^{t B} .
\end{aligned}
$$

Since the energy form $E(t)=\|\dot{u}\|^{2}+\|\omega u\|^{2}$ is conserved, it is reasonable to let

$$
K=D(\sqrt{-A}) \oplus H:=\binom{D(\sqrt{-A})}{H}
$$

with inner product

$$
\left\langle\begin{array}{l|l}
f & \tilde{f} \\
g & \tilde{g}
\end{array}\right\rangle=(g, \tilde{g})+(\omega f, \omega \tilde{f})
$$

For simplicity we assume $\operatorname{Nul}(\sqrt{-A})=\operatorname{Nul}(\omega)=\{0\}$ in which case $K$ becomes a Hilbert space and $e^{t B}$ is a unitary evolution on $K$. Indeed,

$$
\begin{aligned}
\left\|e^{t B}\binom{f}{g}\right\|_{K}^{2} & =\|\cos \omega t g-\omega \sin \omega t f\|^{2}+\|\omega(\cos +\omega f)+\sin \omega t g\|^{2} \\
& =\|\cos \omega t g\|^{2}+\|\omega \sin \omega t f\|^{2}+\|\omega \cos \omega t f\|^{2}+\|\sin \omega t g\|^{2} \\
& =\|\omega f\|^{2}+\|g\|^{2}
\end{aligned}
$$

From Eq. (46.4), it easily follows that $\left.\frac{d}{d t}\right|_{0} e^{t B}\binom{f}{g}$ exists iff $g \in D(\omega)$ and $f \in D\left(-\omega^{2}\right)=D(A)$. Therefore we define $D(B):=D(A) \oplus D(\omega)=$ $D(A) \oplus D(\sqrt{-A})$ and

$$
B=\left(\begin{array}{ll}
0 & I \\
A & 0
\end{array}\right): D(B) \rightarrow \begin{gathered}
D(\omega) \\
\oplus \\
H
\end{gathered}=K
$$

Since $B$ is the infinitesimal generator of a unitary semigroup, it follows that $B^{* K}=-B$, i.e. $B$ is skew adjoint. This may be checked directly as well as follows.

Alternate Proof that $B^{* K}=-B$. For

$$
\begin{aligned}
\binom{u}{v}, & \binom{\tilde{u}}{\tilde{v}} \in D(B)=D(A) \oplus D(\omega), \\
\left\langle B\binom{u}{v},\binom{\tilde{u}}{\tilde{v}}\right\rangle & =\left\langle\binom{ v}{A u},\binom{\tilde{u}}{\tilde{v}}\right\rangle=(A u, \tilde{v})+(\omega v, \omega \tilde{u}) \\
& =(A u, \tilde{v})-(A v, \tilde{u})=(u, A \tilde{v})-(v, A \tilde{u})
\end{aligned}
$$

and similarly

$$
\begin{aligned}
\left\langle\binom{ u}{v}, B\binom{\tilde{u}}{\tilde{v}}\right\rangle & =\left\langle\binom{ u}{v},\binom{\tilde{v}}{A \tilde{u}}\right\rangle=(\omega u, \omega \tilde{v})+(v, A \tilde{u}) \\
& =(-A u, \tilde{v})+(v, A \tilde{u})=-\left\langle B\binom{u}{v},\binom{\tilde{u}}{\tilde{v}}\right\rangle
\end{aligned}
$$

which shows $-B \subset B^{*}$. Conversely if $\binom{\tilde{u}}{\tilde{v}} \in D\left(B^{*}\right)$ and $B^{*}\binom{\tilde{u}}{\tilde{v}}=\binom{f}{g}$, then

$$
\left\langle B\binom{u}{v} \left\lvert\, \begin{array}{l}
\tilde{u}  \tag{46.5}\\
\tilde{v}
\end{array}\right.\right\rangle=\left\langle\left.\begin{array}{l}
u \\
v
\end{array} \right\rvert\, \begin{array}{l}
f \\
g
\end{array}\right\rangle=(v, g)+(\omega u, \omega f)
$$

$(A u, \tilde{v})+(\omega v, \omega \tilde{u})$ for all $u \in D(A), v \in D(\omega)$. Take $u=0$ implies $(\omega v, \omega \tilde{u})=(v, g)$ for all $v \in D(\omega)$ which then implies $\omega \tilde{u} \in D\left(\omega^{*}\right)=D(\omega)$ and hence $-A \tilde{u}=\omega^{2} \tilde{u}=g$. (Note $\tilde{u} \in D(A)$.) Taking $v=0$ in Eq. (46.5) implies $(A u, \tilde{v})=(\omega u, \omega f)=(-A u, f)$. Since

$$
\overline{\operatorname{Ran}(A)}=\operatorname{Nul}(A)^{\perp}=\{0\}^{\perp}=H,
$$

we find that $f=-\tilde{v} \in D(\omega)$ since $f \in D(\omega)$. Therefore $D\left(B^{*}\right) \subset D(B)$ and for $(\tilde{u}, \tilde{v}) \in D\left(B^{*}\right)$ we have

$$
B^{*}\binom{\tilde{u}}{\tilde{v}}=-\binom{-\tilde{v}}{-A \tilde{u}}=-B\binom{\tilde{u}}{\tilde{v}} .
$$

### 46.2 Du Hamel's Principle

Consider

$$
\begin{equation*}
\ddot{u}=A u+f(t) \text { with } u(0)=g \text { and } \dot{u}(0)=h . \tag{46.6}
\end{equation*}
$$

Eq. (46.6) implies, with $v=\dot{u}$, that

$$
\frac{d}{d t}\binom{u}{v}=\binom{v}{\ddot{u}}=\binom{v}{A u+f}=\left(\begin{array}{ll}
0 & I \\
A & 0
\end{array}\right)\binom{u}{v}+\binom{0}{f} .
$$

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$$
\binom{u}{v}(t)=e^{t\left(\begin{array}{ll}
0 & I \\
A & 0
\end{array}\right)}\binom{g}{h}+\int_{0}^{t} e^{(t-\tau)}\left(\begin{array}{cc}
0 & I \\
A & 0
\end{array}\right)\binom{0}{f(\tau)} d \tau
$$

hence

$$
u(t)=\cos (t \sqrt{-A}) g+\frac{\sin (t \sqrt{-A})}{\sqrt{-A}} h+\int_{0}^{t} \frac{\sin ((t-\tau) \sqrt{-A})}{\sqrt{-A}} f(\tau) d \tau
$$

Theorem 46.5. Suppose $f(t) \in D(\sqrt{-A})$ for all $t$ and that $f(t)$ is continuous relative to $\|f\|_{A}:=\|f\|+\|\sqrt{-A} f\|$. Then

$$
u(t):=\int_{0}^{t} \frac{\sin ((t-\tau) \sqrt{-A})}{\sqrt{-A}} f(\tau) d \tau
$$

solves $\ddot{u}=A u+f$ with $u(0)=0, \dot{u}(0)=0$.
Proof. $\dot{u}(t)=\int_{0}^{t} \cos ((t-\tau) \sqrt{-A}) f(\tau) d \tau$.

$$
\begin{aligned}
\ddot{u}(t) & =f(t)-\int_{0}^{t} \sin ((t-\tau) \sqrt{-A}) \sqrt{-A} f(\tau) d \tau \\
& =f(t)-A \int_{0}^{t} \frac{\sin ((t-\tau) \sqrt{-A})}{\sqrt{-A}} f(\tau) d \tau .
\end{aligned}
$$

So $\dot{u}=A u+f$. Note $u(0)=0=\dot{u}(0)$.
Alternate. Let $\omega:=\sqrt{-A}$, then

$$
\begin{aligned}
u(t) & =\int_{0}^{t} \frac{\sin ((t-\tau) \omega)}{\omega} f(\tau) d \tau \\
& =\int_{0}^{t} \frac{\sin \omega t \cos \omega \tau-\sin \omega \tau \cos \omega t}{\omega} f(\tau) d \tau
\end{aligned}
$$

and hence

$$
\begin{aligned}
\dot{u}(t) & =\frac{\sin \omega t \cos \omega t-\sin \omega t \cos \omega t}{\omega} f(t) \\
& +\int_{0}^{t}(\cos \omega t \cos \omega \tau+\sin \omega \tau \sin \omega t) f(\tau) d \tau \\
& =\int_{0}^{t}(\cos \omega t \cos \omega \tau+\sin \omega \tau \sin \omega t) f(\tau) d \tau .
\end{aligned}
$$

Similarly,

## $47.1 n=1$ Case

## Wave Equation on $\mathbb{R}^{n}$

(Ref Courant \& Hilbert Vol II, Chap VI §12.)
We now consider the wave equation

$$
\begin{equation*}
u_{t t}-\triangle u=0 \text { with } u(0, x)=f(x) \text { and } u_{t}(0, x)=g(x) \text { for } x \in \mathbb{R}^{n} \tag{47.1}
\end{equation*}
$$

According to Section 46, the solution (in the $L^{2}-$ sense) is given by

$$
\begin{equation*}
u(t, \cdot)=\left(\cos (t \sqrt{-\triangle}) f+\frac{\sin (t \sqrt{-\triangle})}{\sqrt{-\triangle}} g\right. \tag{47.2}
\end{equation*}
$$

To work out the results in Eq. (47.2) we must diagonalize $\Delta$. This is of course done using the Fourier transform. Let $\mathcal{F}$ denote the Fourier transform in the $x$ - variables only. Then

$$
\begin{aligned}
& \ddot{\hat{u}}(t, k)+|k|^{2} \hat{u}(t, k)=0 \text { with } \\
& \hat{u}(0, k)=\hat{f}(k) \text { and } \dot{\hat{u}}(t, k)=\hat{g}(k) .
\end{aligned}
$$

Therefore

$$
\hat{u}(t, k)=\cos (t|k|) \hat{f}(k)+\frac{\sin (t|k|)}{|k|} \hat{g}(k)
$$

and so

$$
u(t, x)=\mathcal{F}^{-1}\left[\cos (t|k|) \hat{f}(k)+\frac{\sin (t|k|)}{|k|} \hat{g}(k)\right](x)
$$

i.e.

$$
\begin{align*}
& \frac{\sin (t \sqrt{-\triangle})}{\sqrt{-\triangle}} g=\mathcal{F}^{-1}\left[\frac{\sin (t|k|)}{|k|} \hat{g}(k)\right] \text { and }  \tag{47.3}\\
& \cos (t \sqrt{-\triangle}) f=\mathcal{F}^{-1}[\cos (t|k|) \hat{f}(k)]=\frac{d}{d t} \mathcal{F}^{-1}\left[\frac{\sin (t|k|)}{|k|} \hat{g}(k)\right] \tag{47.4}
\end{align*}
$$

As we see from Eq. (47.4) it suffices to compute:

$$
\begin{align*}
\frac{\sin (t \sqrt{-\triangle})}{\sqrt{-\triangle}} g & =\mathcal{F}^{-1}\left(\frac{\sin (t|\xi|)}{|\xi|} \hat{g}(\xi)\right)=\lim _{M \rightarrow \infty} \mathcal{F}^{-1}\left(1_{|\xi| \leq M} \frac{\sin (t|\xi|)}{|\xi|} \hat{g}(\xi)\right) \\
& =\lim _{M \rightarrow \infty} \mathcal{F}^{-1}\left(1_{|\xi| \leq M} \frac{\sin (t|\xi|)}{|\xi|}\right) \star g . \tag{47.5}
\end{align*}
$$

This inverse Fourier transform will be computed in Proposition 47.2 below using the following lemma.

Lemma 47.1. Let $C_{M}$ denote the contour shown in Figure 47.1, then for $\lambda \neq 0$ we have

$$
\lim _{M \rightarrow \infty} \int_{C_{M}} \frac{e^{i \lambda \xi}}{\xi} d \xi=2 \pi i 1_{\lambda>0}
$$

Proof. First assume that $\lambda>0$ and let $\Gamma_{M}$ denote the contour shown in Figure 47.1. Then

$$
\left|\int_{\Gamma_{M}} \frac{e^{i \lambda \xi}}{\xi} d \xi\right| \leq \int_{0}^{\pi}\left|e^{i \lambda M e^{i \theta}}\right| d \theta=2 \pi \int_{0}^{\pi} d \theta e^{-\lambda M \sin \theta} \rightarrow 0 \text { as } M \rightarrow \infty .
$$

Therefore

$$
\lim _{M \rightarrow \infty} \int_{C_{M}} \frac{e^{i \lambda \xi}}{\xi} d \xi=\lim _{M \rightarrow \infty} \int_{C_{M}+\Gamma_{M}} \frac{e^{i \lambda \xi}}{\xi} d \xi=2 \pi i \operatorname{res}_{\xi=0}\left(\frac{e^{i \lambda \xi}}{\xi}\right)=2 \pi i .
$$

Fig. 47.1. A couple of contours in $\mathbb{C}$.

$$
\lim _{M \rightarrow \infty} \int_{C_{M}} \frac{e^{i \lambda \xi}}{\xi} d \xi=\lim _{M \rightarrow \infty} \int_{C_{M}+\tilde{\Gamma}_{M}} \frac{e^{i \lambda \xi}}{\xi} d \xi
$$

and the later integral is 0 since the integrand is holomorphic inside the contour $C_{M}+\tilde{\Gamma}_{M}$.

Proposition 47.2. $\lim _{M \rightarrow \infty} \mathcal{F}^{-1}\left(1_{|\xi| \leq M} \frac{\sin (t|\xi|)}{|\xi|}\right)(x)=\operatorname{sgn}(t) \frac{\sqrt{\pi}}{\sqrt{2}} 1_{|x|<|t|}$.

## Proof. Let

$$
I_{M}=\sqrt{2 \pi} \mathcal{F}^{-1}\left(1_{|\xi| \leq M} \frac{\sin (t|\xi|)}{|\xi|}\right)(x)=\int_{|\xi| \leq M} \frac{\sin (t \xi)}{\xi} e^{i \xi \cdot x} d \xi
$$

Then by deforming the contour we may write

$$
\begin{aligned}
I_{M} & =\int_{C_{M}} \frac{\sin t \xi}{\xi} e^{i \xi \cdot x} d \xi=\frac{1}{2 i} \int_{C_{M}} \frac{e^{i t \xi}-e^{-i t \xi}}{\xi} e^{i \xi \cdot x} d \xi \\
& =\frac{1}{2 i} \int_{C_{M}} \frac{e^{i(x+t) \xi}-e^{i(x-t) \xi}}{\xi} d \xi
\end{aligned}
$$

By Lemma 47.1 we conclude that

$$
\lim _{M \rightarrow \infty} I_{M}=\frac{1}{2 i} 2 \pi i\left(1_{(x+t)>0}-1_{(x-t)>0}\right)=\pi \operatorname{sgn}(t) 1_{|x|<|t|}
$$

(For the last equality, suppose $t>0$. Then $x-t>0$ implies $x+t>0$ so we get 0 and if $x<-t$, i.e. $x+t<0$ then $x-t<0$ and we get 0 again. If $|x|<t$ the first term is 1 while the second is zero. Similar arguments work when $t<0$ as well.)

Theorem 47.3. For $n=1$,

$$
\begin{align*}
\frac{\sin (t \sqrt{-\triangle})}{\sqrt{-\triangle}} g(x) & =\frac{1}{2} \int_{x-t}^{x+t} g(y) d \lambda(y) \text { and }  \tag{47.6}\\
\cos (t \sqrt{-\triangle}) g(x) & =\frac{1}{2}[g(x+t)+g(x-t)] \tag{47.7}
\end{align*}
$$

In particular

$$
\begin{equation*}
u(t, x)=\frac{1}{2}(f(x+t)+f(x-t))+\frac{1}{2} \int_{x-t}^{x+t} g(y) d y \tag{47.8}
\end{equation*}
$$

is the solution to the wave equation (47.2).

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which is equivalent to Eq. (47.8)

### 47.2 Solution for $n=3$

Given a function $f: \mathbb{R}^{n} \rightarrow \mathbb{R}$ and $t \in \mathbb{R}$ let

$$
\bar{f}(x ; t):=\int_{S^{2}} f(x+t \omega) d \sigma(\omega)=f \int_{|y|=|t|} f(x+y) d \sigma(y) .
$$

Theorem 47.4. For $f \in L^{2}\left(\mathbb{R}^{3}\right)$,

$$
\frac{\sin (\sqrt{-\Delta} t)}{\sqrt{-\Delta}} f=\mathcal{F}^{-1}\left[\frac{\sin |\xi| t}{|\xi|} \hat{f}(\xi)\right](x)=t \bar{f}(x ; t)
$$

and

$$
\cos (\sqrt{-\Delta} t) g=\frac{d}{d t}[t \bar{f}(x ; t)]
$$

In particular the solution to the wave equation (47.1) for $n=3$ is given by

$$
\begin{aligned}
u(t, x) & =\frac{\partial}{\partial t}(t \bar{f}(x ; t))+t \bar{g}(x ; t) \\
& =\frac{1}{4 \pi} \int_{|\omega|=1}(t g(x+t \omega)+f(x+t \omega)+t \nabla f(x+t \omega) \cdot \omega) d \sigma(\omega) .
\end{aligned}
$$

Proof. Let $g_{M}:=\mathcal{F}^{-1}\left[\frac{\sin |\xi| t}{|\xi|} 1_{|\xi| \leq M}\right]$, then by symmetry and passing to spherical coordinates,

$$
\begin{aligned}
(2 \pi)^{3 / 2} g_{M}(x) & =\int_{|\xi| \leq M} \frac{\sin |\xi| t}{|\xi|} e^{i \xi \cdot x} d \xi=\int_{|\xi| \leq M} \frac{\sin |\xi| t}{|\xi|} e^{i|x| \xi_{3}} d \xi \\
& =\int_{0}^{M} d \rho \rho^{2} \int_{0}^{2 \pi} d \theta \int_{0}^{\pi} d \phi \frac{\sin \rho t}{\rho} e^{i \rho|x| \cos \phi} \sin \phi \\
& =\left.2 \pi \int_{0}^{M} d \rho \sin \rho t \frac{e^{i \rho|x| \cos \phi}}{-i|x|}\right|_{0} ^{\pi} \\
& =2 \pi \int_{0}^{M} d \rho \sin \rho t \frac{e^{i \rho|x|}-e^{-i \rho|x|}}{i|x|}=\frac{4 \pi}{|x|} \int_{0}^{M} \sin \rho t \sin \rho|x| d \rho .
\end{aligned}
$$

Using

$$
\sin A \sin B=\frac{1}{2}[\cos (A-B)-\cos (A+B)]
$$

in this last equality, shows

$$
\begin{aligned}
g_{M}(x) & =(2 \pi)^{-3 / 2} \frac{2 \pi}{|x|} \int_{0}^{M}[\cos ((t-|x|) \rho)-\cos ((t+|x|) \rho)] d \rho \\
& =(2 \pi)^{-3 / 2} \frac{\pi}{|x|} h_{M}(|x|)
\end{aligned}
$$

where

$$
h_{M}(r):=\int_{-M}^{M}[\cos ((t-r) \alpha)-\cos ((t+r) \alpha)] d \alpha,
$$

an odd function in $r$. Since

$$
\mathcal{F}^{-1}\left[\frac{\sin |\xi| t}{|\xi|} \hat{f}(\xi)\right]=\lim _{M \rightarrow \infty} \mathcal{F}^{-1}\left(\hat{g}_{M}(\xi) \hat{f}(\xi)\right)=\lim _{M \rightarrow \infty}\left(g_{M} \star f\right)(x)
$$

we need to compute $g_{M} \star f$. To this end

$$
\begin{aligned}
& g_{M} \star f(x)=\left(\frac{1}{2 \pi}\right)^{3} \pi \int_{\mathbb{R}^{3}} \frac{1}{|y|} h_{M}(|y|) f(x-y) d y \\
& =\left(\frac{1}{2 \pi}\right)^{3} \pi \int_{0}^{\infty} d \rho \frac{h_{M}(\rho)}{\rho} \int_{|y|=\rho} f(x-y) d \sigma(y) \\
& =\left(\frac{1}{2 \pi}\right)^{3} \pi \int_{0}^{\infty} d \rho \frac{h_{M}(\rho)}{\rho} 4 \pi \rho^{2} f_{|y|=\rho} f(x-y) d \sigma(y) \\
& =\frac{1}{2 \pi} \int_{0}^{\infty} d \rho h_{M}(\rho) \rho \bar{f}(x ; \rho)=\frac{1}{4 \pi} \int_{-\infty}^{\infty} d \rho h_{M}(\rho) \rho \bar{f}(x ; \rho)
\end{aligned}
$$

where the last equality is a consequence of the fact that $h_{M}(\rho) \rho \bar{f}(x ; \rho)$ is an even function of $\rho$. Continuing to work on this expression suing $\rho \rightarrow \rho f(x ; \rho)$ is odd implies

$$
\begin{aligned}
g_{M} \star f(x) & =\frac{1}{4 \pi} \int_{-\infty}^{\infty} d \rho \int_{-M}^{M}[\cos ((t-\rho) \alpha)-\cos ((t+\rho) \alpha)] d \alpha \rho \bar{f}(x ; \rho) \\
& =\frac{1}{2 \pi} \int_{-\infty}^{\infty} d \rho \int_{-M}^{M} \cos ((t-\rho) \alpha) \rho \bar{f}(x ; \rho) d \alpha \\
& =\frac{1}{2 \pi} \operatorname{Re} \int_{-M}^{M} d \rho \int_{-\infty}^{\infty} d \alpha e^{i(t-\rho) \alpha} \rho \bar{f}(x ; \rho) d \alpha \rightarrow t \bar{f}(x ; t) \text { as } M \rightarrow \infty
\end{aligned}
$$

using the 1 - dimensional Fourier inversion formula.

### 47.2.1 Alternate Proof of Theorem 47.4

Lemma 47.5. $\lim _{M \rightarrow \infty} \int_{-M}^{M} \cos (\rho \lambda) d \rho=2 \pi \delta(\lambda)$.
Proof.

$$
\int_{-M}^{M} \cos (\rho \lambda) d \rho=\int_{-M}^{M} e^{i \rho \lambda} d \rho
$$

so that

$$
\int_{\mathbb{R}} \phi(\lambda)\left[\int_{-M}^{M} e^{i \rho \lambda} d \rho\right] d \lambda \rightarrow \int_{\mathbb{R}} d \rho \int_{\mathbb{R}} d \lambda \phi(\lambda) e^{i \lambda \rho}=2 \pi \varphi(0)
$$

by the Fourier inversion formula.
Proof. of Theorem 47.4 again.

$$
\begin{aligned}
& \int \frac{\sin t|\xi|}{|\xi|} e^{i \xi \cdot x} d \xi=\int \frac{\sin t \rho}{\rho} e^{i \rho|x| \cos \theta} \sin \theta d \theta d \varphi \rho^{2} d \rho \\
& =\left.2 \pi \int \frac{\sin t \rho}{\rho} \frac{e^{i \rho|x| \lambda}}{i \rho|x|}\right|_{\lambda=-1} ^{1} d \rho \\
& =\frac{4 \pi}{|x|} \int_{0}^{\infty} \sin t \rho \sin \rho|x| d \rho \\
& =\frac{2 \pi}{|x|} \int_{0}^{\infty}[\cos (\rho(t-|x|))-\cos (\rho(t+|x|))] d \rho \\
& =\frac{4 \pi}{|x|} \int_{-\infty}^{\infty}[\cos (\rho(t-|x|))-\cos (\rho(t+|x|))] d \rho \\
& =\frac{8 \pi^{2}}{|x|}(\delta(t-|x|)-\delta(t+|x|))
\end{aligned}
$$

Therefore

$$
\begin{aligned}
& \mathcal{F}^{-1}\left(\frac{\sin t|\xi|}{|\xi|}\right) * g(x) \\
& =\left(\frac{1}{2 \pi}\right)^{3} 2 \pi^{2} \int_{\mathbb{R}^{3}} \frac{(\delta(t-|y|)-\delta(t+|y|))}{|y|} g(x-y) d \lambda(y) \\
& =\frac{1}{4 \pi} \int_{0}^{\infty}(\delta(t-\rho)-\delta(t+\rho)) g(x+\rho \omega) \frac{\rho^{2}}{\rho} d \rho d \sigma(\omega) \\
& =1_{t>0} t \bar{g}(x ; t)-1_{t<0}(-t) \bar{g}(x ;-t) \\
& =t \bar{g}(x ; t)
\end{aligned}
$$

### 47.3 Du Hamel's Principle

The solution to

$$
u_{t t}=\Delta u+f \text { with } u(0, x)=0 \text { and } u_{t}(0, x)=0
$$

$$
\begin{equation*}
u(t, x)=\frac{1}{4 \pi} \int_{B(x, t)} \frac{f(t-|y-x|, y)}{|y-x|} d y=\frac{1}{4 \pi} \int_{|z|<t} \frac{f(t-|z|, x+z)}{|z|} d z \tag{47.9}
\end{equation*}
$$

Indeed, by Du Hamel's principle,

$$
\begin{aligned}
u(t, x) & =\int_{0}^{t} \frac{\sin ((t-\tau) \sqrt{-\triangle)}}{\sqrt{-\triangle}} f(\tau, x) d \tau=\int_{0}^{t} \frac{\sin (\tau \sqrt{-\triangle)}}{\sqrt{-\triangle}} f(t-\tau, x) d \tau \\
& =\int_{0}^{t} \tau \bar{f}(t-\tau, x ; \tau) d \tau=\frac{1}{4 \pi} \int_{0}^{t} d \tau t^{2} \int_{|\omega|=1} \frac{f(t-\tau, x+\tau \omega)}{\tau} d \sigma(\omega) \\
& =\frac{1}{4 \pi} \int_{B(t, x)} \frac{f(t-|y-x|, y)}{|y-x|} d y(\text { let } y=x+z) \\
& =\frac{1}{4 \pi} \int_{|z|<t} \frac{f(t-|z|, x+z)}{|z|} d z
\end{aligned}
$$

Thinking of $u(t, x)$ as pressure (47.9) says that the pressure at $x$ at time $t$ is the "average" of the disturbance at time $t-|y-x|$ at location $y$.

### 47.4 Spherical Means

Let $n \geq 2$ and suppose $u$ solves $u_{t t}=\Delta u$. Since $\Delta$ is invariant under rotations, i.e. for $R \in O(n)$ we have $\triangle(u \circ R)=(\triangle u) \circ R$, it follows that $u \circ R$ is also a solution to the wave equation. Indeed,

$$
(u(t, \cdot) \circ R)_{t t}=u_{t t}(t, \cdot) \circ R=\Delta u(t, \cdot) \circ R=\triangle(u(t, \cdot) \circ R) .
$$

By the linearity of the wave equation, this also implies, with $d R$ denoting normalized Haar measure on $O(n)$, that

$$
U(t,|x|):=\int_{O(n)}(u(t, R x) \circ R) d R
$$

must be a radial solution of the Wave equation. This implies
$U_{t t}=\triangle_{x} U(t,|x|)=\frac{1}{r^{n-1}} \partial_{r}\left(r^{n-1} \partial_{r} U(t, r)\right)_{r=|x|}=\left[\partial_{r}^{2} U(t, r)+\frac{n-1}{r} \partial_{r} U(t, r)\right]$
Now

$$
U(t,|x|)=\int_{0(n)} u(t, R x) d R=\int_{B(0,|x|)} u(t, y) d \sigma(y) .
$$

Using the translation invariance of $\Delta$ the same argument as above gives the following theorem.

Theorem 47.6. Suppose $u_{t t}=\triangle u$ and $x \in \mathbb{R}^{n}$ and let

$$
\begin{aligned}
U(t, r) & :=\bar{u}(t, x ; r):=\int_{\partial B(x, r)} u(t, y) d \sigma(y) \\
& =\int_{\partial B(0,1)} u(t, x+r \omega) d \sigma(\omega) .
\end{aligned}
$$

Then $U$ solves

$$
U_{t t}=\frac{1}{r^{n-1}} \partial_{r}\left(r^{n-1} U_{r}\right)
$$

with

$$
\begin{aligned}
U(0, r) & =f_{\partial B(0,1)} u(0, x+r \omega) d \sigma(\omega)=\bar{f}(x ; r) \\
U_{t}(0, r) & =\bar{g}(x ; r)
\end{aligned}
$$

Proof. This has already been proved, nevertheless, let us give another proof which does not rely on using integration over $O(n)$. To this hence we compute

$$
\begin{aligned}
\partial_{r} U(t, r) & =\partial_{r} \int_{\partial B(0,1)} u(t, x+r \omega) d \sigma(\omega) \\
& =\int_{\partial B(0,1)} \nabla u(t, x+r \omega) \cdot \omega d \sigma(\omega) \\
& =\frac{1}{\sigma\left(S^{n-1}\right) r^{n-1}} \int_{|y|=r} \nabla u(t, x+y) \cdot \hat{y} d \sigma(y) \\
& =\frac{1}{\sigma\left(S^{n-1}\right) r^{n-1}} \int_{|y| \leq r} \Delta u(t, x+y) d y \\
& =\frac{1}{\sigma\left(S^{n-1}\right) r^{n-1}} \int_{0}^{r} d \rho \int_{|y|=\rho} \Delta u(t, x+y) d \sigma(y)
\end{aligned}
$$

so that

$$
\begin{aligned}
\frac{1}{r^{n-1}} \partial_{r}\left(r^{n-1} U_{r}\right) & =\frac{1}{r^{n-1}} \partial_{r}\left[\frac{1}{\sigma\left(S^{n-1}\right)} \int_{0}^{r} d \rho \int_{|y|=\rho} \Delta u(t, x+y) d \sigma(y)\right] \\
& =\frac{1}{\sigma\left(S^{n-1}\right) r^{n-1}} \int_{|y|=r} \Delta u(t, x+y) d \sigma(y) \\
& =\int_{|y|=r} \Delta u(t, x+y) d \sigma(y) \\
& =\int_{|y|=r} u_{t t}(t, x+y) d \sigma(y)=U_{t t}
\end{aligned}
$$

We can now use the above result to solve the wave equation. For simplicity, assume $n=3$ and let $V(t, r)=r \bar{u}(t, x ; r)=r U(t, r)$. Then for $r>0$ we have

$$
\begin{aligned}
V_{r r} & =2 U_{r}+r U_{r r}=r\left(U_{r r}+\frac{2}{r} U_{r}\right) \\
& =r U_{t t}=V_{t t}
\end{aligned}
$$

This is also valid for $r<0$ because $V(t, r)$ is odd in $r$. Indeed for $r<0$, let $v(t, r)=V(t,-r)$, then $V_{r r}(t, r)=V_{r r}(t,-r)=V_{t t}(t,-r)=V_{t t}(t, r)$. By our solution to the one dimensional wave equation we find

$$
V(t, r)=\frac{1}{2}(V(0, t+r)+V(0, r-t))+\frac{1}{2} \int_{r-t}^{r+t} V_{t}(0, y) d y
$$

Now suppose that $u(0, x)=0$ and $u_{t}(0, x)=g(x)$, in which case

$$
V(0, r)=0 \text { and } V_{t}(0, r)=r \bar{g}(x, r)
$$

and the previous equation becomes
Then

$$
V(t, r)=\frac{1}{2} \int_{r-t}^{r+t} y \bar{g}(x, y) d y
$$

and noting that

$$
\left.\frac{\partial}{\partial r}\right|_{0} V(t, r)=\bar{u}(t, x ; 0)=u(t, x)
$$

we learn

$$
u(t, x)=\frac{1}{2}[t \bar{g}(x ; t)-(-t) \bar{g}(x ;-t)]=t \bar{g}(x ; t)
$$

as before.

### 47.5 Energy methods

Theorem 47.7 (Uniqueness on Bounded Domains). Let $\Omega$ be a bounded domain such that $\Omega$ is a submanifold with $C^{2}$ - boundary and consider the boundary value problem

$$
\begin{array}{ll}
u_{t t}-\triangle u=h & \text { on } \Omega_{T} \\
u=f & \text { on }(\partial \Omega \times[0, T]) \cup(\Omega \times\{t=0\}) \\
u_{t}=g & \text { on } \Omega \times\{t=0\}
\end{array}
$$

If $u \in C^{2}\left(\bar{\Omega}_{T}\right)$ then $u$ is unique

Proof. As usual, using the linearity of the equation, it suffices to consider the special case where $f=0, g=0$ and $h=0$ and to show this implies $u \equiv 0$. Let

$$
E_{\Omega}(t)=\frac{1}{2} \int_{\Omega}\left[\dot{u}(t, x)^{2}+|\nabla u(t, x)|^{2}\right] d x .
$$

Clearly by assumption, $E_{\Omega}(0)=0$ while the usual computation shows

$$
\begin{aligned}
\dot{E}_{\Omega}(t) & =(\dot{u}, \ddot{u})_{L^{2}(\Omega)}+(\nabla u(t), \nabla \dot{u}(t))_{L^{2}(\Omega)} \\
& =(\dot{u}, \Delta u)_{L^{2}(\Omega)}+(\nabla u(t), \nabla \dot{u}(t))_{L^{2}(\Omega)} \\
& =-(\nabla \dot{u}(t), \nabla u(t))_{L^{2}(\Omega)}+\int_{\partial \Omega} \dot{u}(t, x) \frac{\partial u(t, x)}{\partial n} d \sigma(x) \\
& +(\nabla u(t), \nabla \dot{u}(t))_{L^{2}(\Omega)} \\
& =0
\end{aligned}
$$

wherein we have used $u(t, x)=0$ implies $\dot{u}(t, x)=0$ for $x \in \partial \Omega$.
From this we conclude that $E_{\Omega}(t)=0$ and therefore $\dot{u}(t, x)=0$ and hence $u \equiv 0$.

The following proposition is expected to hold given the finite speed of propagation we have seen exhibited above for solutions to the wave equation. Proposition 47.8 (Local Energy). Let $x \in \mathbb{R}^{n}, T>0, u_{t t}=\Delta u$ and define

$$
e(t):=E_{B(x, T-t)}(u ; t):=\frac{1}{2} \int_{B(x, T-t)}\left[|\dot{u}(t, y)|^{2}+|\nabla u(t, y)|^{2}\right] d y .
$$

Then $e(t)$ is decreasing for $0 \leq t \leq T$.

## Proof. First recall that

$$
\frac{d}{d r} \int_{B(x, r)} f d x=\frac{d}{d r} \int_{0}^{r} d \rho \int_{|y-x|=\rho} f(y) d \sigma(y)=\int_{\partial B(x, r)} f d \sigma .
$$

Hence

$$
\begin{aligned}
\dot{e}(t) & =\frac{d}{d t} \int_{B(x, R-t)}\left\{|\dot{u}(t, y)|^{2}+|\nabla u(t, y)|^{2}\right\} d y \\
& =-\frac{1}{2} \int_{\partial B(x, R-t)}\left(|\dot{u}|^{2}+|\nabla u|^{2}\right) d \sigma+\int_{B(x, R-t)}[\dot{u} \ddot{u}+\nabla u \cdot \nabla \dot{u}] d m \\
& =-\frac{1}{2} \int_{\partial B(x, R-t)}\left(|\dot{u}|^{2}+|\nabla u|^{2}\right) d \sigma+\int_{B(x, R-t)}[\dot{u} \Delta u+\nabla u \cdot \nabla \dot{u}] d m \\
& =-\frac{1}{2} \int_{\partial B(x, R-t)}\left(|\dot{u}|^{2}+|\nabla u|^{2}\right) d \sigma+2 \int_{\partial B(x, R-t)} \dot{u} \frac{\partial u}{\partial n} d \sigma \\
& =\frac{1}{2} \int_{\partial B(x, R-t)}\left\{2 \dot{u}(\nabla u \cdot n)-\left(|\dot{u}|^{2}+|\nabla u|^{2}\right)\right\} d \sigma \leq 0
\end{aligned}
$$

wherein we have used the elementary estimate,

$$
2(\nabla u \cdot n) \dot{u} \leq 2|\nabla u||\dot{u}| \leq\left(|\dot{u}|^{2}+|\nabla u|^{2}\right) .
$$

Therefore $e(t) \leq e(0)=0$ for all $t$ i.e. $e(t):=0$.
Corollary 47.9 (Uniqueness of Solutions). Suppose that $u$ is a classical solution to the wave equation with $u(0, \cdot)=0=u_{t}(0, \cdot)$. Then $u \equiv 0$.

Proof. Proposition 47.8 shows

$$
\frac{1}{2} \int_{B(x, T-t)}\left[|\dot{u}(t, y)|^{2}+|\nabla u(t, y)|^{2}\right] d y=E_{B(x, T)}(0)=0
$$

for all $0 \leq t<T$ and $x \in \mathbb{R}^{n}$. This then implies that $\dot{u}(t, y)=0$ for all $y \in \mathbb{R}^{n}$ and $0 \leq t \leq T$ and hence $u \equiv 0$.

Remark 47.10. This result also applies to certain class of weak type solutions in $x$ by first convolving $u$ with an approximate (spatial) delta function, say $u_{\epsilon}(t, x)=u(t, \cdot) * \delta_{\epsilon}(x)$. Then $u_{\epsilon}$ satisfies the hypothesis of Corollary 47.9 and hence is 0 . Now let $\epsilon \downarrow 0$ to find $u \equiv 0$.

Remark 47.11. Proposition 47.8 also exhibits the finite speed of propagation of the wave equation.

### 47.6 Wave Equation in Higher Dimensions

### 47.6.1 Solution derived from the heat kernel

Let

$$
p_{t}^{n}(x):=\frac{1}{(2 \pi t)^{n / 2}} e^{-\frac{1}{2 t}|x|^{2}}
$$

and simply write $p_{t}$ for $p_{t}^{1}$. Then

$$
2 \int_{0}^{\infty} \cos \omega t p_{\lambda}(t) d t=\int_{\mathbb{R}} e^{i t \omega} p_{\lambda}(t) d t=\left.e^{-\lambda \partial_{t}^{2} / 2} e^{i t \omega}\right|_{t=0}=e^{-\lambda \omega^{2} / 2}
$$

Taking $\omega=\sqrt{-\Delta}$ and writing $u(t, x):=\cos (\sqrt{-\Delta} t) g(x)$ the previous identity gives

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$$
\cos (\sqrt{-\Delta} t) g(x)=c_{n} \partial_{t}\left(\frac{1}{t} \partial_{t}\right)^{k-1}\left[t^{2 k-1} \bar{g}(x ; t)\right]
$$

$$
=e^{\lambda \Delta / 2} g(x)=\int_{\mathbb{R}^{n}} p_{\lambda}^{n}(y) g(x-y) d y
$$

$$
=\int_{\mathbb{R}^{n}} \frac{1}{(2 \pi \lambda)^{n / 2}} e^{-\frac{1}{2 \lambda}|y|^{2}} g(x-y) d y
$$

$$
=\frac{1}{(2 \pi \lambda)^{n / 2}} \int_{0}^{\infty} d \rho e^{-\frac{1}{2 \lambda} \rho^{2}} \int_{|y|=\rho} g(x-y) d \sigma(y)
$$

$$
=\frac{\sigma\left(S^{n-1}\right)}{(2 \pi \lambda)^{n / 2}} \int_{0}^{\infty} d \rho e^{-\frac{1}{2 \lambda} \rho^{2}} \rho^{n-1} \bar{g}(x ; \rho),
$$

and so

$$
\begin{aligned}
\int_{0}^{\infty} u(t, x) e^{-\frac{1}{2 \lambda} t^{2}} d t & =\sqrt{\frac{\pi \lambda}{2}} \frac{\sigma\left(S^{n-1}\right)}{(2 \pi \lambda)^{n / 2}} \int_{0}^{\infty} d \rho e^{-\frac{1}{2 \lambda} \rho^{2}} \rho^{n-1} \bar{g}(x ; \rho) \\
& =\sqrt{\frac{\pi}{2}} \frac{\sigma\left(S^{n-1}\right)}{(2 \pi)^{n / 2}} \lambda^{-(n-1) / 2} \int_{0}^{\infty} e^{-\frac{1}{2 \lambda} t^{2}} t^{n-1} \bar{g}(x ; t) d t
\end{aligned}
$$

Suppose $n=2 k+1$ and let $c_{n}:=\sqrt{\frac{\pi}{2}} \frac{\sigma\left(S^{n-1}\right)}{(2 \pi)^{n / 2}}$, then the above equation reads

$$
\begin{aligned}
\int_{0}^{\infty} u(t, x) e^{-\frac{1}{2 \lambda} t^{2}} d t & =c_{n} \lambda^{-k} \int_{0}^{\infty} e^{-\frac{1}{2 \lambda} t^{2}} t^{2 k} \bar{g}(x ; t) d t \\
& =c_{n} \int_{0}^{\infty}\left(-\frac{1}{t} \partial_{t}\right)^{k} e^{-\frac{1}{2 \lambda} t^{2}} t^{2 k} \bar{g}(x ; t) d t \\
& \stackrel{\text { I.B.P. }}{=} c_{n} \int_{0}^{\infty} e^{-\frac{1}{2 \lambda} t^{2}}\left(\partial_{t} M_{t^{-1}}\right)^{k}\left[t^{2 k} \bar{g}(x ; t)\right] d t
\end{aligned}
$$

By the injectivity of the Laplace transform (after making the substitution $t \rightarrow \sqrt{t}$, this implies

$$
\begin{aligned}
\cos (\sqrt{-\Delta} t) g(x) & =u(t, x)=c_{n}\left(\partial_{t} M_{t^{-1}}\right)^{k}\left[t^{2 k} \bar{g}(x ; t)\right] \\
& =c_{n}\left(\partial_{t} M_{t^{-1}} \partial_{t} M_{t^{-1} \ldots \partial_{t} M_{t^{-1}}}\right)\left[t^{2 k} \bar{g}(x ; t)\right] \\
& =c_{n} \partial_{t}(\overbrace{M_{t^{-1}} \partial_{t} M_{t^{-1}}^{k-1} \ldots M_{t^{-1}} \partial_{t}}^{\text {times }})\left[t^{2 k-1} \bar{g}(x ; t)\right] \\
& =c_{n} \partial_{t}\left(\frac{1}{t} \partial_{t}\right)^{k-1}\left[t^{2 k-1} \bar{g}(x ; t)\right] .
\end{aligned}
$$

Hence we have derived the following theorem.
Theorem 47.12. Suppose $n=2 k+1$ is odd and let $c_{n}:=\sqrt{\frac{\pi}{2}} \frac{\sigma\left(S^{n-1}\right)}{(2 \pi)^{n / 2}}$, then
and
$\frac{\sin (\sqrt{-\Delta} t)}{\sqrt{-\Delta}} f(x)=\int_{0}^{t} \cos (\sqrt{-\Delta} \tau) f(x) d \tau=c_{n}\left(\frac{1}{t} \partial_{t}\right)^{k-1}\left[t^{2 k-1} \bar{g}(x ; t)\right]$.
Proof. For the last equality we have used

$$
\left(\frac{1}{t} \partial_{t}\right)^{k-1} t^{2 k-1}=\text { const. } * t^{2 k-1-2(k-1)}=\text { const. } *
$$

so that $\left(\frac{1}{t} \partial_{t}\right)^{k-1}\left[t^{2 k-1} \bar{g}(x ; t)\right]=O(t)$ and in particular is 0 at $t=0$.

### 47.6.2 Solution derived from the Poisson kernel

Suppose we want to write

$$
e^{-|x|}=\int_{0}^{\infty} \phi(s) p_{s}(x) d s
$$

Since

$$
\int_{\mathbb{R}} e^{-|x|} e^{i \lambda x} d x=2 \operatorname{Re} \int_{0}^{\infty} e^{-x} e^{i \lambda x} d x=2 \operatorname{Re}\left(\frac{1}{1-i \lambda}\right)=\frac{2}{1+\lambda^{2}}
$$

and

$$
\int_{\mathbb{R}} p_{s}(x) e^{i \lambda x} d x=\left.e^{s \partial_{x}^{2} / 2} e^{i \lambda x}\right|_{x=0}=e^{-s \lambda^{2} / 2}
$$

$\phi$ must satisfy

$$
\int_{0}^{\infty} \phi(s) e^{-s \lambda^{2} / 2} d s=\frac{2}{1+\lambda^{2}}=\int_{0}^{\infty} e^{-s\left(1+\lambda^{2}\right) / 2} d s=\int_{0}^{\infty} e^{-s / 2} e^{-s \lambda^{2} / 2} d s
$$

from which it follows that $\phi(s)=e^{-s / 2}$. Thus we have derived the formula

$$
\begin{equation*}
e^{-|x|}=\int_{0}^{\infty}(2 \pi s)^{-1 / 2} e^{-s / 2} e^{-\frac{1}{2 s} x^{2}} d s \tag{47.10}
\end{equation*}
$$

Let : $H \rightarrow H$ such that $A=A^{*}$ and $A \leq 0$. By the spectral theorem, we may "substitute" $x=t \sqrt{-A}$ into Eq. (47.10) to learn

$$
e^{-t \sqrt{-A}}=\int_{0}^{\infty}(2 \pi s)^{-1 / 2} e^{-s / 2} e^{\frac{t^{2}}{2 s} A} d s
$$

and in particular taking $A=\Delta$ one finds

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$$
\begin{aligned}
\frac{1}{\sqrt{-\Delta}} e^{-t \sqrt{-\Delta}} f(x) & =\left.\frac{-1}{\sqrt{-\Delta}} e^{-\tau \sqrt{-\Delta}} f(x)\right|_{\tau=t} ^{\infty} \\
& =\int_{t}^{\infty} e^{-\tau \sqrt{-\Delta}} f(x) d \tau \\
& =\int_{\mathbb{R}^{n}} \int_{t}^{\infty} d \tau Q_{\tau}(x-y) f(y) d y .
\end{aligned}
$$

Now

$$
\begin{aligned}
\int_{t}^{\infty} Q_{\tau}(x-y) d \tau & =c_{n} \int_{t}^{\infty} \frac{\tau}{\left(\tau^{2}+|x|^{2}\right)^{\frac{n+1}{2}}} d \tau=\left.\frac{c_{n}}{1-n}\left(\tau^{2}+|x|^{2}\right)^{\frac{1-n}{2}}\right|_{\tau=t} ^{\infty} \\
& =\frac{c_{n}}{n-1}\left(t^{2}+|x|^{2}\right)^{-\frac{n-1}{2}}
\end{aligned}
$$

and hence

$$
\frac{1}{\sqrt{-\Delta}} e^{-t \sqrt{-\Delta}} f(x)=\int_{\mathbb{R}^{n}} \frac{c_{n}}{n-1}\left(t^{2}+|y|^{2}\right)^{-\frac{n-1}{2}} f(x-y) d y
$$

and by analytic continuation,

$$
\begin{aligned}
\frac{1}{\sqrt{-\Delta}} e^{(i t-\epsilon) \sqrt{-\Delta}} f(x) & =\frac{1}{\sqrt{-\Delta}} e^{-(\epsilon-i t) \sqrt{-\Delta}} f(x) \\
& =\frac{c_{n}}{n-1} \int_{\mathbb{R}^{n}}\left((\epsilon-i t)^{2}+|y|^{2}\right)^{-\frac{n-1}{2}} f(x-y) d y \\
& =\frac{c_{n}}{n-1} \int_{\mathbb{R}^{n}}\left(|y|^{2}-(t-i \epsilon)^{2}\right)^{-\frac{n-1}{2}} f(x-y) d y
\end{aligned}
$$

and hence

$$
\frac{1}{\sqrt{-\Delta}} \sin (t \sqrt{-\Delta}) f(x)=c_{n}^{\prime} \lim _{\epsilon \downarrow 0} \int_{\mathbb{R}^{n}} \operatorname{Im}\left(|y|^{2}-(t-i \epsilon)^{2}\right)^{-\frac{n-1}{2}} f(x-y) d y
$$

Now if $|y|>|t|$ then

$$
\lim _{\epsilon \downarrow 0}\left(|y|^{2}-(t-i \epsilon)^{2}\right)^{-\frac{n-1}{2}}=\left(|y|^{2}-t^{2}\right)^{-\frac{n-1}{2}}
$$

is real so

$$
\lim _{\epsilon \downarrow 0} \operatorname{Im}\left(|y|^{2}-(t-i \epsilon)^{2}\right)^{-\frac{n-1}{2}}=0 \text { if }|y|>|t| .
$$

Similarly if $n$ is odd $\lim _{\epsilon \downarrow 0}\left(|y|^{2}-(t-i \epsilon)^{2}\right)^{-\frac{n-1}{2}}=\left(|y|^{2}-t^{2}\right)^{-\frac{n-1}{2}} \in \mathbb{R}$ and so

$$
\lim _{\epsilon \downarrow 0} \operatorname{Im}\left(|y|^{2}-(t-i \epsilon)^{2}\right)^{-\frac{n-1}{2}}
$$

is a distribution concentrated on the sphere $|y|=|t|$ which is the sharp propagation again. See Taylor Vol. 1., p. 221-225 for more on this approach. Let us examine here the special case $n=3$,
$\operatorname{Im}\left(\frac{1}{|y|^{2}-(t-i \epsilon)^{2}}\right)=\operatorname{Im}\left(\frac{1}{|y|^{2}-t^{2}+\epsilon^{2}+2 i \epsilon t}\right)=\frac{-2 \epsilon t}{\left(|y|^{2}-t^{2}+\epsilon^{2}\right)^{2}+4 \epsilon^{2} t^{2}}$
So

$$
\begin{aligned}
I & :=\lim _{\epsilon \downharpoonright 0} \int_{\mathbb{R}^{n}} \operatorname{Im}\left(\frac{1}{|y|^{2}-(t-i \epsilon)^{2}}\right) f(x-y) d y \\
& =\lim _{\epsilon \downarrow 0} \int_{\mathbb{R}^{n}} \frac{-2 \epsilon t}{\left(|y|^{2}-t^{2}+\epsilon^{2}\right)^{2}+4 \epsilon^{2} t^{2}} f(x-y) d y \\
& =4 \pi \lim _{\epsilon \downarrow 0} \int_{0}^{\infty} \rho^{2} \frac{-2 \epsilon t}{\left(\rho^{2}-t^{2}+\epsilon^{2}\right)^{2}+4 \epsilon^{2} t^{2}} \bar{f}(x ; \rho) d \rho \\
& =c t \lim _{\epsilon \downarrow 0} \int_{0}^{\infty} \rho^{2} \frac{\epsilon}{\left(\rho^{2}-t^{2}+\epsilon^{2}\right)^{2}+4 \epsilon^{2} t^{2}} \bar{f}(x ; \rho) d \rho
\end{aligned}
$$

Make the change of variables $\rho=t+\epsilon s$ above to find

$$
\begin{aligned}
I & =c t \lim _{\epsilon \downarrow 0} \int_{-t / \epsilon}^{\infty} \frac{(t+\epsilon s)^{2} \epsilon^{2}}{\left(2 \epsilon s t+\epsilon^{2} s^{2}+\epsilon^{2}\right)^{2}+4 \epsilon^{2} t^{2}} \bar{f}(x ; t+\epsilon s) d s \\
& =c t \lim _{\epsilon \downarrow 0} \int_{-t / \epsilon}^{\infty} \frac{(t+\epsilon s)^{2}}{\left(2 s t+\epsilon s^{2}+\epsilon\right)^{2}+4 t^{2}} \bar{f}(x ; t+\epsilon s) d s \\
& =c t \bar{f}(x ; t) \int_{-\infty}^{\infty} \frac{t^{2}}{4 t^{2} s^{2}+4 t^{2}} d s=\frac{c}{4} t \bar{f}(x ; t) \int_{-\infty}^{\infty} \frac{1}{s^{2}+1} d s \\
& =\frac{c}{4} \pi t \bar{f}(x ; t)
\end{aligned}
$$

which up to an overall constant is the result that we have seen before.

### 47.7 Explain Method of descent $n=2$

$$
u(t, x)=\frac{1}{2} \int_{B(x, t)} \frac{t g(y)+t^{2} h(y)+t \nabla g(y) \cdot(y-x)}{\left(t^{2}-|y-x|^{2}\right)^{1 / 2}} d y .
$$

See constant coefficient PDE notes for more details on this.

Sobolev Theory

## Sobolev Spaces

Definition 48.1. For $p \in[1, \infty], k \in \mathbb{N}$ and $\Omega$ an open subset of $\mathbb{R}^{d}$, let

$$
\begin{aligned}
W_{l o c}^{k, p}(\Omega) & :=\left\{f \in L^{p}(\Omega): \partial^{\alpha} f \in L_{l o c}^{p}(\Omega) \text { (weakly) for all }|\alpha| \leq k\right\}, \\
W^{k, p}(\Omega) & :=\left\{f \in L^{p}(\Omega): \partial^{\alpha} f \in L^{p}(\Omega) \text { (weakly) for all }|\alpha| \leq k\right\}
\end{aligned}
$$

$$
\begin{equation*}
\|f\|_{W^{k, p}(\Omega)}:=\left(\sum_{|\alpha| \leq k}\left\|\partial^{\alpha} f\right\|_{L^{p}(\Omega)}^{p}\right)^{1 / p} \text { if } p<\infty \tag{48.1}
\end{equation*}
$$

and

$$
\begin{equation*}
\|f\|_{W^{k, p}(\Omega)}=\sum_{|\alpha| \leq k}\left\|\partial^{\alpha} f\right\|_{L^{\infty}(\Omega)} \text { if } p=\infty \tag{48.2}
\end{equation*}
$$

In the special case of $p=2$, we write $W_{l o c}^{k, 2}(\Omega)=: H_{l o c}^{k}(\Omega)$ and $W^{k, 2}(\Omega)=$ : $H^{k}(\Omega)$ in which case $\|\cdot\|_{W^{k, 2}(\Omega)}=\|\cdot\|_{H^{k}(\Omega)}$ is a Hilbertian norm associated to the inner product

$$
\begin{equation*}
(f, g)_{H^{k}(\Omega)}=\sum_{|\alpha| \leq k} \int_{\Omega} \partial^{\alpha} f \cdot \overline{\partial^{\alpha} g} d m . \tag{48.3}
\end{equation*}
$$

Theorem 48.2. The function, $\|\cdot\|_{W^{k, p}(\Omega)}$, is a norm which makes $W^{k, p}(\Omega)$ into a Banach space.

Proof. Let $f, g \in W^{k, p}(\Omega)$, then the triangle inequality for the $p$ - norms on $L^{p}(\Omega)$ and $l^{p}(\{\alpha:|\alpha| \leq k\})$ implies

$$
\begin{aligned}
\|f+g\|_{W^{k, p}(\Omega)} & =\left(\sum_{|\alpha| \leq k}\left\|\partial^{\alpha} f+\partial^{\alpha} g\right\|_{L^{p}(\Omega)}^{p}\right)^{1 / p} \\
& \leq\left(\sum_{|\alpha| \leq k}\left[\left\|\partial^{\alpha} f\right\|_{L^{p}(\Omega)}+\left\|\partial^{\alpha} g\right\|_{L^{p}(\Omega)}\right]^{p}\right)^{1 / p} \\
& \leq\left(\sum_{|\alpha| \leq k}\left\|\partial^{\alpha} f\right\|_{L^{p}(\Omega)}^{p}\right)^{1 / p}+\left(\sum_{|\alpha| \leq k}\left\|\partial^{\alpha} g\right\|_{L^{p}(\Omega)}^{p}\right)^{1 / p} \\
& =\|f\|_{W^{k, p}(\Omega)}+\|g\|_{W^{k, p}(\Omega)} .
\end{aligned}
$$

This shows $\|\cdot\|_{W^{k, p}(\Omega)}$ defined in Eq. (48.1) is a norm. We now show completeness.

If $\left\{f_{n}\right\}_{n=1}^{\infty} \subset W^{k, p}(\Omega)$ is a Cauchy sequence, then $\left\{\partial^{\alpha} f_{n}\right\}_{n=1}^{\infty}$ is a Cauchy sequence in $L^{p}(\Omega)$ for all $|\alpha| \leq k$. By the completeness of $L^{p}(\Omega)$, there exists $g_{\alpha} \in L^{p}(\Omega)$ such that $g_{\alpha}=L^{p}-\lim _{n \rightarrow \infty} \partial^{\alpha} f_{n}$ for all $|\alpha| \leq k$. Therefore, for all $\phi \in C_{c}^{\infty}(\Omega)$,

$$
\left\langle f, \partial^{\alpha} \phi\right\rangle=\lim _{n \rightarrow \infty}\left\langle f_{n}, \partial^{\alpha} \phi\right\rangle=(-1)^{|\alpha|} \lim _{n \rightarrow \infty}\left\langle\partial^{\alpha} f_{n}, \phi\right\rangle=(-1)^{|\alpha|} \lim _{n \rightarrow \infty}\left\langle g_{\alpha}, \phi\right\rangle .
$$

This shows $\partial^{\alpha} f$ exists weakly and $g_{\alpha}=\partial^{\alpha} f$ a.e. This shows $f \in W^{k, p}(\Omega)$ and that $f_{n} \rightarrow f \in W^{k, p}(\Omega)$ as $n \rightarrow \infty$.

Example 48.3. Let $u(x):=|x|^{-\alpha}$ for $x \in \mathbb{R}^{d}$ and $\alpha \in \mathbb{R}$. Then

$$
\begin{align*}
\int_{B(0, R)}|u(x)|^{p} d x & =\sigma\left(S^{d-1}\right) \int_{0}^{R} \frac{1}{r^{\alpha p}} r^{d-1} d r=\sigma\left(S^{d-1}\right) \int_{0}^{R} r^{d-\alpha p-1} d r \\
& =\sigma\left(S^{d-1}\right) \cdot\left\{\begin{array}{cc}
\frac{R^{d-\alpha p}}{d-\alpha p} & \text { if } \quad d-\alpha p>0 \\
\infty & \text { otherwise }
\end{array}\right. \tag{48.4}
\end{align*}
$$

and hence $u \in L_{l o c}^{p}\left(\mathbb{R}^{d}\right)$ iff $\alpha<d / p$. Now $\nabla u(x)=-\alpha|x|^{-\alpha-1} \hat{x}$ where $\hat{x}:=x /|x|$. Hence if $\nabla u(x)$ is to exist in $L_{l o c}^{p}\left(\mathbb{R}^{d}\right)$ it is given by $-\alpha|x|^{-\alpha-1} \hat{x}$ which is in $L_{\text {loc }}^{p}\left(\mathbb{R}^{d}\right)$ iff $\alpha+1<d / p$, i.e. if $\alpha<d / p-1=\frac{d-p}{p}$. Let us not check that $u \in W_{l o c}^{1, p}\left(\mathbb{R}^{d}\right)$ provided $\alpha<d / p-1$. To do this suppose $\phi \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ and $\epsilon>0$, then

$$
\begin{aligned}
-\left\langle u, \partial_{i} \phi\right\rangle & =-\lim _{\epsilon \downarrow 0} \int_{|x|>\epsilon} u(x) \partial_{i} \phi(x) d x \\
& =\lim _{\epsilon \downarrow 0}\left\{\int_{|x|>\epsilon} \partial_{i} u(x) \phi(x) d x+\int_{|x|=\epsilon} u(x) \phi(x) \frac{x_{i}}{\epsilon} d \sigma(x)\right\}
\end{aligned}
$$

$$
\left|\int_{|x|=\epsilon} u(x) \phi(x) \frac{x_{i}}{\epsilon} d \sigma(x)\right| \leq\|\phi\|_{\infty} \sigma\left(S^{d-1}\right) \epsilon^{d-1-\alpha} \rightarrow 0 \text { as } \epsilon \downarrow 0
$$

and $\partial_{i} u(x)=-\alpha|x|^{-\alpha-1} \hat{x} \cdot e_{i}$ is locally integrable we conclude that

$$
-\left\langle u, \partial_{i} \phi\right\rangle=\int_{\mathbb{R}^{d}} \partial_{i} u(x) \phi(x) d x
$$

showing that the weak derivative $\partial_{i} u$ exists and is given by the usual pointwise derivative.

### 48.1 Mollifications

Proposition 48.4 (Mollification). Let $\Omega$ be an open subset of $\mathbb{R}^{d}, k \in$ $\mathbb{N}_{0}:=\mathbb{N} \cup\{0\}, p \in[1, \infty)$ and $u \in W_{l o c}^{k, p}(\Omega)$. Then there exists $u_{n} \in C_{c}^{\infty}(\Omega)$ such that $u_{n} \rightarrow u$ in $W_{l o c}^{k, p}(\Omega)$.

Proof. Apply Proposition 29.12 with polynomials, $p_{\alpha}(\xi)=\xi^{\alpha}$, for $|\alpha| \leq k$.

Proposition 48.5. $C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ is dense in $W^{k, p}\left(\mathbb{R}^{d}\right)$ for all $1 \leq p<\infty$.
Proof. The proof is similar to the proof of Proposition 48.4 using Exercise 29.32 in place of Proposition 29.12.

Proposition 48.6. Let $\Omega$ be an open subset of $\mathbb{R}^{d}, k \in \mathbb{N}_{0}:=\mathbb{N} \cup\{0\}$ and $p \geq 1$, then

1. for any $\alpha$ with $|\alpha| \leq k, \partial^{\alpha}: W^{k, p}(\Omega) \rightarrow W^{k-|\alpha|, p}(\Omega)$ is a contraction.
2. For any open subset $V \subset \Omega$, the restriction map $\left.u \rightarrow u\right|_{V}$ is bounded from $W^{k, p}(\Omega) \rightarrow W^{k, p}(V)$.
3. For any $f \in C^{k}(\Omega)$ and $u \in W_{\text {loc }}^{k, p}(\Omega)$, the $f u \in W_{\text {loc }}^{k, p}(\Omega)$ and for $|\alpha| \leq k$,

$$
\begin{equation*}
\partial^{\alpha}(f u)=\sum_{\beta \leq \alpha}\binom{\alpha}{\beta} \partial^{\beta} f \cdot \partial^{\alpha-\beta} u \tag{48.5}
\end{equation*}
$$

where $\binom{\alpha}{\beta}:=\frac{\alpha!}{\beta!(\alpha-\beta)!}$.
4. For any $f \in B C^{k}(\Omega)$ and $u \in W_{l o c}^{k, p}(\Omega)$, the $f u \in W_{l o c}^{k, p}(\Omega)$ and for $|\alpha| \leq k$ Eq. (48.5) still holds. Moreover, the linear map $u \in W^{k, p}(\Omega) \rightarrow$ $f u \in W^{k, p}(\Omega)$ is a bounded operator.
Proof. 1. Let $\phi \in C_{c}^{\infty}(\Omega)$ and $u \in W^{k, p}(\Omega)$, then for $\beta$ with $|\beta| \leq k-|\alpha|$, $\left\langle\partial^{\alpha} u, \partial^{\beta} \phi\right\rangle=(-1)^{|\alpha|}\left\langle u, \partial^{\alpha} \partial^{\beta} \phi\right\rangle=(-1)^{|\alpha|}\left\langle u, \partial^{\alpha+\beta} \phi\right\rangle=(-1)^{|\beta|}\left\langle\partial^{\alpha+\beta} u, \phi\right\rangle$
from which it follows that $\partial^{\beta}\left(\partial^{\alpha} u\right)$ exists weakly and $\partial^{\beta}\left(\partial^{\alpha} u\right)=\partial^{\alpha+\beta} u$. This shows that $\partial^{\alpha} u \in W^{k-|\alpha|, p}(\Omega)$ and it should be clear that $\left\|\partial^{\alpha} u\right\|_{W^{k-|\alpha|, p}(\Omega)} \leq$ $\|u\|_{W^{k, p}(\Omega)}$.

Item 2. is trivial.
3-4. Given $u \in W_{l o c}^{k, p}(\Omega)$, by Proposition 48.4 there exists $u_{n} \in C_{c}^{\infty}(\Omega)$ such that $u_{n} \rightarrow u$ in $W_{l o c}^{k, p}(\Omega)$. From the results in Appendix A.1, $f u_{n} \in$ $C_{c}^{k}(\Omega) \subset W^{k, p}(\Omega)$ and

$$
\begin{equation*}
\partial^{\alpha}\left(f u_{n}\right)=\sum_{\beta \leq \alpha}\binom{\alpha}{\beta} \partial^{\beta} f \cdot \partial^{\alpha-\beta} u_{n} \tag{48.6}
\end{equation*}
$$

holds. Given $V \subset_{o} \Omega$ such that $\bar{V}$ is compactly contained in $\Omega$, we may use the above equation to find the estimate

$$
\begin{aligned}
\left\|\partial^{\alpha}\left(f u_{n}\right)\right\|_{L^{p}(V)} & \leq \sum_{\beta \leq \alpha}\binom{\alpha}{\beta}\left\|\partial^{\beta} f\right\|_{L^{\infty}(V)}\left\|\partial^{\alpha-\beta} u_{n}\right\|_{L^{p}(V)} \\
& \leq C_{\alpha}(f, V) \sum_{\beta \leq \alpha}\left\|\partial^{\alpha-\beta} u_{n}\right\|_{L^{p}(V)} \leq C_{\alpha}(f, V)\left\|u_{n}\right\|_{W^{k, p}(V)}
\end{aligned}
$$

wherein the last equality we have used Exercise 48.36 below. Summing this equation on $|\alpha| \leq k$ shows

$$
\begin{equation*}
\left\|f u_{n}\right\|_{W^{k, p}(V)} \leq C(f, V)\left\|u_{n}\right\|_{W^{k, p}(V)} \text { for all } n \tag{48.7}
\end{equation*}
$$

where $C(f, V):=\sum_{|\alpha| \leq k} C_{\alpha}(f, V)$. By replacing $u_{n}$ by $u_{n}-u_{m}$ in the above inequality it follows that $\left\{f u_{n}\right\}_{n=1}^{\infty}$ is convergent in $W^{k, p}(V)$ and since $V$ was arbitrary $f u_{n} \rightarrow f u$ in $W_{l o c}^{k, p}(\Omega)$. Moreover, we may pass to the limit in Eq. (48.6) and in Eq. (48.7) to see that Eq. (48.5) holds and that

$$
\|f u\|_{W^{k, p}(V)} \leq C(f, V)\|u\|_{W^{k, p}(V)} \leq C(f, V)\|u\|_{W^{k, p}(\Omega)}
$$

Moreover if $f \in B C(\Omega)$ then constant $C(f, V)$ may be chosen to be independent of $V$ and therefore, if $u \in W^{k, p}(\Omega)$ then $f u \in W^{k, p}(\Omega)$.

Alternative direct proof of 4 . We will prove this by induction on $|\alpha|$. If $\alpha=e_{i}$ then, using Lemma 29.9,

$$
\begin{aligned}
-\left\langle f u, \partial_{i} \phi\right\rangle & =-\left\langle u, f \partial_{i} \phi\right\rangle=-\left\langle u, \partial_{i}[f \phi]-\partial_{i} f \cdot \phi\right\rangle \\
& =\left\langle\partial_{i} u, f \phi\right\rangle+\left\langle u, \partial_{i} f \cdot \phi\right\rangle=\left\langle f \partial_{i} u+\partial_{i} f \cdot u, \phi\right\rangle
\end{aligned}
$$

showing $\partial_{i}(f u)$ exists weakly and is equal to $\partial_{i}(f u)=f \partial_{i} u+\partial_{i} f \cdot u \in L^{p}(\Omega)$. Supposing the result has been proved for all $\alpha$ such that $|\alpha| \leq m$ with $m \in$ $[1, k)$. Let $\gamma=\alpha+e_{i}$ with $|\alpha|=m$, then by what we have just proved each summand in Eq. (48.5) satisfies $\partial_{i}\left[\partial^{\beta} f \cdot \partial^{\alpha-\beta} u\right]$ exists weakly and

$$
\partial_{i}\left[\partial^{\beta} f \cdot \partial^{\alpha-\beta} u\right]=\partial^{\beta+e_{i}} f \cdot \partial^{\alpha-\beta} u+\partial^{\beta_{i}} f \cdot \partial^{\alpha-\beta+e} u \in L^{p}(\Omega)
$$

Therefore $\partial^{\gamma}(f u)=\partial_{i} \partial^{\alpha}(f u)$ exists weakly in $L^{p}(\Omega)$ and
$\partial^{\gamma}(f u)=\sum_{\beta \leq \alpha}\binom{\alpha}{\beta}\left[\partial^{\beta+e_{i}} f \cdot \partial^{\alpha-\beta} u+\partial^{\beta} f \cdot \partial^{\alpha-\beta+e_{i}} u\right]=\sum_{\beta \leq \gamma}\binom{\gamma}{\beta}\left[\partial^{\beta} f \cdot \partial^{\gamma-\beta} u\right]$
For the last equality see the combinatorics in Appendix A.1.
Theorem 48.7. Let $\Omega$ be an open subset of $\mathbb{R}^{d}, k \in \mathbb{N}_{0}:=\mathbb{N} \cup\{0\}$ and $p \in$ $1, \infty)$. Then $C^{\infty}(\Omega) \cap W^{k, p}(\Omega)$ is dense in $W^{k, p}(\Omega)$.

Proof. Let $\Omega_{n}:=\{x \in \Omega: \operatorname{dist}(x, \Omega)>1 / n\} \cap B(0, n)$, then

$$
\bar{\Omega}_{n} \subset\{x \in \Omega: \operatorname{dist}(x, \Omega) \geq 1 / n\} \cap \overline{B(0, n)} \subset \Omega_{n+1},
$$

$\bar{\Omega}_{n}$ is compact for every $n$ and $\Omega_{n} \uparrow \Omega$ as $n \rightarrow \infty$. Let $V_{0}=\Omega_{3}, V_{j}:=\Omega_{j+3} \backslash \bar{\Omega}_{j}$ for $j \geq 1, K_{0}:=\bar{\Omega}_{2}$ and $K_{j}:=\bar{\Omega}_{j+2} \backslash \Omega_{j+1}$ for $j \geq 1$ as in figure 48.1. Then


Fig. 48.1. Decomposing $\Omega$ into compact pieces. The compact sets $K_{0}, K_{1}$ and $K_{2}$ are the shaded annular regions while $V_{0}, V_{1}$ and $V_{2}$ are the indicated open annular regions.
$K_{n} \sqsubset \sqsubset V_{n}$ for all $n$ and $\cup K_{n}=\Omega$. Choose $\phi_{n} \in C_{c}^{\infty}\left(V_{n},[0,1]\right)$ such that $\phi_{n}=1$ on $K_{n}$ and set $\psi_{0}=\phi_{0}$ and

$$
\psi_{j}=\left(1-\psi_{1}-\cdots-\psi_{j-1}\right) \phi_{j}=\phi_{j} \prod_{k=1}^{j-1}\left(1-\phi_{k}\right)
$$

for $j \geq 1$. Then $\psi_{j} \in C_{c}^{\infty}\left(V_{n},[0,1]\right)$,

$$
1-\sum_{k=0}^{n} \psi_{k}=\prod_{k=1}^{n}\left(1-\phi_{k}\right) \rightarrow 0 \text { as } n \rightarrow \infty
$$

so that $\sum_{k=0}^{\infty} \psi_{k}=1$ on $\Omega$ with the sum being locally finite.
Let $\epsilon>0$ be given. By Proposition 48.6, $u_{n}:=\psi_{n} u \in W^{k, p}(\Omega)$ with $\operatorname{supp}\left(u_{n}\right) \sqsubset \sqsubset V_{n}$. By Proposition 48.4, we may find $v_{n} \in C_{c}^{\infty}\left(V_{n}\right)$ such that $\left\|u_{n}-v_{n}\right\|_{W^{k, p}(\Omega)} \leq \epsilon / 2^{n+1}$ for all $n$. Let $v:=\sum_{n=1}^{\infty} v_{n}$, then $v \in C^{\infty}(\Omega)$ because the sum is locally finite. Since

$$
\sum_{n=0}^{\infty}\left\|u_{n}-v_{n}\right\|_{W^{k, p}(\Omega)} \leq \sum_{n=0}^{\infty} \epsilon / 2^{n+1}=\epsilon<\infty,
$$

the sum $\sum_{n=0}^{\infty}\left(u_{n}-v_{n}\right)$ converges in $W^{k, p}(\Omega)$. The sum, $\sum_{n=0}^{\infty}\left(u_{n}-v_{n}\right)$, also converges pointwise to $u-v$ and hence $u-v=\sum_{n=0}^{\infty}\left(u_{n}-v_{n}\right)$ is in $W^{k, p}(\Omega)$. Therefore $v \in W^{k, p}(\Omega) \cap C^{\infty}(\Omega)$ and

$$
\|u-v\| \leq \sum_{n=0}^{\infty}\left\|u_{n}-v_{n}\right\|_{W^{k, p}(\Omega)} \leq \epsilon
$$

Notation 48.8 Given a closed subset $F \subset \mathbb{R}^{d}$, let $C^{\infty}(F)$ denote those $u \in$ $C(F)$ that extend to a $C^{\infty}$ - function on an open neighborhood of $F$.

Remark 48.9. It is easy to prove that $u \in C^{\infty}(F)$ iff there exists $U \in C^{\infty}\left(\mathbb{R}^{d}\right)$ such that $u=\left.U\right|_{F}$. Indeed, suppose $\Omega$ is an open neighborhood of $F, f \in$ $C^{\infty}(\Omega)$ and $u=\left.f\right|_{F} \in C^{\infty}(F)$. Using a partition of unity argument (making use of the open sets $V_{i}$ constructed in the proof of Theorem 48.7), one may show there exists $\phi \in C^{\infty}(\Omega,[0,1])$ such that $\operatorname{supp}(\phi) \sqsubset \Omega$ and $\phi=1$ on a neighborhood of $F$. Then $U:=\phi f$ is the desired function.
Theorem 48.10 (Density of $W^{k, p}(\Omega) \cap C^{\infty}(\bar{\Omega})$ in $W^{k, p}(\Omega)$ ). Let $\Omega \subset$ $\mathbb{R}^{d}$ be a manifold with $C^{0}$ - boundary, then for $k \in \mathbb{N}_{0}$ and $p \in[1, \infty)$, $W^{k, p}\left(\Omega^{o}\right) \cap C^{\infty}(\bar{\Omega})$ is dense in $W^{k, p}\left(\Omega^{0}\right)$. This may alternatively be stated by assuming $\Omega \subset \mathbb{R}^{d}$ is an open set such that $\bar{\Omega}^{o}=\Omega$ and $\bar{\Omega}$ is a manifold with $C^{0}$ - boundary, then $W^{k, p}(\Omega) \cap C^{\infty}(\bar{\Omega})$ is dense in $W^{k, p}(\Omega)$.

Before going into the proof, let us point out that some restriction on the boundary of $\Omega$ is needed for assertion in Theorem 48.10 to be valid. For example, suppose

$$
\Omega_{0}:=\left\{x \in \mathbb{R}^{2}: 1<|x|<2\right\} \text { and } \Omega:=\Omega_{0} \backslash\{(1,2) \times\{0\}\}
$$

and $\theta: \Omega \rightarrow(0,2 \pi)$ is defined so that $x_{1}=|x| \cos \theta(x)$ and $x_{2}=|x| \sin \theta(x)$, see Figure 48.2. Then $\theta \in B C^{\infty}(\Omega) \subset W^{k, \infty}(\Omega)$ for all $k \in \mathbb{N}_{0}$ yet $\theta$ can not be approximated by functions from $C^{\infty}(\bar{\Omega}) \subset B C^{\infty}\left(\Omega_{0}\right)$ in $W^{1, p}(\Omega)$. Indeed, if this were possible, it would follows that $\theta \in W^{1, p}\left(\Omega_{0}\right)$. However, $\theta$ is


Fig. 48.2 The region $\Omega_{0}$ along with a vertical in $\Omega$.
not continuous (and hence not absolutely continuous) on the lines $\left\{x_{1}=\rho\right\} \cap \Omega$ for all $\rho \in(1,2)$ and so by Theorem $29.30, \theta \notin W^{1, p}\left(\Omega_{0}\right)$.

The following is a warm-up to the proof of Theorem 48.10
Proposition 48.11 (Warm-up). Let $f: \mathbb{R}^{d-1} \rightarrow \mathbb{R}$ be a continuous function and $\Omega:=\left\{x \in \mathbb{R}^{d}: x_{d}>f\left(x_{1}, \ldots, x_{d-1}\right)\right\}$ and $C^{\infty}(\bar{\Omega})$ denote those $u \in C(\bar{\Omega})$ which are restrictions of $C^{\infty}$ - functions defined on an open neighborhood of $\bar{\Omega}$. Then for $p \in[1, \infty), C^{\infty}(\bar{\Omega}) \cap W^{k, p}(\Omega)$ is dense in $W^{k, p}(\Omega)$.

Proof. By Theorem 48.7, it suffices to show than any $u \in C^{\infty}(\Omega) \cap$ $W^{k, p}(\Omega)$ may be approximated by elements of $C^{\infty}(\bar{\Omega}) \cap W^{k, p}(\Omega)$. For $s>0$ let $u_{s}(x):=u\left(x+s e_{d}\right)$ which is defined for $x \in \Omega-s e_{d}$. Since

$$
\begin{aligned}
& \bar{\Omega}=\left\{x \in \mathbb{R}^{d}: x_{d} \geq f\left(x_{1}, \ldots, x_{d-1}\right)\right\} \\
& \\
& \quad \subset\left\{x \in \mathbb{R}^{d}: x_{d}+s>f\left(x_{1}, \ldots, x_{d-1}\right)\right\}=\Omega-s e_{d}
\end{aligned}
$$

and $\partial^{\alpha} u_{s}=\left(\partial^{\alpha} u\right)_{s}$ for all $\alpha$,

$$
u_{s} \in W^{k, p}\left(\Omega-s e_{d}\right) \cap C^{\infty}\left(\Omega-s e_{d}\right) \subset C^{\infty}(\bar{\Omega}) \cap W^{k, p}(\Omega)
$$

These observations along with the strong continuity of translations in $L^{p}$ (see Proposition 11.13), implies $\lim _{s \downarrow 0}\left\|u-u_{s}\right\|_{W^{k, p}(\Omega)}=0$.

### 48.1.1 Proof of Theorem 48.10

Proof. By Theorem 48.7, it suffices to show than any $u \in C^{\infty}(\Omega) \cap W^{k, p}(\Omega)$ may be approximated by elements of $C^{\infty}(\bar{\Omega}) \cap W^{k, p}(\Omega)$. To understand the main ideas of the proof, suppose that $\Omega$ is the triangular region in Figure 48.3 and suppose that we have used a partition of unity relative to the cover shown so that $u=u_{1}+u_{2}+u_{3}$ with $\operatorname{supp}\left(u_{i}\right) \subset B_{i}$. Now concentrating on


Fig. 48.3. Splitting and moving a function in $C^{\infty}(\Omega)$ so that the result is in $C^{\infty}(\bar{\Omega})$.
$u_{1}$ whose support is depicted as the grey shaded area in Figure 48.3. We now simply translate $u_{1}$ in the direction $v$ shown in Figure 48.3. That is for any small $s>0$, let $w_{s}(x):=u_{1}(x+s v)$, then $v_{s}$ lives on the translated grey area as seen in Figure 48.3. The function $w_{s}$ extended to be zero off its domain of definition is an element of $C^{\infty}(\bar{\Omega})$ moreover it is easily seen, using the same methods as in the proof of Proposition 48.11, that $w_{s} \rightarrow u_{1}$ in $W^{k, p}(\Omega)$

The formal proof follows along these same lines. To do this choose an at most countable locally finite cover $\left\{V_{i}\right\}_{i=0}^{\infty}$ of $\bar{\Omega}$ such that $\bar{V}_{0} \subset \Omega$ and for each $i \geq 1$, after making an affine change of coordinates, $V_{i}=(-\epsilon, \epsilon)^{d}$ for some $\epsilon>0$ and

$$
V_{i} \cap \bar{\Omega}=\left\{(y, z) \in V_{i}: \epsilon>z>f_{i}(y)\right\}
$$

where $f_{i}:(-\epsilon, \epsilon)^{d-1} \rightarrow(-\epsilon, \epsilon)$, see Figure 48.4 below. Let $\left\{\eta_{i}\right\}_{i=0}^{\infty}$ be a par-


Fig. 48.4. The shaded area depicts the support of $u_{i}=u \eta_{i}$.
tition of unity subordinated to $\left\{V_{i}\right\}$ and let $u_{i}:=u \eta_{i} \in C^{\infty}\left(V_{i} \cap \Omega\right)$. Given $\delta>0$, we choose $s$ so small that $w_{i}(x):=u_{i}\left(x+s e_{d}\right)$ (extended to be zero off its domain of definition) may be viewed as an element of $C^{\infty}(\bar{\Omega})$ and such
that $\left\|u_{i}-w_{i}\right\|_{W^{k, p}(\Omega)}<\delta / 2^{i}$. For $i=0$ we set $w_{0}:=u_{0}=u \eta_{0}$. Then, since $\left\{V_{i}\right\}_{i=1}^{\infty}$ is a locally finite cover of $\bar{\Omega}$, it follows that $w:=\sum_{i=0}^{\infty} w_{i} \in C^{\infty}(\bar{\Omega})$ and further we have

$$
\sum_{i=0}^{\infty}\left\|u_{i}-w_{i}\right\|_{W^{k, p}(\Omega)} \leq \sum_{i=1}^{\infty} \delta / 2^{i}=\delta .
$$

This shows

$$
u-w=\sum_{i=0}^{\infty}\left(u_{i}-w_{i}\right) \in W^{k, p}(\Omega)
$$

and $\|u-w\|_{W^{k, p}(\Omega)}<\delta$. Hence $w \in C^{\infty}(\bar{\Omega}) \cap W^{k, p}(\Omega)$ is a $\delta$ - approximation of $u$ and since $\delta>0$ arbitrary the proof is complete.

### 48.2 Difference quotients

Recall from Notation 29.14 that for $h \neq 0$

$$
\partial_{i}^{h} u(x):=\frac{u\left(x+h e^{i}\right)-u(x)}{h} .
$$

Remark 48.12 (Adjoints of Finite Differences). For $u \in L^{p}$ and $g \in L^{q}$,

$$
\begin{aligned}
\int_{\mathbb{R}^{d}} \partial_{i}^{h} u(x) g(x) d x & =\int_{\mathbb{R}^{d}} \frac{u\left(x+h e_{i}\right)-u(x)}{h} g(x) d x \\
& =-\int_{\mathbb{R}^{d}} u(x) \frac{g\left(x-h e_{i}\right)-g(x)}{-h} d x \\
& =-\int_{\mathbb{R}^{d}} u(x) \partial_{i}^{-h} g(x) d x .
\end{aligned}
$$

We summarize this identity by $\left(\partial_{i}^{h}\right)^{*}=-\partial_{i}^{-h}$.
Theorem 48.13. Suppose $k \in \mathbb{N}_{0}, \Omega$ is an open subset of $\mathbb{R}^{d}$ and $V$ is an open precompact subset of $\Omega$.

1. If $1 \leq p<\infty, u \in W^{k, p}(\Omega)$ and $\partial_{i} u \in W^{k, p}(\Omega)$, then

$$
\begin{equation*}
\left\|\partial_{i}^{h} u\right\|_{W^{k, p}(V)} \leq\left\|\partial_{i} u\right\|_{W^{k, p}(\Omega)} \tag{48.8}
\end{equation*}
$$

for all $0<|h|<\frac{1}{2} \operatorname{dist}\left(V, \Omega^{c}\right)$.
2. Suppose that $1<p \leq \infty, u \in W^{k, p}(\Omega)$ and assume there exists a constant $C(V)<\infty$ such that

$$
\left\|\partial_{i}^{h} u\right\|_{W^{k, p}(V)} \leq C(V) \text { for all } 0<|h|<\frac{1}{2} \operatorname{dist}\left(V, \Omega^{c}\right)
$$

Then $\partial_{i} u \in W^{k, p}(V)$ and $\left\|\partial_{i} u\right\|_{W^{k, p}(V)} \leq C(V)$. Moreover if $C:=$ $\sup _{V \subset \subset \Omega} C(V)<\infty$ then in fact $\partial_{i} u \in W^{k, p}(\Omega)$ and there is a constant $c<\infty$ such that

$$
\left\|\partial_{i} u\right\|_{W^{k, p}(\Omega)} \leq c\left(C+\|u\|_{L^{p}(\Omega)}\right)
$$

Proof. 1. Let $|\alpha| \leq k$, then

$$
\left\|\partial^{\alpha} \partial_{i}^{h} u\right\|_{L^{p}(V)}=\left\|\partial_{i}^{h} \partial^{\alpha} u\right\|_{L^{p}(V)} \leq\left\|\partial_{i} \partial^{\alpha} u\right\|_{L^{p}(\Omega)}
$$

wherein we have used Theorem 29.22 for the last inequality. Eq. (48.8) now easily follows.
2. If $\left\|\partial_{i}^{h} u\right\|_{W^{k, p}(V)} \leq C(V)$ then for all $|\alpha| \leq k$,

$$
\left\|\partial_{i}^{h} \partial^{\alpha} u\right\|_{L^{p}(V)}=\left\|\partial^{\alpha} \partial_{i}^{h} u\right\|_{L^{p}(V)} \leq C(V)
$$

So by Theorem 29.22, $\partial_{i} \partial^{\alpha} u \in L^{p}(V)$ and $\left\|\partial_{i} \partial^{\alpha} u\right\|_{L^{p}(V)} \leq C(V)$. From this we conclude that $\left\|\partial^{\beta} u\right\|_{L^{p}(V)} \leq C(V)$ for all $0<|\beta| \leq k+1$ and hence $\|u\|_{W^{k+1, p}(V)} \leq c\left[C(V)+\|u\|_{L^{p}(V)}\right]$ for some constant $c$.
Notation 48.14 Given a multi-index $\alpha$ and $h \neq 0$, let

$$
\partial_{h}^{\alpha}:=\prod_{i=1}^{d}\left(\partial_{i}^{h}\right)^{\alpha_{i}}
$$

The following theorem is a generalization of Theorem 48.13.
Theorem 48.15. Suppose $k \in \mathbb{N}_{0}, \Omega$ is an open subset of $\mathbb{R}^{d}, V$ is an open precompact subset of $\Omega$ and $u \in W^{k, p}(\Omega)$.

1. If $1 \leq p<\infty$ and $|\alpha| \leq k$, then $\left\|\partial_{h}^{\alpha} u\right\|_{W^{k-|\alpha|}(V)} \leq\|u\|_{W^{k, p}(\Omega)}$ for $h$ small.
2. If $1<p \leq \infty$ and $\left\|\partial_{h}^{\alpha} u\right\|_{W^{k, p}(V)} \leq C$ for all $|\alpha| \leq j$ and $h$ near 0 , then $u \in W^{k+j, p}(V)$ and $\left\|\partial^{\alpha} u\right\|_{W^{k, p}(V)} \leq C$ for all $|\alpha| \leq j$.

Proof. Since $\partial_{h}^{\alpha}=\prod_{i} \partial_{h}^{\alpha_{i}}$, item 1. follows from Item 1. of Theorem 48.13 and induction on $|\alpha|$.

For Item 2., suppose first that $k=0$ so that $u \in L^{p}(\Omega)$ and $\left\|\partial_{h}^{\alpha} u\right\|_{L^{p}(V)} \leq$ $C$ for $|\alpha| \leq j$. Then by Proposition 29.16, there exists $\left\{h_{l}\right\}_{l=1}^{\infty} \subset \mathbb{R} \backslash\{0\}$ and $v \in L^{p}(V)$ such that $h_{l} \rightarrow 0$ and $\lim _{l \rightarrow \infty}\left\langle\partial_{h_{e}}^{\alpha} u, \phi\right\rangle=\langle v, \phi\rangle$ for all $\phi \in C_{c}^{\infty}(V)$. Using Remark 48.12,

$$
\langle v, \phi\rangle=\lim _{l \rightarrow \infty}\left\langle\partial_{h_{\ell}}^{\alpha} u, \phi\right\rangle=(-1)^{|\alpha|} \lim _{l \rightarrow \infty}\left\langle u, \partial_{-h_{\ell}}^{\alpha} \phi\right\rangle=(-1)^{|\alpha|}\left\langle u, \partial^{\alpha} \phi\right\rangle
$$

which shows $\partial^{\alpha} u=v \in L^{p}(V)$. Moreover, since weak convergence decreases norms,

$$
\left\|\partial^{\alpha} u\right\|_{L^{p}(V)}=\|v\|_{L^{p}(V)} \leq C .
$$

For the general case if $k \in \mathbb{N}, u \in W^{k, p}(\Omega)$ such that $\left\|\partial_{h}^{\alpha} u\right\|_{W^{k, p}(V)} \leq C$, then (for $p \in(1, \infty)$, the case $p=\infty$ is similar and left to the reader)

$$
\sum_{|\beta| \leq k}\left\|\partial_{h}^{\alpha} \partial^{\beta} u\right\|_{L^{p}(V)}^{p}=\sum_{|\beta| \leq k}\left\|\partial^{\beta} \partial_{h}^{\alpha} u\right\|_{L^{p}(V)}^{p}=\left\|\partial_{h}^{\alpha} u\right\|_{W^{k, p}(V)}^{p} \leq C^{p}
$$

As above this implies $\partial^{\alpha} \partial^{\beta} u \in L^{p}(V)$ for all $|\alpha| \leq j$ and $|\beta| \leq k$ and that

$$
\left\|\partial^{\alpha} u\right\|_{W^{k, p}(V)}^{p}=\sum_{|\beta| \leq k}\left\|\partial^{\alpha} \partial^{\beta} u\right\|_{L^{p}(V)}^{p} \leq C^{p}
$$

### 48.3 Sobolev Spaces on Compact Manifolds

Theorem 48.16 (Change of Variables). Suppose that $U$ and $V$ are open subsets of $\mathbb{R}^{d}, T \in C^{k}(U, V)$ be a $C^{k}$ - diffeomorphism such that $\left\|\partial^{\alpha} T\right\|_{B C(U)}<\infty$ for all $1 \leq|\alpha| \leq k$ and $\epsilon:=\inf _{U}\left|\operatorname{det} T^{\prime}\right|>0$. Then the map $T^{*}: W^{k, p}(V) \rightarrow W^{k, p}(U)$ defined by $u \in W^{k, p}(V) \rightarrow T^{*} u:=u \circ T \in$ $W^{k, p}(U)$ is well defined and is bounded.

Proof. For $u \in W^{k, p}(V) \cap C^{\infty}(V)$, repeated use of the chain and product rule implies,

$$
\begin{aligned}
(u \circ T)^{\prime} & =\left(u^{\prime} \circ T\right) T^{\prime} \\
(u \circ T)^{\prime \prime} & =\left(u^{\prime} \circ T\right)^{\prime} T^{\prime}+\left(u^{\prime} \circ T\right) T^{\prime \prime}=\left(u^{\prime \prime} \circ T\right) T^{\prime} \otimes T^{\prime}+\left(u^{\prime} \circ T\right) T^{\prime \prime} \\
(u \circ T)^{(3)} & =\left(u^{(3)} \circ T\right) T^{\prime} \otimes T^{\prime} \otimes T^{\prime}+\left(u^{\prime \prime} \circ T\right)\left(T^{\prime} \otimes T^{\prime}\right)^{\prime} \\
& +\left(u^{\prime \prime} \circ T\right) T^{\prime} \otimes T^{\prime \prime}+\left(u^{\prime} \circ T\right) T^{(3)}
\end{aligned}
$$

$(u \circ T)^{(l)}=\left(u^{(l)} \circ T\right) \overbrace{T \otimes \cdots \otimes T}^{l \text { times }}+\sum_{j=1}^{l-1}\left(u^{(j)} \circ T\right) p_{j}\left(T^{\prime}, T^{\prime \prime}, \ldots, T^{(l+1-j)}\right)$.

This equation and the boundedness assumptions on $T^{(j)}$ for $1 \leq j \leq k$ implies there is a finite constant $K$ such that

$$
\left|(u \circ T)^{(l)}\right| \leq K \sum_{j=1}^{l}\left|u^{(j)} \circ T\right| \text { for all } 1 \leq l \leq k
$$

By Hölder's inequality for sums we conclude there is a constant $K_{p}$ such that

$$
\sum_{|\alpha| \leq k}\left|\partial^{\alpha}(u \circ T)\right|^{p} \leq K_{p} \sum_{|\alpha| \leq k}\left|\partial^{\alpha} u\right|^{p} \circ T
$$

and therefore

$$
\|u \circ T\|_{W^{k, p}(U)}^{p} \leq K_{p} \sum_{|\alpha| \leq k} \int_{U}\left|\partial^{\alpha} u\right|^{p}(T(x)) d x
$$

Making the change of variables, $y=T(x)$ and using

$$
d y=\left|\operatorname{det} T^{\prime}(x)\right| d x \geq \epsilon d x
$$

we find

$$
\begin{align*}
\|u \circ T\|_{W^{k, p}(U)}^{p} & \leq K_{p} \sum_{|\alpha| \leq k} \int_{U}\left|\partial^{\alpha} u\right|^{p}(T(x)) d x \\
& \leq \frac{K_{p}}{\epsilon} \sum_{|\alpha| \leq k} \int_{V}\left|\partial^{\alpha} u\right|^{p}(y) d y=\frac{K_{p}}{\epsilon}\|u\|_{W^{k, p}(V)}^{p} \tag{48.10}
\end{align*}
$$

This shows that $T^{*}: W^{k, p}(V) \cap C^{\infty}(V) \rightarrow W^{k, p}(U) \cap C^{\infty}(U)$ is a bounded operator. For general $u \in W^{k, p}(V)$, we may choose $u_{n} \in W^{k, p}(V) \cap C^{\infty}(V)$ such that $u_{n} \rightarrow u$ in $W^{k, p}(V)$. Since $T^{*}$ is bounded, it follows that $T^{*} u_{n}$ is Cauchy in $W^{k, p}(U)$ and hence convergent. Finally, using the change of variables theorem again we know,

$$
\left\|T^{*} u-T^{*} u_{n}\right\|_{L^{p}(V)}^{p} \leq \epsilon^{-1}\left\|u-u_{n}\right\|_{L^{p}(U)}^{p} \rightarrow 0 \text { as } n \rightarrow \infty
$$

and therefore $T^{*} u=\lim _{n \rightarrow \infty} T^{*} u_{n}$ and by continuity Eq. (48.10) still holds for $u \in W^{k, p}(V)$.

Let $M$ be a compact $C^{k}$ - manifolds without boundary, i.e. $M$ is a compact Hausdorff space with a collection of charts $x$ in an "atlas" $\mathcal{A}$ such that $x$ : $D(x) \subset_{o} M \rightarrow R(x) \subset_{o} \mathbb{R}^{d}$ is a homeomorphism such that

$$
\left.x \circ y^{-1} \in C^{k}(y(D(x) \cap D(y))), x(D(x) \cap D(y))\right) \text { for all } x, y \in \mathcal{A}
$$

Definition 48.17. Let $\left\{x_{i}\right\}_{i=1}^{N} \subset \mathcal{A}$ such that $M=\cup_{i=1}^{N} D\left(x_{i}\right)$ and let $\left\{\phi_{i}\right\}_{i=1}^{N}$ be a partition of unity subordinate do the cover $\left\{D\left(x_{i}\right)\right\}_{i=1}^{N}$. We now define $u \in W^{k, p}(M)$ if $u: M \rightarrow \mathbb{C}$ is a function such that

$$
\begin{equation*}
\|u\|_{W^{k, p}(M)}:=\sum_{i=1}^{N}\left\|\left(\phi_{i} u\right) \circ x_{i}^{-1}\right\|_{W^{k, p}\left(R\left(x_{i}\right)\right)}<\infty . \tag{48.11}
\end{equation*}
$$

Since $\|\cdot\|_{W^{k, p}\left(R\left(x_{i}\right)\right)}$ is a norm for all $i$, it easily verified that $\|\cdot\|_{W^{k, p}(M)}$ is a norm on $W^{k, p}(M)$.

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Proposition 48.18. If $f \in C^{k}(M)$ and $u \in W^{k, p}(M)$ then $f u \in W^{k, p}(M)$ and

$$
\begin{equation*}
\|f u\|_{W^{k, p}(M)} \leq C\|u\|_{W^{k, p}(M)} \tag{48.12}
\end{equation*}
$$

where $C$ is a finite constant not depending on $u$. Recall that $f: M \rightarrow \mathbb{R}$ is said to be $C^{j}$ with $j \leq k$ if $f \circ x^{-1} \in C^{j}(R(x), \mathbb{R})$ for all $x \in \mathcal{A}$.

Proof. Since $\left[f \circ x_{i}^{-1}\right]$ has bounded derivatives on $\operatorname{supp}\left(\phi_{i} \circ x_{i}^{-1}\right)$, it follows from Proposition 48.6 that there is a constant $C_{i}<\infty$ such that

$$
\begin{aligned}
\left\|\left(\phi_{i} f u\right) \circ x_{i}^{-1}\right\|_{W^{k, p}\left(R\left(x_{i}\right)\right)} & =\left\|\left[f \circ x_{i}^{-1}\right]\left(\phi_{i} u\right) \circ x_{i}^{-1}\right\|_{W^{k, p}\left(R\left(x_{i}\right)\right)} \\
& \leq C_{i}\left\|\left(\phi_{i} u\right) \circ x_{i}^{-1}\right\|_{W^{k, p}\left(R\left(x_{i}\right)\right)}
\end{aligned}
$$

and summing this equation on $i$ shows Eq. (48.12) holds with $C:=\max _{i} C_{i}$.

Theorem 48.19. If $\left\{y_{j}\right\}_{j=1}^{K} \subset \mathcal{A}$ such that $M=\cup_{j=1}^{K} D\left(y_{j}\right)$ and $\left\{\psi_{j}\right\}_{j=1}^{K}$ is a partition of unity subordinate to the cover $\left\{D\left(y_{j}\right)\right\}_{j=1}^{K}$, then the norm

$$
\begin{equation*}
|u|_{W^{k, p}(M)}:=\sum_{j=1}^{K}\left\|\left(\psi_{j} u\right) \circ y_{j}^{-1}\right\|_{W^{k, p}\left(R\left(y_{j}\right)\right)} \tag{48.13}
\end{equation*}
$$

is equivalent to the norm in Eq. (48.11). That is to say the space $W^{k, p}(M)$ along with its topology is well defined independent of the choice of charts and partitions of unity used in defining the norm on $W^{k, p}(M)$.

Proof. Since $|\cdot|_{W^{k, p}(M)}$ is a norm,

$$
\begin{align*}
|u|_{W^{k, p}(M)} & =\left|\sum_{i=1}^{N} \phi_{i} u\right|_{W^{k, p}(M)} \leq \sum_{i=1}^{N}\left|\phi_{i} u\right|_{W^{k, p}(M)} \\
& =\sum_{j=1}^{K}\left\|\sum_{i=1}^{N}\left(\psi_{j} \phi_{i} u\right) \circ y_{j}^{-1}\right\|_{W^{k, p}\left(R\left(y_{j}\right)\right)} \\
& \leq \sum_{j=1}^{K} \sum_{i=1}^{N}\left\|\left(\psi_{j} \phi_{i} u\right) \circ y_{j}^{-1}\right\|_{W^{k, p}\left(R\left(y_{j}\right)\right)} \tag{48.14}
\end{align*}
$$

and since $x_{i} \circ y_{j}^{-1}$ and $y_{j} \circ x_{i}^{-1}$ are $C^{k}$ diffeomorphism and the sets $y_{j}\left(\operatorname{supp}\left(\phi_{i}\right) \cap \operatorname{supp}\left(\psi_{j}\right)\right)$ and $x_{i}\left(\operatorname{supp}\left(\phi_{i}\right) \cap \operatorname{supp}\left(\psi_{j}\right)\right)$ are compact, an application of Theorem 48.16 and Proposition 48.6 shows there are finite constants $C_{i j}$ such that

$$
\begin{aligned}
\left\|\left(\psi_{j} \phi_{i} u\right) \circ y_{j}^{-1}\right\|_{W^{k, p}\left(R\left(y_{j}\right)\right)} & \leq C_{i j}\left\|\left(\psi_{j} \phi_{i} u\right) \circ x_{i}^{-1}\right\|_{W^{k, p}\left(R\left(x_{i}\right)\right)} \\
& \leq C_{i j}\left\|\phi_{i} u \circ x_{i}^{-1}\right\|_{W^{k, p}\left(R\left(x_{i}\right)\right)}
\end{aligned}
$$

which combined with Eq. (48.14) implies

$$
|u|_{W^{k, p}(M)} \leq \sum_{j=1}^{K} \sum_{i=1}^{N} C_{i j}\left\|\phi_{i} u \circ x_{i}^{-1}\right\|_{W^{k, p}\left(R\left(x_{i}\right)\right)} \leq C\|u\|_{W^{k, p}(M)}
$$

where $C:=\max _{i} \sum_{j=1}^{K} C_{i j}<\infty$. Analogously, one shows there is a constant $K<\infty$ such that $\|u\|_{W^{k, p}(M)} \leq K|u|_{W^{k, p}(M)}$.

Lemma 48.20. Suppose $x \in \mathcal{A}(M)$ and $U \subset o M$ such that $U \subset \bar{U} \subset D(x)$, then there is a constant $C<\infty$ such that

$$
\begin{equation*}
\left\|u \circ x^{-1}\right\|_{W^{k, p}(x(U))} \leq C\|u\|_{W^{k, p}(M)} \text { for all } u \in W^{k, p}(M) \tag{48.15}
\end{equation*}
$$

Conversely a function $u: M \rightarrow \mathbb{C}$ with $\operatorname{supp}(u) \subset U$ is in $W^{k, p}(M)$ iff $\left\|u \circ x^{-1}\right\|_{W^{k, p}(x(U))}<\infty$ and in any case there is a finite constant such that

$$
\begin{equation*}
\|u\|_{W^{k, p}(M)} \leq C\left\|u \circ x^{-1}\right\|_{W^{k, p}(x(U))} \tag{48.16}
\end{equation*}
$$

Proof. Choose charts $y_{1}:=x, y_{2}, \ldots, y_{K} \in \mathcal{A}$ such that $\left\{D\left(y_{i}\right)\right\}_{j=1}^{K}$ is an open cover of $M$ and choose a partition of unity $\left\{\psi_{j}\right\}_{j=1}^{K}$ subordinate to the cover $\left\{D\left(y_{j}\right)\right\}_{j=1}^{K}$ such that $\psi_{1}=1$ on a neighborhood of $\bar{U}$. To construct such a partition of unity choose $U_{j} \subset_{o} M$ such that $U_{j} \subset \bar{U}_{j} \subset D\left(y_{j}\right), \bar{U} \subset U_{1}$ and $\cup_{j=1}^{K} U_{j}=M$ and for each $j$ let $\eta_{j} \in C_{c}^{k}\left(D\left(y_{j}\right),[0,1]\right)$ such that $\eta_{j}=1$ on a neighborhood of $\bar{U}_{j}$. Then define $\psi_{j}:=\eta_{j}\left(1-\eta_{0}\right) \cdots\left(1-\eta_{j-1}\right)$ where by convention $\eta_{0} \equiv 0$. Then $\left\{\psi_{j}\right\}_{j=1}^{K}$ is the desired partition, indeed by induction one shows

$$
1-\sum_{j=1}^{l} \psi_{j}=\left(1-\eta_{1}\right) \cdots\left(1-\eta_{l}\right)
$$

and in particular

$$
1-\sum_{j=1}^{K} \psi_{j}=\left(1-\eta_{1}\right) \cdots\left(1-\eta_{K}\right)=0
$$

Using Theorem 48.19, it follows that

$$
\begin{aligned}
\left\|u \circ x^{-1}\right\|_{W^{k, p}(x(U))} & =\left\|\left(\psi_{1} u\right) \circ x^{-1}\right\|_{W^{k, p}(x(U))} \\
& \leq\left\|\left(\psi_{1} u\right) \circ x^{-1}\right\|_{W^{k, p}\left(R\left(y_{1}\right)\right)} \\
& \leq \sum_{j=1}^{K}\left\|\left(\psi_{j} u\right) \circ y_{j}^{-1}\right\|_{W^{k, p}\left(R\left(y_{j}\right)\right)} \\
& =|u|_{W^{k, p}(M)} \leq C\|u\|_{W^{k, p}(M)}
\end{aligned}
$$

which proves Eq. (48.15).
Using Theorems 48.19 and 48.16 there are constants $C_{j}$ for $j=$ $0,1,2 \ldots, N$ such that

$$
\begin{aligned}
&\|u\|_{W^{k, p}(M)} \leq C_{0} \sum_{j=1}^{K}\left\|\left(\psi_{j} u\right) \circ y_{j}^{-1}\right\|_{W^{k, p}\left(R\left(y_{j}\right)\right)} \\
&=C_{0} \sum_{j=1}^{K}\left\|\left(\psi_{j} u\right) \circ y_{1}^{-1} \circ y_{1} \circ y_{j}^{-1}\right\|_{W^{k, p}\left(R\left(y_{j}\right)\right)} \\
& \leq C_{0} \sum_{j=1}^{K} C_{j}\left\|\left(\psi_{j} u\right) \circ x^{-1}\right\|_{W^{k, p}\left(R\left(y_{1}\right)\right)} \\
&=C_{0} \sum_{j=1}^{K} C_{j}\left\|\psi_{j} \circ x^{-1} \cdot u \circ x^{-1}\right\|_{W^{k, p}\left(R\left(y_{1}\right)\right)}
\end{aligned}
$$

This inequality along with $K$ - applications of Proposition 48.6 proves Eq. (48.16).

Theorem 48.21. The space $\left(W^{k, p}(M),\|\cdot\|_{W^{k, p}(M)}\right)$ is a Banach space.
Proof. Let $\left\{x_{i}\right\}_{i=1}^{N} \subset \mathcal{A}$ and $\left\{\phi_{i}\right\}_{i=1}^{N}$ be as in Definition 48.17 and choose $U_{i} \subset_{o} M$ such that $\operatorname{supp}\left(\phi_{i}\right) \subset U_{i} \subset \bar{U}_{i} \subset D\left(x_{i}\right)$. If $\left\{u_{n}\right\}_{n=1}^{\infty} \subset W^{k, p}(M)$ is a Cauchy sequence, then by Lemma 48.20, $\left\{u_{n} \circ x_{i}^{-1}\right\}_{n=1}^{\infty} \subset W^{k, p}\left(x_{i}\left(U_{i}\right)\right)$ is a Cauchy sequence for all $i$. Since $W^{k, p}\left(x_{i}\left(U_{i}\right)\right)$ is complete, there exists $v_{i} \in W^{k, p}\left(x_{i}\left(U_{i}\right)\right)$ such that $u_{n} \circ x_{i}^{-1} \rightarrow \tilde{v}_{i}$ in $W^{k, p}\left(x_{i}\left(U_{i}\right)\right)$. For each $i$ let $v_{i}:=\phi_{i}\left(\tilde{v}_{i} \circ x_{i}\right)$ and notice by Lemma 48.20 that

$$
\left\|v_{i}\right\|_{W^{k, p}(M)} \leq C\left\|v_{i} \circ x_{i}^{-1}\right\|_{W^{k, p}\left(x_{i}\left(U_{i}\right)\right)}=C\left\|\tilde{v}_{i}\right\|_{W^{k, p}\left(x_{i}\left(U_{i}\right)\right)}<\infty
$$

so that $u:=\sum_{i=1}^{N} v_{i} \in W^{k, p}(M)$. Since $\operatorname{supp}\left(v_{i}-\phi_{i} u_{n}\right) \subset U_{i}$, it follows that

$$
\begin{aligned}
\left\|u-u_{n}\right\|_{W^{k, p}(M)} & =\left\|\sum_{i=1}^{N} v_{i}-\sum_{i=1}^{N} \phi_{i} u_{n}\right\|_{W^{k, p}(M)} \\
& \leq \sum_{i=1}^{N}\left\|v_{i}-\phi_{i} u_{n}\right\|_{W^{k, p}(M)} \\
& \leq C \sum_{i=1}^{N}\left\|\left[\phi_{i}\left(\tilde{v}_{i} \circ x_{i}-u_{n}\right)\right] \circ x_{i}^{-1}\right\|_{W^{k, p}\left(x_{i}\left(U_{i}\right)\right)} \\
& =C \sum_{i=1}^{N}\left\|\left[\phi_{i} \circ x_{i}^{-1}\left(\tilde{v}_{i}-u_{n} \circ x_{i}^{-1}\right)\right]\right\|_{W^{k, p}\left(x_{i}\left(U_{i}\right)\right)} \\
& \leq C \sum_{i=1}^{N} C_{i}\left\|\tilde{v}_{i}-u_{n} \circ x_{i}^{-1}\right\|_{W^{k, p}\left(x_{i}\left(U_{i}\right)\right)} \rightarrow 0 \text { as } n \rightarrow \infty
\end{aligned}
$$

wherein the last inequality we have used Proposition 48.6 again.

### 48.4 Trace Theorems

For many more general results on this subject matter, see E. Stein [17, Chapter VI].
Notation 48.22 Let $\mathbb{H}^{d}:=\left\{x \in \mathbb{R}^{d}: x_{d}>0\right\}$ be the open upper half space inside of $\mathbb{R}^{d}$ and if $D>0$ let

$$
\mathbb{H}_{D}^{d}:=\left\{x \in \mathbb{H}^{d}: 0<x_{d}<D\right\} .
$$

Lemma 48.23. Suppose $k \geq 1$ and $D>0$.

1. If $p \in[1, \infty)$ and $C_{c}^{k}\left(\overline{\mathbb{H}^{d}}\right)$, then for all $\alpha \in \mathbb{N}_{0}^{d-1} \times\{0\} \subset \mathbb{N}_{0}^{d}$ with $|\alpha| \leq$ $k-1$,

$$
\begin{equation*}
\left\|\partial^{\alpha} u\right\|_{L^{p}\left(\partial \mathbb{H}^{d}\right)} \leq D^{-1 / p}\left\|\partial^{\alpha} u\right\|_{L^{p}\left(\mathbb{H}_{D}^{d}\right)}+\frac{D^{\frac{p-1}{p}}}{p^{1 / p}}\left\|\partial^{\alpha+e_{d}} u\right\|_{L^{p}\left(\mathbb{H}_{D}^{d}\right)} . \tag{48.17}
\end{equation*}
$$

In particular there there is a constant $C=C(p, k, D, d)$ such that

$$
\begin{equation*}
\|u\|_{W^{k-1, p}\left(\partial \mathbb{H}^{d}\right)} \leq C(p, D, k, d)\|u\|_{W^{k, p}\left(\mathbb{H}^{d}\right)} \tag{48.18}
\end{equation*}
$$

2. For $p=\infty$ and $u \in W^{k, \infty}\left(\mathbb{H}^{d}\right)$, there is a continuous version $\tilde{u}$ of $u$. The function $\tilde{u}$ is in $B C^{k-1}\left(\mathbb{H}^{d}\right)$ and has the property that $\partial^{\alpha} \tilde{u}$ extends to a continuous function $v_{\alpha} \in B C$ ( $\left.\overline{\mathbb{H}}\right)$ for all $|\alpha|<k$ and the function $\left.\tilde{u}\right|_{\partial \mathbb{H}^{d}}$ is in $B C^{k-1}(\partial \bar{H})$ and

$$
\|\tilde{u}\|_{B C^{k-1}\left(\partial \mathbb{H}^{d}\right)} \leq\|u\|_{W^{k, \infty}\left(\mathbb{H}_{D}^{d}\right)} \text { for any } D>0 \text {. }
$$

Proof. 1. Write $x \in \overline{\mathbb{H}^{d}}$ as $x=(y, z) \in \mathbb{R}^{d-1} \times[0, \infty)$ and suppose $\alpha \in \mathbb{N}_{0}^{d-1} \times\{0\} \subset \mathbb{N}_{0}^{d}$ with $|\alpha| \leq k-1$. Then by the fundamental theorem of calculus,

$$
\begin{equation*}
\partial^{\alpha} u(y, 0)=\partial^{\alpha} u(y, z)-\int_{0}^{z} \partial^{\alpha} u_{t}(y, t) d t \tag{48.19}
\end{equation*}
$$

which implies

$$
\left|\partial^{\alpha} u(y, 0)\right| 1_{[0, D]}(z) \leq\left|\partial^{\alpha} u(y, z)\right| 1_{[0, D]}(z)+1_{[0, D]}(z) \int_{0}^{z}\left|\partial^{\alpha} u_{t}(y, t)\right| d t
$$

Taking the $L^{p}\left(\mathbb{H}_{D}^{d}\right)$ - norm of this last equation implies

$$
\left\|\partial^{\alpha} u\right\|_{L^{p}\left(\partial \mathbb{H}_{D}^{d}\right)} \cdot D^{1 / p} \leq\left\|\partial^{\alpha} u\right\|_{L^{p}\left(\mathbb{H}_{D}^{d}\right)}+B
$$

where

$$
\begin{aligned}
B^{p} & =\int_{\mathbb{H}_{D}^{d}}\left(1_{[0, D]}(z) \int_{0}^{z}\left|\partial^{\alpha} u_{t}(y, t)\right| d t\right)^{p} d y d z \\
& \leq \int_{\mathbb{H}_{D}^{d}}\left(1_{[0, D]}(z) \cdot z^{p / q} \int_{0}^{z}\left|\partial^{\alpha} u_{t}(y, t)\right|^{p} d t\right) d y d z \\
& \leq \int_{\mathbb{H}_{D}^{d}}\left(1_{[0, D]}(z) \cdot z^{p / q} \int_{0}^{D}\left|\partial^{\alpha} u_{t}(y, t)\right|^{p} d t\right) d y d z \\
& =\left\|\partial^{\alpha+e_{d}} u\right\|_{L^{p}\left(\mathbb{H}_{D}^{d}\right)}^{p} \frac{D^{p / q+1}}{p / q+1}=\frac{D^{p}}{p}\left\|\partial^{\alpha+e_{d}} u\right\|_{L^{p}\left(\mathbb{H}_{D}^{d}\right)}^{p} .
\end{aligned}
$$

Putting these two equations together shows

$$
\left\|\partial^{\alpha} u\right\|_{L^{p}\left(\partial \mathbb{H}_{D}^{d}\right)} \leq D^{-1 / p}\left[\left\|\partial^{\alpha} u\right\|_{L^{p}\left(\mathbb{H}_{D}^{d}\right)}+\frac{D}{p^{1 / p}}\left\|\partial^{\alpha+e_{d}} u\right\|_{L^{p}\left(\mathbb{H}_{D}^{d}\right)}\right]
$$

which is the same as Eq. (48.17).
Suppose that $p=\infty$ and $u \in W^{k, \infty}\left(\mathbb{H}^{d}\right)$. By Proposition 29.29, we know that $\partial^{\alpha} u$ has a Lipschitz continuous version $v_{\alpha}$ on $\mathbb{H}^{d}$ for each $|\alpha|<k$. Being Lipschitz, each $v_{\alpha}$ has a unique extension to a continuous function $\overline{\mathbb{H}^{d}}$. Let

$$
\eta \in C_{c}^{\infty}\left(B(0,1) \cap\left(-\mathbb{H}^{d}\right),[0, \infty)\right)
$$

be chosen so that $\int_{\mathbb{R}^{n}} \eta(x) d x=1$ and $\eta_{m}(x)=m^{n} \eta(m x)$ and let $u_{m}:=$ $u * \eta_{m}=v_{o} * \eta_{m}$. Since $\operatorname{supp}(\eta) \subset\left(-\mathbb{H}^{d}\right)$ as in Figure 48.5, $u_{m} \in C^{\infty}\left(\overline{\mathbb{H}^{d}}\right)$,


Fig. 48.5. The support of $\eta$.
$\partial^{\alpha} u_{m}=\partial^{\alpha} u * \eta_{m}=v_{\alpha} * \eta_{m}$ for all $|\alpha|<k$ and

$$
\left\|u_{m}-u_{n}\right\|_{B C^{k-1}\left(\overline{\mathbb{H}^{d}} \cap B(0, R)\right)} \rightarrow 0 \text { as } m, n \rightarrow \infty
$$

Therefore $\partial^{\alpha} u_{m} \rightarrow v_{\alpha}$ uniformly on compact subsets of $\overline{\mathbb{H}^{d}}$ for all $|\alpha|<k$. Hence $v_{0} \in C^{k-1}\left(\mathbb{H}^{d}\right)$ and $\partial^{\alpha} v_{0}=v_{\alpha}$ extends to $\overline{\mathbb{H}^{d}}$ for all $|\alpha|<k$ and $\left.v_{0}\right|_{\partial \mathbb{H}^{d}} \in B C^{k-1}(\partial \overline{\mathcal{H}})$.

Theorem 48.24 (Trace Theorem). Suppose $k \geq 1$ and $\Omega \subset_{o} \mathbb{R}^{d}$ such that $\bar{\Omega}$ is a compact manifold with $C^{k}$ - boundary. Then there exists a unique linear map $T: W^{k, p}(\Omega) \rightarrow W^{k-1, p}(\partial \Omega)$ such that $T u=\left.u\right|_{\partial \Omega}$ for all $u \in C^{k}(\bar{\Omega})$.

Proof. Choose a covering $\left\{V_{i}\right\}_{i=0}^{N}$ of $\bar{\Omega}$ such that $\bar{V}_{0} \subset \Omega$ and for each $i \geq 1$, there is $C^{k}$ - diffeomorphism $x_{i}: V_{i} \rightarrow R\left(x_{i}\right) \subset_{o} \mathbb{R}^{d}$ such that

$$
\begin{aligned}
x_{i}\left(\partial \Omega \cap V_{i}\right) & =R\left(x_{i}\right) \cap \mathrm{bd}\left(\mathbb{H}^{d}\right) \text { and } \\
x_{i}\left(\Omega \cap V_{i}\right) & =R\left(x_{i}\right) \cap \mathbb{H}^{d}
\end{aligned}
$$

as in Figure 48.6. Further choose $\phi_{i} \in C_{c}^{\infty}\left(V_{i},[0,1]\right)$ such that $\sum_{i=0}^{N} \phi_{i}=1$


Fig. 48.6. Covering $\Omega$ (the shaded region) as described in the text.
on a neighborhood of $\bar{\Omega}$ and set $y_{i}:=\left.x_{i}\right|_{\partial \Omega \cap V_{i}}$ for $i \geq 1$. Given $u \in C^{k}(\bar{\Omega})$ if $p<\infty$ and $u \in W^{k, \infty}(\Omega)$ if $p=\infty$, we compute

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$$
\begin{aligned}
\left\|\left.u\right|_{\partial \bar{\Omega}}\right\|_{W^{k-1, p}(\partial \bar{\Omega})}= & \sum_{i=1}^{N}\left\|\left.\left(\phi_{i} u\right)\right|_{\partial \bar{\Omega}} \circ y_{i}^{-1}\right\|_{W^{k-1, p}\left(R\left(x_{i}\right) \cap \mathrm{hd}\left(\mathbb{H}^{d}\right)\right)} \\
= & \sum_{i=1}^{N}\left\|\left.\left[\left(\phi_{i} u\right) \circ x_{i}^{-1}\right]\right|_{\mathrm{bd}\left(\mathbb{H}^{d}\right)}\right\|_{W^{k-1, p}\left(R\left(x_{i}\right) \cap \mathrm{Bd}\left(\mathbb{H}^{d}\right)\right)} \\
\leq & \sum_{i=1}^{N} C_{i}\left\|\left[\left(\phi_{i} u\right) \circ x_{i}^{-1}\right]\right\|_{W^{k, p}\left(R\left(x_{i}\right)\right)} \\
\leq & \max C_{i} \cdot \sum_{i=1}^{N}\left\|\left[\left(\phi_{i} u\right) \circ x_{i}^{-1}\right]\right\|_{W^{k, p}\left(R\left(x_{i}\right) \cap \mathbb{H}^{d}\right)} \\
& +\left\|\left[\left(\phi_{0} u\right) \circ x_{0}^{-1}\right]\right\|_{W^{k, p}\left(R\left(x_{0}\right)\right)} \\
\leq & C\|u\|_{W^{k, p}(\Omega)}
\end{aligned}
$$

where $C=\max \left\{1, C_{1}, \ldots, C_{N}\right\}$. The proof is complete if $p=\infty$ and follows by the B.L.T. Theorem 2.68 and the fact that $C^{k}(\bar{\Omega})$ is dense inside $W^{k, p}(\Omega)$ if $p<\infty$.
Notation 48.25 In the sequel will often abuse notation and simply write $\left.u\right|_{\partial \bar{\Omega}}$ for the "function" $T u \in W^{k-1, p}(\partial \bar{\Omega})$.

Proposition 48.26 (Integration by parts). Suppose $\Omega \subset_{o} \mathbb{R}^{d}$ such that $\bar{\Omega}$ is a compact manifold with $C^{1}$ - boundary, $p \in[1, \infty]$ and $q=\frac{p}{p-1}$ is the conjugate exponent. Then for $u \in W^{1, p}(\Omega)$ and $v \in W^{1, q}(\Omega)$,

$$
\begin{equation*}
\int_{\Omega} \partial_{i} u \cdot v d m=-\int_{\Omega} u \cdot \partial_{i} v d m+\left.\left.\int_{\partial \bar{\Omega}} u\right|_{\partial \bar{\Omega}} \cdot v\right|_{\partial \bar{\Omega}} n_{i} d \sigma \tag{48.20}
\end{equation*}
$$

where $n: \partial \bar{\Omega} \rightarrow \mathbb{R}^{d}$ is unit outward pointing norm to $\partial \bar{\Omega}$.
Proof. Equation 48.20 holds for $u, v \in C^{2}(\bar{\Omega})$ and therefore for $(u, v) \in$ $W^{k, p}(\Omega) \times W^{k, q}(\Omega)$ since both sides of the equality are continuous in $(u, v) \in$ $W^{k, p}(\Omega) \times W^{k, q}(\Omega)$ as the reader should verify.

Warning BRUCE: We might need $p \in(1, \infty)$ here. To fix this, I think if $p=1$ one should replace $u$ by $u_{M}:=\psi_{M}(u)$ where $\psi_{M}(x)=\int_{0}^{x} \alpha(y / m) d y$ and $\alpha \in C_{c}\left(\mathbb{R},[0,1)\right.$ such that $\alpha=1$ on $[-1,1]$. Then $u_{M} \in W^{1, \infty}(\Omega)$ and $u_{M} \rightarrow u$ in $W^{1,1}(\Omega)$ and hence the argument given above goes through. We need only approximate $v \in W^{1, q}(\Omega)$ with $q<\infty$ now. We should then pass to the limit as $M \rightarrow \infty$.
Definition 48.27. Let $W_{0}^{k, p}(\Omega):={\overline{C_{c}^{\infty}(\Omega)}}^{W^{k, p}(\Omega)}$ be the closure of $C_{c}^{\infty}(\Omega)$ inside $W^{k, p}(\Omega)$.

Remark 48.28. Notice that if $T: W^{k, p}(\Omega) \rightarrow W^{k-1, p}(\partial \bar{\Omega})$ is the trace operator in Theorem 48.24, then $T\left(W_{0}^{k, p}(\Omega)\right)=\{0\} \subset W^{k-1, p}(\partial \bar{\Omega})$ since $T u=\left.u\right|_{\partial \bar{\Omega}}=0$ for all $u \in C_{c}^{\infty}(\Omega)$.

Corollary 48.29. Suppose $\Omega \subset_{o} \mathbb{R}^{d}$ such that $\bar{\Omega}$ is a compact manifold with $C^{1}$ - boundary, $p \in[1, \infty)$ and $T: W^{1, p}(\Omega) \rightarrow L^{p}(\partial \Omega)$ is the trace operator of Theorem 48.24. Then $W_{0}^{1, p}(\Omega)=\operatorname{Nul}(T)$.

Proof. It has already been observed in Remark 48.28 that $W_{0}^{1, p}(\Omega) \subset$ $\operatorname{Nul}(T)$. Suppose $u \in \operatorname{Nul}(T)$ and $\operatorname{supp}(u)$ is compactly contained in $\Omega$. The mollification $u_{\epsilon}(x)$ defined in Proposition 48.4 will be in $C_{c}^{\infty}(\Omega)$ for $\epsilon>0$ sufficiently small and by Proposition 48.4, $u_{\epsilon} \rightarrow u$ in $W^{1, p}(\Omega)$. Thus $u \in$ $W_{0}^{1, p}(\Omega)$. So to finish the proof that $\operatorname{Nul}(T) \subset W_{0}^{1, p}(\Omega)$, it suffices to show every $u \in W_{0}^{1, p}(\Omega)$ may be approximated by $v \in W_{0}^{1, p}(\Omega)$ such that $\operatorname{supp}(v)$ is compactly contained in $\Omega$. Two proofs of this last assertion will now be given.

Proof 1. For $u \in \operatorname{Nul}(T) \subset W^{1, p}(\Omega)$ define

$$
\tilde{u}(x)=\left\{\begin{array}{cc}
u(x) & \text { for } x \in \bar{\Omega} \\
0 & \text { for } x \notin \bar{\Omega} .
\end{array}\right.
$$

Then clearly $\tilde{u} \in L^{p}\left(\mathbb{R}^{d}\right)$ and moreover by Proposition 48.26 , for $v \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$,

$$
\int_{\mathbb{R}^{d}} \tilde{u} \cdot \partial_{i} v d m=\int_{\Omega} u \cdot \partial_{i} v d m=-\int_{\Omega} \partial_{i} u \cdot v d m
$$

from which it follows that $\partial_{i} \tilde{u}$ exists weakly in $L^{p}\left(\mathbb{R}^{d}\right)$ and $\partial_{i} \tilde{u}=1_{\Omega} \partial_{i} u$ a.e.. Thus $\tilde{u} \in W^{1, p}\left(\mathbb{R}^{d}\right)$ with $\|\tilde{u}\|_{W^{1, p}\left(\mathbb{R}^{d}\right)}=\|u\|_{W^{1, p}(\Omega)}$ and $\operatorname{supp}(\tilde{u}) \subset \bar{\Omega}$. (The reader should compare this result with Proposition 48.30 below.)

Choose $V \in C_{c}^{1}\left(\mathbb{R}^{d}, \mathbb{R}^{d}\right)$ such that $V(x) \cdot n(x)>0$ for all $x \in \partial \bar{\Omega}$ and define

$$
\tilde{u}_{\epsilon}(x)=T_{\epsilon} \tilde{u}(x):=\tilde{u} \circ e^{\epsilon V}(x) .
$$

Notice that $\operatorname{supp}\left(\tilde{u}_{\epsilon}\right) \subset e^{-\epsilon V}(\bar{\Omega}) \sqsubset \sqsubset \Omega$ for all $\epsilon$ sufficiently small. By the change of variables Theorem 48.16, we know that $\tilde{u}_{\epsilon} \in W^{1, p}(\Omega)$ and since $\operatorname{supp}\left(\tilde{u}_{\epsilon}\right)$ is a compact subset of $\Omega$, it follows from the first paragraph that $\tilde{u}_{\epsilon} \in W_{0}^{1, p}(\Omega)$.

To so finish this proof, it only remains to show $\tilde{u}_{\epsilon} \rightarrow u$ in $W^{1, p}(\Omega)$ as $\epsilon \downarrow 0$. Looking at the proof of Theorem 48.16, the reader may show there are constants $\delta>0$ and $C<\infty$ such that

$$
\begin{equation*}
\left\|T_{\epsilon} v\right\|_{W^{1, p}\left(\mathbb{R}^{d}\right)} \leq C\|v\|_{W^{1, p}\left(\mathbb{R}^{d}\right)} \text { for all } v \in W^{1, p}\left(\mathbb{R}^{d}\right) \tag{48.21}
\end{equation*}
$$

By direct computation along with the dominated convergence it may be shown that

$$
\begin{equation*}
T_{\epsilon} v \rightarrow v \text { in } W^{1, p}\left(\mathbb{R}^{d}\right) \text { for all } v \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right) \tag{48.22}
\end{equation*}
$$

As is now standard, Eqs. (48.21) and (48.22) along with the density of $C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ in $W^{1, p}\left(\mathbb{R}^{d}\right)$ allows us to conclude $T_{\epsilon} v \rightarrow v$ in $W^{1, p}\left(\mathbb{R}^{d}\right)$ for all $v \in W^{1, p}\left(\mathbb{R}^{d}\right)$ which completes the proof that $\tilde{u}_{\epsilon} \rightarrow u$ in $W^{1, p}(\Omega)$ as $\epsilon \rightarrow 0$.

Proof 2. As in the first proof it suffices to show that any $u \in W_{0}^{1, p}(\Omega)$ may be approximated by $v \in W^{1, p}(\Omega)$ with $\operatorname{supp}(v) \sqsubset \Omega$. As above extend $u$

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to $\Omega^{c}$ by 0 so that $\tilde{u} \in W^{1, p}\left(\mathbb{R}^{d}\right)$. Using the notation in the proof of 48.24 , it suffices to show $u_{i}:=\phi_{i} \tilde{u} \in W^{1, p}\left(\mathbb{R}^{d}\right)$ may be approximated by $u_{i} \in W^{1, p}(\Omega)$ with $\operatorname{supp}\left(u_{i}\right) \sqsubset \Omega$. Using the change of variables Theorem 48.16, the problem may be reduced to working with $w_{i}=u_{i} \circ x_{i}^{-1}$ on $B=R\left(x_{i}\right)$. But in this case we need only define $w_{i}^{\epsilon}(y):=w_{i}^{\epsilon}\left(y-\epsilon e_{d}\right)$ for $\epsilon>0$ sufficiently small. Then $\operatorname{supp}\left(w_{i}^{\epsilon}\right) \subset \mathbb{H}^{d} \cap B$ and as we have already seen $w_{i}^{\epsilon} \rightarrow w_{i}$ in $W^{1, p}\left(\mathbb{H}^{d}\right)$. Thus $u_{i}^{\epsilon}:=w_{i}^{\epsilon} \circ x_{i} \in W^{1, p}(\Omega), u_{i}^{\epsilon} \rightarrow u_{i}$ as $\epsilon \downarrow 0$ with $\operatorname{supp}\left(u_{i}\right) \sqsubset \Omega$.

### 48.5 Extension Theorems

Proposition 48.30. Let $k \in \mathbb{N}_{0}, p \in[1, \infty]$ and suppose $\Omega$ is any open subset of $\mathbb{R}^{d}$. Then the extension by zero map,

$$
u \in W_{0}^{k, p}(\Omega) \rightarrow 1_{\Omega} u \in W^{k, p}\left(\mathbb{R}^{d}\right)
$$

is a contraction. Recall $W_{0}^{k, p}(\Omega)$ was defined in Definition 48.27) above.
Proof. The result holds for $u \in C_{c}^{\infty}(\Omega)$ and hence for all $u \in W_{0}^{k, p}(\Omega)$.
Lemma 48.31. Let $R>0, B:=B(0, R) \subset \mathbb{R}^{d}, B^{ \pm}:=\left\{x \in B: \pm x_{d}>0\right\}$ and $\Gamma:=\left\{x \in B: x_{d}=0\right\}$. Suppose that $u \in C^{k}(B \backslash \Gamma) \cap C(B)$ and for each $|\alpha| \leq k, \partial^{\alpha} u$ extends to a continuous function $v_{\alpha}$ on $B$. Then $u \in C^{k}(B)$ and $\partial^{\alpha} u=v_{\alpha}$ for all $|\alpha| \leq k$.

Proof. For $x \in \Gamma$ and $i<d$, then by continuity, the fundamental theorem of calculus and the dominated convergence theorem,

$$
\begin{aligned}
u\left(x+\Delta e_{i}\right)-u(x) & =\lim _{\substack{y \rightarrow x \\
y \in B \backslash \Gamma}}\left[u\left(y+\Delta e_{i}\right)-u(y)\right]=\lim _{\substack{y \rightarrow x \\
y \in B \backslash \Gamma}} \int_{0}^{\Delta} \partial_{i} u\left(y+s e_{i}\right) d s \\
& =\lim _{\substack{y \rightarrow x \\
y \in B \backslash \Gamma}} \int_{0}^{\Delta} v_{e_{i}}\left(y+s e_{i}\right) d s=\int_{0}^{\Delta} v_{e_{i}}\left(x+s e_{i}\right) d s
\end{aligned}
$$

and similarly, for $i=d$,

$$
\begin{aligned}
u\left(x+\Delta e_{d}\right)-u(x) & =\lim _{\substack{y \rightarrow x \\
y \in B^{\operatorname{sgn}(\Delta) \backslash \Gamma}}}\left[u\left(y+\Delta e_{d}\right)-u(y)\right] \\
& =\lim _{\substack{y \rightarrow x \\
y \in B^{\operatorname{sgn}(\Delta) \backslash \Gamma}}} \int_{0}^{\Delta} \partial_{d} u\left(y+s e_{d}\right) d s \\
& =\lim _{\substack{y \rightarrow x \\
y \in B^{\operatorname{sgn}(\Delta) \backslash \Gamma}}} \int_{0}^{\Delta} v_{e_{d}}\left(y+s e_{d}\right) d s=\int_{0}^{\Delta} v_{e_{d}}\left(x+s e_{d}\right) d s
\end{aligned}
$$

These two equations show, for each $i, \partial_{i} u(x)$ exits and $\partial_{i} u(x)=v_{e_{i}}(x)$. Hence we have shown $u \in C^{1}(B)$.

Suppose it has been proven for some $l \geq 1$ that $\partial^{\alpha} u(x)$ exists and is given by $v_{\alpha}(x)$ for all $|\alpha| \leq l<k$. Then applying the results of the previous paragraph to $\partial^{\alpha} u(x)$ with $|\alpha|=l$ shows that $\partial_{i} \partial^{\alpha} u(x)$ exits and is given by $v_{\alpha+e_{i}}(x)$ for all $i$ and $x \in B$ and from this we conclude that $\partial^{\alpha} u(x)$ exists and is given by $v_{\alpha}(x)$ for all $|\alpha| \leq l+1$. So by induction we conclude $\partial^{\alpha} u(x)$ exists and is given by $v_{\alpha}(x)$ for all $|\alpha| \leq k$, i.e. $u \in C^{k}(B)$.
Lemma 48.32. Given any $k+1$ distinct points, $\left\{c_{i}\right\}_{i=0}^{k}$, in $\mathbb{R} \backslash\{0\}$, the $(k+1) \times(k+1)$ matrix $C$ with entries $C_{i j}:=\left(c_{i}\right)^{j}$ is invertible.

Proof. Let $a \in \mathbb{R}^{k+1}$ and define $p(x):=\sum_{j=0}^{k} a_{j} x^{j}$. If $a \in \operatorname{Nul}(C)$, then

$$
0=\sum_{j=0}^{k}\left(c_{i}\right)^{j} a_{j}=p\left(c_{i}\right) \text { for } i=0,1, \ldots, k .
$$

Since $\operatorname{deg}(p) \leq k$ and the above equation says that $p$ has $k+1$ distinct roots, we conclude that $a \in \operatorname{Nul}(C)$ implies $p \equiv 0$ which implies $a=0$. Therefore $\operatorname{Nul}(C)=\{0\}$ and $C$ is invertible.

Lemma 48.33. Let $B, B^{ \pm}$and $\Gamma$ be as in Lemma 48.31 and $\left\{c_{i}\right\}_{i=0}^{k}$, be $k+1$ distinct points in $(\infty,-1]$ for example $c_{i}=-(i+1)$ will work. Also let $a \in \mathbb{R}^{k+1}$ be the unique solution (see Lemma 48.32 to $C^{t r} a=\mathbf{1}$ where $\mathbf{1}$ denotes the vector of all ones in $\mathbb{R}^{k+1}$, i.e. a satisfies

$$
\begin{equation*}
1=\sum_{j=0}^{k}\left(c_{i}\right)^{j} a_{i} \text { for } j=0,1,2 \ldots, k \tag{48.23}
\end{equation*}
$$

For $u \in C_{c}^{k}\left(\overline{\mathbb{H}^{d}}\right)^{1}$ with $\operatorname{supp}(u) \subset B \cap \overline{\mathbb{H}^{d}}$ and $x=(y, z) \in \mathbb{R}^{d}$ define

$$
\tilde{u}(x)=\tilde{u}(y, z)=\left\{\begin{array}{cc}
u(y, z) & \text { if } z \geq 0  \tag{48.24}\\
\sum_{i=0}^{k} a_{i} u\left(y, c_{i} z\right) & \text { if } z \leq 0 .
\end{array}\right.
$$

Then $\tilde{u} \in C_{c}^{k}\left(\mathbb{R}^{d}\right)$ with $\operatorname{supp}(\tilde{u}) \subset B$ and moreover there exists a constant $M$ independent of $u$ such that

$$
\begin{equation*}
\|\tilde{u}\|_{W^{k, p}(B)} \leq M\|u\|_{W^{k, p}\left(B^{+}\right)} \tag{48.25}
\end{equation*}
$$

Proof. By Eq. (48.23) with $j=0$,

$$
\sum_{i=0}^{k} a_{i} u\left(y, c_{i} 0\right)=u(y, 0) \sum_{i=0}^{k} a_{i}=u(y, 0)
$$

${ }^{1}$ Or more generally, one may assume $u \in C^{k}\left(\mathbb{H}^{d}\right) \cap C_{c}\left(\overline{\mathbb{H}^{d}}\right)$ such that each $\partial^{\alpha} u$ for $|\alpha| \leq k$ extends to a continuous function on $\overline{\mathbb{H}^{d}}$.

This shows that $\tilde{u}$ in Eq. (48.24) is well defined and that $\tilde{u} \in C\left(\mathbb{H}^{d}\right)$. Let $K^{-}:=\{(y, z):(y,-z) \in \operatorname{supp}(u)\}$. Since $c_{i} \in(\infty,-1]$, if $x=(y, z) \notin K^{-}$ and $z<0$ then $\left(y, c_{i} z\right) \notin \operatorname{supp}(u)$ and therefore $\tilde{u}(x)=0$ and therefore $\operatorname{supp}(\tilde{u})$ is compactly contained inside of $B$. Similarly if $\alpha \in \mathbb{N}_{0}^{d}$ with $|\alpha| \leq k$, Eq. (48.23) with $j=\alpha_{d}$ implies

$$
v_{\alpha}(x):=\left\{\begin{array}{cc}
\left(\partial^{\alpha} u\right)(y, z) & \text { if } z \geq 0 \\
\sum_{i=0}^{k} a_{i} c_{i}^{\alpha_{d}}\left(\partial^{\alpha} u\right)\left(y, c_{i} z\right) & \text { if } z \leq 0 .
\end{array}\right.
$$

is well defined and $v_{\alpha} \in C\left(\mathbb{R}^{d}\right)$. Differentiating Eq. (48.24) shows $\partial^{\alpha} \tilde{u}(x)=$ $v_{\alpha}(x)$ for $x \in B \backslash \Gamma$ and therefore we may conclude from Lemma 48.31 that $\tilde{u} \in C_{c}^{k}(B) \subset C^{k}\left(\mathbb{R}^{d}\right)$ and $\partial^{\alpha} \tilde{u}=v_{\alpha}$ for all $|\alpha| \leq k$.

We now verify Eq. (48.25) as follows. For $|\alpha| \leq k$,

$$
\begin{aligned}
\left\|\partial^{\alpha} \tilde{u}\right\|_{L^{p}\left(B^{-}\right)}^{p} & =\int_{\mathbb{R}^{d}} 1_{z<0}\left|\sum_{i=0}^{k} a_{i} c_{i}^{\alpha_{d}}\left(\partial^{\alpha} u\right)\left(y, c_{i} z\right)\right|^{p} d y d z \\
& \leq C \int_{\mathbb{R}^{d}} 1_{z<0} \sum_{i=0}^{k}\left|\left(\partial^{\alpha} u\right)\left(y, c_{i} z\right)\right|^{p} d y d z \\
& =C \int_{\mathbb{R}^{d}} 1_{z>0} \sum_{i=0}^{k} \frac{1}{\left|c_{i}\right|}\left|\left(\partial^{\alpha} u\right)(y, z)\right|^{p} d y d z \\
& =C\left(\sum_{i=0}^{k} \frac{1}{\left|c_{i}\right|}\right)\left\|\partial^{\alpha} u\right\|_{L^{p}\left(B^{+}\right)}^{p}
\end{aligned}
$$

where $C:=\left(\sum_{i=0}^{k}\left|a_{i} c_{i}^{\alpha_{d}}\right|^{q}\right)^{p / q}$. Summing this equation on $|\alpha| \leq k$ shows there exists a constant $M^{\prime}$ such that $\|\tilde{u}\|_{W^{k, p}\left(B^{-}\right)} \leq M^{\prime}\|u\|_{W^{k, p}\left(B^{+}\right)}$and hence Eq. (48.25) holds with $M=M^{\prime}+1$.

Corollary 48.34. Let $k \geq 1, B, B^{ \pm},\left\{c_{i}\right\}_{i=0}^{k}$ be as in Lemma 48.33 and suppose that $u \in W^{k, \infty}\left(\mathbb{H}^{d}\right)$ with $\operatorname{supp}(u) \subset B \cap \overline{\mathbb{H}^{d}}$. By item 2. of Lemma 48.23, by modifying $u$ on a null set we may assume that $u \in B C^{k-1}\left(\mathbb{H}^{d}\right)$ with $\partial^{\alpha} u \in C\left(\overline{\mathbb{H}^{d}}\right)$ for all $|\alpha|<k$. Then the function $\tilde{u}$ defined in $E q$. (48.24) in $W^{k, \infty}\left(\mathbb{R}^{d}\right)$, with $\operatorname{supp}(\tilde{u}) \subset B$ and Eq. (48.25) is still valid.

Proof. By Lemma 48.33, $\tilde{u} \in C^{k-1}\left(\mathbb{R}^{d}\right)$. Let $\phi \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ and $|\alpha|=k-1$ and $i \in\{1,2, \ldots, d\}$, then by standard integration by parts

$$
\begin{aligned}
\left\langle\tilde{u}, \partial^{\alpha} \partial_{i} \phi\right\rangle & =(-1)^{|\alpha|}\left\langle\partial^{\alpha} \tilde{u}, \partial_{i} \phi\right\rangle \\
& =(-1)^{|\alpha|}\left\langle 1_{B^{+}} \partial^{\alpha} \tilde{u}, \partial_{i} \phi\right\rangle+(-1)^{|\alpha|}\left\langle 1_{B^{-}} \partial^{\alpha} \tilde{u}, \partial_{i} \phi\right\rangle .
\end{aligned}
$$

Making use of Proposition 48.26 and the change of variables theorem,

$$
\begin{aligned}
\left\langle\tilde{u}, \partial^{\alpha} \partial_{i} \phi\right\rangle & =-(-1)^{|\alpha|}\left\langle 1_{B^{+}} \partial_{i} \partial^{\alpha} \tilde{u}, \phi\right\rangle-(-1)^{|\alpha|}\left\langle 1_{B^{-}} \partial_{i} \partial^{\alpha} \tilde{u}, \phi\right\rangle \\
& =-(-1)^{|\alpha|}\left\langle\partial_{i} \partial^{\alpha} \tilde{u}, \phi\right\rangle
\end{aligned}
$$

wherein we have used the fact that $\left.\partial^{\alpha} \tilde{u}\right|_{\partial B^{+}}=\left.\partial^{\alpha} \tilde{u}\right|_{\partial B^{-}}$for all $|\alpha|<k$ to see that the boundary terms from the integrals cancel. Hence it follows that $\partial^{\alpha} \partial_{i} \tilde{u}$ exists weakly and is given by the expected formula, namely by differentiating Eq. (48.24) away from $\Gamma$ and piecing the results together. The verification of Eq. (48.25) is as before.

Theorem 48.35 (Extension Theorem). Suppose $k \geq 1$ and $\Omega \subset_{o} \mathbb{R}^{d}$ such that $\bar{\Omega}$ is a compact manifold with $C^{k}$ - boundary. Given $U \subset_{o} \mathbb{R}^{d}$ such that $\bar{\Omega} \subset U$, there exists a bounded linear (extension) operator $E: W^{k, p}(\Omega) \rightarrow$ $W^{k, p}\left(\mathbb{R}^{d}\right)$ such that

1. $E u=u$ a.e. in $\Omega$ and
2. $\operatorname{supp}(E u) \subset U$.

Proof. As in the proof of Theorem 48.24, choose a covering $\left\{V_{i}\right\}_{i=0}^{N}$ of $\bar{\Omega}$ such that $\bar{V}_{0} \subset \Omega, \cup_{i=0}^{N} \bar{V}_{i} \subset U$ and for each $i \geq 1$, there is $C^{k}$ - diffeomorphism $x_{i}: V_{i} \rightarrow R\left(x_{i}\right) \subset_{o} \mathbb{R}^{d}$ such that

$$
x_{i}\left(\partial \Omega \cap V_{i}\right)=R\left(x_{i}\right) \cap \operatorname{bd}\left(\mathbb{H}^{d}\right) \text { and } x_{i}\left(\Omega \cap V_{i}\right)=R\left(x_{i}\right) \cap \mathbb{H}^{d}=B^{+}
$$

where $B^{+}$is as in Lemma 48.33 and Corollary 48.34, refer to Figure 48.6. Further choose $\phi_{i} \in C_{c}^{\infty}\left(V_{i},[0,1]\right)$ such that $\sum_{i=0}^{N} \phi_{i}=1$ on a neighborhood of $\bar{\Omega}$ and set $y_{i}:=\left.x_{i}\right|_{\partial \Omega \cap V_{i}}$ for $i \geq 1$. Given $u \in C^{k}(\bar{\Omega})$ if $p<\infty(u \in$ $W^{k, \infty}(\Omega)$ if $\left.p=\infty\right)$ and $i \geq 1$, the function $v_{i}:=\left(\phi_{i} u\right) \circ x_{i}^{-1}$ may be viewed as a function in $C^{k}\left(\mathbb{H}^{d}\right) \cap C_{c}\left(\overline{\mathbb{H}^{d}}\right)\left(W^{k, \infty}\left(\mathbb{H}^{d}\right)\right)$ with $\operatorname{supp}(u) \subset B$. Let $\tilde{v}_{i} \in C_{c}^{k}(B)\left(W^{k, \infty}(B)\right)$ be defined as in Eq. (48.24) above and define $\tilde{u}:=\phi_{0} u+\sum_{i=1}^{N} \tilde{v}_{i} \circ x_{i} \in C_{c}^{k}\left(\mathbb{R}^{d}\right)\left(W^{k, \infty}\left(\mathbb{R}^{d}\right)\right)$. Notice that $\tilde{u}=u$ on $\bar{\Omega}$, $\operatorname{supp}(u) \subset U$ and by Lemma 48.20,

$$
\begin{aligned}
\|\tilde{u}\|_{W^{k, p}\left(\mathbb{R}^{d}\right)} & \leq\left\|\phi_{0} u\right\|_{W^{k, p}\left(\mathbb{R}^{d}\right)}+\sum_{i=1}^{N}\left\|\tilde{v}_{i} \circ x_{i}\right\|_{W^{k, p}\left(\mathbb{R}^{d}\right)} \\
& \leq\left\|\phi_{0} u\right\|_{W^{k, p}(\Omega)}+\sum_{i=1}^{N}\left\|\tilde{v}_{i}\right\|_{W^{k, p}\left(R\left(x_{i}\right)\right)} \\
& \leq C\left(\phi_{0}\right)\|u\|_{W^{k, p}(\Omega)}+\sum_{i=1}^{N}\left\|v_{i}\right\|_{W^{k, p}\left(B^{+}\right)} \\
& =C\left(\phi_{0}\right)\|u\|_{W^{k, p}(\Omega)}+\sum_{i=1}^{N}\left\|\left(\phi_{i} u\right) \circ x_{i}^{-1}\right\|_{W^{k, p}\left(B^{+}\right)} \\
& \leq C\left(\phi_{0}\right)\|u\|_{W^{k, p}(\Omega)}+\sum_{i=1}^{N} C_{i}\|u\|_{W^{k, p}(\Omega)} .
\end{aligned}
$$

This completes the proof for $p=\infty$ and shows for $p<\infty$ that the map $u \in$ $C^{k}(\bar{\Omega}) \rightarrow E u:=\tilde{u} \in C_{c}^{k}(U)$ is bounded as map from $W^{k, p}(\Omega)$ to $W^{k, p}(U)$. As usual, we now extend $E$ using the B.L.T. Theorem 2.68 to a bounded linear map from $W^{k, p}(\Omega)$ to $W^{k, p}(U)$. So for general $u \in W^{k, p}(\Omega), E u=W^{k, p}(U)$ $-\lim _{n \rightarrow \infty} \tilde{u}_{n}$ where $u_{n} \in C^{k}(\bar{\Omega})$ and $u=W^{k, p}(\Omega)-\lim _{n \rightarrow \infty} u_{n}$. By passing to a subsequence if necessary, we may assume that $\tilde{u}_{n}$ converges a.e. to $E u$ from which it follows that $E u=u$ a.e. on $\Omega$ and $\operatorname{supp}(E u) \subset U$.

### 48.6 Exercises

Exercise 48.36. Show the norm in Eq. (48.1) is equivalent to the norm

$$
|f|_{W^{k, p}(\Omega)}:=\sum_{|\alpha| \leq k}\left\|\partial^{\alpha} f\right\|_{L^{p}(\Omega)}
$$

Solution 48.37. 48.36This is a consequence of the fact that all norms on $l^{p}(\{\alpha:|\alpha| \leq k\})$ are equivalent. To be more explicit, let $a_{\alpha}=\left\|\partial^{\alpha} f\right\|_{L^{p}(\Omega)}$, then

$$
\sum_{|\alpha| \leq k}\left|a_{\alpha}\right| \leq\left(\sum_{|\alpha| \leq k}\left|a_{\alpha}\right|^{p}\right)^{1 / p}\left(\sum_{|\alpha| \leq k} 1^{q}\right)^{1 / q}
$$

while
$\left(\sum_{|\alpha| \leq k}\left|a_{\alpha}\right|^{p}\right)^{1 / p} \leq\left(\sum_{|\alpha| \leq k}^{p}\left[\sum_{|\beta| \leq k}\left|a_{\beta}\right|\right]^{p}\right)^{1 / p} \leq[\#\{\alpha:|\alpha| \leq k\}]^{1 / p} \sum_{|\beta| \leq k}\left|a_{\beta}\right|$.

## Sobolev Inequalities

### 49.1 Morrey's Inequality

Notation 49.1 Let $S^{d-1}$ be the sphere of radius one centered at zero inside $\mathbb{R}^{d}$. For a set $\Gamma \subset S^{d-1}, x \in \mathbb{R}^{d}$, and $r \in(0, \infty)$, let

$$
\Gamma_{x, r} \equiv\{x+s \omega: \omega \in \Gamma \text { such that } 0 \leq s \leq r\} .
$$

So $\Gamma_{x, r}=x+\Gamma_{0, r}$ where $\Gamma_{0, r}$ is a cone based on $\Gamma$, see Figure 49.1 below.


Fig. 49.1. The cone $\Gamma_{0, r}$.

Notation 49.3 If $\Omega \subset \mathbb{R}^{d}$ is a measurable set and $f: \mathbb{R}^{d} \rightarrow \mathbb{C}$ is a measurable function let

$$
f_{\Omega}:=f_{\Omega} f(x) d x:=\frac{1}{m(\Omega)} \int_{\Omega} f(x) d x .
$$

By Theorem 9.35,

$$
\begin{equation*}
\int_{\Gamma_{x, r}} f(y) d y=\int_{\Gamma_{0, r}} f(x+y) d y=\int_{0}^{r} d t t^{d-1} \int_{\Gamma} f(x+t \omega) d \sigma(\omega) \tag{49.1}
\end{equation*}
$$

and letting $f=1$ in this equation implies

$$
\begin{equation*}
m\left(\Gamma_{x, r}\right)=|\Gamma| r^{d} / d \tag{49.2}
\end{equation*}
$$

Lemma 49.4. Let $\Gamma \subset S^{d-1}$ be a measurable set such that $|\Gamma|>0$. For $u \in C^{1}\left(\bar{\Gamma}_{x, r}\right)$,

$$
\begin{equation*}
f_{\Gamma_{x, r}}|u(y)-u(x)| d y \leq \frac{1}{|\Gamma|} \int_{\Gamma_{x, r}} \frac{|\nabla u(y)|}{|x-y|^{d-1}} d y . \tag{49.3}
\end{equation*}
$$

Proof. Write $y=x+s \omega$ with $\omega \in S^{d-1}$, then by the fundamental theorem of calculus,

$$
u(x+s \omega)-u(x)=\int_{0}^{s} \nabla u(x+t \omega) \cdot \omega d t
$$

and therefore,

$$
\begin{aligned}
\int_{\Gamma}|u(x+s \omega)-u(x)| d \sigma(\omega) & \leq \int_{0}^{s} \int_{\Gamma}|\nabla u(x+t \omega)| d \sigma(\omega) d t \\
& =\int_{0}^{s} t^{d-1} d t \int_{\Gamma} \frac{|\nabla u(x+t \omega)|}{|x+t \omega-x|^{d-1}} d \sigma(\omega) \\
& =\int_{\Gamma_{x, s}} \frac{|\nabla u(y)|}{|y-x|^{d-1}} d y \leq \int_{\Gamma_{x, r}} \frac{|\nabla u(y)|}{|x-y|^{d-1}} d y
\end{aligned}
$$

wherein the second equality we have used Eq. (49.1). Multiplying this inequality by $s^{d-1}$ and integrating on $s \in[0, r]$ gives

$$
\int_{\Gamma_{x, r}}|u(y)-u(x)| d y \leq \frac{r^{d}}{d} \int_{\Gamma_{x, r}} \frac{|\nabla u(y)|}{|x-y|^{d-1}} d y=\frac{m\left(\Gamma_{x, r}\right)}{|\Gamma|} \int_{\Gamma_{x, r}} \frac{|\nabla u(y)|}{|x-y|^{d-1}} d y
$$

which proves Eq. (49.3).
Corollary 49.5. Suppose $d<p \leq \infty, \Gamma \in \mathcal{B}_{S^{d-1}}$ such that $|\Gamma|>0, r \in$ $(0, \infty)$ and $u \in C^{1}\left(\bar{\Gamma}_{x, r}\right)$. Then

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$$
|u(x)| \leq C(|\Gamma|, r, d, p)\|u\|_{W^{1, p}\left(\Gamma_{x, r}\right)}
$$

where

$$
C(|\Gamma|, r, d, p):=\frac{1}{|\Gamma|^{1 / p}} \max \left(\frac{d^{-1 / p}}{r},\left(\frac{p-1}{p-d}\right)^{1-1 / p}\right) \cdot r^{1-d / p}
$$

Proof. For $y \in \Gamma_{x, r}$,

$$
|u(x)| \leq|u(y)|+|u(y)-u(x)|
$$

and hence using Eq. (49.3) and Hölder's inequality,

$$
\begin{align*}
|u(x)| \leq & f_{\Gamma_{x, r}}|u(y)| d y+\frac{1}{|\Gamma|} \int_{\Gamma_{x, r}} \frac{|\nabla u(y)|}{|x-y|^{d-1}} d y \\
\leq & \frac{1}{m\left(\Gamma_{x, r}\right)}\|u\|_{L^{p}\left(\Gamma_{x, r}\right)}\|1\|_{L^{p}\left(\Gamma_{x, r}\right)} \\
& \quad+\frac{1}{|\Gamma|}\|\nabla u\|_{L^{p}\left(\Gamma_{x, r}\right)}\left\|\frac{1}{|x-\cdot|^{d-1}}\right\|_{L^{q}\left(\Gamma_{x, r}\right)} \tag{49.5}
\end{align*}
$$

where $q=\frac{p}{p-1}$ as before. Now

$$
\begin{aligned}
\left\|\frac{1}{|\cdot|^{d-1}}\right\|_{L^{q}\left(\Gamma_{0, r}\right)}^{q} & =\int_{0}^{r} d t t^{d-1} \int_{\Gamma}\left(t^{d-1}\right)^{-q} d \sigma(\omega) \\
& =|\Gamma| \int_{0}^{r} d t\left(t^{d-1}\right)^{1-\frac{p}{p-1}}=|\Gamma| \int_{0}^{r} d t t^{-\frac{d-1}{p-1}}
\end{aligned}
$$

and since

$$
1-\frac{d-1}{p-1}=\frac{p-d}{p-1}
$$

we find

$$
\begin{equation*}
\left\|\frac{1}{|\cdot|^{d-1}}\right\|_{L^{q}\left(\Gamma_{0, r}\right)}=\left(\frac{p-1}{p-d}|\Gamma| r^{\frac{p-d}{p-1}}\right)^{1 / q}=\left(\frac{p-1}{p-d}|\Gamma|\right)^{\frac{p-1}{p}} r^{1-\frac{d}{p}} \tag{49.6}
\end{equation*}
$$

Combining Eqs. (49.5), Eq. (49.6) along with the identity,

$$
\begin{equation*}
\frac{1}{m\left(\Gamma_{x, r}\right)}\|1\|_{L^{q}\left(\Gamma_{x, r}\right)}=\frac{1}{m\left(\Gamma_{x, r}\right)} m\left(\Gamma_{x, r}\right)^{1 / q}=\left(|\Gamma| r^{d} / d\right)^{-1 / p} \tag{49.7}
\end{equation*}
$$

shows
$|u(x)| \leq\|u\|_{L^{p}\left(\Gamma_{x, r}\right)}\left(|\Gamma| r^{d} / d\right)^{-1 / p}+\frac{1}{|\Gamma|}\|\nabla u\|_{L^{p}\left(\Gamma_{x, r}\right)}\left(\frac{p-1}{p-d}|\Gamma|\right)^{1-1 / p} r^{1-d / p}$ $=\frac{1}{|\Gamma|^{1 / p}}\left[\|u\|_{L^{p}\left(\Gamma_{x, r}\right)} \frac{d^{-1 / p}}{r}+\|\nabla u\|_{L^{p}\left(\Gamma_{x, r}\right)}\left(\frac{p-1}{p-d}\right)^{1-1 / p}\right] r^{1-d / p}$. $\leq \frac{1}{|\Gamma|^{1 / p}} \max \left(\frac{d^{-1 / p}}{r},\left(\frac{p-1}{p-d}\right)^{1-1 / p}\right)\|u\|_{W^{1, p}\left(\Gamma_{x, r}\right)} \cdot r^{1-d / p}$.

Corollary 49.6. For $d \in \mathbb{N}$ and $p \in(d, \infty]$ there are constants $\alpha=\alpha_{d}$ and $\beta=\beta_{d}$ such that if $u \in C^{1}\left(\mathbb{R}^{d}\right)$ then for all $x, y \in \mathbb{R}^{d}$,

$$
\begin{equation*}
|u(y)-u(x)| \leq 2 \beta \alpha^{1 / p}\left(\frac{p-1}{p-d}\right)^{\frac{p-1}{p}}\|\nabla u\|_{L^{p}(B(x, r) \cap B(y, r))} \cdot|x-y|^{\left(1-\frac{d}{p}\right)} \tag{49.8}
\end{equation*}
$$ where $r:=|x-y|$.

Proof. Let $r:=|x-y|, V:=B_{x}(r) \cap B_{y}(r)$ and $\Gamma, \Lambda \subset S^{d-1}$ be chosen so that $x+r \Gamma=\partial B_{x}(r) \cap B_{y}(r)$ and $y+r \Lambda=\partial B_{y}(r) \cap B_{x}(r)$, i.e.

$$
\Gamma=\frac{1}{r}\left(\partial B_{x}(r) \cap B_{y}(r)-x\right) \text { and } \Lambda=\frac{1}{r}\left(\partial B_{y}(r) \cap B_{x}(r)-y\right)=-\Gamma
$$

Also let $W=\Gamma_{x, r} \cap \Lambda_{y, r}$, see Figure 49.2 below. By a scaling,


Fig. 49.2. The geometry of two intersecting balls of radius $r:=|x-y|$. Here $W=\Gamma_{x, r} \cap \Lambda_{y, r}$ and $V=B(x, r) \cap B(y, r)$.

$$
\beta_{d}:=\frac{\left|\Gamma_{x, r} \cap \Lambda_{y, r}\right|}{\left|\Gamma_{x, r}\right|}=\frac{\left|\Gamma_{x, 1} \cap \Lambda_{y, 1}\right|}{\left|\Gamma_{x, 1}\right|} \in(0,1)
$$

is a constant only depending on $d$, i.e. we have $\left|\Gamma_{x, r}\right|=\left|\Lambda_{y, r}\right|=\beta|W|$. Integrating the inequality

$$
|u(x)-u(y)| \leq|u(x)-u(z)|+|u(z)-u(y)|
$$

over $z \in W$ gives

$$
\begin{aligned}
|u(x)-u(y)| & \leq \int_{W}|u(x)-u(z)| d z+\int_{W}|u(z)-u(y)| d z \\
& =\frac{\beta}{\left|\Gamma_{x, r}\right|}\left(\int_{W}|u(x)-u(z)| d z+\int_{W}|u(z)-u(y)| d z\right) \\
& \leq \frac{\beta}{\left|\Gamma_{x, r}\right|}\left(\int_{\Gamma_{x, r}}|u(x)-u(z)| d z+\int_{\Lambda_{y, r}}|u(z)-u(y)| d z\right) .
\end{aligned}
$$

Hence by Lemma 49.4, Hölder's inequality and translation and rotation invariance of Lebesgue measure,

$$
\begin{align*}
|u(x)-u(y)| & \leq \frac{\beta}{|\Gamma|}\left(\int_{\Gamma_{x, r}} \frac{|\nabla u(z)|}{|x-z|^{d-1}} d z+\int_{\Lambda_{y, r}} \frac{|\nabla u(z)|}{|z-y|^{d-1}} d z\right) \\
& \leq \frac{\beta}{|\Gamma|}\binom{\|\nabla u\|_{L^{p}\left(\Gamma_{x, r}\right)}\left\|\frac{1}{\mid x-\cdot d^{d-1}}\right\|_{L^{q}\left(\Gamma_{x, r}\right)}}{+\|\nabla u\|_{L^{p}\left(\Lambda_{y, r}\right)}\left\|\frac{1}{|y-\cdot|^{d-1}}\right\|_{L^{q}\left(\Lambda_{y, r}\right)}} \\
& \leq \frac{2 \beta}{|\Gamma|}\|\nabla u\|_{L^{p}(V)}\left\|\frac{1}{|\cdot|^{d-1}}\right\|_{L^{q}\left(\Gamma_{0, r}\right)} \tag{49.9}
\end{align*}
$$

where $q=\frac{p}{p-1}$ is the conjugate exponent to $p$. Combining Eqs. (49.9) and (49.6) gives Eq. (49.8) with $\alpha:=|\Gamma|^{-1}$.

Theorem 49.7 (Morrey's Inequality). If $d<p \leq \infty, u \in W^{1, p}\left(\mathbb{R}^{d}\right)$, then there exists a unique version $u^{*}$ of $u$ (i.e. $u^{*}=u$ a.e.) such that $u^{*}$ is continuous. Moreover $u^{*} \in C^{0,1-\frac{d}{p}}\left(\mathbb{R}^{d}\right)$ and

$$
\begin{equation*}
\left\|u^{*}\right\|_{C^{0,1-\frac{d}{p}}\left(\mathbb{R}^{d}\right)} \leq C\|u\|_{W^{1, p}\left(\mathbb{R}^{d}\right)} \tag{49.10}
\end{equation*}
$$

where $C=C(p, d)$ is a universal constant. Moreover, the estimates in Eqs. (49.3), (49.4) and (49.8) still hold when $u$ is replaced by $u^{*}$.

Proof. For $p<\infty$ and $u \in C_{c}^{1}\left(\mathbb{R}^{d}\right)$, Corollaries 49.5 and 49.6 imply

$$
\|u\|_{B C\left(\mathbb{R}^{d}\right)} \leq C\|u\|_{W^{1, p}\left(\mathbb{R}^{d}\right)} \text { and } \frac{|u(y)-u(x)|}{|x-y|^{1-\frac{d}{p}}} \leq C\|\nabla u\|_{L^{p}\left(\mathbb{R}^{d}\right)}
$$

which implies $[u]_{1-\frac{d}{p}} \leq C\|\nabla u\|_{L^{p}\left(\mathbb{R}^{d}\right)} \leq C\|u\|_{W^{1, p}\left(\mathbb{R}^{d}\right)}$ and hence

$$
\begin{equation*}
\|u\|_{C^{0,1-\frac{d}{p}}\left(\mathbb{R}^{d}\right)} \leq C\|u\|_{W^{1, p}\left(\mathbb{R}^{d}\right)} \tag{49.11}
\end{equation*}
$$

Now suppose $u \in W^{1, p}\left(\mathbb{R}^{d}\right)$, choose (using Exercise 29.32) $u_{n} \in C_{c}^{1}\left(\mathbb{R}^{d}\right)$ such that $u_{n} \rightarrow u$ in $W^{1, p}\left(\mathbb{R}^{d}\right)$. Then by Eq. (49.11), $\left\|u_{n}-u_{m}\right\|_{C^{0,1-\frac{d}{p}\left(\mathbb{R}^{d}\right)}} \rightarrow 0$ as $m, n \rightarrow \infty$ and therefore there exists $u^{*} \in C^{0,1-\frac{d}{p}}\left(\mathbb{R}^{d}\right)$ such that $u_{n} \rightarrow u^{*}$ in $C^{0,1-\frac{d}{p}}\left(\mathbb{R}^{d}\right)$. Clearly $u^{*}=u$ a.e. and Eq. (49.10) holds.

If $p=\infty$ and $u \in W^{1, \infty}\left(\mathbb{R}^{d}\right)$, then by Proposition 29.29 there is a version $u^{*}$ of $u$ which is Lipschitz continuous. Now in both cases, $p<\infty$ and $p=\infty$, the sequence $u_{m}:=u * \eta_{m}=u^{*} * \eta_{m} \in C^{\infty}\left(\mathbb{R}^{d}\right)$ and $u_{m} \rightarrow u^{*}$ uniformly on compact subsets of $\mathbb{R}^{d}$. Using Eq. (49.3) with $u$ replaced by $u_{m}$ along with a (by now) standard limiting argument shows that Eq. (49.3) still holds with $u$ replaced by $u^{*}$. The proofs of Eqs. (49.4) and (49.8) only relied on Eq. (49.3) and hence go through without change. Similarly the argument in the first paragraph only relied on Eqs. (49.4) and (49.8) and hence Eq. (49.10) is also valid for $p=\infty$.
Corollary 49.8 (Morrey's Inequality). Suppose $\Omega \subset_{o} \mathbb{R}^{d}$ such that $\bar{\Omega}$ is compact $C^{1}$-manifold with boundary and $d<p \leq \infty$. Then for $u \in W^{1, p}(\Omega)$, there exists a unique version $u^{*}$ of $u$ such that $u^{*} \in C^{0,1-\frac{d}{p}}\left(\mathbb{R}^{d}\right)$ and we further have

$$
\begin{equation*}
\left\|u^{*}\right\|_{C^{0,1-\frac{d}{p}}(\Omega)} \leq C\|u\|_{W^{1, p}(\Omega)} \tag{49.12}
\end{equation*}
$$

where $C=C(p, d, \Omega)$.
Proof. Let $U$ be a precompact open subset of $\mathbb{R}^{d}$ and $E: W^{1, p}(\Omega) \rightarrow$ $W^{1, p}\left(\mathbb{R}^{d}\right)$ be an extension operator as in Theorem 48.35. For $u \in W^{1, p}(\Omega)$ with $d<p \leq \infty$, Theorem 49.7 implies there is a version $U^{*} \in C^{0,1-\frac{d}{p}}\left(\mathbb{R}^{d}\right)$ of $E u$. Letting $u^{*}:=\left.U^{*}\right|_{\Omega}$, we have and moreover,

$$
\left\|u^{*}\right\|_{C^{0,1-\frac{d}{p}}(\Omega)} \leq\left\|U^{*}\right\|_{C^{0,1-\frac{d}{p}}\left(\mathbb{R}^{d}\right)} \leq C\|E u\|_{W^{1, p}\left(\mathbb{R}^{d}\right)} \leq C\|u\|_{W^{1, p}(\Omega)} .
$$

The following example shows that $L^{\infty}\left(\mathbb{R}^{d}\right) \nsubseteq W^{1, d}\left(\mathbb{R}^{d}\right)$, i.e. $W^{1, d}\left(\mathbb{R}^{d}\right)$ contains unbounded elements. Therefore Theorem 49.7 and Corollary 49.8 are not valid for $p=d$. It turns out that for $p=d, W^{1, d}\left(\mathbb{R}^{d}\right)$ embeds into $B M O\left(\mathbb{R}^{d}\right)$ - the space of functions with "bounded mean oscillation."
Example 49.9. Let $u(x)=\psi(x) \log \log \left(1+\frac{1}{|x|}\right)$ where $\psi \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ is chosen so that $\psi(x)=1$ for $|x| \leq 1$. Then $u \notin L^{\infty}\left(\mathbb{R}^{d}\right)$ while $u \in W^{1, d}\left(\mathbb{R}^{d}\right)$. Let us check this claim. Using Theorem 9.35, one easily shows $u \in L^{p}\left(\mathbb{R}^{d}\right)$. A short computation shows, for $|x|<1$, that

$$
\begin{aligned}
\nabla u(x) & =\frac{1}{\log \left(1+\frac{1}{|x|}\right)} \frac{1}{1+\frac{1}{|x|}} \nabla \frac{1}{|x|} \\
& =\frac{1}{1+\frac{1}{|x|}} \frac{1}{\log \left(1+\frac{1}{|x|}\right)}\left(-\frac{1}{|x|} \hat{x}\right)
\end{aligned}
$$

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where $\hat{x}=x /|x|$ and so again by Theorem 9.35,

$$
\begin{aligned}
\int_{\mathbb{R}^{d}}|\nabla u(x)|^{d} d x & \geq \int_{|x|<1}\left(\frac{1}{|x|^{2}+|x|} \frac{1}{\log \left(1+\frac{1}{|x|}\right)}\right)^{d} d x \\
& \geq \sigma\left(S^{d-1}\right) \int_{0}^{1}\left(\frac{2}{r \log \left(1+\frac{1}{r}\right)}\right)^{d} r^{d-1} d r=\infty .
\end{aligned}
$$

### 49.2 Rademacher's Theorem

Theorem 49.10. Suppose that $u \in W_{\text {loc }}^{1, p}(\Omega)$ for some $d<p \leq \infty$. Then $u$ is differentiable almost everywhere and $w-\partial_{i} u=\partial_{i} u$ a.e. on $\Omega$.

Proof. We clearly may assume that $p<\infty$. For $v \in W_{l o c}^{1, p}(\Omega)$ and $x, y \in \Omega$ such that $\overline{B(x, r) \cap B(y, r)} \subset \Omega$ where $r:=|x-y|$, the estimate in Corollary 49.6, gives

$$
\begin{array}{r}
|v(y)-v(x)| \leq C\|\nabla u\|_{L^{p}(B(x, r) \cap B(y, r))} \cdot|x-y|^{\left(1-\frac{d}{p}\right)} \\
=C\|\nabla v\|_{L^{p}(B(x, r) \cap B(y, r))} \cdot r^{\left(1-\frac{d}{p}\right)} . \tag{49.13}
\end{array}
$$

Let $u$ now denote the unique continuous version of $u \in W_{l o c}^{1, p}(\Omega)$. The by the Lebesgue differentiation Theorem 20.12, there exists an exceptional set $E \subset \Omega$ such that $m(E)=0$ and

$$
\lim _{r \downarrow 0} f_{B(x, r)}|\nabla u(y)-\nabla u(x)|^{p} d y=0 \text { for } x \in \Omega \backslash E .
$$

Fix a point $x \in \Omega \backslash E$ and let $v(y):=u(y)-u(x)-\nabla u(x) \cdot(y-x)$ and notice that $\nabla v(y)=\nabla u(y)-\nabla u(x)$. Applying Eq. (49.13) to $v$ then implies

$$
\begin{aligned}
\mid u(y) & -u(x)-\nabla u(x) \cdot(y-x) \mid \\
& \leq C\|\nabla u(\cdot)-\nabla u(x)\|_{L^{p}(B(x, r) \cap B(y, r))} \cdot r^{\left(1-\frac{d}{p}\right)} \\
& \leq C\left(\int_{B(x, r)}|\nabla u(y)-\nabla u(x)|^{p} d y\right)^{1 / p} \cdot r^{\left(1-\frac{d}{p}\right)} \\
& =C \sigma\left(S^{d-1}\right)^{1 / p} r^{d / p}\left({\left.\underset{B(x, r)}{ }|\nabla u(y)-\nabla u(x)|^{p} d y\right)^{1 / p} \cdot r^{\left(1-\frac{d}{p}\right)}} \quad=C \sigma\left(S^{d-1}\right)^{1 / p}\left({\left.\underset{B(x, r)}{ }|\nabla u(y)-\nabla u(x)|^{p} d y\right)^{1 / p} \cdot|x-y|}^{f} \quad l\right.\right.
\end{aligned}
$$

which shows $u$ is differentiable at $x$ and $\nabla u(x)=w-\nabla u(x)$.

Theorem 49.11 (Rademacher's Theorem). Let u be locally Lipschitz continuous on $\Omega \subset_{o} \mathbb{R}^{d}$. Then $u$ is differentiable almost everywhere and $w$ $\partial_{i} u=\partial_{i} u$ a.e. on $\Omega$.

Proof. By Proposition $29.29 \partial_{i}^{(w)} u$ exists weakly and is in $\partial_{i} u \in L^{\infty}\left(\mathbb{R}^{d}\right)$ for $i=1,2, \ldots, d$. The result now follows from Theorem 49.10.

### 49.3 Gagliardo-Nirenberg-Sobolev Inequality

In this section our goal is to prove an inequality of the form:

$$
\begin{equation*}
\|u\|_{L^{q}} \leq C\|\nabla u\|_{L^{p}\left(\mathbb{R}^{d}\right)} \text { for } u \in C_{c}^{1}\left(\mathbb{R}^{d}\right) \tag{49.14}
\end{equation*}
$$

For $\lambda>0$, let $u_{\lambda}(x)=u(\lambda x)$. Then

$$
\left\|u_{\lambda}\right\|_{L^{q}}^{q}=\int_{\mathbb{R}^{d}}|u(\lambda x)|^{q} d x=\int_{\mathbb{R}^{d}}|u(y)|^{q} \frac{d y}{\lambda^{d}}
$$

and hence $\left\|u_{\lambda}\right\|_{L^{q}}=\lambda^{-d / q}\|u\|_{L^{q}}$. Moreover, $\nabla u_{\lambda}(x)=\lambda(\nabla u)(\lambda x)$ and thus

$$
\left\|\nabla u_{\lambda}\right\|_{L^{p}}=\lambda\left\|(\nabla u)_{\lambda}\right\|_{L^{p}}=\lambda \lambda^{-d / p}\|\nabla u\|_{L^{p}}
$$

If (49.14) is to hold for all $u \in C_{c}^{1}\left(\mathbb{R}^{d}\right)$ then we must have

$$
\lambda^{-d / q}\|u\|_{L^{q}}=\left\|u_{\lambda}\right\|_{L^{q}} \leq C\left\|\nabla u_{\lambda}\right\|_{L^{p}\left(\mathbb{R}^{d}\right)}=C \lambda^{1-d / p}\|\nabla u\|_{L^{p}} \text { for all } \lambda>0
$$

which is only possible if

$$
\begin{equation*}
1-d / p+d / q=0, \text { i.e. } 1 / p=1 / d+1 / q \tag{49.15}
\end{equation*}
$$

Notation 49.12 For $p \in[1, d]$, let $p^{*}:=\frac{d p}{d-p}$ with the convention that $p^{*}=$ $\infty$ if $p=d$. That is $p^{*}=q$ where $q$ solves Eq. (49.15).
Theorem 49.13. Let $p=1$ so $1^{*}=\frac{d}{d-1}$, then

$$
\begin{equation*}
\|u\|_{1^{*}}=\|u\|_{\frac{d}{d-1}} \leq \prod_{i=1}^{d}\left(\int_{\mathbb{R}^{d}}\left|\partial_{i} u(x)\right| d x\right)^{\frac{1}{d}} \leq d^{-\frac{1}{2}}\|\nabla u\|_{1} \tag{49.16}
\end{equation*}
$$

for all $u \in W^{1,1}\left(\mathbb{R}^{d}\right)$.
Proof. Since there exists $u_{n} \in C_{c}^{1}\left(\mathbb{R}^{d}\right)$ such that $u_{n} \rightarrow u$ in $W^{1,1}\left(\mathbb{R}^{d}\right)$, a simple limiting argument shows that it suffices to prove Eq. (49.16) for $u \in C_{c}^{1}\left(\mathbb{R}^{d}\right)$. To help the reader understand the proof, let us give the proof for $d \leq 3$ first and with the constant $d^{-1 / 2}$ being replaced by 1 . After that the general induction argument will be given. (The adventurous reader may skip directly to the paragraph containing Eq. (49.17).

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$\left(d=1, p^{*}=\infty\right)$ By the fundamental theorem of calculus,

$$
|u(x)|=\left|\int_{-\infty}^{x} u^{\prime}(y) d y\right| \leq \int_{-\infty}^{x}\left|u^{\prime}(y)\right| d y \leq \int_{\mathbb{R}}\left|u^{\prime}(x)\right| d x
$$

Therefore $\|u\|_{L^{\infty}} \leq\left\|u^{\prime}\right\|_{L^{1}}{ }^{1}$, proving the $d=1$ case.
$\left(d=2, p^{*}=2\right)$ Applying the same argument as above to $y_{1} \rightarrow u\left(y_{1}, x_{2}\right)$ and $y_{2} \rightarrow u\left(x_{1}, y_{2}\right)$, ,

$$
\begin{aligned}
& \left|u\left(x_{1}, x_{2}\right)\right| \leq \int_{-\infty}^{\infty}\left|\partial_{1} u\left(y_{1}, x_{2}\right)\right| d y_{1} \leq \int_{-\infty}^{\infty}\left|\nabla u\left(y_{1}, x_{2}\right)\right| d y_{1} \text { and } \\
& \left|u\left(x_{1}, x_{2}\right)\right| \leq \int_{-\infty}^{\infty}\left|\partial_{2} u\left(x_{1}, y_{2}\right)\right| d y_{2} \leq \int_{-\infty}^{\infty}\left|\nabla u\left(x_{1}, y_{2}\right)\right| d y_{2}
\end{aligned}
$$

and therefore

$$
\left|u\left(x_{1}, x_{2}\right)\right|^{2} \leq \int_{-\infty}^{\infty}\left|\partial_{1} u\left(y_{1}, x_{2}\right)\right| d y_{1} \cdot \int_{-\infty}^{\infty}\left|\partial_{2} u\left(x_{1}, y_{2}\right)\right| d y_{2}
$$

Integrating this equation relative to $x_{1}$ and $x_{2}$ gives

$$
\begin{aligned}
\|u\|_{L^{2}}^{2} & =\int_{\mathbb{R}^{2}}|u(x)|^{2} d x \leq\left(\int_{-\infty}^{\infty}\left|\partial_{1} u(x)\right| d x\right)\left(\int_{-\infty}^{\infty}\left|\partial_{2} u(x)\right| d x\right) \\
& \leq\left(\int_{-\infty}^{\infty}|\nabla u(x)| d x\right)^{2}
\end{aligned}
$$

which proves the $d=2$ case.
$\left(d=3, p^{*}=3 / 2\right)$ Let $x^{1}=\left(y_{1}, x_{2}, x_{3}\right), x^{2}=\left(x_{1}, y_{2}, x_{3}\right)$, and $x^{3}=$ $\left(x_{1}, x_{2}, y_{3}\right)$. Then as above,

$$
|u(x)| \leq \int_{-\infty}^{\infty}\left|\partial_{i} u\left(x^{i}\right)\right| d y_{i} \text { for } i=1,2,3
$$

${ }^{1}$ Actually we may do better here by observing

$$
\begin{aligned}
|u(x)| & =\frac{1}{2}\left|\int_{-\infty}^{x} u^{\prime}(y) d y-\int_{x}^{\infty} u^{\prime}(y) d y\right| \\
& \leq \frac{1}{2} \int_{\mathbb{R}}\left|u^{\prime}(x)\right| d x
\end{aligned}
$$

and this leads to an improvement in Eq. (49.17) to

$$
\begin{aligned}
\|u\|_{1^{*}} & \leq \frac{1}{2} \prod_{i=1}^{d}\left(\int_{\mathbb{R}^{d}}\left|\partial_{i} u(x)\right| d x\right)^{\frac{1}{d}} \\
& \leq \frac{1}{2} d^{-\frac{1}{2}}\|\nabla u\|_{1}
\end{aligned}
$$

and hence

$$
|u(x)|^{\frac{3}{2}} \leq \prod_{i=1}^{3}\left(\int_{-\infty}^{\infty}\left|\partial_{i} u\left(x^{i}\right)\right| d y_{i}\right)^{\frac{1}{2}}
$$

Integrating this equation on $x_{1}$ gives,

$$
\begin{aligned}
\int_{\mathbb{R}}|u(x)|^{\frac{3}{2}} d x_{1} & \leq\left(\int_{-\infty}^{\infty}\left|\partial_{1} u\left(x^{1}\right)\right| d y_{1}\right)^{\frac{1}{2}} \int \prod_{i=2}^{3}\left(\int_{-\infty}^{\infty}\left|\partial_{i} u\left(x^{i}\right)\right| d y_{i}\right)^{\frac{1}{2}} d x_{1} \\
& \leq\left(\int_{-\infty}^{\infty}\left|\partial_{1} u(x)\right| d x_{1}\right)^{\frac{1}{2}} \prod_{i=2}^{3}\left(\int_{-\infty}^{\infty}\left|\partial_{i} u\left(x^{i}\right)\right| d x_{1} d y_{i}\right)^{\frac{1}{2}}
\end{aligned}
$$

wherein the second equality we have used the Hölder's inequality with $p=$ $q=2$. Integrating this result on $x_{2}$ and using Hölder's inequality gives

$$
\begin{aligned}
& \int_{\mathbb{R}^{2}}|u(x)|^{\frac{3}{2}} d x_{1} d x_{2} \\
& \leq\left(\int_{\mathbb{R}^{2}}\left|\partial_{2} u(x)\right| d x_{1} d x_{2}\right)^{\frac{1}{2}} \int_{\mathbb{R}} d x_{2}\left(\int_{-\infty}^{\infty}\left|\partial_{1} u(x)\right| d x_{1}\right)^{\frac{1}{2}} \times \\
& \quad\left(\int_{\mathbb{R}^{2}}\left|\partial_{3} u\left(x^{3}\right)\right| d x_{1} d y_{3}\right)^{\frac{1}{2}} \\
& \leq\left(\int_{\mathbb{R}^{2}}\left|\partial_{2} u(x)\right| d x_{1} d x_{2}\right)^{\frac{1}{2}}\left(\int_{\mathbb{R}^{2}}\left|\partial_{1} u(x)\right| d x_{1} d x_{2}\right)^{\frac{1}{2}}\left(\int_{\mathbb{R}^{3}}\left|\partial_{3} u(x)\right| d x\right)^{\frac{1}{2}} .
\end{aligned}
$$

One more integration of $x_{3}$ and application of Hölder's inequality, implies

$$
\int_{\mathbb{R}^{3}}|u(x)|^{\frac{3}{2}} d x \leq \prod_{i=1}^{3}\left(\int_{\mathbb{R}^{3}}\left|\partial_{i} u(x)\right| d x\right)^{\frac{1}{2}} \leq\left(\int_{\mathbb{R}^{3}}|\nabla u(x)| d x\right)^{\frac{3}{2}}
$$

proving the $d=3$ case.
For general $d\left(p^{*}=\frac{d}{d-1}\right)$, as above let $x^{i}=\left(x_{1}, \ldots, x_{i-1}, y_{i}, x_{i+1} \ldots, x_{d}\right)$. Then

$$
|u(x)| \leq \int_{-\infty}^{\infty}\left|\partial_{i} u\left(x^{i}\right)\right| d y_{i}
$$

and

$$
\begin{equation*}
|u(x)|^{\frac{d}{d-1}} \leq \prod_{i=1}^{d}\left(\int_{-\infty}^{\infty}\left|\partial_{i} u\left(x^{i}\right)\right| d y_{i}\right)^{\frac{1}{d-1}} \tag{49.17}
\end{equation*}
$$

Integrating this equation relative to $x_{1}$ and making use of Hölder's inequality in the form

$$
\begin{equation*}
\left\|\prod_{i=2}^{d} f_{i}\right\|_{1} \leq \prod_{i=2}^{d}\left\|f_{i}\right\|_{d-1} \tag{49.18}
\end{equation*}
$$

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$$
\begin{aligned}
\int_{\mathbb{R}}|u(x)|^{\frac{d}{d-1}} d x_{1} \leq & \left(\int_{\mathbb{R}} \partial_{1} u(x) d x_{1}\right)^{\frac{1}{d-1}} \int_{\mathbb{R}} d x_{1} \prod_{i=2}^{d}\left(\int_{\mathbb{R}}\left|\partial_{i} u\left(x^{i}\right)\right| d y_{i}\right)^{\frac{1}{d-1}} \\
\leq & \left(\int_{\mathbb{R}} \partial_{1} u(x) d x_{1}\right)^{\frac{1}{d-1}} \prod_{i=2}^{d}\left(\int_{\mathbb{R}^{2}}\left|\partial_{i} u\left(x^{i}\right)\right| d x_{1} d y_{i}\right)^{\frac{1}{d-1}} \\
= & \left(\int_{\mathbb{R}} \partial_{1} u(x) d x_{1}\right)^{\frac{1}{d-1}}\left(\int_{\mathbb{R}^{2}}\left|\partial_{2} u(x)\right| d x_{1} d x_{2}\right)^{\frac{1}{d-1}} \times \\
& \prod_{i=3}^{d}\left(\int_{\mathbb{R}^{2}}\left|\partial_{i} u\left(x^{i}\right)\right| d x_{1} d y_{i}\right)^{\frac{1}{d-1}} .
\end{aligned}
$$

Integrating this equation on $x_{2}$ and using Eq. (49.18) once again implies,

$$
\begin{aligned}
& \int_{\mathbb{R}^{2}}|u(x)|^{\frac{d}{d-1}} d x_{1} d x_{2} \\
& \leq\left(\int_{\mathbb{R}^{2}}\left|\partial_{2} u(x)\right| d x_{1} d x_{2}\right)^{\frac{1}{d-1}} \int_{\mathbb{R}} d x_{2}\left(\int_{\mathbb{R}} \partial_{1} u(x) d x_{1}\right)^{\frac{1}{d-1}} \\
& \quad \times \prod_{i=3}^{d}\left(\int_{\mathbb{R}^{2}}\left|\partial_{i} u\left(x^{i}\right)\right| d x_{1} d y_{i}\right)^{\frac{1}{d-1}} \\
& \leq\left(\int_{\mathbb{R}^{2}}\left|\partial_{2} u(x)\right| d x_{1} d x_{2}\right)^{\frac{1}{d-1}}\left(\int_{\mathbb{R}^{2}}\left|\partial_{1} u(x)\right| d x_{1} d x_{2}\right)^{\frac{1}{d-1}} \\
& \quad \times \prod_{i=3}^{d}\left(\int_{\mathbb{R}^{3}}\left|\partial_{i} u\left(x^{i}\right)\right| d x_{1} d x_{2} d y_{i}\right)^{\frac{1}{d-1}} .
\end{aligned}
$$

Continuing this way inductively, one shows

$$
\begin{aligned}
& \int_{\mathbb{R}^{k}}|u(x)|^{\frac{d}{d-1}} d x_{1} d x_{2} \ldots d x_{k} \\
& \leq \prod_{i=1}^{k}\left(\int_{\mathbb{R}^{k}}\left|\partial_{i} u(x)\right| d x_{1} d x_{2} \ldots d x_{k}\right)^{\frac{1}{d-1}} \\
& \times \prod_{i=k+1}^{d}\left(\int_{\mathbb{R}^{3}}\left|\partial_{i} u\left(x^{i}\right)\right| d x_{1} d x_{2} \ldots d x_{k} d y_{k+1}\right)^{\frac{1}{d-1}}
\end{aligned}
$$

and in particular when $k=d$,

$$
\begin{align*}
\int_{\mathbb{R}^{d}}|u(x)|^{\frac{d}{d-1}} d x & \leq\left(\frac{1}{2}\right)^{\frac{d}{d-1}} \prod_{i=1}^{d}\left(\int_{\mathbb{R}^{d}}\left|\partial_{i} u(x)\right| d x_{1} d x_{2} \ldots d x_{d}\right)^{\frac{1}{d-1}}  \tag{49.19}\\
& \leq \prod_{i=1}^{d}\left(\int_{\mathbb{R}^{d}}|\nabla u(x)| d x\right)^{\frac{1}{d-1}}=\left(\int_{\mathbb{R}^{d}}|\nabla u(x)| d x\right)^{\frac{d}{d-1}} \tag{49.22}
\end{align*}
$$

For this choice of $s, s 1^{*}=p^{*}=(s-1) q$ and Eq. (49.21) becomes

$$
\left[\int_{\mathbb{R}^{d}}|u|^{p^{*}} d m\right]^{1 / 1^{*}} \leq \frac{s}{\sqrt{d}}\|\nabla u\|_{L^{p}} \cdot\left[\int_{\mathbb{R}^{d}}|u|^{p^{*}} d m\right]^{1 / q}
$$

Since

$$
\begin{aligned}
\frac{1}{1^{*}}-\frac{1}{q} & =\frac{d-1}{d}-\frac{p-1}{p}=\frac{p(d-1)-d(p-1)}{d p} \\
& =\frac{d-p}{p d}=\frac{1}{p^{*}}
\end{aligned}
$$

Eq. (49.22) implies Eq. (49.20).
Corollary 49.15. Suppose $\Omega \subset \mathbb{R}^{d}$ is bounded open set with $C^{1}$-boundary, then for all $p \in[1, d)$ and $1 \leq q \leq p^{*}$ there exists $C=C(\Omega, p, q)$ such that

$$
\|u\|_{L^{q}(\Omega)} \leq C\|u\|_{W^{1, p}(\Omega)}
$$

Proof. Let $U$ be a precompact open subset of $\mathbb{R}^{d}$ such that $\bar{\Omega} \subset U$ and $E: W^{1, p}(\Omega) \rightarrow W^{1, p}\left(\mathbb{R}^{d}\right)$ be an extension operator as in Theorem 48.35. Then for $u \in C^{1}(\bar{\Omega}) \cap W^{1, p}(\Omega)$,

$$
\|u\|_{L^{p^{*}}(\Omega)} \leq C\|E u\|_{L^{p^{*}}\left(\mathbb{R}^{d}\right)} \leq C\|\nabla(E u)\|_{L^{p}\left(\mathbb{R}^{d}\right)} \leq C\|u\|_{W^{1, p}(\Omega)}
$$ i.e.

$$
\begin{equation*}
\|u\|_{L^{p^{*}}(\Omega)} \leq C\|u\|_{W^{1, p}(\Omega)} \tag{49.23}
\end{equation*}
$$

Since $C^{1}(\bar{\Omega})$ is dense in $W^{1, p}(\Omega)$, Eq. (49.23) holds for all $u \in W^{1, p}(\Omega)$. Finally for all $1 \leq q<p^{*}$,

$$
\begin{aligned}
\|u\|_{L^{q}(\Omega)} & \leq\|u\|_{L^{P^{*}}(\Omega)} \cdot\|1\|_{L^{r}(\Omega)}=\|u\|_{L^{p^{*}}}(\lambda(\Omega))^{\frac{1}{r}} \\
& \leq C(\lambda(\Omega))^{\frac{1}{r}}\|u\|_{W^{1, p}(\Omega)}
\end{aligned}
$$

where $\frac{1}{r}+\frac{1}{p^{*}}=\frac{1}{q}$.

### 49.4 Sobolev Embedding Theorems Summary

Let us summarize what we have proved up to this point in the following theorem.

Theorem 49.16. Let $p \in[1, \infty]$ and $u \in W^{1, p}\left(\mathbb{R}^{d}\right)$. Then

1. Morrey's Inequality. If $p>d$, then $W^{1, p} \hookrightarrow C^{0,1-\frac{d}{p}}$ and

$$
\left\|u^{*}\right\|_{C^{0,1-\frac{d}{p}}\left(\mathbb{R}^{d}\right)} \leq C\|u\|_{W^{1, p}\left(\mathbb{R}^{d}\right)}
$$

2. When $p=d$ there is an $L^{\infty}$ - like space called $B M O$ (which is not defined in these notes) such that $W^{1, p} \hookrightarrow B M O$.
3. GNS Inequality. If $1 \leq p<d$, then $W^{1, p} \hookrightarrow L^{p^{*}}$

$$
\|u\|_{L^{p^{*}}} \leq d^{-1 / 2} \frac{p(d-1)}{d-p}\|\nabla u\|_{L^{p}}
$$

where $p^{*}=\frac{d p}{d-p}$ or equivalently $\frac{1}{p^{*}}=\frac{1}{p}-\frac{1}{d}$.
Our next goal is write out the embedding theorems for $W^{k, p}(\Omega)$ for general $k$ and $p$.

Notation 49.17 Given a number $s \geq 0$, let

$$
s_{+}=\left\{\begin{array}{cc}
s & \text { if } n \notin \mathbb{N}_{0} \\
s+\delta & \text { if } n \in \mathbb{N}_{0}
\end{array}\right.
$$

where $\delta>0$ is some arbitrarily small number. When $s=k+\alpha$ with $k \in \mathbb{N}_{0}$ and $0 \leq \alpha<1$ we will write $C^{k, \alpha}(\Omega)$ simply as $C^{s}(\Omega)$. Warning, although $C^{k, 1}(\Omega) \subset C^{k+1}(\Omega)$ it is not true that $C^{k, 1}(\Omega)=C^{k+1}(\Omega)$.

Theorem 49.18 (Sobolev Embedding Theorems). Suppose $\Omega=\mathbb{R}^{d}$ or $\Omega \subset \mathbb{R}^{d}$ is bounded open set with $C^{1}$-boundary, $p \in[1, \infty), k, l \in \mathbb{N}$ with $l \leq k$.

1. If $p<d / l$ then $W^{k, p}(\Omega) \hookrightarrow W^{k-l, q}(\Omega)$ provided $q:=\frac{d p}{d-p l}$, i.e. $q$ solves

$$
\frac{1}{q}=\frac{1}{p}-\frac{l}{d}>0
$$

and there is a constant $C<\infty$ such that

$$
\|u\|_{W^{k-l, q}(\Omega)} \leq C\|u\|_{W^{k, p}(\Omega)} \text { for all } u \in W^{k, p}(\Omega)
$$

2. If $p>d / k$, then $W^{k, p}(\Omega) \hookrightarrow C^{k-(d / p)+}(\Omega)$ and there is a constant $C<$ $\infty$ such that

$$
\|u\|_{C^{k-(d / p)+}(\Omega)} \leq C\|u\|_{W^{k, p}(\Omega)} \text { for all } u \in W^{k, p}(\Omega)
$$

Proof. 1. $(p<d / l)$ If $u \in W^{k, p}(\Omega)$, then $\partial^{\alpha} u \in W^{1, p}(\Omega)$ for all $|\alpha| \leq$ $k-1$. Hence by Corollary 49.15, $\partial^{\alpha} u \in L^{p^{*}}(\Omega)$ for all $|\alpha| \leq k-1$ and therefore $W^{k, p}(\Omega) \hookrightarrow W^{k-1, p^{*}}(\Omega)$ and there exists a constant $C_{1}$ such that

$$
\begin{equation*}
\|u\|_{W^{k-1, p_{1}}(\Omega)} \leq C\|u\|_{W^{k, p}(\Omega)} \text { for all } u \in W^{k, p}(\Omega) . \tag{49.24}
\end{equation*}
$$

Define $p_{j}$ inductively by, $p_{1}:=p^{*}$ and $p_{j}:=p_{j-1}^{*}$. Since $\frac{1}{p_{j}}=\frac{1}{p_{j-1}}-\frac{1}{d}$ it is easily checked that $\frac{1}{p_{l}}=\frac{1}{p}-\frac{l}{d}>0$ since $p<d / l$. Hence using Eq. (49.24) repeatedly we learn that the following inclusion maps are all bounded:

$$
W^{k, p}(\Omega) \hookrightarrow W^{k-1, p_{1}}(\Omega) \hookrightarrow W^{k-2, p_{2}}(\Omega) \ldots \hookrightarrow W^{k-l, p_{l}}(\Omega)
$$

This proves the first item of the theorem. The following lemmas will be used in the proof of item 2 .

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Lemma 49.19. Suppose $j \in \mathbb{N}$ and $p \geq d$ and $j>d / p$ (i.e. $j \geq 1$ if $p>d$ and $j \geq 2$ if $p=d$ ) then

$$
W^{j, p}(\Omega) \hookrightarrow C^{j-(d / p)_{+}}(\Omega)
$$

and there is a constant $C<\infty$ such that

$$
\begin{equation*}
\|u\|_{C^{j-(d / p)+(\Omega)}} \leq C\|u\|_{W^{j, p}(\Omega)} . \tag{49.25}
\end{equation*}
$$

Proof. By the usual methods, it suffices to show that the estimate in Eq. (49.25) holds for all $u \in C^{\infty}(\bar{\Omega})$.

For $p>d$ and $|\alpha| \leq j-1$,

$$
\left\|\partial^{\alpha} u\right\|_{C^{0,1-d / p}(\Omega)} \leq C\left\|\partial^{\alpha} u\right\|_{W^{1, p}(\Omega)} \leq C\|u\|_{W^{j, p}(\Omega)}
$$

and hence

$$
\|u\|_{C^{j-d / p}(\Omega)}:=\|u\|_{C^{j-1,1-d / p}(\Omega)} \leq C\|u\|_{W^{j, p}(\Omega)}
$$

which is Eq. (49.25).
When $p=d$ (so now $j \geq 2$ ), choose $q \in(1, d)$ be close to $d$ so that $j>d / q$ and $q^{*}=\frac{q d}{d-q}>d$. Then

$$
W^{j, d}(\Omega) \hookrightarrow W^{j, q}(\Omega) \hookrightarrow W^{j-1, q^{*}}(\Omega) \hookrightarrow C^{j-2,1-d / q^{*}}(\Omega) .
$$

Since $d / q^{*} \downarrow 0$ as $q \uparrow d$, we conclude that $W^{j, d}(\Omega) \hookrightarrow C^{j-2, \alpha}(\Omega)$ for any $\alpha \in(0,1)$ which we summarizes by writing

$$
W^{j, d}(\Omega) \hookrightarrow C^{j-(d / d)}(\Omega) .
$$

Proof. Continuation of the proof of Theorem 49.18. Item 2., $(p>d / k)$. If $p \geq d$, the result follows from Lemma 49.19. So nos suppose that $d>p>d / k$ and choose the largest $l$ such that $1 \leq l<k$ and $d / l>p$ and let $q=\frac{d p}{d-p l}$, i.e. $q$ solves $q \geq d$ and

$$
\frac{1}{q}=\frac{1}{p}-\frac{l}{d} \text { or } \frac{d}{q}=\frac{d}{p}-l
$$

Then
$W^{k, p}(\Omega) \hookrightarrow W^{k-l, q}(\Omega) \hookrightarrow C^{k-l-(d / q)_{+}}(\Omega)=C^{k-l-\left(\frac{d}{p}-l\right)_{+}}(\Omega)=C^{k-\left(\frac{d}{p}\right)_{+}}(\Omega)$
as desired.
Remark 49.20 (Rule of thumb.). Assign the "degrees of regularity" $k-(d / p)_{+}$ to the space $W^{k, p}$ and $k+\alpha$ to the space $C^{k, \alpha}$. If

$$
X, Y \in\left\{W^{k, p}: k \in \mathbb{N}_{0}, p \in[1, \infty]\right\} \cup\left\{C^{k, \alpha}: k \in \mathbb{N}_{0}, \alpha \in[0,1]\right\}
$$

with $\operatorname{deg}_{\text {reg }}(X) \geq \operatorname{deg}_{\text {reg }}(Y)$, then $X \hookrightarrow Y$.

Example 49.21. 1. $W^{k, p} \hookrightarrow W^{k-\ell, q}$ iff $k-\frac{d}{p} \geq k-\ell-\frac{d}{q}$ iff $\ell \geq \frac{d}{p}-\frac{d}{q}$ iff $\frac{1}{q} \geq \frac{1}{p}-\frac{\ell}{d}$.
2. $W^{k, p} \subset C^{0, \alpha}$ iff $k-\left(\frac{d}{p}\right)_{+} \geq \alpha$.

### 49.5 Compactness Theorems

Lemma 49.22. Suppose $K_{m}: X \rightarrow Y$ are compact operators and $\| K-$ $K_{m} \|_{L(X, Y)} \rightarrow 0$ as $n \rightarrow \infty$ then $K$ is compact.

Proof. Let $\left\{x_{n}\right\}_{n=1}^{\infty} \subset X$ be given such that $\left\|x_{n}\right\| \leq 1$. By Cantor's diagonalization scheme we may choose $\left\{x_{n}^{\prime}\right\} \subset\left\{x_{n}\right\}$ such that $y_{m}:=\lim _{n \rightarrow \infty} K_{m} x_{n}^{\prime} \in$ $Y$ exists for all $m$. Hence

$$
\begin{aligned}
\left\|K x_{n}^{\prime}-K x_{\ell}^{\prime}\right\| & =\left\|K\left(x_{n}^{\prime}-x_{\ell}^{\prime}\right)\right\| \leq\left\|K\left(x_{n}^{\prime}-x_{\ell}^{\prime}\right)\right\| \\
& \leq\left\|K-K_{m}\right\|\left\|x_{n}^{\prime}-x_{\ell}^{\prime}\right\|+\left\|K_{m}\left(x_{n}^{\prime}-x_{\ell}^{\prime}\right)\right\| \\
& \leq\left\|K-K_{m}\right\|+\left\|K_{m}\left(x_{n}^{\prime}-x_{\ell}^{\prime}\right)\right\|
\end{aligned}
$$

and therefore,

$$
\limsup _{l, n \rightarrow \infty}\left\|K x_{n}^{\prime}-K x_{\ell}^{\prime}\right\| \leq\left\|K-K_{m}\right\| \rightarrow 0 \text { as } m \rightarrow \infty
$$

Lemma 49.23. Let $\eta \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right), C_{\eta} f=\eta * f, \Omega \subset \mathbb{R}^{d}$ be a bounded open set with $C^{1}$-boundary, $U$ be an open precompact subset of $\mathbb{R}^{d}$ such that $\bar{\Omega} \subset U$ and $E: W^{1,1}(\Omega) \rightarrow W^{1,1}\left(\mathbb{R}^{d}\right)$ be an extension operator as in Theorem 48.35 . Then to every bounded sequence $\left\{\tilde{u}_{n}\right\}_{n=1}^{\infty} \subset W^{1,1}(\Omega)$ there has a subsequence $\left\{u_{n}^{\prime}\right\}_{n=1}^{\infty}$ such that $C_{\eta} E u_{n}^{\prime}$ is uniformly convergent to a function in $C_{c}\left(\mathbb{R}^{d}\right)$.

Proof. Let $u_{n}:=E \tilde{u}_{n}$ and $C:=\sup \left\|u_{n}\right\|_{W^{1,1}\left(\mathbb{R}^{d}\right)}$ which is finite by assumption. So $\left\{u_{n}\right\}_{n=1}^{\infty} \subset W^{1,1}\left(\mathbb{R}^{d}\right)$ is a bounded sequence such that $\operatorname{supp}\left(u_{n}\right) \subset U \subset \bar{U} \sqsubset \sqsubset \mathbb{R}^{d}$ for all $n$. Since $\eta$ is compactly supported there exists a precompact open set $V$ such that $\bar{U} \subset V$ and $v_{n}:=\eta * u_{n} \in C_{c}^{\infty}(V) \subset$ $C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ for all $n$. Since,

$$
\begin{aligned}
\left\|v_{n}\right\|_{L^{\infty}} & \leq\|\eta\|_{L^{\infty}}\left\|u_{n}\right\|_{L^{1}} \leq\|\eta\|_{L^{\infty}}\left\|u_{n}\right\|_{L^{1}} \leq C\|\eta\|_{L^{\infty}} \text { and } \\
\left\|D v_{n}\right\|_{L^{\infty}} & =\left\|\eta * D u_{n}\right\|_{L^{\infty}} \leq\|\eta\|_{L^{\infty}}\left\|D u_{n}\right\|_{L^{1}} \leq C\|\eta\|_{L^{\infty}},
\end{aligned}
$$

it follows by the Arzela-Ascoli theorem that $\left\{v_{n}\right\}_{n=1}^{\infty}$ has a uniformly convergent subsequence.
Lemma 49.24. Let $\eta \in C_{c}^{\infty}(B(0,1),[0, \infty))$ such that $\int_{\mathbb{R}^{d}} \eta d m=1, \eta_{m}(x)=$ $m^{n} \eta(m x)$ and $K_{m} u=\left.\left(C_{\eta_{m}} E u\right)\right|_{\Omega}$. Then for all $p \in[1, d)$ and $q \in\left[1, p^{*}\right)$,

$$
\lim _{m \rightarrow \infty}\left\|K_{m}-i\right\|_{B\left(W^{1, p}(\Omega), L^{q}(\Omega)\right)}=0
$$

where $i: W^{1, p}(\Omega) \rightarrow L^{q}(\Omega)$ is the inclusion map.

Proof. For $u \in C_{c}^{1}(U)$ let $v_{m}:=\eta_{m} * u-u$, then

$$
\begin{aligned}
\left|v_{m}(x)\right| & \leq\left|\eta_{m} * u(x)-u(x)\right|=\left|\int_{\mathbb{R}^{d}} \eta_{m}(y)(u(x-y)-u(x)) d y\right| \\
& =\left|\int_{\mathbb{R}^{d}} \eta(y)\left[u\left(x-\frac{y}{m}\right)-u(x)\right] d y\right| \\
& \leq \frac{1}{m} \int_{\mathbb{R}^{d}} d y|y| \eta(y) \int_{0}^{1} d t\left|\nabla u\left(x-t \frac{y}{m}\right)\right|
\end{aligned}
$$

and so by Minikowski's inequality for integrals,

$$
\begin{align*}
\left\|v_{m}\right\|_{L^{r}} & \leq \frac{1}{m} \int_{\mathbb{R}^{d}} d y|y| \eta(y) \int_{0}^{1} d t\left\|\nabla u\left(\cdot-t \frac{y}{m}\right)\right\|_{L^{r}} \\
& \leq \frac{1}{m}\left(\int_{\mathbb{R}^{d}}|y| \eta(y) d y\right)\|\nabla u\|_{L^{r}} \leq \frac{1}{m}\|u\|_{W^{1, r}\left(\mathbb{R}^{d}\right)} \tag{49.26}
\end{align*}
$$

By the interpolation inequality in Corollary 10.25, Theorem 49.14 and Eq. (49.26) with $r=1$,

$$
\begin{aligned}
\left\|v_{m}\right\|_{L^{d}} & \leq\left\|v_{m}\right\|_{L^{1}}^{\lambda}\left\|v_{m}\right\|_{L^{p^{*}}}^{1-\lambda} \\
& \leq \frac{1}{m^{\lambda}}\|u\|_{W^{1,1}\left(\mathbb{R}^{d}\right)}^{\lambda}\left[d^{-1 / 2} \frac{p(d-1)}{d-p}\left\|\nabla v_{m}\right\|_{L^{p}}\right]^{1-\lambda} \\
& \leq C m^{-\lambda}\|u\|_{W^{1,1}\left(\mathbb{R}^{d}\right)}^{\lambda}\left\|v_{m}\right\|_{W^{1, p}\left(\mathbb{R}^{d}\right)}^{1-\lambda} \\
& \leq C m^{-\lambda}\|u\|_{W^{1,1}\left(\mathbb{R}^{d}\right)}^{\lambda}\left\|v_{m}\right\|_{W^{1, p}\left(\mathbb{R}^{d}\right)}^{1-\lambda} \\
& \leq C(p,|U|) m^{-\lambda}\|u\|_{W^{1, p}\left(\mathbb{R}^{d}\right)}^{\lambda}\left\|v_{m}\right\|_{W^{1, p}\left(\mathbb{R}^{d}\right)}^{1-\lambda}
\end{aligned}
$$

where $\lambda \in(0,1)$ is determined by

$$
\frac{1}{q}=\frac{\lambda}{1}+\frac{1-\lambda}{p^{*}}=\lambda\left(1-\frac{1}{p^{*}}\right)+\frac{1}{p^{*}}
$$

Now using Proposition 11.12,

$$
\begin{aligned}
\left\|v_{m}\right\|_{W^{1, p}\left(\mathbb{R}^{d}\right)} & =\left\|\eta_{m} * u-u\right\|_{W^{1, p}\left(\mathbb{R}^{d}\right)} \\
& \leq\left\|\eta_{m} * u\right\|_{W^{1, p}\left(\mathbb{R}^{d}\right)}+\|u\|_{W^{1, p}\left(\mathbb{R}^{d}\right)} \leq 2\|u\|_{W^{1, p}\left(\mathbb{R}^{d}\right)} .
\end{aligned}
$$

Putting this all together shows

$$
\begin{aligned}
\left\|K_{m} u-u\right\|_{L^{q}(\Omega)} & \leq\left\|K_{m} u-E u\right\|_{L^{q}}=\left\|v_{m}\right\|_{L^{q}} \\
& \leq C(p,|U|) m^{-\lambda}\|u\|_{W^{1, p}\left(\mathbb{R}^{d}\right)}^{\lambda}\left(2\|u\|_{W^{1, p}\left(\mathbb{R}^{d}\right)}\right)^{1-\lambda} \\
& \leq C(p,|U|) m^{-\lambda}\|E u\|_{W^{1, p}\left(\mathbb{R}^{d}\right)} \\
& \leq C(p,|U|) m^{-\lambda}\|u\|_{W^{1, p}(\Omega)}
\end{aligned}
$$

from which it follows that

$$
\left\|K_{m}-i\right\|_{B\left(W^{1, p}(\Omega), L^{q}(\Omega)\right)} \leq C m^{-\lambda} \rightarrow 0 \text { as } m \rightarrow \infty
$$

Theorem 49.25 (Rellich - Kondrachov Compactness Theorem). Suppose $\Omega \subset \mathbb{R}^{d}$ is a precompact open subset with $C^{1}$-boundary, $p \in[1, d)$ and $1 \leq q<p^{*}$ then $W^{1, p}(\Omega)$ is compactly embedded in $L^{q}(\Omega)$.

Proof. If $\left\{u_{n}\right\}_{n=1}^{\infty}$ is contained in the unit ball in $W^{1, p}(\Omega)$, then by Lemma $49.23\left\{K_{m} u_{n}\right\}_{n=1}^{\infty}$ has a uniformly convergent subsequence and hence is convergent in $L^{q}(\Omega)$. This shows $K_{m}: W^{1, p}(\Omega) \rightarrow L^{q}(\Omega)$ is compact for every $m$. By Lemma 49.24, $K_{m} \rightarrow i$ in the $L\left(W^{1, p}(\Omega), L^{q}(\Omega)\right)$ - norm and so by Lemma $49.22 i: W^{1, p}(\Omega) \rightarrow L^{q}(\Omega)$ is compact.

Corollary 49.26. The inclusion of $W^{k, p}(\Omega)$ into $W^{k-\ell, q}(\Omega)$ is compact provided $l \geq 1$ and $\frac{1}{q}>\frac{1}{p}-\frac{l}{d}=\frac{d-p l}{d p}>0$, i.e. $q<\frac{d p}{d-p l}$.

Proof. Case (i) Suppose $\ell=1, q \in\left[1, p^{*}\right)$ and $\left\{u_{n}\right\}_{n=1}^{\infty} \subset W^{k, p}(\Omega)$ is bounded. Then $\left\{\partial^{\alpha} u_{n}\right\}_{n=1}^{\infty} \subset W^{1, p}(\Omega)$ is bounded for all $|\alpha| \leq k-1$ and therefore there exist a subsequence $\left\{\tilde{u}_{n}\right\}_{n=1}^{\infty} \subset\left\{u_{n}\right\}_{n=1}^{\infty}$ such that $\partial^{\alpha} \tilde{u}_{n}$ is convergent in $L^{q}(\Omega)$ for all $|\alpha| \leq k-1$. This shows that $\left\{\tilde{u}_{n}\right\}$ is $W^{k-1, q}(\Omega)$ - convergent and so proves this case.

Case (ii) $\ell>1$. Let $\tilde{p}$ be defined so that $\frac{1}{\tilde{p}}=\frac{1}{p}-\frac{\ell-1}{d}$. Then

$$
W^{k, p}(\Omega) \subset W^{k-\ell+1, \tilde{p}}(\Omega) \subset \subset W^{k-\ell, q}(\Omega)
$$

and therefore $W^{k, p}(\Omega) \subset \subset W^{k-\ell, q}(\Omega)$.
Example 49.27. It is necessary to assume that The inclusion of $L^{2}([0,1]) \hookrightarrow$ $L^{1}([0,1])$ is continuous (in fact a contraction) but not compact. To see this, take $\left\{u_{n}\right\}_{n=1}^{\infty}$ to be the Haar basis for $L^{2}$. Then $u_{n} \rightarrow 0$ weakly in both $L^{2}$ and $L^{1}$ so if $\left\{u_{n}\right\}_{n=1}^{\infty}$ were to have a convergent subsequence the limit would have to be $0 \in L^{1}$. On the other hand, since $\left|u_{n}\right|=1,\left\|u_{n}\right\|_{2}=\left\|u_{n}\right\|_{1}=1$ and any subsequential limit would have to have norm one and in particular not be 0.

Lemma 49.28. Let $\Omega$ be a precompact open set such that $\bar{\Omega}$ is a manifold with $C^{1}$ - boundary. Then for all $p \in[1, \infty), W^{1, p}(\Omega)$ is compactly embedded in $L^{p}(\Omega)$. Moreover if $p>d$ and $0 \leq \beta<1-\frac{d}{p}$, then $W^{1, p}(\Omega)$ is compactly embedded in $C^{0, \beta}(\Omega)$. In particular, $W^{1, p}(\Omega) \subset \subset L^{\infty}(\Omega)$ for all $d<p \leq \infty$.

Proof. Case 1, $p \in[1, d)$. By Theorem 49.25, $W^{1, p}(\Omega) \subset \subset L^{q}(\Omega)$ for all $1 \leq q<p^{*}$. Since $p^{*}>p$ we may choose $q=p$ to learn $W^{1, p}(\Omega) \subset \subset L^{p}(\Omega)$. Case 2, $p \in[d, \infty)$. For any $p_{0} \in[1, d)$, we have

$$
W^{1, p}(\Omega) \hookrightarrow W^{1, p_{0}}(\Omega) \subset \subset L^{p_{0}^{*}}(\Omega)
$$

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Since $p_{0}^{*}=\frac{p_{0} d}{d-p_{o}} \uparrow \infty$ as $p_{0} \uparrow d$, we see that $W^{1, p}(\Omega) \subset \subset L^{q}(\Omega)$ for all $q<\infty$. Moreover by Morrey's inequality (Corollary 49.8) and Proposition5.13 we have $W^{1, p}(\Omega) \hookrightarrow C^{0,1-\frac{d}{p}}(\Omega) \subset \subset C^{0, \beta}(\Omega)$ which completes the proof.

Remark 49.29. Similar proofs may be given to show $W^{k, p} \subset \subset C^{k-\frac{d}{p}-\delta}$ for all $\delta>0$ provided $k-\frac{d}{p}>0$ and $k-\frac{d}{p}-\delta>0$.

Lemma 49.30 (Poincaré Lemma). Assume $1 \leq p \leq \infty, \Omega$ is a precompact open connected subset of $\mathbb{R}^{d}$ such that $\bar{\Omega}$ is a manifold with $C^{1}$-boundary. Then exist $C=C(\Omega, \rho)$ such that

$$
\begin{equation*}
\left\|u-u_{\Omega}\right\|_{L^{p}(\Omega)} \leq C\|\nabla u\|_{L^{p}(\Omega)} \text { for all } u \in W^{1, p}(\Omega), \tag{49.27}
\end{equation*}
$$

where $u_{\Omega}:=f_{\Omega} u d m$ is the average of $u$ on $\Omega$ as in Notation 49.3.
Proof. For sake of contradiction suppose there is no $C<\infty$ such that Eq. (49.27) holds. Then there exists a sequence $\left\{u_{n}\right\}_{n=1}^{\infty} \subset W^{1, p}(\Omega)$ such that

$$
\left\|u_{n}-\left(u_{n}\right)_{\Omega}\right\|_{L^{p}(\Omega)}>n\left\|\nabla u_{n}\right\|_{L^{p}(\Omega)} \text { for all } n .
$$

Let

$$
u_{n}:=\frac{u_{n}-\left(u_{n}\right)_{\Omega}}{\left\|u_{n}-\left(u_{n}\right)_{\Omega}\right\|_{L^{p}(\Omega)}} .
$$

Then $u_{n} \in W^{1, p}(\Omega),\left(u_{n}\right)_{\Omega}=0,\left\|u_{n}\right\|_{L^{p}(\Omega)}=1$ and $1=\left\|u_{n}\right\|_{L^{p}(\Omega)}>$ $n\left\|\nabla u_{n}\right\|_{L^{p}(\Omega)}$ for all $n$. Therefore $\left\|\nabla u_{n}\right\|_{L^{p}(\Omega)}<\frac{1}{n}$ and in particular $\sup \left\|u_{n}\right\|_{W^{1, p}(\Omega)}<\infty$ and hence by passing to a subsequence if necessary ${ }_{\text {there exists }}^{n} u \in L^{p}(\Omega)$ such that $u_{n} \rightarrow u$ in $L^{p}(\Omega)$. Since $\nabla u_{n} \rightarrow 0$ in $L^{p}(\Omega)$, it follows that $u_{n}$ is convergent in $W^{1, p}(\Omega)$ and hence $u \in W^{1, p}(\Omega)$ and $\nabla u=\lim _{n \rightarrow \infty} \nabla u_{n}=0$ in $L^{p}(\Omega)$. Since $\nabla u=0, u \in W^{k, p}(\Omega)$ for all $k \in \mathbb{N}$ and hence $u \in C^{\infty}(\Omega)$ and $\nabla u=0$ and $\Omega$ is connected implies $u$ is constant. Since $0=\lim _{n \rightarrow \infty}\left(u_{n}\right)_{\Omega}=u_{\Omega}$ we must have $u \equiv 0$ which is clearly impossible since $\|u\|_{L^{p}(\Omega)}=\lim _{n \rightarrow \infty}\left\|u_{n}\right\|_{L^{p}(\Omega)}=1$.

Theorem 49.31 (Poincaré Lemma). Let $\Omega$ be a precompact open subset of $\mathbb{R}^{d}$ and $p \in[1, \infty)$. Then

$$
\begin{equation*}
\|u\|_{L^{p}} \leq \operatorname{diam}(\Omega)\|\nabla u\|_{L^{p}} \text { for all } u \in W_{0}^{1, p}(\Omega) \tag{49.28}
\end{equation*}
$$

Proof. Let $\operatorname{diam}(\Omega)=M$. By translating $\Omega$ if necessary we may assume $\Omega \subset[-M, M]^{d}$. For $1 \leq p<\infty$ we may assume $u \in C_{c}^{\infty}(\Omega)$ since $C_{c}^{\infty}(\Omega)$ is dense in $W_{0}^{1, p}(\Omega)$. Then by the fundamental theorem of calculus,

$$
\begin{aligned}
|u(x)| & =\frac{1}{2}\left|\int_{-M}^{x_{1}} \partial_{1} u\left(y, x_{2}, \ldots, x_{d}\right) d y-\int_{x_{1}}^{M} \partial_{1} u\left(y, x_{2}, \ldots, x_{d}\right) d y\right| \\
& \leq \frac{1}{2} \int_{-M}^{M}\left|\partial_{1} u\left(y, x_{2}, \ldots, x_{d}\right)\right| d y=M \int_{-M}^{M}\left|\partial_{1} u\left(y, x_{2}, \ldots, x_{d}\right)\right| \frac{d y}{2 M}
\end{aligned}
$$

and hence by Jensen's inequality,

$$
|u(x)|^{p} \leq M^{p} \int_{-M}^{M}\left|\partial_{1} u\left(y, x_{2}, \ldots, x_{d}\right)\right|^{p} \frac{d y}{2 M}
$$

Integrating this equation over $x$ implies,

$$
\|u\|_{L^{p}}^{p} \leq M^{p} \int_{\Omega}\left|\partial_{1} u(x)\right|^{p} d x \leq M^{p} \int_{\Omega}|\nabla u(x)|^{p} d x
$$

which gives Eq. (49.28).

### 49.6 Fourier Transform Method

See $L^{2}$ - Sobolev spaces for another proof of the following theorem.
Theorem 49.32. Suppose $s>t \geq 0,\left\{u_{n}\right\}_{n=1}^{\infty}$ is a bounded sequence (say by 1) in $H^{s}\left(\mathbb{R}^{d}\right)$ such that $K=\overline{U_{n} \operatorname{supp}\left(u_{n}\right)} \sqsubset \sqsubset \mathbb{R}^{d}$. Then there exist a subsequence $\left\{v_{n}\right\}_{n=1}^{\infty} \subset\left\{u_{n}\right\}_{n=1}^{\infty}$ which is convergent in $H^{t}\left(\mathbb{R}^{d}\right)$.

## Proof. Since

$$
\begin{aligned}
\left|\partial_{\xi}^{\alpha} \hat{u}_{n}(\xi)\right| & =\left|\partial_{\xi}^{\alpha} \int_{\mathbb{R}^{d}} e^{-i \xi \cdot x} u_{n}(x) d x\right|=\left|\int_{\mathbb{R}^{d}}(-i x)^{\alpha} e^{-i \xi \cdot x} u_{n}(x) d x\right| \\
& \leq\left\|x^{\alpha}\right\|_{L^{2}(K)}\left\|u_{n}\right\|_{L^{2}} \leq C_{\alpha}\left\|u_{n}\right\|_{H^{s}\left(\mathbb{R}^{d}\right)} \leq C_{\alpha}
\end{aligned}
$$

$\hat{u}_{n}$ and all of it's derivatives are uniformly bounded. By the Arzela-Ascoli theorem and Cantor's Diagonalization argument, there exists a subsequence $\left\{v_{n}\right\}_{n=1}^{\infty} \subset\left\{u_{n}\right\}_{n=1}^{\infty}$ such that $\hat{v}_{n}$ and all of its derivatives converge uniformly on compact subsets in $\xi$-space. If $\hat{v}(\xi):=\lim _{n \rightarrow \infty} \hat{v}_{n}(\xi)$, then by the dominated convergence theorem,

$$
\begin{aligned}
\int_{|\xi| \leq R}\left(1+|\xi|^{2}\right)^{s}|\hat{v}(\xi)|^{2} d \xi & =\lim _{n \rightarrow \infty} \int_{|\xi| \leq R}\left(1+|\xi|^{2}\right)^{s}\left|\hat{v}_{n}(\xi)\right|^{2} d \xi \\
& \leq \limsup _{n \rightarrow \infty}\left\|v_{n}\right\|_{H^{s}\left(\mathbb{R}^{d}\right)}^{2} \leq 1
\end{aligned}
$$

Since $R$ is arbitrary this implies $\hat{v} \in L^{2}\left(\left(1+|\xi|^{2}\right)^{s} d \xi\right)$ and $\|v\|_{H^{s}\left(\mathbb{R}^{d}\right)} \leq 1$. Set $g_{n}:=v-v_{n}$ while $v=\mathcal{F}^{-1} \hat{v}$. Then $\left\{g_{n}\right\}_{n=1}^{\infty} \subset H^{s}\left(\mathbb{R}^{d}\right)$ and we wish to show $g_{n} \rightarrow 0$ in $H^{t}\left(\mathbb{R}^{d}\right)$. Let $d \mu_{t}(\xi)=\left(1+|\xi|^{2}\right)^{t} d \xi$, then for any $R<\infty$,

$$
\begin{aligned}
& \left\|g_{n}\right\|_{H^{t}}^{2}=\int\left|\hat{g}(\xi)-\hat{g}_{n}(\xi)\right|^{2} d \mu_{t}(\xi) \\
& =\int_{|\xi| \leq R}\left|\hat{g}(\xi)-\hat{g}_{n}(\xi)\right|^{2} d \mu_{t}(\xi)+\int_{|\xi| \geq R}\left|\hat{g}(\xi)-\hat{g}_{n}(\xi)\right|^{2} d \mu_{t}(\xi)
\end{aligned}
$$

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$$
\begin{aligned}
|\Gamma| u(x) & =(-1)^{m} \int_{\gamma} d \omega \int_{0}^{\rho} \partial_{t}^{m}[\tau(t) u(x+t \omega)] \frac{t^{m-1}}{(m-1)!} d t \\
& =\frac{(-1)^{m}}{(m-1)!} \int_{\gamma} d \omega \int_{0}^{\rho} t^{m-d} \partial_{t}^{m}[\tau(t) u(x+t \omega)] t^{d-1} d t \\
& =\frac{(-1)^{m}}{(m-1)!} \int_{\gamma} d \omega \int_{0}^{\rho} t^{m-d} \sum_{k=0}^{m}\binom{m}{k}\left[\tau^{(m-k)}(t)\left(\partial_{\omega}^{k} u\right)(x+t \omega)\right] t^{d-1} d t \\
& =\frac{(-1)^{m}}{(m-1)!} \int_{\gamma} d \omega \int_{0}^{\rho} t^{m-d} \sum_{k=0}^{m}\binom{m}{k} \rho^{k-m}\left[g^{(m-k)}(t)\left(\partial_{\omega}^{k} u\right)(x+t \omega)\right] t^{d-1} \\
& =\frac{(-1)^{m}}{(m-1)!} \sum_{k=0}^{m}\binom{m}{k} \rho^{k-m} \int_{\Gamma_{x, \rho}}|y-x|^{m-d}\left[g^{(m-k)}(|y-x|)\left(\partial_{\frac{k}{k-x}} u\right)(y)\right]
\end{aligned}
$$

and hence
$u(x)=\frac{(-1)^{m}}{|\Gamma|(m-1)!} \sum_{k=0}^{m}\binom{m}{k} \rho^{k-m} \int_{\Gamma_{x, \rho}}|y-x|^{m-d}\left[g^{(m-k)}(|y-x|)\left(\partial^{k-x} u\right)(y)\right]$ and hence by the Hölder's inequality,

$$
|u(x)| \leq C(g) \frac{(-1)^{m}}{|\Gamma|(m-1)!} \sum_{k=0}^{m}\binom{m}{k} \rho^{k-m}\left[\int_{\Gamma_{x, \rho}}|y-x|^{q(m-d)} d y\right]^{1 / q}\left[\int_{\Gamma_{x, \rho}} \left\lvert\,\left(\partial_{\frac{k}{k-t}}\right.\right.\right.
$$

From the same computation as in Eq. (48.4) we find

$$
\begin{aligned}
\int_{\Gamma_{x, \rho}}|y-x|^{q(m-d)} d y & =\sigma(\Gamma) \int_{0}^{\rho} r^{q(m-d)} r^{d-1} d r=\sigma(\Gamma) \frac{\rho^{q(m-d)+d}}{q(m-d)+d} \\
& =\sigma(\Gamma) \frac{\rho^{\frac{p m-d}{p-1}}}{p m-d}(p-1)
\end{aligned}
$$

provided that $p m-d>0$ (i.e. $m>d / p$ ) wherein we have used

$$
q(m-d)+d=\frac{p}{p-1}(m-d)+d=\frac{p(m-d)+d(p-1)}{p-1}=\frac{p m-d}{p-1}
$$

This gives the estimate
$\left[\int_{\Gamma_{x, \rho}}|y-x|^{q(m-d)} d y\right]^{1 / q} \leq\left[\frac{\sigma(\Gamma)(p-1)}{p m-d}\right]^{\frac{p-1}{p}} \rho^{\frac{p m-d}{p}}=\left[\frac{\sigma(\Gamma)(p-1)}{p m-d}\right]^{\frac{p-1}{p}} \rho^{m}$
Thus we have obtained the estimate that

$$
\begin{aligned}
|u(x)| \leq & \frac{C(g)}{|\Gamma|(m-1)!}\left[\frac{\sigma(\Gamma)(p-1)}{p m-d}\right]^{\frac{p-1}{p}} \times \\
& \rho^{m-d / p} \sum_{k=0}^{m}\binom{m}{k} \rho^{k-m}\left\|\partial_{y-x}^{k} u\right\|_{L^{p}\left(\Gamma_{x, p}\right)} .
\end{aligned}
$$

### 49.8 Exercises

Exercise 49.33. Let $a_{i} \geq 0$ and $p_{i} \in[1, \infty)$ for $i=1,2, \ldots, d$ satisfy $\sum_{i=1}^{d} p_{i}^{-1}=1$, then

$$
\prod_{i=1}^{d} a_{i} \leq \sum_{i=1}^{d} \frac{1}{p_{i}} a_{i}^{p_{i}} .
$$

Hint: This may be proved by induction on $d$ making use of Lemma 1.27. Alternatively see Example 10.11, where this is already proved using Jensen's inequality.

Solution 49.34 (49.33). We may assume that $a_{i}>0$, in which case

$$
\prod_{i=1}^{d} a_{i}=e^{\sum_{i=1}^{d} \ln a_{i}}=e^{\sum_{i=1}^{d} \frac{1}{p_{i}} \ln a_{i}^{p_{i}}} \leq \sum_{i=1}^{d} \frac{1}{p_{i}} e^{\ln a_{i}^{p_{i}}}=\sum_{i=1}^{d} \frac{1}{p_{i}} a_{i}^{p_{i}} .
$$

This was already done in Example 10.11

Variable Coefficient Equations

## $2^{\text {nd }}$ order differential operators

Notations 50.1 Let $\Omega$ be a precompact open subset of $\mathbb{R}^{d}, A_{i j}=A_{j i}, A_{i}, A_{0} \in$ $B C^{\infty}(\bar{\Omega})$ for $i, j=1, \ldots, d$,

$$
p(x, \xi):=-\sum_{i, j=1}^{d} A_{i j} \xi_{i} \xi_{j}+\sum_{i=1}^{d} A_{i} \xi_{i}+A_{0}
$$

and

$$
L=p(x, \partial)=-\sum_{i, j=1}^{d} A_{i j} \partial_{i} \partial_{j}+\sum_{i=1}^{d} A_{i} \partial_{i}+A_{0}
$$

We also let

$$
L^{\dagger}=-\sum_{i, j=1}^{d} \partial_{i} \partial_{j} M_{A_{i j}}-\sum_{i=1}^{d} \partial_{i} M_{A_{i}}+A_{0}
$$

Remark 50.2. The operators $L$ and $L^{\dagger}$ have the following properties.

1. The operator $L^{\dagger}$ is the formal adjoint of $L$, i.e.

$$
(L u, v)=\left(u, L^{\dagger} v\right) \text { for all } u, v \in \mathcal{D}(\Omega)=C_{c}^{\infty}(\Omega)
$$

2. We may view $L$ as an operator on $\mathcal{D}^{\prime}(\Omega)$ via the formula $u \in \mathcal{D}^{\prime}(\Omega) \rightarrow$ $L u \in \mathcal{D}^{\prime}(\Omega)$ where

$$
\langle L u, \phi\rangle:=\left\langle u, L^{\dagger} \phi\right\rangle \text { for all } \phi \in C_{c}^{\infty}(\Omega) .
$$

3. The restriction of $L$ to $H^{k+2}(\Omega)$ gives a bounded linear transformation

$$
L: H^{k+2}(\Omega) \rightarrow H^{k}(\Omega) \text { for } k \in \mathbb{N}_{0}
$$

Indeed, $L$ may be written as

$$
L=-\sum_{i, j=1}^{d} M_{A_{i j}} \partial_{i} \partial_{j}+\sum_{i=1}^{d} M_{A_{i}} \partial_{i}+M_{A_{0}}
$$

Now $\partial_{i}: H^{k}(\Omega) \rightarrow H^{k+1}(\Omega)$ is bounded and $M_{\psi}: H^{k}(\Omega) \rightarrow H^{k}(\Omega)$ is bounded where $\psi \in B C^{\infty}(\bar{\Omega})$. Therefore, for $k \in \mathbb{N}_{0}, L: H^{k+2}(\Omega) \rightarrow$ $H^{k}(\Omega)$ is bounded.

Definition 50.3. For $u \in \mathcal{D}^{\prime}(\Omega)$, let

$$
\|u\|_{H^{-1}(\Omega)}:=\sup _{0 \neq \phi \in \mathcal{D}(\Omega)} \frac{|\langle u, \phi\rangle|}{\|\phi\|_{H_{0}^{1}(\Omega)}}
$$

and

$$
H^{-1}(\Omega):=\left\{u \in \mathcal{D}^{\prime}(\Omega):\|u\|_{H^{-1}(\Omega)}<\infty\right\}
$$

Example 50.4. Let $\Omega=\mathbb{R}^{d}$ and $S \subset \Omega$ be the unit sphere in $\mathbb{R}^{d}$. Then define $\sigma \in \mathcal{D}^{\prime}(\Omega)$ by

$$
\langle\sigma, \phi\rangle:=\int_{S} \phi d \sigma
$$

Let us shows that $\sigma \in H^{-1}(\Omega)$. For this let $T: H^{1}(\Omega) \rightarrow L^{2}(S, d \sigma)$ denote the trace operator, i.e. the unique bounded linear operator such that $T \phi=\left.\phi\right|_{S}$ for all $\phi \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$. Since $T$ is bounded,

$$
|\langle\sigma, \phi\rangle| \leq \sigma(S)^{1 / 2}\|T \phi\|_{L^{2}(S)} \leq \sigma(S)^{1 / 2}\|T\|_{L\left(H^{1}(\Omega), L^{2}(S)\right)}\|\phi\|_{H^{-1}(\Omega)}
$$

This shows $\sigma \in H^{-1}(\Omega)$ and $\|\sigma\|_{H^{-1}(\Omega)} \leq \sigma(S)^{1 / 2}\|T\|_{L\left(H^{1}(\Omega), L^{2}(S)\right)}$.
Lemma 50.5. Suppose $\Omega$ is an open subset of $\mathbb{R}^{d}$ such that $\bar{\Omega}$ is a manifold with $C^{0}$ - boundary and $\Omega=\bar{\Omega}^{o}$, then the map $\left.u \in\left[H_{0}^{1}(\Omega)\right]^{*} \rightarrow u\right|_{\mathcal{D}(\Omega)} \in$ $H^{-1}(\Omega)$ is a unitary map of Hilbert spaces.

Proof. By definition $C_{c}^{\infty}(\Omega)$ is dense in $H_{0}^{1}(\Omega)$, and hence it follows that the map $\left.u \in\left[H_{0}^{1}(\Omega)\right]^{*} \rightarrow u\right|_{\mathcal{D}(\Omega)} \in H^{-1}(\Omega)$ is isometric. If $u \in H^{-1}(\Omega)$, it has a unique extension to $H_{0}^{1}(\Omega)={\overline{C_{c}^{\infty}(\Omega)}}^{H^{1}(\Omega)}$ and this provides the inverse map.

If we identify $L^{2}(\Omega)=H^{0}(\Omega)$ with elements of $\mathcal{D}^{\prime}(\Omega)$ via $u \rightarrow(u, \cdot)_{L^{2}(\Omega)}$, then

$$
\mathcal{D}^{\prime}(\Omega) \supset H^{-1}(\Omega) \supset H^{0}(\Omega)=L^{2}(\Omega) \supset H^{1}(\Omega) \supset H^{2}(\Omega) \supset \ldots
$$

Proposition 50.6. The following mapping properties hold:

1. If $\chi \in B C^{1}(\bar{\Omega})$. Then $M_{\chi}: H^{-1}(\Omega) \rightarrow H^{-1}(\Omega)$ is a bounded operator.
2. If $V=\sum_{i=1}^{d} M_{A_{i}} \partial_{i}+M_{A_{0}}$ with $A_{i}, A_{0} \in B C^{1}(\bar{\Omega})$, then $V: L^{2}(\Omega) \rightarrow$ $H^{-1}(\Omega)$ is a bounded operator.
3. The map $L: \mathcal{D}^{\prime}(\Omega) \rightarrow \mathcal{D}^{\prime}(\Omega)$ restricts to a bounded linear map from $H^{1}(\Omega)$ to $H^{-1}(\Omega)$. Also
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Proof. Let us begin by showing $M_{\chi}: H_{0}^{1}(\Omega) \rightarrow H_{0}^{1}(\Omega)$ is a bounded linear map. In order to do this choose $\chi_{n} \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ such that $\chi_{n} \rightarrow \chi$ in $B C^{1}(\bar{\Omega})$. Then for $\phi \in C_{c}^{\infty}(\Omega), \chi_{n} \phi \in C_{c}^{\infty}(\Omega) \subset H_{0}^{1}(\Omega)$ and there is a constant $K<\infty$ such that

$$
\left\|\chi_{n} \phi\right\|_{H_{0}^{1}(\Omega)}^{2} \leq K\left\|\chi_{n}\right\|_{B C^{1}(\bar{\Omega})}\|\phi\|_{H_{0}^{1}(\Omega)}^{2} .
$$

By density this estimate holds for all $\phi \in H_{0}^{1}(\Omega)$ and by replacing $\chi_{n}$ by $\chi_{n}-\chi_{m}$ we also learn that

$$
\left\|\left(\chi_{n}-\chi_{m}\right) \phi\right\|_{H_{0}^{1}(\Omega)}^{2} \leq K\left\|\chi_{n}-\chi_{m}\right\|_{B C^{1}(\bar{\Omega})}\|\phi\|_{H_{0}^{1}(\Omega)}^{2} \rightarrow 0 \text { as } m, n \rightarrow \infty .
$$

By completeness of $H_{0}^{1}(\Omega)$ it follows that $\chi \phi \in H_{0}^{1}(\Omega)$ for all $\phi \in H_{0}^{1}(\Omega)$ and

$$
\left\|\chi_{n} \phi\right\|_{H_{0}^{1}(\Omega)}^{2} \leq K\|\chi\|_{B C^{1}(\bar{\Omega})}\|\phi\|_{H_{0}^{1}(\Omega)}^{2}
$$

1. If $u \in H^{-1}(\Omega)$ and $\phi \in H_{0}^{1}(\Omega)$, then by definition, $\left\langle M_{\chi} u, \phi\right\rangle=\langle u, \chi \phi\rangle$ and therefore,

$$
\begin{aligned}
\left|\left\langle M_{\chi} u, \phi\right\rangle\right| & =|\langle u, \chi \phi\rangle| \leq\|u\|_{H^{-1}(\Omega)}\|\chi \phi\|_{H_{0}^{1}(\Omega)} \\
& \leq K\|\chi\|_{B C^{1}(\bar{\Omega})}\|u\|_{H^{-1}(\Omega)}\|\phi\|_{H_{0}^{1}(\Omega)}
\end{aligned}
$$

which implies $M_{\chi} u \in H^{-1}(\Omega)$ and

$$
\left\|M_{\chi} u\right\|_{H^{-1}(\Omega)} \leq K\|\chi\|_{B C^{1}(\bar{\Omega})}\|u\|_{H^{-1}(\Omega)}
$$

2. For $u \in L^{2}(\Omega)$ and $\phi \in C_{c}^{\infty}(\Omega)$

$$
\begin{aligned}
\left|\left\langle\partial_{i} u, \phi\right\rangle\right| & =\left|\left\langle u, \partial_{i} \phi\right\rangle\right| \leq\|u\|_{L^{2}(\Omega)} \cdot\left\|\partial_{i} \phi\right\|_{L^{2}(\Omega)} \\
& \leq\|u\|_{L^{2}(\Omega)}\|\phi\|_{H_{0}^{1}(\Omega)}
\end{aligned}
$$

and therefore $\left\|\partial_{i} u\right\|_{H^{-1}(\Omega)} \leq\|u\|_{L^{2}(\Omega)}$. For general $V=\sum_{i=1}^{d} M_{A_{i}} \partial_{i}+$ $M_{A_{0}}$, we have

$$
\begin{aligned}
\|A u\|_{H^{-1}(\Omega)} & \leq \sum_{i=1}^{d} K\left\|A_{i}\right\|_{B C^{1}(\bar{\Omega})}\left\|\partial_{i} u\right\|_{H^{-1}(\Omega)}+\left\|A_{0}\right\|_{\infty}\|u\|_{L^{2}(\Omega)} \\
& \leq\left[\sum_{i=1}^{d} K\left\|A_{i}\right\|_{B C^{1}(\bar{\Omega})}+\left\|A_{0}\right\|_{\infty}\right]\|u\|_{L^{2}(\Omega)}
\end{aligned}
$$

3. Since $V: H^{1}(\Omega) \rightarrow L^{2}(\Omega)$ and $i: L^{2}(\Omega) \hookrightarrow H^{-1}(\Omega)$ are both bounded maps, to prove $L=-\sum_{i, j=1}^{d} M_{A_{i j}} \partial_{i} \partial_{j}+V$ is bounded from $H^{1}(\Omega) \rightarrow$ $H^{-1}(\Omega)$ it suffices to show $M_{A_{i j}} \partial_{i} \partial_{j}: H^{1}(\Omega) \rightarrow H^{-1}(\Omega)$ is a bounded. But $M_{A_{i j}} \partial_{i} \partial_{j}: H^{1}(\Omega) \rightarrow H^{-1}(\Omega)$ is bounded since it is the composition of the following bounded maps:

$$
H^{1}(\Omega) \xrightarrow{\partial_{j}} L^{2}(\Omega) \xrightarrow{\partial_{i}} H^{-1}(\Omega) \xrightarrow{M_{A_{i j}}} H^{-1}(\Omega)
$$

Lemma 50.7. Suppose $\chi \in B C^{\infty}(\bar{\Omega})$ then

1. $\left[L, M_{\chi}\right]=V$ is a first order differential operator acting on $\mathcal{D}^{\prime}(\Omega)$ which necessarily satisfies $V: H^{k}(\Omega) \rightarrow H^{k-1}(\Omega)$ for $k=0,1,2, \ldots$ etc.
2. If $u \in H^{k}(\Omega)$, then

$$
\left[L, M_{\chi}\right] u \in H^{k-1}(\Omega) \text { for } k=0,1,2, \ldots
$$

and

$$
\left\|\left[L, M_{\chi}\right] u\right\|_{H^{k-1}(\Omega)} \leq C_{k}(\chi)\|u\|_{H^{k}(\Omega)}
$$

Proof. On smooth functions $u \in C^{\infty}(\Omega)$,

$$
L(\chi u)=\chi L u-2 \sum_{i, j=1}^{d} A_{i j} \partial_{i} \chi \cdot \partial_{j} u+\left(\sum_{i=1}^{d} A_{i} \partial_{i} \chi-\sum_{i, j=1}^{d} A_{i j} \partial_{i} \partial_{j} \chi\right) \cdot u
$$

and therefore

$$
\begin{aligned}
{\left[L, M_{\chi}\right] u } & =-2 \sum_{i, j=1}^{d} A_{i j} \partial_{i} \chi \cdot \partial_{j} u+\left(\sum_{i=1}^{d} A_{i} \partial_{i} \chi-\sum_{i, j=1}^{d} A_{i j} \partial_{i} \partial_{j} \chi\right) \cdot u \\
& =: V u
\end{aligned}
$$

Similarly,

$$
\begin{align*}
L^{\dagger}(\chi u)= & -\sum_{i, j=1}^{d} \partial_{i} \partial_{j}\left[\chi A_{i j} u\right]-\sum_{i=1}^{d} \partial_{i}\left(\chi A_{i} u\right)+A_{0} \chi u \\
= & \chi L^{\dagger} u-2 \sum_{i, j=1}^{d} \partial_{i} \chi \cdot \partial_{j}\left[A_{i j} u\right] \\
& -\sum_{i=1}^{d} A_{i} \partial_{i} \chi \cdot u-\left(\sum_{i, j=1}^{d} A_{i j} \partial_{i} \partial_{j} \chi\right) u \tag{50.1}
\end{align*}
$$

Noting that

$$
\begin{aligned}
V^{\dagger} u & =2 \sum_{i, j=1}^{d} \partial_{j}\left[\partial_{i} \chi \cdot A_{i j} u\right]+\left(\sum_{i=1}^{d} A_{i} \partial_{i} \chi-\sum_{i, j=1}^{d} A_{i j} \partial_{i} \partial_{j} \chi\right) \cdot u \\
& =2 \sum_{i, j=1}^{d} \partial_{i} \chi \cdot \partial_{j}\left[A_{i j} u\right]+\left(\sum_{i=1}^{d} A_{i} \partial_{i} \chi+\sum_{i, j=1}^{d} A_{i j} \partial_{i} \partial_{j} \chi\right) \cdot u
\end{aligned}
$$

Eq. (50.1) may be written as

$$
\left[L^{\dagger}, M_{\chi}\right]=-V^{\dagger} .
$$

Now suppose $k=0$, then in this case for $\phi \in \mathcal{D}(\Omega)$,

$$
\begin{aligned}
\left|\left\langle\left[L, M_{\chi}\right] u, \phi\right\rangle\right| & =\left|\left\langle u,\left[M_{\chi}, L^{\dagger}\right] \phi\right\rangle\right|=\left|\left\langle u, V^{\dagger} \phi\right\rangle\right| \\
& \leq\|u\|_{L^{2}(\Omega)}\left\|V^{\dagger} \phi\right\|_{L^{2}(\Omega)} \leq C\|u\|_{L^{2}(\Omega)}\|\phi\|_{H_{0}^{1}(\Omega)} .
\end{aligned}
$$

This implies $\left\|\left[L, M_{\chi}\right] u\right\|_{H^{-1}(\Omega)} \leq C\|u\|_{L^{2}}$ and in particular $\left[L, M_{\chi}\right] u \in$ $H^{-1}(\Omega)$. For $k>0,\left[L, M_{\chi}\right] u=V u$ with $V$ as above and therefore by Proposition48.6, there exists $C<\infty$ such that $\|V u\|_{H^{k-1}(\Omega)} \leq C\|u\|_{H^{k}(\Omega)}$

Definition 50.8. The operator $L$ is uniformly elliptic on $\Omega$ if there exists $\epsilon>0$ such that $\left(A_{i j}(x)\right) \geq \epsilon I$ for all $x \in \Omega$, i.e. $A_{i j}(x) \xi_{i} \xi_{j} \geq \epsilon|\xi|^{2}$ for all $x \in \Omega$ and $\xi \in \mathbb{R}^{d}$.

Suppose now that $L$ is uniformly elliptic. Let us outline the results to be proved below.

### 50.1 Outline of future results

1. We consider $L$ with Dirichlet boundary conditions meaning we will view $L$ as a mapping from $H_{0}^{1}(\Omega) \rightarrow H^{-1}(\Omega)=\left[H_{0}^{1}(\Omega)\right]^{*}$. Proposition 51.13 below states there exists $C=C(L)<\infty$ such that $(L+C): H_{0}^{1}(\Omega) \rightarrow$ $H^{-1}(\Omega)$ is an isomorphism of Hilbert spaces. The proof uses the Dirichlet form

$$
\mathcal{E}(u, v):=\langle L u, v\rangle \text { for } u, v \in H_{0}^{1}(\Omega) .
$$

Notice for $v \in \mathcal{D}(\Omega)$ and $u \in H_{0}^{1}(\Omega)$,

$$
\begin{aligned}
\mathcal{E}(u, v) & =\langle L u, v\rangle=\left\langle u, L^{\dagger} v\right\rangle \\
& =\int_{\Omega} u\left(-\partial_{i} \partial_{j}\left(A_{i j} v\right)-\partial_{i}\left(A_{i} v\right)+A_{0} v\right) d m \\
& =\int_{\Omega}\left[\partial_{i} u \cdot \partial_{j}\left(A_{i j} v\right)-u \partial_{i}\left(A_{i} v\right)+u A_{0} v\right] d m \\
& =\int_{\Omega}\left[A_{i j} \partial_{i} u \cdot \partial_{j} v+\left(A_{i}+\partial_{j} A_{i j}\right) \partial_{i} u \cdot v+A_{0} u v\right] d m .
\end{aligned}
$$

Since the last expression is continuous for $(u, v) \in H_{0}^{1}(\Omega) \times H_{0}^{1}(\Omega)$, we have shown

$$
\mathcal{E}(u, v)=\int_{\Omega}\left[A_{i j} \partial_{i} u \cdot \partial_{j} v+\left(A_{i}+\partial_{j} A_{i j}\right) \partial_{i} u \cdot v+A_{0} u v\right] d m
$$

for all $u, v \in H_{0}^{1}(\Omega)$.
2. To implement other boundary conditions, we will need to consider $L$ acting on subspaces of $H^{2}(\Omega)$ which are determined by the boundary conditions. Rather than describe the general case here, let us consider an example where $L=-\Delta$ and the boundary condition is $\frac{\partial u}{\partial n}=\rho u$ on $\partial \Omega$ where $\partial_{n} u=\nabla u \cdot n, n$ is the outward normal on $\partial \Omega$ and $\rho: \partial \Omega \rightarrow \mathbb{R}$ is a smooth function. In this case, let

$$
D:=\left\{u \in H^{2}(\Omega): \frac{\partial u}{\partial n}=\rho u \text { on } \partial \Omega\right\} .
$$

We will eventually see that $D$ is a dense subspace of $H^{1}(\Omega)$. For $u \in D$ and $v \in H^{1}(\Omega)$,

$$
\begin{align*}
(-\Delta u, v) & =\int_{\Omega} \nabla u \cdot \nabla v d m-\int_{\partial \Omega} v \partial_{n} u d \sigma \\
& =\int_{\Omega} \nabla u \cdot \nabla v d m-\int_{\partial \Omega} \rho u v d \sigma=: \mathcal{E}(u, v) . \tag{50.2}
\end{align*}
$$

The latter expression extends by continuity to all $u \in H^{1}(\Omega)$. Given $\mathcal{E}$ as in Eq. (50.2) let $-\Delta_{\mathcal{E}}: H^{1}(\Omega) \rightarrow\left[H^{1}(\Omega)\right]^{*}$ be defined by $-\Delta_{\mathcal{E}} u:=$ $\mathcal{E}(u, \cdot)$ so that $-\Delta_{\mathcal{E}} u$ is an extension of $-\Delta u$ as a linear functional on $H_{0}^{1}(\Omega)$ to one on $H^{1}(\Omega) \supset H_{0}^{1}(\Omega)$. It will be shown below that there exists $C<\infty$ such that $\left(-\Delta_{\mathcal{E}}+C\right): H^{1}(\Omega) \rightarrow\left[H^{1}(\Omega)\right]^{*}$ is an isomorphism of Hilbert spaces.
3. The Dirichlet form $\mathcal{E}$ in Eq. (50.2) may be rewritten in a way as to avoid the surface integral term. To do this, extend the normal vector field $n$ along $\partial \Omega$ to a smooth vector field on $\bar{\Omega}$. Then by integration by parts,

$$
\begin{aligned}
\int_{\partial \Omega} \rho u v d \sigma & =\int_{\partial \Omega} n_{i}^{2} \rho u v d \sigma=\int_{\Omega} \partial_{i}\left[n_{i} \rho u v\right] d m \\
& =\left[\int_{\Omega} \nabla \cdot(\rho n) u v+\rho n_{i} \partial_{i} u \cdot v+\rho n_{i} u \cdot \partial_{i} v\right] d m
\end{aligned}
$$

In this way we see that the Dirichlet form $\mathcal{E}$ in Eq. (50.2) may be written as

$$
\begin{equation*}
\mathcal{E}(u, v)=\int_{\Omega}\left[\nabla u \cdot \nabla v+a_{i 0} \partial_{i} u \cdot v+a_{0 i} u \partial_{i} v+a_{00} u v\right] d m \tag{50.3}
\end{equation*}
$$

with $a_{00}=\nabla \cdot(\rho n), a_{i 0}=\rho n_{i}=a_{0 i}$. This should motivate the next section where we consider generalizations of the form $\mathcal{E}$ in Eq. (50.3).

## Dirichlet Forms

In this section $\Omega$ will be an open subset of $\mathbb{R}^{d}$.

### 51.1 Basics

Notation 51.1 (Dirichlet Forms) For $\alpha, \beta \in \mathbb{N}_{0}^{d}$ with $|\alpha|,|\beta| \leq 1$, suppose $a_{\alpha, \beta} \in B C^{\infty}(\bar{\Omega})$ and $\rho \in B C^{\infty}(\bar{\Omega})$ with $\rho>0$, let

$$
\begin{equation*}
\mathcal{E}(u, v)=\sum_{|\alpha|,|\beta| \leq 1} \int_{\Omega} a_{\alpha \beta} \partial^{\alpha} u \cdot \partial^{\beta} v d \mu \tag{51.1}
\end{equation*}
$$

where $d \mu:=\rho d m$. We will also write $(u, v):=\int_{\Omega} u v d \mu$ and $L^{2}$ for $L^{2}(\Omega, \mu)$. In the sequel we will often write $a_{i, \beta}$ for $a_{e_{i}, \beta}, a_{\alpha, j}$ for $a_{\alpha, e_{j}}$ and $a_{i j}$ for $a_{e_{i}, e_{j}}$.
Proposition 51.2. Let $\mathcal{E}$ be as in Notation 51.1 then

$$
|\mathcal{E}(u, v)| \leq C\|u\|_{H^{1}}\|v\|_{H^{1}} \text { for all } u, v \in H^{1}
$$

where $C$ is a constant depending on $d$ and upper bounds for $\left\{\left\|a_{\alpha \beta}\right\|_{B C(\bar{\Omega})}:|\alpha|,|\beta|\right.$
Proof. To simplify notation in the proof, let $\|\cdot\|$ denote the $L^{2}(\Omega, \mu)-$ norm. Then

$$
\begin{aligned}
|\mathcal{E}(u, v)| & \leq C \sum_{i j}\left\{\left\|\partial_{i} u\right\|\left\|\partial_{j} v\right\|+\left\|\partial_{i} u\right\|\|v\|+\|u\|\left\|\partial_{i} v\right\|+\|u\|\|v\|\right\} \\
& \leq C\|u\|_{H^{1}} \cdot\|v\|_{H^{1}}
\end{aligned}
$$

Notation 51.3 Given $\mathcal{E}$ as in Notation 51.1, let $\mathcal{L}_{\mathcal{E}}$ and $\mathcal{L}_{\mathcal{E}}^{\dagger}$ be the bounded linear operators from $H^{1}(\Omega)$ to $\left[H^{1}(\Omega)\right]^{*}$ defined by

$$
\mathcal{L}_{\mathcal{E}} u:=\mathcal{E}(u, \cdot) \text { and } \mathcal{L}_{\mathcal{E}}^{\dagger} u:=\mathcal{E}(\cdot, u)
$$

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It follows directly from the definitions that $\left\langle\mathcal{L}_{\mathcal{E}} u, v\right\rangle=\left\langle u, \mathcal{L}_{\mathcal{E}}^{\dagger} v\right\rangle$ for all $u, v \in H^{1}(\Omega)$. The Einstein summation convention will be used below when convenient.

Proposition 51.4. Suppose $\Omega$ is a precompact open subset of $\mathbb{R}^{d}$ such that $\bar{\Omega}$ is a manifold with $C^{2}$ - boundary, Then for all $u \in H^{2}(\Omega)$ and $v \in H^{1}(\Omega)$,

$$
\begin{equation*}
\left\langle\mathcal{L}_{\mathcal{E}} u, v\right\rangle=\mathcal{E}(u, v)=(L u, v)+\int_{\partial \Omega} B u \cdot v \rho d \sigma \tag{51.2}
\end{equation*}
$$

and for all $u \in H^{1}(\Omega)$ and $v \in H^{2}(\Omega)$,

$$
\begin{equation*}
\left\langle u, \mathcal{L}_{\mathcal{E}}^{\dagger} v\right\rangle=\mathcal{E}(u, v)=\left(u, L^{\dagger} v\right)+\int_{\partial \Omega} u \cdot B^{\dagger} v \rho d \sigma \tag{51.3}
\end{equation*}
$$

where

$$
\begin{align*}
B & =n_{j} a_{i j} \partial_{i}+n_{j} a_{0 j}=n \cdot a \nabla+n \cdot a_{0, \cdot},  \tag{51.4}\\
B^{\dagger} & =n_{i}\left[a_{i j} \partial_{j}+a_{i 0}\right]=a n \cdot \nabla+n \cdot a \cdot, 0  \tag{51.5}\\
L u & :=\rho^{-1} \sum_{|\alpha|,|\beta| \leq 1}(-1)^{|\beta|} \partial^{\beta}\left[\rho a_{\alpha \beta} \partial^{\alpha} u\right] \tag{51.6}
\end{align*}
$$

and

$$
\begin{equation*}
L^{\dagger} v:=\rho^{-1} \sum_{|\alpha|,|\beta| \leq 1}(-1)^{|\alpha|} \partial^{\alpha}\left[\rho a_{\alpha \beta} \partial^{\beta} v\right] \tag{51.7}
\end{equation*}
$$

We may also write $L, L^{\dagger}$ as

$$
\begin{align*}
L & =-a_{i j} \partial_{j} \partial_{i}+\left(a_{i 0}-a_{0 j}-\rho^{-1} \partial_{j}\left[\rho a_{i j}\right]\right) \partial_{i}+\left(a_{00}-\rho^{-1} \partial_{j}\left[\rho a_{0 j}\right]\right)  \tag{51.8}\\
L^{\dagger} & =-a_{i j} \partial_{i} \partial_{j}+\left(a_{0 j}-a_{j 0}-\rho^{-1} \partial_{i}\left[\rho a_{i j}\right]\right) \partial_{j}+\left(a_{00}-\rho^{-1} \partial_{i}\left[\rho a_{i 0}\right]\right) \tag{51.9}
\end{align*}
$$

Proof. Suppose $u \in H^{2}(\Omega)$ and $v \in H^{1}(\Omega)$, then by integration by parts,

$$
\begin{aligned}
\mathcal{E}(u, v) & =\sum_{|\alpha|,|\beta| \leq 1} \int_{\Omega}(-1)^{|\beta|} \rho^{-1} \partial^{\beta}\left[\rho a_{\alpha \beta} \partial^{\alpha} u\right] \cdot v d \mu+\sum_{|\alpha| \leq 1} \sum_{j=1}^{d} \int_{\partial \Omega} n_{j}\left[a_{\alpha j} \partial^{\alpha} u\right] \cdot v \rho \\
& =(L u, v)+\int_{\partial \Omega} B u \cdot v \rho d \sigma
\end{aligned} \text { where }
$$

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$$
\begin{aligned}
L u= & \rho^{-1} \sum_{|\alpha|,|\beta| \leq 1}(-1)^{|\beta|} \partial^{\beta}\left[\rho a_{\alpha \beta} \partial^{\alpha} u\right]=-\rho^{-1} \sum_{|\alpha| \leq 1} \sum_{j=1}^{d} \partial_{j}\left(\rho a_{\alpha j} \partial^{\alpha} u\right)+\sum_{|\alpha| \leq 1} a_{\alpha 0} \delta \\
= & -\rho^{-1} \sum_{i, j=1}^{d} \partial_{j}\left(\rho a_{i j} \partial_{i} u\right)-\rho^{-1} \sum_{j=1}^{d} \partial_{j}\left(\rho a_{0 j} u\right)+\sum_{i=1}^{d} a_{i 0} \partial_{i} u+a_{00} u \\
= & -\sum_{i, j=1}^{d} a_{i j} \partial_{j} \partial_{i} u-\rho^{-1} \sum_{i, j=1}^{d}\left(\partial_{j}\left[\rho a_{i j}\right]\right) \partial_{i} u-\sum_{j=1}^{d} a_{0 j} \partial_{j} u \\
& -\rho^{-1} \sum_{j=1}^{d}\left(\partial_{j}\left[\rho a_{0 j}\right]\right) u+\sum_{i=1}^{d} a_{i 0} \partial_{i} u+a_{00} u
\end{aligned}
$$

and

$$
B u=\sum_{|\alpha| \leq 1} \sum_{j=1}^{d} n_{j}\left(a_{\alpha j} \partial^{\alpha} u\right)=\sum_{i, j=1}^{d} n_{j} a_{i j} \partial_{i} u+\sum_{j=1}^{d} n_{j} a_{0 j} u .
$$

Similarly for $u \in H^{1}(\Omega)$ and $v \in H^{2}(\Omega)$,

$$
\begin{aligned}
\mathcal{E}(u, v) & =\sum_{|\alpha|,|\beta| \leq 1} \int_{\Omega} u \cdot(-1)^{|\alpha|} \rho^{-1} \partial^{\alpha}\left[\rho a_{\alpha \beta} \partial^{\beta} v\right] d \mu+\sum_{i=1}^{d} \sum_{|\beta| \leq 1} \int_{\Omega} u \cdot n_{i}\left[a_{i \beta} \partial^{\beta} v\right] \\
& =\left(u, L^{\dagger} v\right)+\int_{\partial \Omega} u \cdot B^{\dagger} v \rho d \sigma
\end{aligned}
$$

where $B^{\dagger} v=n_{i}\left[a_{i j} \partial_{j}+a_{i 0}\right]$ and

$$
\begin{aligned}
L^{\dagger} v= & -\rho^{-1} \partial_{i}\left(\rho a_{i j} \partial_{j} v\right)+a_{0 j} \partial_{j} v-\rho^{-1} \partial_{i}\left(\rho a_{i 0} v\right)+a_{00} v \\
= & -a_{i j} \partial_{i} \partial_{j} v-\rho^{-1}\left(\partial_{i}\left[\rho a_{i j}\right]\right) \partial_{j} v \\
& +a_{0 j} \partial_{j} v-a_{i 0} \partial_{i} v-\rho^{-1}\left(\partial_{i}\left[\rho a_{i 0}\right]\right) v+a_{00} v \\
= & {\left[-a_{i j} \partial_{i} \partial_{j}+\left(a_{0 j}-a_{j 0}-\rho^{-1} \partial_{i}\left[\rho a_{i j}\right]\right) \partial_{j}+a_{00}-\rho^{-1} \partial_{i}\left[\rho a_{i 0}\right]\right] v . }
\end{aligned}
$$

Proposition 51.4 shows that to the Dirichlet form $\mathcal{E}$ there is an associated second order elliptic operator $L$ along with boundary conditions $B$ as in Eqs. (51.6) and (51.4). The next proposition shows how to reverse this procedure and associate a Dirichlet form $\mathcal{E}$ to a second order elliptic operator $L$ with boundary conditions.

Proposition 51.5 (Following Folland p. 240.). Let $A_{j}, A_{0} \rho \in B C^{\infty}(\Omega)$ and $A_{i j}=A_{j i} \in B C^{\infty}(\Omega)$ with $\left(A_{i j}\right)>0$ and $\rho>0$ and let

$$
\begin{equation*}
L=-A_{i j} \partial_{i} \partial_{j}+A_{i} \partial_{i}+A_{0} \tag{51.10}
\end{equation*}
$$

and $(u, v):=\int_{\Omega}$ uvpdm. Also suppose $\alpha: \partial \Omega \rightarrow \mathbb{R}$ and $V: \partial \Omega \rightarrow \mathbb{R}^{d}$ are smooth functions such that $V(x) \cdot n(x)>0$ for all $x \in \partial \Omega$ and let $B_{0} u:=$
$V \cdot \nabla u+\alpha u$. Then there exists a Dirichlet form $\mathcal{E}$ as in Notation 51.1 and $\beta \in$ $C^{\infty}(\partial \Omega \rightarrow(0, \infty))$ such that Eq. (51.2) holds with $B u=\beta B_{0} u$. In particular if $u \in H^{2}(\Omega)$, then $B u=0$ iff $B_{0} u=0$ on $\partial \Omega$.

Proof. Since mixed partial derivatives commute on $H^{2}(\Omega)$, the term $a_{i j} \partial_{j} \partial_{i}$ in Eq. (51.8) may be written as

$$
\frac{1}{2}\left(a_{i j}+a_{j i}\right) \partial_{j} \partial_{i}
$$

With this in mind we must find coefficients $\left\{a_{\alpha, \beta}:|\alpha|,|\beta| \leq 1\right\}$ as in Notation 51.1, such that

$$
\begin{align*}
A_{i j} & =\frac{1}{2}\left(a_{i j}+a_{j i}\right)  \tag{51.11}\\
A_{i} & =\left(a_{i 0}-a_{0 j}-\rho^{-1} \partial_{j}\left[\rho a_{i j}\right]\right),  \tag{51.12}\\
A_{0} & =a_{00}-\rho^{-1} \partial_{j}\left[\rho a_{0 j}\right],  \tag{51.13}\\
a^{\operatorname{tr}} n & =\beta V \text { and }  \tag{51.14}\\
n_{i} a_{0 i} & =\beta \alpha . \tag{51.15}
\end{align*}
$$

Eq. (51.11) will be satisfied if

$$
a_{i j}=A_{i j}+c_{i j}
$$

where $c_{i j}=-c_{j i}$ are any functions in $B C^{\infty}(\Omega)$. Dotting Eq. (51.14) with $n$ shows that

$$
\begin{equation*}
\beta=\frac{a^{\text {tr }} n \cdot n}{V \cdot n}=\frac{n \cdot a n}{V \cdot n}=\frac{n \cdot A n}{V \cdot n} \tag{51.16}
\end{equation*}
$$

and Eq. (51.14) may now be written as

$$
\begin{equation*}
w:=A n-\frac{n \cdot A n}{V \cdot n} V=c n \tag{51.17}
\end{equation*}
$$

which means we have to choose $c=\left(c_{i j}\right)$ so that Eq. (51.17) holds. This is easily done, since $w \cdot n=0$ by construction we may define $c \xi:=w(n \cdot \xi)-n(w \cdot \xi)$ for $\xi \in \mathbb{R}^{d}$. Then $c$ is skew symmetric and $c n=w$ as desired. Since $c_{i j}$ are smooth functions on $\partial \Omega$, a partition of unity argument shows that $c_{i j}=-c_{j i}$ may be extended to element of $C^{\infty}(\bar{\Omega})$. (These extensions are highly nonunique but it does not matter.) With these choices, Eq. (51.11) and Eq. (51.14) now hold with $\beta$ as in Eq. (51.16). We now choose $a_{0 i} \in C^{\infty}(\bar{\Omega})$ such that $a_{0 i}=\beta \alpha n_{i}$ on $\partial \Omega$. Once these choices are made, it should be clear that Eqs. (51.13) and (51.14) may be solved uniquely for the functions $a_{0 j}$ and $a_{00}$.

### 51.2 Weak Solutions for Elliptic Operators

For the rest of this subsection we will assume $\rho=1$. This can be done here by absorbing $\rho$ into the coefficient $a_{\alpha \beta}$.

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Definition 51.6. The Dirichlet for $\mathcal{E}$ is uniformly elliptic on $\Omega$ if there exists $\epsilon>0$ such that $\left(a_{i j}(x)\right) \geq \epsilon I$ for all $x \in \Omega$, i.e. $a_{i j}(x) \xi_{i} \xi_{j} \geq \epsilon|\xi|^{2}$ for all $x \in \Omega$ and $\xi \in \mathbb{R}^{d}$.
Assumption 4 For the remainder of this chapter, it will be assumed that $\mathcal{E}$ is uniformly elliptic on $\Omega$.
Lemma 51.7. If $\xi^{2} \leq A \xi+B$ then $\xi^{2} \leq A^{2}+2 B$.
Proof. $\xi^{2} \leq \frac{1}{2} A^{2}+\frac{1}{2} \xi^{2}+B$. Therefore $\frac{1}{2} \xi^{2} \leq \frac{1}{2} A^{2}+B$ or $\xi^{2} \leq A^{2}+2 B$.
Theorem 51.8. Keeping the notation and assumptions of Proposition 51.2 along with Assumption 4, then

$$
\begin{equation*}
\mathcal{E}(u, u)+C_{\epsilon}\|u\|_{L^{2}(\Omega)} \geq \frac{\epsilon}{2}\|u\|_{H^{1}(\Omega)}, \tag{51.18}
\end{equation*}
$$

where $C_{\epsilon}=\frac{2 C^{2}}{\epsilon}+C+\frac{\epsilon}{2}$.
Proof. To simplify notation in the proof, let $\|\cdot\|$ denote the $L^{2}(\Omega)$ - norm. Since

$$
\begin{aligned}
& \int_{\Omega} a_{i j} \partial_{i} u \cdot \partial_{j} u d m \geq \epsilon \int_{\Omega}|\nabla u|^{2} d x=\epsilon\|\nabla u\|_{L^{2}}^{2} \\
& \quad \mathcal{E}(u, u) \geq \epsilon\|\nabla u\|_{L^{2}}^{2}-C\left(\|\nabla u\|\|u\|+\|u\|^{2}\right)
\end{aligned}
$$

and so

$$
\|\nabla u\|^{2} \leq \frac{C}{\epsilon}\|u\|\|\nabla u\|+\left(\frac{1}{\epsilon} \mathcal{E}(u, u)+\frac{C}{\epsilon}\|u\|^{2}\right)
$$

Therefore by Lemma 51.7 with $A=\frac{C}{\epsilon}\|u\|, B=\left(\frac{1}{\epsilon} \mathcal{E}(u, u)+\frac{C}{\epsilon}\|u\|^{2}\right)$ and $\xi=\|\nabla u\|$,

$$
\begin{aligned}
&\|\nabla u\|^{2} \leq \frac{C^{2}}{\epsilon^{2}}\|u\|^{2}+\frac{2}{\epsilon}\left(\mathcal{E}(u, u)+\frac{C}{\epsilon}\|u\|^{2}\right) \\
&=\frac{2}{\epsilon} \mathcal{E}(u, u)+\left(\frac{C^{2}}{\epsilon^{2}}+\frac{2 C}{\epsilon}\right)\|u\|^{2}
\end{aligned}
$$

Hence

$$
\frac{\epsilon}{2}\|\nabla u\|^{2} \leq \mathcal{E}(u, u)+\left(\frac{2 C^{2}}{\epsilon}+C\right)\|u\|^{2}
$$

which, after adding $\frac{\epsilon}{2}\|u\|^{2}$ to both sides of this equation, gives Eq. (51.18)
The following theorem is an immediate consequence of Theorem 51.8 and the Lax-Milgram Theorem 53.9.
Corollary 51.9. The quadratic form

$$
Q(u, v):=\mathcal{E}(u, v)+C_{\epsilon}(u, v)
$$

satisfies the assumptions of the Lax Milgram Theorem 53.9 on $H^{1}(\Omega)$ or any closed subspace $X$ of $H^{1}(\Omega)$.

Theorem 51.10 (Weak Solutions). Let $\mathcal{E}$ be as in Notation 51.1 and $C_{\epsilon}$ be as in Theorem 51.8,

$$
Q(u, v):=\mathcal{E}(u, v)+C_{\epsilon}(u, v) \text { for } u, v \in H^{1}(\Omega)
$$

and $X$ be a closed subspace of $H^{1}(\Omega)$. Then the maps $\mathcal{L}: X \rightarrow X^{*}$ and $\mathcal{L}^{\dagger}: X \rightarrow X^{*}$ defined by

$$
\begin{aligned}
\mathcal{L} v & :=Q(v, \cdot)=\left(\mathcal{L}_{\mathcal{E}}+C\right) v \text { and } \\
\mathcal{L}^{\dagger} v & :=Q(\cdot, v)=\left(\mathcal{L}_{\mathcal{E}}^{\dagger}+C\right) v
\end{aligned}
$$

are linear isomorphisms of Hilbert spaces satisfying

$$
\left\|\mathcal{L}^{-1}\right\|_{L\left(X^{*}, X\right)} \leq \frac{2}{\epsilon} \text { and }\left\|\left(\mathcal{L}^{\dagger}\right)^{-1}\right\|_{L\left(X^{*}, X\right)} \leq \frac{2}{\epsilon}
$$

In particular for $f \in X^{*}$, there exist a unique solution $u \in X$ to $\mathcal{L} u=f$ and this solution satisfies the estimate

$$
\|u\|_{H^{1}(\Omega)} \leq \frac{2}{\epsilon}\|f\|_{X^{*}}
$$

Remark 51.11. If $X \supset H_{0}^{1}(\Omega)$ and $u \in X$ then for $\phi \in C_{c}^{\infty}(\Omega) \subset X$,

$$
\langle\mathcal{L} u, \phi\rangle=Q(u, \phi)=\left(u,\left(L^{\dagger}+C\right) \phi\right)=\langle(L+C) u, \phi\rangle
$$

That is to say $\left.\mathcal{L} u\right|_{C_{c}^{\infty}(\Omega)}=(L+C) u$. In particular any solution $u \in X$ to $\mathcal{L} u=f \in X^{*}$ solves

$$
(L+C) u=\left.f\right|_{C_{c}^{\infty}(\Omega)} \in \mathcal{D}^{\prime}(\Omega)
$$

Remark 51.12. Suppose that $\Gamma \subset \partial \Omega$ is a measurable set such that $\sigma(\Gamma)>0$ and $X_{\Gamma}:=\left\{u \in H^{1}(\Omega):\left.1_{\Gamma} u\right|_{\partial \Omega}=0\right\}$. If $u \in H^{2}(\Omega)$ solves $\mathcal{L} u=f$ for some $f \in L^{2}(\Omega) \subset X^{*}$, then by Proposition 51.4,

$$
(f, u):=\langle\mathcal{L} u, v\rangle=\mathcal{E}(u, v)+C(u, v)=((L+C) u, v)+\int_{\partial \Omega} B u \cdot v d \sigma
$$

for all $v \in X_{\Gamma} \subset H^{1}(\Omega)$. Taking $v \in \mathcal{D}(\Omega) \subset X_{\Gamma}$ in Eq. (51.19) shows $(L+C) u=f$ a.e. and

$$
\int_{\partial \Omega} B u \cdot v d \sigma=0 \text { for all } v \in X_{\Gamma}
$$

Therefore we may conclude, $u$ solves

$$
\begin{aligned}
(L+C) u & =f \text { a.e. with } \\
B u(x) & =0 \text { for } \sigma-\text { a.e. } x \in \partial \Omega \backslash \Gamma \text { and } \\
u(x) & =0 \text { for } \sigma-\text { a.e. } x \in \Gamma .
\end{aligned}
$$

The following proposition records the important special case of Theorem 51.10 when $X=H_{0}^{1}(\Omega)$ and hence $X^{*}=H^{-1}(\Omega)$. The point to note here is that $\mathcal{L} u=(L+C) u$ when $X=H_{0}^{1}(\Omega)$, i.e. $\mathcal{L} u$ equals $[(L+C) u]$ extended by continuity to a linear functional on $X^{*}=\left[H_{0}^{1}(\Omega]^{*}\right.$.

Proposition 51.13. Assume $L$ is elliptic as above. Then there exist $C>0$ sufficiently large such that $(L+C): H_{0}^{1}(\Omega) \rightarrow H^{-1}(\Omega)$ is bijective with bounded inverse. Moreover

$$
\left\|(L+C)^{-1}\right\|_{L\left(H_{0}^{1}(\Omega), H^{-1}(\Omega)\right)} \leq 2 / \epsilon
$$

or equivalently

$$
\|u\|_{H_{0}^{1}(\Omega)} \leq \frac{2}{\epsilon}\|(L+C) u\|_{H^{-1}(\Omega)} \text { for all } u \in H_{0}^{1}(\Omega)
$$

Our next goal, see Theorem 52.15, is to prove the elliptic regularity result, namely if $X=H_{0}^{1}(\Omega)$ or $X=H^{1}(\Omega)$ and $u \in X$ satisfies $\mathcal{L} u \in H^{k}(\Omega)$, then $u \in H^{k+2}(\Omega) \cap X$.

## Elliptic Regularity

Assume that $\bar{\Omega}$ is a compact manifold with $C^{2}$ - boundary and satisfying $\bar{\Omega}^{o}=\Omega$ and let $\mathcal{E}$ be the Dirichlet form defined in Notation 51.1 and $L$ be as in Eq. (51.6) or Eq. (51.8). We will assume $\mathcal{E}$ or equivalently that $L$ is uniformly elliptic on $\Omega$. This section is devoted to proving the following elliptic regularity theorem.
Theorem 52.1 (Elliptic Regularity Theorem). Suppose $X=H_{0}^{1}(\Omega)$ or $H^{1}(\Omega)$ and $\mathcal{E}$ is as above. If $u \in X$ such that $\mathcal{L}_{\mathcal{E}} u \in H^{k}(\Omega)$ for some $k \in$ $\mathbb{N}_{0} \cup\{-1\}$, then $u \in H^{k+2}(\Omega)$ and

$$
\begin{equation*}
\|u\|_{H^{k+2}}(\Omega) \leq C\left(\left\|\mathcal{L}_{\mathcal{E}} u\right\|_{H^{k}(\Omega)}+\|u\|_{L^{2}(\Omega)}\right) \tag{52.1}
\end{equation*}
$$

### 52.1 Interior Regularity

Theorem 52.2 (Elliptic Interior Regularity). To each $\chi \in C_{c}^{\infty}(\Omega)$ there exist a constant $C=C(\chi)$ such that

$$
\begin{equation*}
\|\chi u\|_{H^{1}(\Omega)} \leq C\left\{\|L u\|_{H^{-1}(\Omega)}+\|u\|_{L^{2}(\Omega)}\right\} \text { for all } u \in H^{1}(\Omega) \tag{52.2}
\end{equation*}
$$

In particular, if $W$ is a precompact open subset of $\Omega$, then

$$
\begin{equation*}
\|u\|_{H^{1}(W)} \leq C\left\{\|L u\|_{H^{-1}(\Omega)}+\|u\|_{L^{2}(\Omega)}\right\} . \tag{52.3}
\end{equation*}
$$

Proof. For $u \in H^{1}(\Omega), \chi u \in H_{0}^{1}(\Omega)$ and hence by Proposition 51.13, Proposition 50.6 and Lemma 50.7,

$$
\begin{aligned}
\|\chi u\|_{H^{1}(\Omega)} & \leq \frac{2}{\epsilon}\left\|\left(L+C_{\epsilon}\right)(\chi u)\right\|_{H^{-1}(\Omega)} \\
& =\frac{2}{\epsilon}\left\|\chi\left(L+C_{\epsilon}\right) u+\left[L, M_{\chi}\right] u\right\|_{H^{-1}(\Omega)} \\
& \leq \frac{2}{\epsilon} C(\chi)\left\{\left\|\left(L+C_{\epsilon}\right) u\right\|_{H^{-1}(\Omega)}+\|u\|_{L^{2}(\Omega)}\right\} \\
& \leq \frac{2}{\epsilon} C(\chi)\left\{\|L u\|_{H^{-1}(\Omega)}+C_{\epsilon}\|u\|_{H^{-1}(\Omega)}+\|u\|_{L^{2}(\Omega)}\right\}
\end{aligned}
$$

from which Eq. (52.2) follows. To prove Eq. (52.3), choose $\chi \in C_{c}^{\infty}(\Omega,[0,1])$ such that $\chi=1$ on a neighborhood of $\bar{W}$ in which case

$$
\|u\|_{H^{1}(W)}=\|\chi u\|_{H^{1}(W)} \leq\|\chi u\|_{H^{1}(\Omega)} \leq C\left\{\|L u\|_{H^{-1}(\Omega)}+\|u\|_{L^{2}(\Omega)}\right\} .
$$

Exercise 52.3. Let $v \in \mathbb{R}^{d}$ with $|v|=1, u \in L^{2}(\Omega)$ and $W$ be an open set such that $\bar{W} \sqsubset \sqsubset \Omega$. For all $h \neq 0$ sufficiently show

$$
\left\|\partial_{v}^{h} u\right\|_{H^{-1}(W)} \leq\|u\|_{L^{2}(\Omega)} .
$$

Notice that $\partial_{v}^{h} u \in L^{2}(W) \subset H^{-1}(W)$.
Solution 52.4 (52.3). Let $W_{1}$ be a precompact open subset of $\Omega$ such that $\bar{W} \subset W_{1} \subset \bar{W}_{1} \subset \Omega$. Then for $\phi \in \mathcal{D}(W)$ and $h$ close to zero,

$$
\begin{aligned}
\left|\left\langle\partial_{v}^{h} u, \phi\right\rangle\right| & =\left|\left\langle u, \partial_{v}^{-h} \phi\right\rangle\right| \leq\|u\|_{L^{2}\left(W_{1}\right)}\left\|\partial_{v}^{-h} \phi\right\|_{L^{2}\left(W_{1}\right)} \\
& \leq\|u\|_{L^{2}\left(W_{1}\right)}\left\|\partial_{v} \phi\right\|_{L^{2}(\Omega)} \text { (Theorem 48.13) } \\
& \leq\|u\|_{L^{2}(\Omega)}\|\phi\|_{H^{1}(\Omega)} .
\end{aligned}
$$

Hence

$$
\begin{aligned}
\left\|\partial_{v}^{h} u\right\|_{H^{-1}(W)} & =\sup \left\{\left|\left\langle\partial_{v}^{h} u, \phi\right\rangle\right|: \phi \in \mathcal{D}(W) \text { with }\|\phi\|_{H^{1}(\Omega)}=1\right\} \\
& \leq\|u\|_{L^{2}(\Omega)} .
\end{aligned}
$$

Theorem 52.5 (Interior Regularity). Suppose $L$ is $2^{\text {nd }}$ order uniformly elliptic operator on $\Omega$ and $u \in H^{1}(\Omega)$ satisfies $L u \in H^{k}(\Omega)^{1}$ for some $k=$ $-1,0,1,2, \ldots$, then $u \in H_{l o c}^{k+2}(\Omega)$. Moreover, if $W \subset \subset \Omega$ then there exists $C=C_{k}(W)<\infty$ such that

$$
\begin{equation*}
\|u\|_{H^{k+2}(W)} \leq C\left(\|L u\|_{H^{k}(\Omega)}+\|u\|_{L^{2}(\Omega)}\right) . \tag{52.4}
\end{equation*}
$$

Proof. The proof is by induction on $k$ with Theorem 52.2 being the case $k=-1$. Suppose that the interior regularity theorem holds for $-1 \leq k \leq k_{0}$. We will now complete the induction proof by showing it holds for $k=k_{0}+1$.

So suppose that $u \in H^{1}(\Omega)$ such that $L u \in H^{k_{0}+1}(\Omega)$ and $W=W_{0} \subset \Omega$ is fixed. Choose open sets $W_{1}, W_{2}$ and $W_{3}$ such that $\bar{W}_{0} \subset W_{1} \subset \bar{W}_{1} \subset$ $W_{2} \subset \bar{W}_{2} \subset W_{3} \subset \bar{W}_{3} \subset \Omega$ as in Figure 52.1. The idea now is to apply the induction hypothesis to the function $\partial_{v}^{h} u$ where $v \in \mathbb{R}^{d}$ and $\partial_{v}^{h}$ is the finite difference operator in Definition 29.14. For the remainder of the proof $h \neq 0$ will be assumed to be sufficiently small so that the following computations make sense. To simplify notation let $D^{h}=\partial_{v}^{h}$.

[^6]\[

$$
\begin{aligned}
-\left\langle u, D^{-h} L^{\dagger} \phi\right\rangle & =\left\langle L D^{h} u, \phi\right\rangle=\left\langle D^{h} L u+L^{h} \tau_{v}^{h} u, \phi\right\rangle \\
& =-\left\langle u, L^{\dagger} D^{-h} \phi\right\rangle+\left\langle\tau_{v}^{h} u,\left(L^{h}\right)^{\dagger} \phi\right\rangle
\end{aligned}
$$
\]

So the real identity which needs to be proved here is that $\left[D^{-h}, L^{\dagger}\right] \phi=$ $-\tau_{v}^{-h}\left(L^{h}\right)^{\dagger} \phi$ for all $\phi \in \mathcal{D}\left(W_{3}\right)$. This can be done as above or it can be inferred (making use of the properties $L^{\dagger}$ is the formal adjoint of $L$ and $-D^{-h}$ is the formal adjoint of $D^{h}$ ) from the computations already done in the previous paragraph with $u$ being a smooth function.

Since $L^{h}$ is a second order differential operator with coefficients which have bounded derivatives to all orders with bounds independent of $h$ small, $\left[L, D^{h}\right] u=L^{h} \tau_{v}^{h} u \in H^{k_{0}}\left(W_{1}\right)$ and there is a constant $C<\infty$ such that

$$
\begin{align*}
\left\|\left[L, D^{h}\right] u\right\|_{H^{k_{0}}\left(W_{1}\right)} & =\left\|L^{h} \tau_{v}^{h} u\right\|_{H^{k_{0}}\left(W_{1}\right)} \\
& \leq C\left\|\tau_{v}^{h} u\right\|_{H^{k_{0}+2}\left(W_{2}\right)} \leq C\|u\|_{H^{k_{0}+2}\left(W_{3}\right)} . \tag{52.7}
\end{align*}
$$

Combining Eqs. $(52.5-52.7)$ implies that $L D^{h} u \in H^{k_{0}}\left(W_{2}\right)$ and

$$
\left\|L D^{h} u\right\|_{H^{k_{0}}\left(W_{1}\right)} \lesssim\|L u\|_{H^{k_{0}+1}\left(W_{2}\right)}+\|u\|_{H^{k_{0}+2}\left(W_{3}\right)}
$$

Therefore by the induction hypothesis, $D^{h} u \in H^{k_{0}+2}\left(W_{0}\right)$ and

$$
\begin{aligned}
\left\|D^{h} u\right\|_{H^{k_{0}+2}\left(W_{0}\right)} & \lesssim\left\|L D^{h} u\right\|_{H^{k_{0}}\left(W_{1}\right)}+\left\|D^{h} u\right\|_{L^{2}\left(W_{1}\right)} \\
& \lesssim\|L u\|_{H^{k_{0}+1}\left(W_{2}\right)}+\|u\|_{H^{k_{0}+2}\left(W_{3}\right)}+\|u\|_{H^{1}\left(W_{2}\right)} \\
& \lesssim\|L u\|_{H^{k_{0}+1}\left(W_{2}\right)}+\|u\|_{H^{k_{0}+2}\left(W_{3}\right)} \\
& \lesssim\|L u\|_{H^{k_{0}+1}(\Omega)}+\|L u\|_{H^{k_{0}}(\Omega)}+\|u\|_{L^{2}(\Omega)} \quad \text { (by induction) } \\
& \lesssim\|L u\|_{H^{k_{0}+1}(\Omega)}+\|u\|_{L^{2}(\Omega)}
\end{aligned}
$$

So by Theorem 48.13, $\partial_{v} u \in H^{k_{0}+2}\left(W_{0}\right)$ for all $v=e_{i}$ with $i=1,2, \ldots, d$ and

$$
\left\|\partial_{i} u\right\|_{H^{k_{0}+2}(W)}=\left\|\partial_{i} u\right\|_{H^{k_{0}+2}\left(W_{0}\right)} \lesssim\|L u\|_{H^{k_{0}+1}(\Omega)}+\|u\|_{L^{2}(\Omega)}
$$

Thus $u \in H^{k_{0}+3}\left(W_{0}\right)$ and Eq. (52.4) holds.
Corollary 52.6. Suppose $L$ is as above and $u \in H^{1}(\Omega)$ such that $L u \in$ $B C^{\infty}(\Omega)$ then $u \in C^{\infty}(\Omega)$.

Proof. Choose $\Omega_{0} \subset \subset \Omega$ so $L u \in B C^{\infty}\left(\bar{\Omega}_{0}\right)$. Therefore $L u \in H^{k}\left(\Omega_{0}\right)$ for all $k=0,1,2, \ldots$. Hence $u \in H_{\text {loc }}^{k+2}\left(\Omega_{0}\right)$ for all $k=0,1,2, \ldots$. Then by Sobolev embedding Theorem 49.18, $u \in C^{\infty}\left(\Omega_{0}\right)$. Since $\Omega_{0}$ is an arbitrary precompact open subset of $\Omega, u \in C^{\infty}(\Omega)$.

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### 52.2 Boundary Regularity Theorem

Example 52.7. Let $\Omega=D(0,1)$ and $u(z)=(1+z) \log (1+z)$. Since $u(z)$ is holomorphic on $\Omega$ it is also harmonic, i.e. $\Delta u=0 \in H^{k}(\Omega)$ for all $k$. However we will now show that while $u \in H^{1}(\Omega)$ it is not in $H^{2}(\Omega)$. Because $u$ is holomorphic,

$$
u_{x}=\frac{\partial u}{\partial z}=1+\log (1+z) \text { and } u_{y}=i \frac{\partial u}{\partial z}=i+i \log (1+z)
$$

from which it is easily shown $u \in H^{1}(\Omega)$. On the other hand,

$$
u_{x x}=\frac{\partial^{2}}{\partial z^{2}} u=\frac{1}{1+z}
$$

and

$$
\begin{aligned}
\int_{\Omega}\left|\frac{1}{1+z}\right|^{2} d x d y & =\int_{\Omega+1}\left|\frac{1}{z}\right|^{2} d x d y \\
& \geq \int_{C}\left|\frac{1}{z}\right|^{2} d x d y=\frac{\pi}{2} \int_{0}^{1} \frac{1}{r^{2}} r d r=\infty
\end{aligned}
$$

where $C$ is the cone in Figure 52.2. This shows $u \notin H^{2}(\Omega)$ and the problems come from the bad behavior of $u$ near $-1 \in \partial \Omega$.


Fig. 52.2. The cone used in showing $u$ not in $H^{2}(\Omega)$.

This example shows that in order to get an elliptic regularity result which is valid all the way up to the boundary, it is necessary to impose some sort of
boundary conditions on the solution which will rule out the bad behavior of the example. Since the Dirichlet form contains boundary information, we will do this by working with $\mathcal{E}$ rather than the operator $L$ on $\mathcal{D}^{\prime}(\Omega)$ associated to $\mathcal{E}$. Having to work with the quadratic form makes life a bit more difficult.

## Notations 52.8 Let

1. $N_{r}:=\left\{x \in \mathbb{H}^{n}:|x|<r\right\}$.
2. $X=H_{0}^{1}\left(N_{r}\right)$ or $X$ be the closed subspace $H^{1}\left(N_{r}\right)$ given by

$$
\begin{equation*}
X=\left\{u \in H^{1}\left(N_{r}\right):\left.u\right|_{\partial \mathbb{H}^{n} \cap \bar{N}_{r}}=0\right\} . \tag{52.8}
\end{equation*}
$$

3. For $s \leq r$ let $X_{s}=\left\{u \in X: u=0\right.$ on $\mathbb{H}^{n} \backslash N_{\rho}$ for some $\left.\rho<s\right\}$.


Fig. 52.3. Nested half balls.

Remark 52.9. 1. If $\phi \in C^{\infty}\left(\overline{\mathbb{H}}^{n}\right)$ and vanishes on $\mathbb{H}^{n} \backslash N_{\rho}$ for some $\rho<r$ then $\phi u \in X_{r}$ for all $u \in X$.
2. If $u \in X_{r}$ then $\partial_{h}^{\alpha} u \in X_{r}$ for all $\alpha$ such that $\alpha_{d}=0$ and $|h|$ sufficiently small.
Lemma 52.10 (Commutator). If $\psi \in C^{\infty}\left(\bar{N}_{r}\right)$ then for $\gamma \in \mathbb{N}_{0}^{d-1} \times\{0\}$ there exists $C_{\gamma}(\psi)<\infty$ such that

$$
\begin{equation*}
\left\|\left[\psi, \partial_{h}^{\gamma}\right] f\right\|_{L^{2}\left(N_{\rho}\right)} \leq C_{\gamma}(\psi) \sum_{\alpha<\gamma}\left\|\partial^{\alpha} f\right\|_{L^{2}\left(N_{r}\right)} \tag{52.9}
\end{equation*}
$$

for all $f \in L^{2}\left(N_{r}\right)$ with $\partial^{\alpha} f \in L^{2}\left(N_{r}\right)$ for $\alpha<\gamma$.
Proof. The proof will be by induction on $|\gamma|$. If $\gamma=e_{i}$ for some $i<d$, then

$$
\begin{aligned}
\partial_{h}^{i}(\psi f)(x) & :=\frac{\psi\left(x+h e_{i}\right) f\left(x+h e_{i}\right)-\psi(x) f(x)}{h} \\
& =\frac{\left[\psi\left(x+h e_{i}\right)-\psi(x)\right] f\left(x+h e_{i}\right)+\psi(x)\left[f\left(x+h e_{i}\right)-f(x)\right]}{h} \\
& =\partial_{h}^{i} \psi(x) f\left(x+h e_{i}\right)+\psi(x) \partial_{h}^{i} f(x)
\end{aligned}
$$

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which gives

$$
\begin{equation*}
\left[\partial_{h}^{i}, \psi\right] f=\left(\partial_{h}^{i} \psi\right) \tau_{h}^{i} f \tag{52.10}
\end{equation*}
$$

This then implies that

$$
\left\|\left[\psi, \partial_{h}^{i}\right] f\right\|_{L^{2}\left(N_{\rho}\right)} \leq C(\psi)\|f\|_{L^{2}\left(N_{r}\right)}
$$

Now suppose $|\gamma|>1$ with $\gamma=e_{i}+\gamma^{\prime}$ so that $\partial_{h}^{\gamma}=\partial_{h}^{\gamma^{\prime}} \partial_{h}^{i}$ with $\left|\gamma^{\prime}\right|=|\gamma|-1$. Then

$$
\left[\psi, \partial_{h}^{\gamma}\right]=\left[\psi, \partial_{h}^{\gamma^{\prime}}\right] \partial_{h}^{i}+\partial_{h}^{\gamma^{\prime}}\left[\psi, \partial_{h}^{i}\right]
$$

and therefore by the induction hypothesis and Theorems 48.13 and 48.15,

$$
\begin{align*}
\left\|\left[\psi, \partial_{h}^{\gamma}\right] f\right\|_{L^{2}} & \leq C_{\gamma^{\prime}}(\psi) \sum_{\alpha<\gamma^{\prime}}\left\|\partial^{\alpha} \partial_{h}^{i} f\right\|_{L^{2}}+\left\|\partial^{\gamma^{\prime}}\left[\psi, \partial_{h}^{i}\right] f\right\|_{L^{2}} \\
& \leq C_{\gamma^{\prime}}(\psi) \sum_{\alpha<\gamma^{\prime}}\left\|\partial^{\alpha+e_{i}} f\right\|_{L^{2}}+\left\|\partial^{\gamma^{\prime}}\left[\left(\partial_{h}^{i} \psi\right) \tau_{h}^{i} f\right]\right\|_{L^{2}} \tag{52.11}
\end{align*}
$$

But

$$
\partial^{\gamma^{\prime}}\left[\left(\partial_{h}^{i} \psi\right) \tau_{h}^{i} f\right]=\sum_{\beta_{1}+\beta_{2}=\gamma^{\prime}} \frac{\gamma^{\prime}!}{\beta_{1}!\beta_{2}!}\left(\partial_{h}^{i} \partial^{\beta_{1}} \psi\right) \tau_{h}^{i} \partial^{\beta_{2}} f
$$

and hence

$$
\begin{equation*}
\left\|\partial^{\gamma^{\prime}}\left[\left(\partial_{h}^{i} \psi\right) \tau_{h}^{i} f\right]\right\| \leq C \sum_{\beta \leq \gamma^{\prime}}\left\|\partial^{\beta} f\right\|_{L^{2}} \tag{52.12}
\end{equation*}
$$

Combining Eqs. (52.11) and (52.12) gives the desired result,

$$
\left\|\left[\psi, \partial_{h}^{\gamma}\right] f\right\|_{L^{2}} \leq C_{\gamma}(\psi) \sum_{\alpha<\gamma}\left\|\partial^{\alpha} f\right\|_{L^{2}}
$$

Lemma 52.11 (Warmup for Proposition 52.12). Let $a_{\alpha \beta} \in B C^{\infty}\left(\mathbb{H}^{d}\right)$ with $\left(a_{i j}\right) \geq \epsilon \delta_{i j}$ for some $\epsilon>0$,

$$
\begin{equation*}
\langle\mathcal{L} u, v\rangle=\mathcal{E}(u, v)=\int_{\mathbb{H}^{d}} \sum_{|\alpha|,|\beta| \leq 1} a_{\alpha \beta} \partial^{\alpha} u \cdot \partial^{\beta} v d x \tag{52.13}
\end{equation*}
$$

$X=H_{0}^{1}\left(\mathbb{H}^{d}\right)$ or $H^{1}\left(\mathbb{H}^{d}\right)$. There exists $C<\infty$ such that if $u \in X$ such that $\mathcal{L} u=: f \in L^{2}\left(\mathbb{H}^{d}\right)$, then $u \in H^{2}\left(\mathbb{R}^{d}\right)$ and

$$
\begin{equation*}
\|u\|_{H^{2}\left(\mathbb{H}^{d}\right)} \leq C\left(\|f\|_{L^{2}\left(\mathbb{H}^{d}\right)}+\|u\|_{X^{*}}\right) \tag{52.14}
\end{equation*}
$$

Proof. If $\mathcal{L} u=f \in X^{*}$ then $(\mathcal{L}+C) u=f+C u$, so by the Lax-Milgram method,

$$
\|u\|_{X} \lesssim\|f+C u\|_{X^{*}} \leq\|f\|_{X^{*}}+C\|u\|_{X^{*}} \lesssim\|\mathcal{L} u\|_{X^{*}}+\|u\|_{X^{*}}
$$

We wish to prove $\partial_{i} u \in H^{1}\left(\mathbb{H}^{d}\right)$ for all $i<d$ and

$$
\left\|\partial_{i} u\right\|_{H^{1}\left(\mathbb{H}^{d}\right)} \lesssim\|\mathcal{L} u\|_{L^{2}\left(\mathbb{H}^{d}\right)}+\|u\|_{X^{*}} .
$$

To do this consider

$$
\begin{aligned}
\left\langle\mathcal{L} \partial_{i}^{h} u, v\right\rangle & =\int_{\mathbb{H}^{d}} \sum_{|\alpha|,|\beta| \leq 1} a_{\alpha \beta} \partial_{i}^{h} \partial^{\alpha} u \cdot \partial^{\beta} v d x \\
& =\int_{\mathbb{H}^{d}} \sum_{|\alpha|,|\beta| \leq 1}\left\{\partial_{i}^{h}\left(a_{\alpha \beta} \partial^{\alpha} u\right)+\left[a_{\alpha \beta}, \partial_{i}^{h}\right] \partial^{\alpha} u\right\} \cdot \partial^{\beta} v d x \\
& =:-\left\langle\mathcal{L} u, \partial_{i}^{-h} v\right\rangle-\int_{\mathbb{H}^{d}} \sum_{|\alpha|,|\beta| \leq 1}\left(\partial_{h}^{i} a_{\alpha \beta}\right) \tau_{h}^{i} \partial^{\alpha} u \cdot \partial^{\beta} v d x \\
& =-\left\langle\mathcal{L} u, \partial_{i}^{-h} v\right\rangle-\mathcal{E}_{\partial_{h}^{i} a}\left(\tau_{h}^{i} u, v=-\left(f, \partial_{i}^{-h} v\right)-\mathcal{E}_{\partial_{h}^{i} a}\left(\tau_{h}^{i} u, v\right)\right. \\
& =\left(\partial_{i}^{h} f, v\right)-\mathcal{E}_{\partial_{h}^{i} a}\left(\tau_{h}^{i} u, v\right)
\end{aligned}
$$

wherein we have made use of Eq. (52.10) in the third equality. From this it follows that

$$
\mathcal{L} \partial_{i}^{h} u=\partial_{i}^{h} \mathcal{L} u-\mathcal{E}_{\partial_{h}^{i} a}\left(\tau_{h}^{i} u, \cdot\right) \in X^{*}
$$

and

$$
\begin{aligned}
\left\|\mathcal{L} \partial_{i}^{h} u\right\|_{X^{*}} & \leq\left\|\partial_{i}^{h} \mathcal{L} u\right\|_{X^{*}}+\left\|\mathcal{E}_{\partial_{h}^{i} a}\left(\tau_{h}^{i} u, \cdot\right)\right\|_{X^{*}} \lesssim\|\mathcal{L} u\|_{L^{2}}+\|u\|_{X} \\
& \lesssim\|\mathcal{L} u\|_{L^{2}}+\|\mathcal{L} u\|_{X^{*}}+\|u\|_{X^{*}} \lesssim\|\mathcal{L} u\|_{L^{2}}+\|u\|_{X^{*}}
\end{aligned}
$$

Therefore,

$$
\begin{aligned}
\left\|\partial_{i}^{h} u\right\|_{X} & \lesssim\left\|\mathcal{L} \partial_{i}^{h} u\right\|_{X^{*}}+\left\|\partial_{i}^{h} u\right\|_{X^{*}} \lesssim\|\mathcal{L} u\|_{L^{2}}+\|u\|_{X^{*}}+\|u\|_{L^{2}} \\
& \lesssim\|\mathcal{L} u\|_{L^{2}}+\|u\|_{X^{*}}
\end{aligned}
$$

Since $h$ is small but arbitrary we conclude that $\partial_{i} u \in X$ and

$$
\left\|\partial_{i} u\right\|_{X} \lesssim\|\mathcal{L} u\|_{L^{2}}+\|u\|_{X^{*}} \text { for all } i<d
$$

Finally if $i=d$, we have that $f=\mathcal{L} u=\sum_{\alpha \neq 2 e_{d}} A_{\alpha} \partial^{\alpha} u+\partial_{d}^{2} u$ which implies (writing $A_{d, d}$ for $A_{2 e_{d}}$ )

$$
\partial_{d}^{2} u=A_{d, d}^{-1}\left(f-\sum_{\alpha \neq 2 e_{d}} A_{\alpha} \partial^{\alpha} u\right) \in L^{2}
$$

because we have shown that $\partial_{i} \partial_{j} u \in L^{2}$ if $\{i, j\} \neq\{d\}$. Moreover we have the estimate that

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Given $\rho<t$, there exists $C<\infty$ such that for all $\gamma \in \mathbb{N}_{0}^{d-1} \times\{0\}$ with $|\gamma| \leq k+1, \partial^{\gamma} u \in H^{1}\left(N_{\rho}\right)$ and

$$
\begin{equation*}
\left\|\partial^{\gamma} u\right\|_{H^{1}\left(N_{\rho}\right)} \leq C\left(\|f\|_{H^{k}\left(N_{t}\right)}+\|u\|_{H^{1}\left(N_{t}\right)}\right) \tag{52.17}
\end{equation*}
$$

Proof. Let $\rho<r<s<t$ and consider the half nested balls as in Figure 52.4 below. The proof will be by induction on $j=|\gamma|$. When $j=0$ the


Fig. 52.4. A collection of nested half balls along with the cutoff function $\phi$.
assertion is trivial. Assume now there exists $j \in[1, k+1] \cap \mathbb{N}$ such that $\partial^{\gamma} u \in H^{1}\left(N_{s}\right)$ for all $|\gamma|<j$ with $\gamma_{d}=0$ and

$$
\left\|\partial^{\gamma} u\right\|_{H^{1}\left(N_{s}\right)} \leq C\left(\|f\|_{H^{k}\left(N_{t}\right)}+\|u\|_{H^{1}\left(N_{t}\right)}\right)
$$

Fix $\phi \in C_{c}^{\infty}\left(\bar{N}_{t}\right)$ such that $\phi=0$ on $\bar{N}_{t} \backslash \bar{N}_{r}$ and $\phi=1$ in a neighborhood of $\bar{N}_{\rho}$. Suppose $\gamma$ is a multi-index such that $|\gamma|=j$ and $\gamma_{d}=0$. Then $\partial_{h}^{\gamma}(\phi u) \in$ $X_{r}$ for $h$ sufficiently small.

With out loss of generality we may assume $\gamma_{1}>0$ and write $\gamma=e_{1}+\gamma^{\prime}$ and $\partial_{h}^{\gamma}=\partial_{h}^{1} \partial_{h}^{\gamma^{\prime}}$. For $v \in X_{r}$,

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$$
\begin{aligned}
& E_{1}:=\sum_{|\alpha|,|\beta| \leq 1} \int_{N_{t}}\left[a_{\alpha \beta}, \partial_{h}^{\gamma}\right] \partial^{\alpha}(\phi u) \cdot \partial^{\beta} v \\
& E_{2}:=\sum_{|\alpha|,|\beta| \leq 1} \int_{N_{t}} \partial_{h}^{\gamma}\left(a_{\alpha \beta}\left[\partial^{\alpha}, \phi\right] u\right) \cdot \partial^{\beta} v \\
& E_{3}:=\sum_{|\alpha|,|\beta| \leq 1}(-1)^{|\gamma|} \int_{N_{t}} a_{\alpha \beta} \partial^{\alpha} u \cdot\left[\phi, \partial^{\beta}\right] \partial_{-h}^{\gamma} v \text { and } \\
& E_{4}:=-\int_{N_{t}} \partial_{h}^{\gamma^{\prime}}[\phi f] \cdot \partial_{-h}^{1} v .
\end{aligned}
$$

To finish the proof we will estimate each of the terms $E_{i}$ for $i=1, \ldots, 4$. Using Lemma 52.10,

$$
\begin{aligned}
\left|E_{1}\right| & \leq \sum_{|\alpha|,|\beta| \leq 1} \int_{N_{t}}\left|\left[a_{\alpha \beta}, \partial_{h}^{\gamma}\right] \partial^{\alpha}(\phi u) \cdot \partial^{\beta} v\right| \\
& \leq\|v\|_{H^{1}\left(N_{r}\right)} \sum_{|\alpha|,|\beta| \leq 1}\left\|\left[a_{\alpha \beta}, \partial_{h}^{\gamma}\right] \partial^{\alpha}(\phi u)\right\|_{L^{2}\left(N_{r}\right)} \\
& \leq\|v\|_{H^{1}\left(N_{r}\right)} \sum_{|\alpha|,|\beta| \leq 1} \sum_{\delta<\gamma} C_{\gamma}\left(a_{\alpha \beta}\right)\left\|\partial^{\delta} \partial^{\alpha}(\phi u)\right\|_{L^{2}\left(N_{r}\right)} \\
& \lesssim\|v\|_{H^{1}\left(N_{r}\right)} \sum_{\delta<\gamma}\left\|\partial^{\delta} u\right\|_{H^{1}\left(N_{r}\right)} \\
& \lesssim\|v\|_{H^{1}\left(N_{s}\right)}\left(\|f\|_{H^{k}\left(N_{s}\right)}+\|u\|_{H^{1}\left(N_{s}\right)}\right) \quad \text { (by induction). }
\end{aligned}
$$

For $E_{2}$,

$$
\begin{aligned}
\left|E_{2}\right| & =\left|\sum_{|\beta| \leq 1,|\alpha|=1} \int_{N_{t}} \partial_{h}^{\gamma}\left(a_{\alpha \beta}\left(\partial^{\alpha} \phi\right) u\right) \cdot \partial^{\beta} v\right| \\
& \leq\|v\|_{H^{1}\left(N_{r}\right)} \sum_{|\beta| \leq 1,|\alpha|=1}\left\|\partial_{h}^{\gamma}\left[a_{\alpha \beta}\left(\partial^{\alpha} \phi\right) u\right]\right\|_{L^{2}\left(N_{r}\right)} \\
& \leq C\|v\|_{H^{1}\left(N_{r}\right)} \sum_{|\beta| \leq 1,|\alpha|=1}\left\|\partial^{\gamma}\left[a_{\alpha \beta}\left(\partial^{\alpha} \phi\right) u\right]\right\|_{L^{2}\left(N_{r}\right)} \\
& \leq C\|v\|_{H^{1}\left(N_{r}\right)} \sum_{\delta \leq \gamma}\left\|\partial^{\delta} u\right\|_{L^{2}\left(N_{r}\right)} \\
& \leq C\|v\|_{H^{1}\left(N_{r}\right)} \sum_{|\delta| \leq j-1, \delta_{n=0}}\left\|\partial^{\delta} u\right\|_{H^{1}\left(N_{r}\right)} \\
& \left.\leq C\|v\|_{H^{1}\left(N_{s}\right)}\|f\|_{H^{k}\left(N_{s}\right)}+\|u\|_{H^{1}\left(N_{s}\right)}\right) \text { (by induction). }
\end{aligned}
$$

For $E_{3}$,

$$
\begin{aligned}
\left|E_{3}\right| & \leq \sum_{|\alpha| \leq 1,|\beta|=1}\left|\int_{N_{t}} a_{\alpha \beta} \partial^{\alpha} u \cdot\left(\partial^{\beta} \phi\right) \partial_{-h}^{\gamma} v\right| \\
& =\sum_{|\alpha| \leq 1,|\beta|=1}\left|\int_{N_{t}} \partial_{h}^{\gamma^{\prime}}\left[a_{\alpha \beta} \partial^{\alpha} u \cdot\left(\partial^{\beta} \phi\right)\right] \cdot \partial_{-h}^{i} v\right| \\
& \leq \sum_{|\alpha| \leq 1,|\beta|=1}\|v\|_{H^{1}\left(N_{r}\right)}\left\|\partial^{\gamma^{\prime}}\left[a_{\alpha \beta} \partial^{\alpha} u \cdot\left(\partial^{\beta} \phi\right)\right]\right\|_{L^{2}\left(N_{r}\right)} \\
& \leq C\|v\|_{H^{1}\left(N_{r}\right)} \sum_{|\alpha| \leq 1} \sum_{\delta \leq \gamma^{\prime}}\left\|\partial^{\delta+\alpha} u\right\|_{L^{2}\left(N_{r}\right)} \\
& \leq C\|v\|_{H^{1}\left(N_{r}\right)} \sum_{\delta \leq \gamma^{\prime}}\left\|\partial^{\delta} u\right\|_{H^{1}\left(N_{r}\right)} \\
& \leq C\|v\|_{H^{1}\left(N_{s}\right)}\left(\|f\|_{H^{k}\left(N_{s}\right)}+\|u\|_{H^{1}\left(N_{s}\right)}\right) \text { (by the induction hypothesis). }
\end{aligned}
$$

Finally for $E_{4}$,

$$
\begin{aligned}
\left|E_{4}\right| & =\left|\int_{N_{t}} \partial_{h}^{\gamma^{\prime}}[\phi f] \cdot \partial_{-h}^{1} v\right| \leq\left\|\partial_{-h}^{1} v\right\|_{L^{2}\left(N_{r}\right)}\left\|\partial_{h}^{\gamma^{\prime}}(\phi f)\right\|_{L^{2}\left(N_{r}\right)} \\
& \leq\|v\|_{H^{1}\left(N_{r}\right)}\left\|\partial^{\gamma^{\prime}}(\phi f)\right\|_{L^{2}\left(N_{r}\right)} \\
& \leq\|v\|_{H^{1}\left(N_{s}\right)}\|\phi f\|_{H^{j-1}\left(N_{r}\right)} \leq C\|v\|_{H^{1}\left(N_{s}\right)}\|f\|_{H^{k}\left(N_{s}\right)} .
\end{aligned}
$$

Putting all of these estimates together proves, whenever $|\gamma|=j$,

$$
\begin{equation*}
\left|Q\left(\partial_{h}^{\gamma}(\phi u), v\right)\right| \leq C\|v\|_{H^{1}\left(N_{s}\right)}\left(\|f\|_{H^{k}\left(N_{s}\right)}+\|u\|_{H^{1}\left(N_{s}\right)}\right) \tag{52.18}
\end{equation*}
$$

for all $v \in X_{s}$. In particular we may take $v=\partial_{h}^{\gamma}(\phi u) \in X_{s}$ in the above inequality to learn

$$
\begin{equation*}
Q\left(\partial_{h}^{\gamma}(\phi u), \partial_{h}^{\gamma}(\phi u)\right) \leq C\left\|\partial_{h}^{\gamma}(\phi u)\right\|_{H^{1}\left(N_{s}\right)}\left(\|f\|_{H^{k}\left(N_{s}\right)}+\|u\|_{H^{1}\left(N_{s}\right)}\right) . \tag{52.19}
\end{equation*}
$$

But by coercivity of $Q$,

$$
\begin{align*}
\left\|\partial_{h}^{\gamma}(\phi u)\right\|_{H^{1}\left(N_{s}\right)}^{2} & \leq C\left[Q\left(\partial_{h}^{\gamma}(\phi u), \partial_{h}^{\gamma}(\phi u)\right)+\left\|\partial_{h}^{\gamma}(\phi u)\right\|_{L^{2}\left(N_{s}\right)}^{2}\right] \\
& \lesssim\left\|\partial_{h}^{\gamma}(\phi u)\right\|_{H^{1}\left(N_{s}\right)}\left(\|f\|_{H^{k}\left(N_{s}\right)}+\|u\|_{H^{1}\left(N_{s}\right)}\right) \\
& +\left\|\partial_{h}^{\gamma}(\phi u)\right\|_{H^{1}\left(N_{s}\right)}\left\|\partial_{h}^{\gamma}(\phi u)\right\|_{L^{2}\left(N_{s}\right)} \\
& \lesssim\left\|\partial_{h}^{\gamma}(\phi u)\right\|_{H^{1}\left(N_{s}\right)}\binom{\|f\|_{H^{k}\left(N_{s}\right)}+\|u\|_{H^{1}\left(N_{s}\right)}}{+\left\|\partial_{h}^{\gamma}(\phi u)\right\|_{L^{2}\left(N_{s}\right)}} \tag{52.20}
\end{align*}
$$

and hence

$$
\begin{equation*}
\left\|\partial_{h}^{\gamma}(\phi u)\right\|_{H^{1}\left(N_{s}\right)} \lesssim\|f\|_{H^{k}\left(N_{s}\right)}+\|u\|_{H^{1}\left(N_{s}\right)}+\left\|\partial_{h}^{\gamma}(\phi u)\right\|_{L^{2}\left(N_{s}\right)} . \tag{52.21}
\end{equation*}
$$

$$
L=\sum_{|\alpha|,|\beta| \leq 1}(-1)^{|\beta|} \partial^{\beta} a_{\alpha \beta} \partial_{\alpha}=\sum_{|\alpha| \leq 2} A_{\alpha} \partial^{\alpha},
$$

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then $L u=f$ in the distributional sense. Writing $\tilde{A}$ for $A_{(0,0, \ldots, 0,2)}$,

$$
f=\tilde{A} \partial_{d}^{2} u+\sum_{|\alpha| \leq 2, \alpha_{d}<2} A_{\alpha} \partial^{\alpha} u
$$

so that

$$
\partial_{d}^{2} u=\frac{1}{\tilde{A}}\left(f-\sum_{|\alpha| \leq 2, \alpha_{d}<2} A_{\alpha} \partial^{\alpha} u\right)
$$

and

$$
\begin{equation*}
\partial^{\gamma} u=\partial^{\gamma^{\prime}} \partial_{d}^{2} u=\partial^{\gamma^{\prime}}\left(\frac{1}{\tilde{A}} f-\sum_{|\alpha| \leq 2, \alpha_{d}<2} \frac{A_{\alpha}}{\tilde{A}} \partial^{\alpha} u\right) . \tag{52.23}
\end{equation*}
$$

Now by the product rule

$$
\begin{equation*}
\sum_{|\alpha| \leq 2, \alpha_{d}<2} \partial^{\gamma^{\prime}}\left(\frac{A_{\alpha}}{\tilde{A}} \partial^{\alpha} u\right):=\sum_{|\alpha| \leq 2, \alpha_{d}<2, \delta \leq \gamma^{\prime}}\binom{\gamma^{\prime}}{\delta} \partial^{\left(\gamma^{\prime}-\delta+\alpha\right)}\left(\frac{A_{\alpha}}{\tilde{A}}\right) \cdot \partial^{(\delta+\alpha)} u \tag{52.24}
\end{equation*}
$$

Since $\left(\gamma^{\prime}+\alpha\right)_{d}<j$, the induction hypothesis (i.e. Eq. (52.22) is valid for $|\gamma|<j$ ) shows the right member of Eq. (52.24) is in $L^{2}\left(N_{\rho}\right)$ and gives the estimate

$$
\begin{aligned}
\left\|\sum_{|\alpha| \leq 2, \alpha_{d}<2} \partial^{\gamma^{\prime}}\left(\frac{A_{\alpha}}{\tilde{A}} \partial^{\alpha} u\right)\right\|_{L^{2}\left(N_{\rho}\right)} & \lesssim \sum_{|\alpha| \leq 2, \alpha_{d}<2, \delta \leq \gamma^{\prime}}\left\|\partial^{(\delta+\alpha)} u\right\|_{L^{2}\left(N_{\rho}\right)} \\
& \lesssim\|f\|_{H^{k}\left(N_{t}\right)}+\|u\|_{H^{1}\left(N_{t}\right)}
\end{aligned}
$$

Combining this with Eq. (52.23) gives $\partial^{\gamma} u \in L^{2}\left(N_{\rho}\right)$ and

$$
\begin{align*}
\left\|\partial^{\gamma} u\right\|_{L^{2}\left(N_{\rho}\right)} & \lesssim\|f\|_{H^{\left|\gamma^{\prime}\right|\left(N_{t}\right)}}+\left\|\sum_{|\alpha| \leq 2, \alpha_{d}<2} \partial^{\gamma^{\prime}}\left(\frac{A_{\alpha}}{\tilde{A}} \partial^{\alpha} u\right)\right\|_{L^{2}\left(N_{\rho}\right)} \\
& \lesssim\|f\|_{H^{k}\left(N_{t}\right)}+\|f\|_{H^{k}\left(N_{t}\right)}+\|u\|_{H^{1}\left(N_{t}\right)} \lesssim\|f\|_{H^{k}\left(N_{t}\right)}+\|u\|_{H^{1}\left(N_{t}\right)} . \tag{52.25}
\end{align*}
$$

The following assumptions an notation will be in force for the remainder of this chapter.

Assumption 5 Let $\Omega$ be a bounded open subset such that $\bar{\Omega}^{o}=\Omega$ and $\bar{\Omega}$ is a $C^{\infty}$ - manifold with boundary, $X$ be either $H_{0}^{1}(\Omega)$ or $H^{1}(\Omega)$ and $\mathcal{E}$ be a Dirichlet form as in Notation 51.1 which is assumed to be elliptic. Also if $W$ is an open subset of $\mathbb{R}^{d}$ let

$$
X_{W}:=\{v \in X: \operatorname{supp}(v) \sqsubset \sqsubset W \cap \bar{\Omega}\} .
$$

Lemma 52.14. For each $p \in \partial \Omega$ there exists precompact open neighborhoods $V$ and $W$ in $\mathbb{R}^{d}$ such that $\bar{V} \subset W$, for each $k \in \mathbb{N}$ there is a constant $C_{k}<\infty$ such that if $u \in X$ and $f \in H^{k}(\Omega)$ satisfies

$$
\begin{equation*}
\mathcal{E}(u, v)=\int_{\Omega} f v d x \text { for all } v \in X_{W} \tag{52.26}
\end{equation*}
$$

then $u \in H^{k+2}(V \cap \Omega)$ and

$$
\begin{equation*}
\|u\|_{H^{k+2}(V \cap \Omega)} \leq C\left(\|f\|_{H^{k}(\Omega)}+\|u\|_{H^{1}(\Omega)}\right) \tag{52.27}
\end{equation*}
$$

Proof. Let $W$ be an open neighborhood of $p$ such that there exists a chart $\psi: W \rightarrow B(0, r)$ with inverse $\phi:=\psi^{-1}: B(0, r) \rightarrow W$ satisfying:

1. The maps $\psi$ and $\phi$ has bounded derivatives to all orders.
2. $\psi(W \cap \Omega)=B(0, r) \cap \mathbb{H}^{d}=N_{r}$ and $\left.\psi(W \cap \operatorname{bd}(\Omega))=B(0, r) \cap \operatorname{bd}\left(\mathbb{H}^{d}\right)\right)$.

Now let $\rho<r$ and define $V:=\phi(B(0, \rho))$, see Figure 52.5.


Fig. 52.5. Flattening out the boundary of $\Omega$ in a neighborhood of $p$.

Suppose that $u \in X$ satisfies Eq. (52.26) and $v \in X_{W}$. Then making the change of variables $x=\phi(y)$,

$$
\int_{\Omega} f v d m=\int_{N_{r}} f(\phi(y)) v(\phi(y)) J(y) d y=\int_{N_{r}} \tilde{f}(y) \tilde{v}(y) d y
$$

Proof. Cover $\partial \Omega$ with $\left\{V_{i}\right\}_{i=1}^{N}$ and $\left\{W_{i}\right\}_{i=1}^{N}$ as in the above Lemma 52.14 such that $\bar{V}_{i} \sqsubset ᄃ W_{i}$. Also choose a precompact open subset $V_{0}$ contained in $\Omega$ such that $\left\{V_{i}\right\}_{i=0}^{N}$ covers $\bar{\Omega}$. Choose $W_{0}$ such that $\bar{V}_{0} \subset W_{0}$ and $\bar{W}_{0} \subset \Omega$. If $\mathcal{L}_{\mathcal{E}} u=: f \in H^{k}(\Omega)$, then by Lemma 52.14 for $i \geq 1$ and Theorem 52.5 for $i=0, u \in H^{k+2}\left(V_{i}\right)$ and there exist $C_{i}<\infty$ such that

$$
\begin{equation*}
\|u\|_{H^{k+2}\left(V_{i} \cap \Omega\right)} \leq C_{i}\left(\|f\|_{H^{k}\left(W_{i} \cap \Omega\right)}+\|u\|_{H^{1}\left(W_{i} \cap \Omega\right)}\right) \tag{52.30}
\end{equation*}
$$

Summing Eq. (52.30) on $i$ implies $u \in H^{k+2}(\Omega)$ and

$$
\begin{equation*}
\|u\|_{H^{k+2}(\Omega)} \leq C\left(\|f\|_{H^{k}(\Omega)}+\|u\|_{X}\right) \tag{52.31}
\end{equation*}
$$

Finally

$$
\begin{aligned}
\|u\|_{X}^{2} & \leq C\left(\mathcal{E}(u, u)+\|u\|_{H^{-1}(\Omega)}^{2}\right) \\
& =C\left((f, u)_{L^{2}(\Omega)}+\|u\|_{H^{-1}(\Omega)}^{2}\right) \\
& \leq C\left(\|f\|_{L^{2}(\Omega)}\|u\|_{L^{2}(\Omega)}+\|u\|_{H^{-1}(\Omega)}^{2}\right) \\
& \leq C\left(\frac{1}{2 \delta}\|f\|_{L^{2}(\Omega)}^{2}+\frac{\delta}{2}\|u\|_{X}^{2}+\|u\|_{H^{-1}(\Omega)}^{2}\right)
\end{aligned}
$$

for any $\delta>0$. Choosing $\delta$ so that $C \delta=1$, we find

$$
\frac{1}{2}\|u\|_{X}^{2} \leq C\left(\frac{1}{2 \delta}\|f\|_{L^{2}(\Omega)}^{2}+\|u\|_{H^{-1}(\Omega)}^{2}\right)
$$

which implies with a new constant $C$ that

$$
\begin{equation*}
\|u\|_{X} \leq C\left(\|f\|_{L^{2}(\Omega)}+\|u\|_{H^{-1}(\Omega)}\right) . \tag{52.32}
\end{equation*}
$$

Combining Eqs. (52.31) and (52.32) implies Eq. (52.29).

## Unbounded operators and quadratic forms

### 53.1 Unbounded operator basics

Definition 53.1. If $X$ and $Y$ are Banach spaces and $D$ is a subspace of $X$, then a linear transformation $T$ from $D$ into $Y$ is called a linear transformation (or operator) from $X$ to $Y$ with domain $D$. We will sometimes wr If $D$ is dense in $X, T$ is said to be densely defined.

Notation 53.2 If $S$ and $T$ are operators from $X$ to $Y$ with domains $D(S)$ and $D(T)$ and if $D(S) \subset D(T)$ and $S x=T x$ for $x \in D(S)$, then we say $T$ is an extension of $S$ and write $S \subset T$.

We note that $X \times Y$ is a Banach space in the norm

$$
\|\langle x, y\rangle\|=\sqrt{\|x\|^{2}+\|y\|^{2}} .
$$

If $H$ and $K$ are Hilbert spaces, then $H \times K$ and $K \times H$ become Hilbert spaces by defining

$$
\left(\langle x, y\rangle,\left\langle x^{\prime}, y^{\prime}\right\rangle\right)_{H \times K}:=\left(x, x^{\prime}\right)_{H}+\left(y, y^{\prime}\right)_{K}
$$

and

$$
\left(\langle y, x\rangle,\left\langle y^{\prime}, x^{\prime}\right\rangle\right)_{K \times H}:=\left(x, x^{\prime}\right)_{H}+\left(y, y^{\prime}\right)_{K} .
$$

Definition 53.3. If $T$ is an operator from $X$ to $Y$ with domain $D$, the graphof $T$ is

$$
\Gamma(T):=\{\langle x, D x\rangle: x \in D(T)\} \subset H \times K .
$$

Note that $\Gamma(T)$ is a subspace of $X \times Y$.
Definition 53.4. An operator $T: X \rightarrow Y$ is closed if $\Gamma(T)$ is closed in $X \times Y$.
Remark 53.5. It is easy to see that $T$ is closed iff for all sequences $x_{n} \in D$ such that there exists $x \in X$ and $y \in Y$ such that $x_{n} \rightarrow x$ and $T x_{n} \rightarrow y$ implies that $x \in \mathcal{D}$ and $T x=y$.

Let $H$ be a Hilbert space with inner product $(\cdot, \cdot)$ and norm $\|v\|:=\sqrt{(v, v)}$. As usual we will write $H^{*}$ for the continuous dual of $H$ and $\overline{H^{*}}$ for the continuous conjugate linear functionals on $H$. Our convention will be that $(\cdot, v) \in H^{*}$ is linear while $(v, \cdot) \in \overline{H^{*}}$ is conjugate linear for all $v \in H$.

Lemma 53.6. Suppose that $T: H \rightarrow K$ is a densely defined operator between two Hilbert spaces $H$ and K. Then

1. $T^{*}$ is always a closed but not necessarily densely defined operator.
2. If $T$ is closable, then $\bar{T}^{*}=T^{*}$.
3. $T$ is closable iff $T^{*}: K \rightarrow H$ is densely defined.
4. If $T$ is closable then $\bar{T}=T^{* *}$.

Proof. Suppose $\left\{v_{n}\right\} \subset D(T)$ is a sequence such that $v_{n} \rightarrow 0$ in $H$ and $T v_{n} \rightarrow k$ in $K$ as $n \rightarrow \infty$. Then for $l \in D\left(T^{*}\right)$, by passing to the limit in the equality, $\left(T v_{n}, l\right)=\left(v_{n}, T^{*} l\right)$ we learn $(k, l)=\left(0, T^{*} l\right)=0$. Hence if $T^{*}$ is densely defined, this implies $k=0$ and hence $T$ is closable. This proves one direction in item 3. To prove the other direction and the remaining items of the Lemma it will be useful to express the graph of $T^{*}$ in terms of the graph of $T$. We do this now.

Recall that $k \in D\left(T^{*}\right)$ and $T^{*} k=h$ iff $(k, T x)_{K}=(h, x)_{H}$ for all $x \in$ $D(T)$. This last condition may be written as $(k, y)_{K}-(h, x)_{H}=0$ for all $\langle x, y\rangle \in \Gamma(T)$.

Let $V: H \times K \rightarrow K \times H$ be the unitary map defined by $V\langle x, y\rangle=\langle-y, x\rangle$. With this notation, we have $\langle k, h\rangle \in \Gamma\left(T^{*}\right)$ iff $\langle k, h\rangle \perp V \Gamma(T)$, i.e.

$$
\begin{equation*}
\Gamma\left(T^{*}\right)=(V \Gamma(T))^{\perp}=V\left(\Gamma(T)^{\perp}\right) \tag{53.1}
\end{equation*}
$$

where the last equality is a consequence of $V$ being unitary. As a consequence of Eq. (53.1), $\Gamma\left(T^{*}\right)$ is always closed and hence $T^{*}$ is always a closed operator, and this proves item 1 . Moreover if $T$ is closable, then

$$
\Gamma\left(T^{*}\right)=V \Gamma(T)^{\perp}=V \overline{\Gamma(T)}^{\perp}=V \Gamma(\bar{T})^{\perp}=\Gamma\left(\bar{T}^{*}\right)
$$

which proves item 2.
Now suppose $T$ is closable and $k \perp \mathcal{D}\left(T^{*}\right)$. Then

$$
\langle k, 0\rangle \in \Gamma\left(T^{*}\right)^{\perp}=V \Gamma(T)^{\perp \perp}=V \overline{\Gamma(T)}=V \Gamma(\bar{T}),
$$

where $\bar{T}$ denotes the closure of $T$. This implies that $\langle 0, k\rangle \in \Gamma(\bar{T})$. But $\bar{T}$ is a well defined operator (by the assumption that $T$ is closable) and hence $k=\bar{T} 0=0$. Hence we have shown $\mathcal{D}\left(T^{*}\right)^{\perp}=\{0\}$ which implies $\mathcal{D}\left(T^{*}\right)$ is dense in $K$. This completes the proof of item 3 .
4. Now assume $T$ is closable so that $T^{*}$ is densely defined. Using the obvious analogue of Eq. (53.1) for $T^{*}$ we learn $\Gamma\left(T^{* *}\right)=U \Gamma\left(T^{*}\right)^{\perp}$ where $U\langle y, x\rangle=\langle-x, y\rangle=-V^{-1}\langle y, x\rangle$. Therefore,

$$
\Gamma\left(T^{* *}\right)=U V\left(\Gamma(T)^{\perp}\right)^{\perp}=-\overline{\Gamma(T)}=\overline{\Gamma(T)}=\Gamma(\bar{T})
$$

and hence $\bar{T}=T^{* *}$.

Lemma 53.7. Suppose that $H$ and $K$ are Hilbert spaces, $T: H \rightarrow K$ is a densely defined operator which has a densely defined adjoint $T^{*}$. Then $\operatorname{Nul}\left(T^{*}\right)=\operatorname{Ran}(T)^{\perp}$ and $\operatorname{Nul}(\bar{T})=\operatorname{Ran}\left(T^{*}\right)^{\perp}$ where $\bar{T}$ denotes the closure of $T$.

Proof. Suppose that $k \in \operatorname{Nul}\left(T^{*}\right)$ and $h \in \mathcal{D}(T)$, then $(k, T h)=$ $\left(T^{*} k, h\right)=0$. Since $h \in \mathcal{D}(T)$ is arbitrary, this proves that $\operatorname{Nul}\left(T^{*}\right) \subset$ $\operatorname{Ran}(T)^{\perp}$. Now suppose that $k \in \operatorname{Ran}(T)^{\perp}$. Then $0=(k, T h)$ for all $h \in D(T)$. This shows that $k \in \mathcal{D}\left(T^{*}\right)$ and that $T^{*} k=0$. The assertion $\operatorname{Nul}(\bar{T})=$ $\operatorname{Ran}\left(T^{*}\right)^{\perp}$ follows by replacing $T$ by $T^{*}$ in the equality, $\operatorname{Nul}\left(T^{*}\right)=\operatorname{Ran}(T)^{\perp}$.

Definition 53.8. A quadratic form $q$ on $H$ is a dense subspace $\mathcal{D}(q) \subset H$ called the domain of $q$ and a sesquilinear form $q: \mathcal{D}(q) \times \mathcal{D}(q) \rightarrow \mathbb{C}$. (Sesquilinear means that $q(\cdot, v)$ is linear while $q(v, \cdot)$ is conjugate linear on $\mathcal{D}(q)$ for all $v \in \mathcal{D}(q)$.) The form $q$ is symmetric if $q(v, w)=\overline{q(w, v)}$ for all $v, w \in \mathcal{D}(q), q$ is positive if $q(v) \geq 0$ (here $q(v)=q(v, v)$ ) for all $v \in \mathcal{D}(q)$, and $q$ is semi-bounded if there exists $M_{0} \in(0, \infty)$ such that $q(v, v) \geq-M_{0}\|v\|^{2}$ for all $v \in \mathcal{D}(q)$.

### 53.2 Lax-Milgram Methods

For the rest of this section $q$ will be a sesquilinear form on $H$ and to simplify notation we will write $X$ for $\mathcal{D}(q)$.

Theorem 53.9 (Lax-Milgram). Let $q: X \times X \rightarrow \mathbb{C}$ be a sesquilinear form and suppose the following added assumptions hold.

1. $X$ is equipped with a Hilbertian inner product $(\cdot, \cdot)_{X}$.
2. The form $q$ is bounded on $X$, i.e. there exists a constant $C<\infty$ such that $|q(v, w)| \leq C\|v\|_{X} \cdot\|w\|_{X}$ for all $v, w \in X$.
3. The form $q$ is coercive, i.e. there exists $\epsilon>0$ such that $|q(v, v)| \geq \epsilon\|v\|_{X}^{2}$ for all $v \in X$.

Then the maps $\mathcal{L}: X \rightarrow \overline{X^{*}}$ and $\mathcal{L}^{\dagger}: X \rightarrow X^{*}$ defined by $\mathcal{L} v:=q(v, \cdot)$ and $\mathcal{L}^{\dagger} v:=q(\cdot, v)$ are linear and (respectively) conjugate linear isomorphisms of Hilbert spaces. Moreover

$$
\left\|\mathcal{L}^{-1}\right\| \leq \epsilon^{-1} \text { and }\left\|\left(\mathcal{L}^{\dagger}\right)^{-1}\right\| \leq \epsilon^{-1} .
$$

Proof. The operator $\mathcal{L}$ is bounded because

$$
\begin{equation*}
\|\mathcal{L} v\|_{X^{*}}=\sup _{w \neq 0} \frac{|q(v, w)|}{\|w\|_{X}} \leq C\|v\|_{X} \tag{53.2}
\end{equation*}
$$

Similarly $\mathcal{L}^{\dagger}$ is bounded with $\left\|\mathcal{L}^{\dagger}\right\| \leq C$.

Let $\beta: X \rightarrow \overline{X^{*}}$ denote the linear Riesz isomorphism defined by $\beta(x)=$ $(x, \cdot)_{X}$ for $x \in X$. Define $R:=\beta^{-1} \mathcal{L}: X \rightarrow X$ so that $\mathcal{L}=\beta R$, i.e.

$$
\mathcal{L} v=q(v, \cdot)=(R v, \cdot)_{X} \text { for all } v \in X
$$

Notice that $R$ is a bounded linear map with operator bound less than $C$ by Eq. (53.2). Since

$$
\left(\mathcal{L}^{\dagger} v\right)(w)=q(w, v)=(R w, v)_{X}=\left(w, R^{*} v\right)_{X} \text { for all } v, w \in X
$$

we see that $\mathcal{L}^{\dagger} v=\left(\cdot, R^{*} v\right)_{X}$, i.e. $R^{*}=\bar{\beta}^{-1} \mathcal{L}^{\dagger}$, where $\bar{\beta}(x):=\overline{(x, \cdot)_{X}}=(\cdot, x)_{X}$. Since $\beta$ and $\bar{\beta}$ are linear and conjugate linear isometric isomorphisms, to finish the proof it suffices to show $R$ is invertible and that $\left\|R^{-1}\right\|_{X} \leq \epsilon^{-1}$.

Since

$$
\begin{equation*}
\left|\left(v, R^{*} v\right)_{X}\right|=\left|(R v, v)_{X}\right|=|q(v, v)| \geq \epsilon\|v\|_{X}^{2} \tag{53.3}
\end{equation*}
$$

one easily concludes that $\operatorname{Nul}(R)=\{0\}=\operatorname{Nul}\left(R^{*}\right)$. By Lemma 53.7, $\overline{\operatorname{Ran}(R)}=\operatorname{Nul}\left(R^{*}\right)^{\perp}=\{0\}^{\perp}=X$ and so we have shown $R: X \rightarrow X$ is injective and has a dense range. From Eq. (53.3) and the Schwarz inequality, $\epsilon\|v\|_{X}^{2} \leq\|R v\|_{X}\|v\|_{X}$, i.e.

$$
\begin{equation*}
\|R v\|_{X} \geq \epsilon\|v\|_{X} \text { for all } v \in X \tag{53.4}
\end{equation*}
$$

This inequality proves the range of $R$ is closed. Indeed if $\left\{v_{n}\right\}$ is a sequence in $X$ such that $R v_{n} \rightarrow w \in X$ as $n \rightarrow \infty$ then Eq. (53.4) implies

$$
\epsilon\left\|v_{n}-v_{m}\right\|_{X} \leq\left\|R v_{n}-R v_{m}\right\|_{X} \rightarrow 0 \text { as } m, n \rightarrow \infty
$$

Thus $v:=\lim _{n \rightarrow \infty} v_{n}$ exists in $X$ and hence $w=R v \in \operatorname{Ran}(R)$ and so $\operatorname{Ran}(R)=\overline{\operatorname{Ran}(R)}^{X}=X$. So $R: X \rightarrow X$ is a bijective map and hence invertible. By replacing $v$ by $R^{-1} v$ in Eq. (53.4) we learn $R^{-1}$ is bounded with operator norm no larger than $\epsilon^{-1}$.
Theorem 53.10. Let $q$ be a bounded coercive sesquilinear form on $X$ as in Theorem 53.9. Further assume that the inclusion map $i: X \rightarrow H$ is bounded and let $L$ and $L^{\dagger}$ be the unbounded linear operators on $H$ defined by:

$$
\begin{aligned}
\mathcal{D}(L) & :=\{v \in X: w \in X \rightarrow q(v, w) \text { is } H \text { - continuous }\}, \\
\mathcal{D}\left(L^{\dagger}\right) & :=\{w \in X: v \in X \rightarrow q(v, w) \text { is } H \text { - continuous }\}
\end{aligned}
$$

and for $v \in \mathcal{D}(L)$ and $w \in \mathcal{D}\left(L^{\dagger}\right)$ define $L v \in H$ and $L^{\dagger} w \in H$ by requiring

$$
q(v, \cdot)=(L v, \cdot) \text { and } q(\cdot, w)=\left(\cdot, L^{\dagger} w\right)
$$

Then $\mathcal{D}(L)$ and $\mathcal{D}\left(L^{\dagger}\right)$ are dense subspaces of $X$ and hence of $H$. The operators $L^{-1}: H \rightarrow \mathcal{D}(L) \subset H$ and $\left(L^{\dagger}\right)^{-1}: H \rightarrow \mathcal{D}\left(L^{\dagger}\right) \subset H$ are bounded when viewed as operators from $H$ to $H$ with norms less than or equal to $\epsilon^{-1}\|i\|_{L(X, H)}^{2}$. Furthermore, $L^{*}=L^{\dagger}$ and $\left(L^{\dagger}\right)^{*}=L$ and in particular both $L$ and $L^{\dagger}=L^{*}$ are closed operators.

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Proof. Let $\alpha: \underline{H \rightarrow \overline{X^{*}}}$ be defined by $\alpha(v)=\left.(v, \cdot)\right|_{X}$. If $(v, \cdot)_{X}$ is perpendicular to $\alpha(H)=\overline{i^{*}\left(H^{*}\right)} \subset \overline{X^{*}}$, then

$$
0=\left((v, \cdot)_{X}, \alpha(w)\right)_{\overline{X^{*}}}=\left((v, \cdot)_{X},(w, \cdot)\right)_{\overline{X^{*}}}=(v, w) \text { for all } w \in H
$$

Taking $w=v$ in this equation shows $v=0$ and hence the orthogonal complement of $\alpha(H)$ in $\overline{X^{*}}$ is $\{0\}$ which implies $\alpha(H)=\overline{i^{*}\left(H^{*}\right)}$ is dense in $\overline{X^{*}}$.

Using the notation in Theorem 53.9, we have $v \in \mathcal{D}(L)$ iff $\mathcal{L} v \in \overline{i^{*}\left(H^{*}\right)}=$ $\alpha(H)$ iff $v \in \mathcal{L}^{-1}(\alpha(H))$ and for $v \in \mathcal{D}(L), \mathcal{L} v=\left.(L v, \cdot)\right|_{X}=\alpha(L v)$. This and a similar computation shows
$\mathcal{D}(L)=\mathcal{L}^{-1}\left(\overline{i^{*}\left(H^{*}\right)}\right)=\mathcal{L}^{-1}(\alpha(H))$ and $\mathcal{D}\left(L^{\dagger}\right):=\left(\mathcal{L}^{\dagger}\right)^{-1}\left(i^{*}\left(H^{*}\right)\right)=\left(\mathcal{L}^{\dagger}\right)^{-1}(\bar{\alpha}(1$
and for $v \in \mathcal{D}(L)$ and $w \in \mathcal{D}\left(L^{\dagger}\right)$ we have $\mathcal{L} v=\left.(L v, \cdot)\right|_{X}=\alpha(L v)$ and $\mathcal{L}^{\dagger} w=$ $\left.\left(\cdot, L^{\dagger} w\right)\right|_{X}=\bar{\alpha}\left(L^{\dagger} w\right)$. The following commutative diagrams summarizes the relationships of $L$ and $\mathcal{L}$ and $L^{\dagger}$ and $\mathcal{L}^{\dagger}$,

$$
\begin{array}{cccc}
X & \xrightarrow{\mathcal{L}} \overline{X^{*}} & X & \xrightarrow{\mathcal{L}^{\dagger}} X^{*} \\
i & \uparrow & \uparrow \text { and } i & \uparrow
\end{array}
$$

where in each diagram $i$ denotes an inclusion map. Because $\mathcal{L}$ and $\mathcal{L}^{\dagger}$ are invertible, $L: D(L) \rightarrow H$ and $L^{\dagger}: D\left(L^{\dagger}\right) \rightarrow H$ are invertible as well. Because both $\mathcal{L}$ and $\mathcal{L}^{\dagger}$ are isomorphisms of $X$ onto $\overline{X^{*}}$ and $X^{*}$ respectively and $\alpha(H)$ is dense in $\overline{X^{*}}$ and $\bar{\alpha}(H)$ is dense in $X^{*}$, the spaces $\mathcal{D}(L)$ and $\mathcal{D}\left(L^{\dagger}\right)$ are dense subspaces of $X$, and hence also of $H$.

For the norm bound assertions let $v \in \mathcal{D}(L) \subset X$ and use the coercivity estimate on $q$ to find

$$
\begin{aligned}
\epsilon\|v\|_{H}^{2} & \leq \epsilon\|i\|_{L(X, H)}^{2}\|v\|_{X}^{2} \leq\|i\|_{L(X, H)}^{2}|q(v, v)|=\|i\|_{L(X, H)}^{2}\left|(L v, v)_{H}\right| \\
& \leq\|i\|_{L(X, H)}^{2}\|L v\|_{H}\|v\|_{H} .
\end{aligned}
$$

Hence $\epsilon\|v\|_{H} \leq\|i\|_{L(X, H)}^{2}\|L v\|_{H}$ for all $v \in \mathcal{D}(L)$. By replacing $v$ by $L^{-1} v$ (for $v \in H$ ) in this last inequality, we find

$$
\left\|L^{-1} v\right\|_{H} \leq \frac{\|i\|_{L(X, H)}^{2}}{\epsilon}\|v\|_{H}, \text { i..e }\left\|L^{-1}\right\|_{B(H)} \leq \epsilon^{-1}\|i\|_{L(X, H)}^{2}
$$

Similarly one shows that $\left\|\left(L^{\dagger}\right)^{-1}\right\|_{B(H)} \leq \epsilon^{-1}\|i\|_{L(X, H)}^{2}$ as well.
For $v \in \mathcal{D}(L)$ and $w \in \mathcal{D}\left(L^{\dagger}\right)$,

$$
\begin{equation*}
(L v, w)=q(v, w)=\left(v, L^{\dagger} w\right) \tag{53.5}
\end{equation*}
$$

which shows $L^{\dagger} \subset L^{*}$. Now suppose that $w \in \mathcal{D}\left(L^{*}\right)$, then

$$
q(v, w)=(L v, w)=\left(v, L^{*} w\right) \text { for all } v \in \mathcal{D}(L)
$$

By continuity if follows that

$$
q(v, w)=\left(v, L^{*} w\right) \text { for all } v \in X
$$

and therefore by the definition of $L^{\dagger}, w \in \mathcal{D}\left(L^{\dagger}\right)$ and $L^{\dagger} w=L^{*} w$, i.e. $L^{*} \subset L^{\dagger}$. Since we have shown $L^{\dagger} \subset L^{*}$ and $L^{*} \subset L^{\dagger}, L^{\dagger}=L^{*}$. A similar argument shows that $\left(L^{\dagger}\right)^{*}=L$. Because the adjoints of operators are always closed, both $L=\left(L^{\dagger}\right)^{*}$ and $L^{\dagger}=L^{*}$ are closed operators.

Corollary 53.11. If $q$ in Theorem 53.10 is further assumed to be symmetric then $L$ is self-adjoint, i.e. $L^{*}=L$.

Proof. This simply follows from Theorem 53.10 upon observing that $L=$ $L^{\dagger}$ when $q$ is symmetric.

### 53.3 Close, symmetric, semi-bounded quadratic forms and self-adjoint operators

Definition 53.12. A symmetric, sesquilinear quadratic form $q: X \times X \rightarrow \mathbb{C}$ is closed if whenever $\left\{v_{n}\right\}_{n=1}^{\infty} \subset X$ is a sequence such that $v_{n} \rightarrow v$ in $H$ and

$$
q\left(v_{n}-v_{m}\right):=q\left(v_{n}-v_{m}, v_{n}-v_{m}\right) \rightarrow 0 \text { as } m, n \rightarrow \infty
$$

implies that $v \in X$ and $\lim _{n \rightarrow \infty} q\left(v-v_{n}\right)=0$. The form $q$ is said to be closable iff for all $\left\{v_{n}\right\} \subset X$ such that $v_{n} \rightarrow 0 \in H$ and $q\left(v_{n}-v_{m}\right) \rightarrow 0$ as $m, n \rightarrow \infty$ implies that $q\left(v_{n}\right) \rightarrow 0$ as $n \rightarrow \infty$.

Example 53.13. Let $H$ and $K$ be Hilbert spaces and $T: H \rightarrow K$ be a densely defined operator. Set $q(v, w):=(T v, T w)_{K}$ for $v, w \in X:=\mathcal{D}(q):=\mathcal{D}(T)$. Then $q$ is a positive symmetric quadratic form on $H$ which is closed iff $T$ is closed and is closable iff $T$ is closable.

For the remainder of this section let $q: X \times X \rightarrow \mathbb{C}$ be a symmetric, sesquilinear quadratic form which is semi-bounded and satisfies $q(v) \geq$ $-M_{0}\|v\|^{2}$ for all $v \in X$ and some $M_{0}<\infty$.

Notation 53.14 For $v, w \in X$ and $M>M_{0}$ let $(v, w)_{M}:=q(v, w)+$ $M(v, w)$. Notice that

$$
\begin{align*}
\|v\|_{M}^{2} & =q(v)+M\|v\|^{2}=q(v)+M_{0}\|v\|^{2}+\left(M-M_{0}\right)\|v\|^{2} \\
& \geq\left(M-M_{0}\right)\|v\|^{2} \tag{53.6}
\end{align*}
$$

from which it follows that $(\cdot, \cdot)_{M}$ is an inner product on $X$ and $i: X \rightarrow H$ is bounded by $\left(M-M_{0}\right)^{-1 / 2}$. Let $H_{M}$ denote the Hilbert space completion of $\left(X,(\cdot, \cdot)_{M}\right)$.

Formally, $H_{M}=\mathcal{C} / \sim$, where $\mathcal{C}$ denotes the collection of $\|\cdot\|_{M^{-}}$ Cauchy sequences in $X$ and $\sim$ is the equivalence relation, $\left\{v_{n}\right\} \sim\left\{u_{n}\right\}$ iff $\lim _{n \rightarrow \infty}\left\|v_{n}-u_{n}\right\|_{M}=0$. For $v \in X$, let $i(v)$ be the equivalence class of the constant sequence with elements $v$. Notice that if $\left\{v_{n}\right\}$ and $\left\{u_{n}\right\}$ are in $\mathcal{C}$, then $\lim _{m, n \rightarrow \infty}\left(v_{n}, u_{m}\right)_{M}$ exists. Indeed, let $C$ be a finite upper bound for $\left\|u_{n}\right\|_{M}$ and $\left\|v_{n}\right\|_{M}$. (Why does this bound exist?) Then

$$
\begin{align*}
\left|\left(v_{n}, u_{m}\right)_{M}-\left(v_{k}, u_{l}\right)_{M}\right| & =\left|\left(v_{n}-v_{k}, u_{m}\right)_{M}+\left(v_{k}, u_{m}-u_{l}\right)_{M}\right| \\
& \leq C\left\{\left\|v_{n}-v_{k}\right\|_{M}+\left\|u_{m}-u_{l}\right\|_{M}\right\} \tag{53.7}
\end{align*}
$$

and this last expression tends to zero as $m, n, k, l \rightarrow \infty$. Therefore, if $\bar{v}$ and $\bar{u}$ denote the equivalence class of $\left\{v_{n}\right\}$ and $\left\{u_{n}\right\}$ in $\mathcal{C}$ respectively, we may define $(\bar{v}, \bar{u})_{M}:=\lim _{m, n \rightarrow \infty}\left(v_{n}, u_{m}\right)_{M}$. It is easily checked that $H_{M}$ with this inner product is a Hilbert space and that $i: X \rightarrow H_{M}$ is an isometry.

Remark 53.15. The reader should verify that all of the norms, $\left\{\|\cdot\|_{M}: M>M_{0}\right\}$, on $X$ are equivalent so that $H_{M}$ is independent of $M>M_{0}$.

Lemma 53.16. The inclusion map $i: X \rightarrow H$ extends by continuity to $a$ continuous linear map $\hat{\imath}$ from $H_{M}$ into $H$. Similarly, the quadratic form $q: X \times$ $X \rightarrow \mathbb{C}$ extends by continuity to a continuous quadratic form $\hat{q}: H_{M} \times H_{M} \rightarrow$ $\mathbb{C}$. Explicitly, if $\bar{v}$ and $\bar{u}$ denote the equivalence class of $\left\{v_{n}\right\}$ and $\left\{u_{n}\right\}$ in $\mathcal{C}$ respectively, then $\hat{\imath}(\bar{v})=H-\lim _{n \rightarrow \infty} v_{n}$ and $\hat{q}(\bar{v}, \bar{u})=\lim _{m, n \rightarrow \infty} q\left(v_{n}, u_{n}\right)$.

Proof. This routine verification is left to the reader.
Lemma 53.17. Let $q$ be as above and $M>M_{0}$ be given.

1. The quadratic form $q$ is closed iff $\left(X,(\cdot, \cdot)_{M}\right)$ is a Hilbert space.
2. The quadratic form $q$ is closable iff the map $\hat{\imath}: H_{M} \rightarrow H$ is injective. In this case we identify $H_{M}$ with $\hat{\imath}\left(H_{M}\right) \subset H$ and therefore we may view $\hat{q}$ as a quadratic form on $H$. The form $\hat{q}$ is called the closure of $q$ and as the notation suggests is a closed quadratic form on $H$.

A more explicit description of $\hat{q}$ is as follows. The domain $\mathcal{D}(\hat{q})$ consists of those $v \in H$ such that there exists $\left\{v_{n}\right\} \subset X$ such that $v_{n} \rightarrow v$ in $H$ and $q\left(v_{n}-v_{m}\right) \rightarrow 0$ as $m, n \rightarrow \infty$. If $v, w \in \mathcal{D}(\hat{q})$ and $v_{n} \rightarrow v$ and $w_{n} \rightarrow w$ as just described, then $\hat{q}(v, w):=\lim _{n \rightarrow \infty} q\left(v_{n}, w_{n}\right)$.

Proof. 1. Suppose $q$ is closed and $\left\{v_{n}\right\}_{n=1}^{\infty} \subset X$ is a $\|\cdot\|_{M}$ - Cauchy sequence. By the inequality in Eq. (53.6), $\left\{v_{n}\right\}_{n=1}^{\infty}$ is $\|\cdot\|_{H}$ - Cauchy and hence $v:=\lim _{n \rightarrow \infty} v_{n}$ exists in $H$. Moreover,

$$
q\left(v_{n}-v_{m}\right)=\left\|v_{n}-v_{m}\right\|_{M}^{2}-M\left\|v_{n}-v_{m}\right\|_{H}^{2} \rightarrow 0
$$

and therefore, because $q$ is closed, $v \in \mathcal{D}(q)=X$ and $\lim _{n \rightarrow \infty} q\left(v-v_{n}\right)=0$ and hence $\lim _{n \rightarrow \infty}\left\|v_{n}-v\right\|_{M}^{2}=0$. The converse direction is simpler and will be left to the reader.
2. The proof that $q$ is closable iff the map $\hat{\imath}: H_{M} \rightarrow H$ is injective will be complete once the reader verifies that the following assertions are equivalent. 1) $\hat{\imath}: H_{1} \rightarrow H$ is injective, 2) $\hat{\imath}(\bar{v})=0$ implies $\left.\bar{v}=0,3\right)$ if $v_{n} \xrightarrow{H} 0$ and $q\left(v_{n}-v_{m}\right) \rightarrow 0$ as $m, n \rightarrow \infty$ implies that $q\left(v_{n}\right) \rightarrow 0$ as $n \rightarrow \infty$.

By construction $H_{M}$ equipped with the inner product $(\cdot, \cdot)_{M}:=\hat{q}(\cdot, \cdot)+$ $M(\cdot, \cdot)$ is complete. So by item 1. it follows that $\hat{q}$ is a closed quadratic form on $H$ if $q$ is closable.
Example 53.18. Suppose $H=L^{2}([-1,1]), \mathcal{D}(q)=C([-1,1])$ and $q(f, g):=$ $f(0) \bar{g}(0)$ for all $f, g \in \mathcal{D}(q)$. The form $q$ is not closable. Indeed, let $f_{n}(x)=$ $\left(1+x^{2}\right)^{-n}$, then $f_{n} \rightarrow 0 \in L^{2}$ as $n \rightarrow \infty$ and $q\left(f_{n}-f_{m}\right)=0$ for all $m, n$ while $q\left(f_{n}-0\right)=q\left(f_{n}\right)=1 \nrightarrow 0$ as $n \rightarrow \infty$. This example also shows the operator $T: H \rightarrow \mathbb{C}$ defined by $\mathcal{D}(T)=C([-1,1])$ with $T f=f(0)$ is not closable.

Let us also compute $T^{*}$ for this example. By definition $\lambda \in D\left(T^{*}\right)$ and $T^{*} \lambda=f$ iff $(f, g)=\lambda \overline{T g}=\lambda \overline{g(0)}$ for all $g \in C([-1,1])$. In particular this implies $(f, g)=0$ for all $g \in C([-1,1])$ such that $g(0)=0$. However these functions are dense in $H$ and therefore we conclude that $f=0$ and hence $\mathcal{D}\left(T^{*}\right)=\{0\}!!$

Exercise 53.19. Keeping the notation in Example 53.18, show $\overline{\Gamma(T)}=H \times \mathbb{C}$ which is clearly not the graph of a linear operator $S: H \rightarrow \mathbb{C}$.
Proposition 53.20. Suppose that $A: H \rightarrow H$ is a densely defined positive symmetric operator, i.e. $(A v, w)=(v, A w)$ for all $v, w \in \mathcal{D}(A)$ and $(v, A v) \geq$ 0 for all $v \in \mathcal{D}(A)$. Define $q_{A}(v, w):=(v, A w)$ for $v, w \in \mathcal{D}(A)$. Then $q_{A}$ is closable and the closure $\hat{q}_{A}$ is a non-negative, symmetric closed quadratic form on $H$.

Proof. Let $(\cdot, \cdot)_{1}=(\cdot, \cdot)+q_{A}(\cdot, \cdot)$ on $\mathcal{D}(A) \times \mathcal{D}(A), v_{n} \in \mathcal{D}(A)$ such that $H$ - $\lim _{n \rightarrow \infty} v_{n}=0$ and

$$
q_{A}\left(v_{n}-v_{m}\right)=\left(A\left(v_{n}-v_{m}\right),\left(v_{n}-v_{m}\right)\right) \rightarrow 0 \text { as } m, n \rightarrow \infty
$$

Then
$\limsup _{n \rightarrow \infty} q_{A}\left(v_{n}\right) \leq \lim _{n \rightarrow \infty}\left\|v_{n}\right\|_{1}^{2}=\lim _{m, n \rightarrow \infty}\left(v_{m}, v_{n}\right)_{1}=\lim _{m, n \rightarrow \infty}\left\{\left(v_{m}, v_{n}\right)+\left(v_{m}, A v_{n}\right)\right\}=$
where the last equality follows by first letting $m \rightarrow \infty$ and then $n \rightarrow \infty$. Notice that the above limits exist because of Eq. (53.7).

Lemma 53.21. Let $A$ be a positive self-adjoint operator on $H$ and define $q_{A}(v, w):=(v, A w)$ for $v, w \in \mathcal{D}(A)=\mathcal{D}\left(q_{A}\right)$. Then $q_{A}$ is closable and the closure of $q_{A}$ is

$$
\hat{q}_{A}(v, w)=(\sqrt{A} v, \sqrt{A} w) \text { for } v, w \in X:=\mathcal{D}\left(\hat{q}_{A}\right)=\mathcal{D}(\sqrt{A})
$$

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Proof. Let $\hat{q}(v, w)=(\sqrt{A} v, \sqrt{A} w)$ for $v, w \in X=\mathcal{D}(\sqrt{A})$. Since $\sqrt{A}$ is self-adjoint and hence closed, it follows from Example 53.13 that $\hat{q}$ is closed. Moreover, $\hat{q}$ extends $q_{A}$ because if $v, w \in \mathcal{D}(A)$, then $v, w \in \mathcal{D}(A)=\mathcal{D}\left((\sqrt{A})^{2}\right)$ and $\hat{q}(v, w)=(\sqrt{A} v, \sqrt{A} w)=(v, A w)=q_{A}(v, w)$. Thus to show $\hat{q}$ is the closure of $q_{A}$ it suffices to show $\mathcal{D}(A)$ is dense in $X=\mathcal{D}(\sqrt{A})$ when equipped with the Hilbertian norm, $\|w\|_{1}^{2}=\|w\|^{2}+\hat{q}(w)$.

Let $v \in \mathcal{D}(\sqrt{A})$ and define $v_{m}:=1_{[0, m]}(A) v$. Then using the spectral theorem along with the dominated convergence theorem one easily shows that $v_{m} \in X=\mathcal{D}(A), \lim _{m \rightarrow \infty} v_{m}=v$ and $\lim _{m \rightarrow \infty} \sqrt{A} v_{m}=\sqrt{A} v$. But this is equivalent to showing that $\lim _{m \rightarrow \infty}\left\|v-v_{m}\right\|_{1}=0$.

Theorem 53.22. Suppose $q: X \times X \rightarrow \mathbb{C}$ is a symmetric, closed, semibounded (say $q(v, v) \geq-M_{0}\|v\|^{2}$ ) sesquilinear form. Let $L: H \rightarrow H$ be the possibly unbounded operator defined by

$$
D(L):=\{v \in X: q(v, \cdot) \text { is } H-\text { continuous }\}
$$

and for $v \in D(L)$ let $L v \in H$ be the unique element such that $q(v, \cdot)=$ $\left.(L v, \cdot)\right|_{X}$. Then

1. $L$ is a densely defined self-adjoint operator on $H$ and $L \geq-M_{0} I$.
2. $D(L)$ is a form core for $q$, i.e. the closure of $D(L)$ is a dense subspace in $\left(X,\|\cdot\|_{M}\right)$. More explicitly, for all $v \in X$ there exists $v_{n} \in D(L)$ such that $v_{n} \rightarrow v$ in $H$ and $q\left(v-v_{n}\right) \rightarrow 0$ as $n \rightarrow \infty$.
3. For and $M \geq M_{0}, D(q)=D(\sqrt{L+M I})$.
4. Letting $q_{L}(v, w):=(L v, w)$ for all $v, w \in D(L)$, we have $q_{L}$ is closable and $\hat{q}_{L}=q$.

Proof. 1. From Lemma 53.17, $\left(X,(\cdot, \cdot)_{X}:=(\cdot, \cdot)_{M}\right)$ is a Hilbert space for any $M>M_{0}$. Applying Theorem 53.10 and Corollary 53.11 with $q$ being $(\cdot, \cdot)_{X}$ gives a self-adjoint operator $L_{M}: H \rightarrow H$ such that

$$
D\left(L_{M}\right):=\left\{v \in X:(v, \cdot)_{X} \text { is } H-\text { continuous }\right\}
$$

and for $v \in D\left(L_{M}\right)$,

$$
\begin{equation*}
\left(L_{M} v, w\right)_{H}=(v, w)_{X}=q(v, w)+M(v, w) \text { for all } w \in X \tag{53.8}
\end{equation*}
$$

Since $(v, \cdot)_{X}$ is $H$ - continuous iff $q(v, \cdot)$ is $H$ - continuous it follows that $D\left(L_{M}\right)=D(L)$ and moreover Eq. (53.8) is equivalent to

$$
\left(\left(L_{M}-M I\right) v, w\right)_{H}=q(v, w) \text { for all } w \in X
$$

Hence it follows that $L:=L_{M}-M I$ and so $L$ is self-adjoint. Since $(L v, v)=$ $q(v, v) \geq-M_{0}\|v\|^{2}$, we see that $L \geq-M_{0} I$.
2. The density of $\mathcal{D}(L)=\mathcal{D}\left(L_{M}\right)$ in $\left(X,(\cdot, \cdot)_{M}\right)$ is a direct consequence of Theorem 53.10.
3. For

$$
v, w \in \mathcal{D}(Q):=\mathcal{D}\left(\sqrt{L_{M}}\right)=\mathcal{D}(\sqrt{L+M I})=\mathcal{D}\left(\sqrt{L+M_{0} I}\right)
$$

let $Q(v, w):=\left(\sqrt{L_{M}} v, \sqrt{L_{M}} w\right)$. For $v, w \in D(L)$ we have

$$
Q(v, w)=\left(L_{M} v, w\right)=(L v, w)+M(v, w)=q(v, w)+M(v, w)=(v, w)_{M}
$$

By Lemma $53.21, Q$ is a closed, non-negative symmetric form on $H$ and $\mathcal{D}(L)=\mathcal{D}\left(L_{M}\right)$ is dense in $(\mathcal{D}(Q), Q)$. Hence if $v \in \mathcal{D}(Q)$ there exists $v_{n} \in \mathcal{D}(L)$ such that $Q\left(v-v_{n}\right) \rightarrow 0$ and this implies $q\left(v_{m}-v_{n}\right) \rightarrow 0$ as $m, n \rightarrow \infty$. Since $q$ is closed, this implies $v \in \mathcal{D}(q)$ and furthermore that $Q(v, w)=(v, w)_{M}$ for all $v, w \in \mathcal{D}(Q)$.

Conversely, by item 2., if $v \in X=\mathcal{D}(q)$, there exists $v_{n} \in \mathcal{D}(L)$ such that $\left\|v-v_{m}\right\|_{M} \rightarrow 0$. From this it follows that $Q\left(v_{m}-v_{n}\right) \rightarrow 0$ as $m, n \rightarrow \infty$ and therefore since $Q$ is closed, $v \in \mathcal{D}(Q)$ and again $Q(v, w)=(v, w)_{M}$ for all $v, w \in \mathcal{D}(q)$. This proves item 3. and also shows that

$$
q(v, w)=(\sqrt{L+M I} v, \sqrt{L+M I} w)-M(v, w) \text { for all } v, w \in X
$$

where $X:=\mathcal{D}\left(\sqrt{L_{M}}\right)$.
4. Since $q_{L} \subset q, q_{L}$ is closable and the closure of $q_{L}$ is still contained in $q$. Since $q_{L}=Q-L(\cdot, \cdot)$ on $D(L)$ and the closure of $\left.Q\right|_{D(L)}=(\cdot, \cdot)_{M}$, it is easy to conclude that the closure of $q_{L}$ is $q$ as well.
Notation 53.23 Let $\mathcal{P}$ denote the collection of positive self-adjoint operators on $H$ and $\mathcal{Q}$ denote the collection of positive and closed symmetric forms on $H$.
Theorem 53.24. The map $A \in \mathcal{P} \rightarrow \hat{q}_{A} \in \mathcal{Q}$ is bijective, where $\hat{q}_{A}(v, w):=$ $(\sqrt{A} v, \sqrt{A} w)$ with $\mathcal{D}\left(\hat{q}_{A}\right)=\mathcal{D}(\sqrt{A})$ is the closure of the quadratic form $q_{A}(v, w):=(A v, w)$ for $v, w \in \mathcal{D}(q):=\mathcal{D}(A)$. The inverse map is given by $q \in \mathcal{Q} \rightarrow A_{q} \in \mathcal{P}$ where $A_{q}$ is uniquely determined by

$$
\begin{aligned}
\mathcal{D}\left(A_{q}\right) & =\{v \in \mathcal{D}(q): q(v, \cdot) \text { is } H-\text { continuous }\} \text { and } \\
\left(A_{q} v, w\right) & =q(v, w) \text { for } v \in \mathcal{D}\left(A_{q}\right) \text { and } w \in \mathcal{D}(q)
\end{aligned}
$$

Proof. From Lemma 53.21, $\hat{q}_{A} \in \mathcal{Q}$ and $\hat{q}_{A}$ is the closure of $q_{A}$. From Theorem $53.22 A_{q} \in \mathcal{P}$ and

$$
q(\cdot, \cdot)=\left(\sqrt{A_{q}} \cdot, \sqrt{A_{q}} \cdot\right)=\hat{q}_{A_{q}} .
$$

So to finish the proof it suffices to show $A \in \mathcal{P} \rightarrow \hat{q}_{A} \in \mathcal{Q}$ is injective. However, again by Theorem 53.22, if $q \in \mathcal{Q}$ and $A \in \mathcal{P}$ such that $q=\hat{q}_{A}$, then $v \in \mathcal{D}\left(A_{q}\right)$ and $A_{q} v=w$ iff

$$
(\sqrt{A} v, \sqrt{A} \cdot)=q(v, \cdot)=\left.\left(A_{q} v, \cdot\right)\right|_{X}
$$

But this implies $\sqrt{A} v \in \mathcal{D}(\sqrt{A})$ and $A_{q} v=\sqrt{A} \sqrt{A} v=A v$. But by the spectral theorem, $D(\sqrt{A} \sqrt{A})=D(A)$ and so we have proved $A_{q}=A$.

### 53.4 Construction of positive self-adjoint operators

The main theorem concerning closed symmetric semi-bounded quadratic forms $q$ is Friederich's extension theorem.

Corollary 53.25 (The Friederich's extension). Suppose that $A: H \rightarrow H$ is a densely defined positive symmetric operator. Then $A$ has a positive selfadjoint extension $\hat{A}$. Moreover, $\hat{A}$ is the only self-adjoint extension of $A$ such that $\mathcal{D}(\hat{A}) \subset \mathcal{D}\left(\hat{q}_{A}\right)$.

Proof. By Proposition 53.20, $q:=\hat{q}_{A}$ exists in $\mathcal{Q}$. By Theorem 53.24, there exists a unique positive self-adjoint operator $B$ on $H$ such that $\hat{q}_{B}=q$. Since for $v \in \mathcal{D}(A), q(v, w)=(A v, w)$ for all $w \in X$, it follows from Eq. (??) and (??) that $v \in \mathcal{D}(B)$ and $B v=A v$. Therefore $A:=B$ is a self-adjoint extension of $A$.

Suppose that $C$ is another self-adjoint extension of $A$ such that $\mathcal{D}(C) \subset X$. Then $\hat{q}_{C}$ is a closed extension of $q_{A}$. Thus $q=\hat{q}_{A} \subset \hat{q}_{C}$, i.e. $\mathcal{D}\left(\hat{q}_{A}\right) \subset \mathcal{D}\left(\hat{q}_{C}\right)$ and $\hat{q}_{A}=\hat{q}_{C}$ on $\mathcal{D}\left(\hat{q}_{A}\right) \times \mathcal{D}\left(\hat{q}_{A}\right)$. For $v \in \mathcal{D}(C)$ and $w \in \mathcal{D}(A)$, we have that

$$
\hat{q}_{C}(v, w)=(C v, w)=(v, C w)=(v, A w)=(v, B w)=q(v, w) .
$$

By continuity it follows that

$$
\hat{q}_{C}(v, w)=(C v, w)=(v, B w)=q(v, w)
$$

for all $w \in \mathcal{D}(B)$. Therefore, $v \in \mathcal{D}\left(B^{*}\right)=\mathcal{D}(B)$ and $B v=B^{*} v=C v$. That is $C \subset B$. Taking adjoints of this equation shows that $B=B^{*} \subset C^{*}=C$. Thus $C=B$.

Corollary 53.26 (von Neumann). Suppose that $D: H \rightarrow K$ is a closed operator, then $A=D^{*} D$ is a positive self-adjoint operator on $H$.

Proof. The operator $D^{*}$ is densely defined by Lemma 53.6. The quadratic form $q(v, w):=(D v, D w)$ for $v, w \in X:=\mathcal{D}(D)$ is closed (Example 53.13) and positive. Hence by Theorem 53.24 there exists an $A \in \mathcal{P}$ such that $q=\hat{q}_{A}$, i.e.

$$
\begin{equation*}
(D v, D w)=(\sqrt{A} v, \sqrt{A} w) \text { for all } v, w \in X=\mathcal{D}(D)=\mathcal{D}(\sqrt{A}) \tag{53.9}
\end{equation*}
$$

Recalling that $v \in D(A) \subset X$ and $A v=g$ happens iff

$$
(D v, D w)=q(v, w)=(g, w) \text { for all } w \in X
$$

and this happens iff $D v \in D\left(D^{*}\right)$ and $D^{*} D v=g$. Thus we have shown $D^{*} D=A$ which is self-adjoint and positive.

### 53.5 Applications to partial differential equations

Let $U \subset \mathbb{R}^{n}$ be an open set, $\rho \in C^{1}(U \rightarrow(0, \infty))$ and for $i, j=1,2, \ldots, n$ let $a_{i j} \in C^{1}(U, \mathbb{R})$. Take $H=L^{2}(U, \rho d x)$ and define

$$
q(f, g):=\int_{U} \sum_{i, j=1}^{n} a_{i j}(x) \partial_{i} f(x) \partial_{j} g(x) \rho(x) d x
$$

for $f, g \in X=C_{c}^{2}(U)$.
Proposition 53.27. Suppose that $a_{i j}=a_{j i}$ and that $\sum_{i, j=1}^{n} a_{i j}(x) \xi_{i} \xi_{j} \geq 0$ for all $\xi \in \mathbb{R}^{n}$. Then $q$ is a symmetric closable quadratic form on $H$. Hence there exists a unique self-adjoint operator $\hat{A}$ on $H$ such that $\hat{q}=\hat{q}_{\hat{A}}$. Moreover $\hat{A}$ is an extension of the operator

$$
A f(x)=-\frac{1}{\rho(x)} \sum_{i, j=1}^{n} \partial_{j}\left(\rho(x) a_{i j}(x) \partial_{i} f(x)\right)
$$

for $f \in \mathcal{D}(A)=C_{c}^{2}(U)$.
Proof. A simple integration by parts argument shows that $q(f, g)=$ $(A f, g)_{H}=(f, A g)_{H}$ for all $f, g \in \mathcal{D}(A)=C_{c}^{2}(U)$. Thus by Proposition 53.20, $q$ is closable. The existence of $\hat{A}$ is a result of Theorem 53.24. In fact $\hat{A}$ is the Friederich's extension of $A$ as in Corollary 53.25.

Given the above proposition and the spectral theorem, we now know that (at least in some weak sense) we may solve the general heat and wave equations: $u_{t}=-A u$ for $t \geq 0$ and $u_{t t}=-A u$ for $t \in \mathbb{R}$. Namely, we will take

$$
u(t, \cdot):=e^{-t \hat{A}} u(0, \cdot)
$$

and

$$
u(t, \cdot)=\cos (t \sqrt{\hat{A}}) u(0, \cdot)+\frac{\sin (t \sqrt{\hat{A}})}{\sqrt{\hat{A}}} u_{t}(0, \cdot)
$$

respectively. In order to get classical solutions to the equations we would have to better understand the operator $\hat{A}$ and in particular its domain and the domains of the powers of $\hat{A}$. This will be one of the topics of the next part of the course dealing with Sobolev spaces.
Remark 53.28. By choosing $\mathcal{D}(A)=C_{c}^{2}(U)$ we are essentially using Dirichlet boundary conditions for $A$ and $\hat{A}$. If $U$ is a bounded region with $C^{2}$-boundary, we could have chosen (for example VERIFY THIS EXAMPLE)

$$
\mathcal{D}(A)=\left\{f \in C^{2}(U) \cap C^{1}(\bar{U}): \text { with } \partial u / \partial n=0 \text { on } \partial U\right\} .
$$

This would correspond to Neumann boundary conditions. Proposition 53.27 would be valid with this domain as well provided we assume that $a_{i, j}$ and $\rho$ are in $C^{1}(\bar{U})$.

For a second application let $H=L^{2}\left(U, \rho d x ; \mathbb{R}^{N}\right)$ and for $j=1,2, \ldots, n$, let $A_{j}: U \rightarrow \mathcal{M}_{N \times N}$ (the $N \times N$ matrices) be a $C^{1}$ function. Set $\mathcal{D}(D):=$ $C_{c}^{1}\left(U \rightarrow \mathbb{R}^{N}\right)$ and for $S \in \mathcal{D}(D)$ let $D S(x)=\sum_{i=1}^{n} A_{i}(x) \partial_{i} S(x)$.

Proposition 53.29 ("Dirac Like Operators"). The operator D on $H$ defined above is closable. Hence $A:=D^{*} \bar{D}$ is a self-adjoint operator on $H$, where $\bar{D}$ is the closure of $D$.

Proof. Again a simple integration by parts argument shows that $\mathcal{D}(D) \subset$ $\mathcal{D}\left(D^{*}\right)$ and that for $S \in \mathcal{D}(D)$,

$$
D^{*} S(x)=\frac{1}{\rho(x)} \sum_{i=1}^{n}-\partial_{i}\left(\rho(x) A_{i}(x) S(x)\right) .
$$

In particular $D^{*}$ is a densely defined operator and hence $D$ is closable. The result now follows from Corollary 53.26.

## $L^{2}-$ operators associated to $\mathcal{E}$

Let $\bar{\Omega}$ be a $C^{2}$ - manifold with boundary such that $\Omega=\bar{\Omega}^{o}, \rho \in C^{\infty}(\bar{\Omega})$ with $\rho>0$ on $\bar{\Omega}, a_{\alpha, \beta} \in B C^{\infty}(\bar{\Omega})$ such that $\left(a_{i j}\right) \geq \epsilon I$ for some $\epsilon>0$ and $\mathcal{E}$ be the elliptic Dirichlet form given by

$$
\mathcal{E}(u, v)=\int_{\Omega} \sum a_{\alpha \beta} \partial^{\alpha} u \partial^{\beta} v d \mu
$$

where $d \mu:=\rho d m$. Let $H=L^{2}(\mu) \cong L^{2}(m)$ and $X=H^{1}(\Omega)$ or $H_{0}^{1}(\Omega)$.
Definition 54.1. Let

$$
\begin{aligned}
& D\left(L_{\mathcal{E}}\right)=\left\{u \in X: \mathcal{L}_{\mathcal{E}} u:=\mathcal{E}(u, \cdot) \in L^{2}(\rho d x)\right\} \\
& D\left(L_{\mathcal{E}}^{\dagger}\right)=\left\{u \in X: \mathcal{L}_{\mathcal{E}}^{\dagger} u:=\mathcal{E}(\cdot, u) \in L^{2}(\rho d x)\right\}
\end{aligned}
$$

and for $u \in D\left(L_{\mathcal{E}}\right)\left(u \in D\left(L_{\mathcal{E}}^{\dagger}\right)\right)$ define $L_{\mathcal{E}} u\left(L_{\mathcal{E}}^{\dagger} u\right)$ to be the unique element in $L^{2}(\mu)$ such that

$$
\begin{aligned}
& \mathcal{E}(u, v)=\left(L_{\mathcal{E}} u, v\right)_{L^{2}(\mu)}=\int_{\Omega} L_{\mathcal{E}} u \cdot v \rho d m \\
& \mathcal{E}(v, u)=\left(v, L_{\mathcal{E}} u\right)_{L^{2}(\mu)}=\int_{\Omega} v \cdot L_{\mathcal{E}}^{\dagger} u \rho d m
\end{aligned}
$$

## for all $v \in X$.

Theorem 54.2. If $X=H_{0}^{1}(\Omega)$ let $B=B^{\dagger} \equiv 0$ and if $X=H^{1}(\Omega)$ let $B$ and $B^{\dagger}$ be given (as in Proposition 51.4) by

$$
\begin{aligned}
& B u:=\sum a_{\alpha j} \partial^{\alpha} u n_{j}=n \cdot a \nabla u+\left(n \cdot a_{0, \cdot}\right) u \\
& B^{\dagger} u:=\sum a_{i \beta} \partial^{\beta} u n_{i}=a n \cdot \nabla u+(n \cdot a \cdot a \cdot, 0) u .
\end{aligned}
$$

Then

$$
\begin{aligned}
& D\left(L_{\mathcal{E}}\right)=\left\{u \in H^{2}(\Omega) \cap X:\left.B u\right|_{\partial \Omega}=0\right\} \\
& D\left(L_{\mathcal{E}}^{\dagger}\right)=\left\{u \in H^{2}(\Omega) \cap X:\left.B^{\dagger} u\right|_{\partial \Omega}=0\right\}
\end{aligned}
$$

and

$$
\begin{aligned}
L_{\mathcal{E}} u & =\frac{1}{\rho} \sum\left((-\partial)^{\beta} \rho a_{\alpha \beta} \partial^{\alpha} u\right)=: L u \\
L_{\mathcal{E}}^{\dagger} u & \left.=\frac{1}{\rho} \sum(-\partial)^{\alpha} \rho a_{\alpha \beta} \partial^{\beta} u\right)=: L^{\dagger} u
\end{aligned}
$$

Moreover $L_{\mathcal{E}}=\left(L_{\mathcal{E}}^{\dagger}\right)^{*}$ and $L_{\mathcal{E}}^{*}=L_{\mathcal{E}}^{\dagger}$.
Proof. By replacing $\mathcal{E}(u, v)$ by $\mathcal{E}(u, v)+C(u, v), L_{\mathcal{E}} \rightarrow L_{\mathcal{E}}+C$ and $L_{\mathcal{E}}^{\dagger} \rightarrow$ $L_{\mathcal{E}}^{\dagger}+C$ for a sufficiently large constant $C$, we may assume that $\mathcal{E}(u, v)$ satisfies

$$
\epsilon\|u\|_{H^{1}}^{2} \leq \mathcal{E}(u, u) \text { for all } u \in H^{1}
$$

Then by Theorem 53.10, $L_{\mathcal{E}}^{*}=L_{\mathcal{E}}^{\dagger},\left(L_{\mathcal{E}}^{\dagger}\right)^{*}=L_{\mathcal{E}}$ and

$$
L_{\mathcal{E}}: D\left(L_{\mathcal{E}}\right) \rightarrow L^{2}(\Omega) \text { and } L_{\mathcal{E}}^{\dagger}: D\left(L_{\mathcal{E}}^{\dagger}\right) \rightarrow L^{2}(\Omega)
$$

are linear isomorphisms. By the elliptic regularity Theorem 52.15 , both $D\left(L_{\mathcal{E}}\right)$ and $D\left(L_{\mathcal{E}}^{\dagger}\right)$ are subspaces of $H^{2}(\Omega)$ and moreover there is a constant $C<\infty$ such that

$$
\begin{equation*}
\|u\|_{H^{2}(\Omega)} \leq C\left\|L_{\mathcal{E}} u\right\|_{L^{2}(\Omega)} \tag{54.1}
\end{equation*}
$$

From Proposition 51.4 (integration by parts), for $u \in H^{2}(\Omega)$ and $v \in X$,

$$
\begin{equation*}
\mathcal{E}(u, v)=(L u, v)_{L^{2}(\mu)}+\left.\left.\int_{\partial \Omega} B u\right|_{\partial \Omega} \cdot v\right|_{\partial \Omega} \rho d \sigma \tag{54.2}
\end{equation*}
$$

while, by definition, if $u \in D\left(L_{\mathcal{E}}\right)$ then

$$
\begin{equation*}
\mathcal{E}(u, v)=\left(L_{\mathcal{E}} u, v\right)_{L^{2}(\mu)} \text { for all } v \in X \tag{54.3}
\end{equation*}
$$

Choosing $v \in H_{0}^{1}(\Omega) \subset X$, comparing Eqs. (54.2) and (54.3) shows that $L u=$ $L_{\mathcal{E}} u$. So for $u \in D\left(L_{\mathcal{E}}\right), L_{\mathcal{E}} u=L u$ and moreover we must have $\left.B u\right|_{\partial \Omega}=0$ as well. Therefore

$$
D\left(L_{\mathcal{E}}\right) \subset H^{2}(\Omega) \cap\left\{u:\left.B u\right|_{\partial \Omega}=0\right\}
$$

Conversely if $u \in H^{2}(\Omega)$ with $\left.B u\right|_{\partial \Omega}=0, \mathcal{E}(u, v)=(L u, v)_{L^{2}(\mu)}$ for all $v \in X$ and therefore by definition of $\mathcal{L}_{\mathcal{E}}, u \in D\left(L_{\mathcal{E}}\right)$ and $L_{\mathcal{E}} u=L u$. The assertions involving $L_{\mathcal{E}}^{\dagger}$ are proved in the same way.

### 54.1 Compact perturbations of the identity and the Fredholm Alternative

Definition 54.3. A bounded operator $F: H \rightarrow B$ is Fredholm iff the $\operatorname{dim} \operatorname{Nul}(F)<\infty$, dim coker $(F)<\infty$ and $\operatorname{Ran}(F)$ is closed in $B$. (Recall: $\operatorname{coker}(F):=B / \operatorname{Ran}(F)$.) The index of $F$ is the integer,

$$
\begin{align*}
\operatorname{index}(F) & =\operatorname{dim} \operatorname{Nul}(F)-\operatorname{dim} \operatorname{coker}(F)  \tag{54.4}\\
& =\operatorname{dim} \operatorname{Nul}(F)-\operatorname{dim} \operatorname{Nul}\left(F^{*}\right) \tag{54.5}
\end{align*}
$$

## $109454 L^{2}$ - operators associated to $\mathcal{E}$

Example 54.4. Suppose that $H$ and $B$ are finite dimensional Hilbert spaces and $F: H \rightarrow B$ is a linear operator. In this case, the rank nullity theorem implies

$$
\begin{aligned}
\operatorname{index}(F) & =\operatorname{dim} \operatorname{Nul}(F)-\operatorname{dim} \operatorname{coker}(F) \\
& =\operatorname{dim} \operatorname{Nul}(F)-[\operatorname{dim} B-\operatorname{dim} \operatorname{Ran}(F)] \\
& =\operatorname{dim} \operatorname{Nul}(F)+\operatorname{dim} \operatorname{Ran}(F)-\operatorname{dim} B \\
& =\operatorname{dim} H-\operatorname{dim} B .
\end{aligned}
$$

Theorem 54.5. If $R: H \rightarrow H$ is finite rank, then $F=I+R$ is Fredholm and index $(F)=0$.

Proof. Let $H_{1}=\operatorname{Nul}(R), H_{2}=\operatorname{Ran}(R), P_{i}: H \rightarrow H_{i}$ be orthogonal projection, $\left\{\psi_{i}\right\}_{i=1}^{n}$ be an orthonormal basis for $\operatorname{Ran}(R)$ and $\phi_{i}:=R^{*} \psi_{i}$ for $i=1,2, \ldots, n$. Then for $h \in H$,

$$
R h=\sum_{i=1}^{n}\left(R h, \psi_{i}\right) \psi_{i}=\sum_{i=1}^{n}\left(h, R^{*} \psi_{i}\right) \psi_{i}=\sum_{i=1}^{n}\left(h, \phi_{i}\right) \psi_{i}
$$

and hence $\left\{\phi_{1}, \ldots, \phi_{n}\right\}^{\perp} \subset \operatorname{Nul}(R)$. Therefore $H_{2}=\operatorname{Nul}(R)^{\perp} \subset \operatorname{span}\left\{\phi_{1}, \ldots, \phi_{n}\right\}$ is finite dimensional. For $h=h_{1}+h_{2} \in H_{1} \oplus H_{2}$,

$$
\begin{align*}
F h & =\left(P_{1}+P_{2}\right)\left(h_{1}+h_{2}+R h_{2}\right)=\left(h_{1}+P_{1} R h_{2}\right)+\left(h_{2}+P_{2} R P_{2} h_{2}\right) \\
& =\left(h_{1}+P_{1} R h_{2}\right)+\left(I_{H_{2}}+P_{2} R P_{2}\right) h_{2} . \tag{54.6}
\end{align*}
$$

From Eq. (54.6) we see that $h=h_{1}+h_{2} \in \operatorname{Nul}(F)$ iff $h_{2} \in \operatorname{Nul}\left(I_{H_{2}}+P_{2} R P_{2}\right)$ and $h_{1}=-P_{1} R h_{2}$ and hence

$$
\begin{equation*}
\operatorname{Nul}(F) \cong \operatorname{Nul}\left(I_{H_{2}}+P_{2} R P_{2}\right) \tag{54.7}
\end{equation*}
$$

It is also easily seen from Eq. (54.6) that

$$
\begin{equation*}
\operatorname{Ran}(F)=H_{1} \oplus \operatorname{Ran}\left(I_{H_{2}}+P_{2} R P_{2}\right) \tag{54.8}
\end{equation*}
$$

Since $H_{2}$ is finite dimensional, $\operatorname{Ran}\left(I_{H_{2}}+P_{2} R P_{2}\right)$ is a closed subspace of $\mathrm{H}_{2}$ and so $\operatorname{Ran}(F)$ is closed. Moreover

$$
\begin{align*}
\operatorname{coker}(F) & =H / \operatorname{Ran}(F)=\left[H_{1} \oplus H_{2}\right] /\left[H_{1} \oplus \operatorname{Ran}\left(I_{H_{2}}+P_{2} R P_{2}\right)\right] \\
& \cong H_{2} / \operatorname{Ran}\left(I_{H_{2}}+P_{2} R P_{2}\right)=\operatorname{coker}\left(I_{H_{2}}+P_{2} R P_{2}\right) \tag{54.9}
\end{align*}
$$

So by Eqs. (54.7), (54.9) and Example 54.4,

$$
\begin{aligned}
\operatorname{index}(F) & =\operatorname{dim} \operatorname{Nul}(F)-\operatorname{dim} \operatorname{coker}(F) \\
& =\operatorname{dim} \operatorname{Nul}\left(I_{H_{2}}+P_{2} R P_{2}\right)-\operatorname{dim} \operatorname{coker}\left(I_{H_{2}}+P_{2} R P_{2}\right) \\
& =\operatorname{index}\left(I_{H_{2}}+P_{2} R P_{2}\right)=0
\end{aligned}
$$

Corollary 54.6. If $K: H \rightarrow H$ compact then $F=I+K$ is Fredholm and index $(F)=0$.

Proof. Choose $R_{1}: H \rightarrow H$ finite rank such that $\varepsilon:=K-R_{1}$ is a bounded operator with operator norm less than one. Then

$$
F=I+K=I+\varepsilon+R_{1}=(I+\varepsilon)\left(I+(I+\varepsilon)^{-1} R_{1}\right)=U(I+R)
$$

where $U:=(I+\varepsilon): H \rightarrow H$ is invertible and $R:=(I+\varepsilon)^{-1} R_{1}: H \rightarrow H$ is finite rank. Therefore, $\operatorname{Ran}(F)=U(\operatorname{Ran}(I+R))$ is closed,

$$
\begin{equation*}
\operatorname{dim} \operatorname{coker}(F)=\operatorname{dim} \operatorname{coker}(I+R)<\infty \tag{54.10}
\end{equation*}
$$

$\operatorname{Nul}(F)=\operatorname{Nul}(I+R)$, and

$$
\begin{equation*}
\operatorname{dim} \operatorname{Nul}(F)=\operatorname{dim} \operatorname{Nul}(I+R) \tag{54.11}
\end{equation*}
$$

From this it follows that $F$ is Fredholm and index $(F)=\operatorname{index}(I+R)=0$.

### 54.2 Solvability of $L u=f$ and properties of the solution

Theorem 54.7. Let $\bar{\Omega} \subset \mathbb{R}^{d}$ be a $C^{\infty}$ - manifold with boundary such that $\Omega=\bar{\Omega}^{o}$. Let $\mathcal{E}$ be an elliptic Dirichlet form, $L:=L_{\mathcal{E}}$ be the associated operator.

1. For $C>0$ sufficiently large, $(L+C): D(L) \rightarrow L^{2}(\Omega)$ is a linear isomorphism and

$$
(L+C)^{-1}: L^{2}(\Omega) \rightarrow D(L) \subset H^{2}(\Omega)
$$

is a bounded operator and $D(L)$ is a closed subspace of $H^{2}(\Omega)$.
2. $(L+C)^{-1}$ as viewed as an operator from $L^{2}(\Omega)$ to $L^{2}(\Omega)$ is compact.
3. If $u \in D(L)$ and $L u \in H^{k}(\Omega)$ then $u \in H^{k+2}(\Omega)$.
4. If $u \in D(L)$ and $L u \in C^{\infty}(\bar{\Omega})$ then $u \in C^{\infty}(\bar{\Omega}):=\bigcap_{k=0}^{\infty} C^{k}(\bar{\Omega})$.
5. If $u \in D(L)$ is an eigenfunction of $L$, i.e. $L u=\lambda u$ for some $\lambda \in \mathbb{C}$, then $u \in C^{\infty}(\bar{\Omega})$.

Proof. 1. It was already shown in the proof of Theorem 54.2 that $(L+C)$ : $D(L) \rightarrow L^{2}(\Omega)$ is bijective. Moreover the bound in Eq. (54.1) shows that $(L+C)^{-1}: L^{2}(\Omega) \rightarrow H^{2}(\Omega)$ is bounded. If $\left\{u_{n}\right\}_{n=1}^{\infty} \subset D(L) \subset H^{2}(\Omega)$ is a sequence such that $u_{n} \rightarrow u \in H^{2}(\Omega)$, then $\left\{(L+C) u_{n}\right\}_{n=1}^{\infty}$ is convergent in $L^{2}(\Omega)$ since $(L+C): H^{2}(\Omega) \rightarrow L^{2}(\Omega)$ is bounded. Because $L$ is a closed operator, it follows that $u \in D(L)$ and so $D(L)$ is a closed subspace of $H^{2}(\Omega)$.
2. This follows from item 1. and the Rellich - Kondrachov Compactness Theorem 49.25 which implies the embedding $H^{2}(\Omega) \hookrightarrow L^{2}(\Omega)$ is compact.
3. If $f \in H^{k}(\Omega)$ and $u \in D(L)$ such that $L u=f \in H^{k}(\Omega)$ then $\mathcal{L}_{\mathcal{E}} u=f$ and hence the elliptic regularity Theorem 53.10 gives the result.
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$$
\begin{aligned}
\operatorname{index}(L) & =\operatorname{dim} \operatorname{coker}(L)-\operatorname{dim} \operatorname{Nul}(L) \\
& =\operatorname{dim} \operatorname{coker}(I-K)-\operatorname{dim} \operatorname{Nul}(I-K) \\
& =\operatorname{index}(I-K)=0 .
\end{aligned}
$$

The same argument applies to $L^{*}$ to show $L^{*}$ is Fredholm and index $\left(L^{*}\right)=0$. Because $\operatorname{Ran}(L)$ is closed and $\operatorname{Ran}(L)^{\perp}=\operatorname{Nul}\left(L^{*}\right), \operatorname{Ran}(L)=\operatorname{Nul}\left(L^{*}\right)^{\perp}$ and

$$
L^{2}(\Omega)=\operatorname{Nul}\left(L^{*}\right)^{\perp} \oplus \operatorname{Nul}\left(L^{*}\right)=\operatorname{Ran}(L) \oplus \operatorname{Nul}\left(L^{*}\right)
$$

Thus $\operatorname{dim} \operatorname{coker}(L)=\operatorname{dim} \operatorname{Nul}\left(L^{*}\right)$ and so

$$
\begin{aligned}
0 & =\operatorname{index}(L)=\operatorname{dim} \operatorname{Nul}(L)-\operatorname{dim} \operatorname{coker}(L) \\
& =\operatorname{dim} \operatorname{Nul}(L)-\operatorname{dim} \operatorname{Nul}\left(L^{*}\right)
\end{aligned}
$$

This proves items 1-4 and finishes the proof of the theorem since item 5 . is a direct consequence of items 3 and 4 .

Example 54.9 (Dirichlet Boundary Conditions). Let $\Delta$ denote the Laplacian with Dirichlet boundary conditions, i.e. $D(\Delta)=H_{0}^{1}(\Omega) \cap H^{2}(\Omega)$. If $u \in D(\Delta)$ then

$$
\begin{equation*}
\int_{\Omega} \nabla u \cdot \nabla v d m=(-\Delta u, v) \text { for all } v \in H_{0}^{1}(\Omega) \tag{54.12}
\end{equation*}
$$

and in particular for $u \in \operatorname{Nul}(\Delta)$ we have

$$
\int_{\Omega}|\nabla u|^{2} d m=(-\Delta u, u)=0 .
$$

By the Poincaré inequality in Theorem 49.31 (or by more direct means) this implies $u=0$ and therefore $\operatorname{Nul}(\Delta)=\{0\}$. It now follows by the Fredholm alternative in Theorem 54.8 that there exists a unique solution $u \in D(\Delta)$ to $\Delta u=f$ for any $f \in L^{2}(\Omega)$.

Example 54.10 (Neuwmann Boundary Conditions). Suppose $\Delta$ is the Laplacian on $\Omega$ with Neuwmann boundary conditions, i.e.

$$
D(\Delta)=\left\{u \in H^{2}(\Omega): \frac{\partial u}{\partial n}=0\right\}
$$

If $u \in D(\Delta)$ then

$$
\begin{equation*}
\int_{\Omega} \nabla u \cdot \nabla v d m=(-\Delta u, v) \text { for all } v \in H^{1}(\Omega) \tag{54.13}
\end{equation*}
$$

so that the Dirichlet form associated to $\Delta$ is symmetric and hence $\Delta=\Delta^{*}$. Moreover if $u \in \operatorname{Nul}(\Delta)$, then

$$
0=(-\Delta u, v)=\int_{\Omega}|\nabla u|^{2} d m
$$

i.e. $\nabla u=0$. As in the proof of the Poincaré Lemma 49.30 (or using the Poincaré Lemma itself), $u$ is constant on each connected component of $\Omega$. Assuming, for simplicity, that $\Omega$ is connected, we have shown

$$
\operatorname{Nul}\left(\Delta_{N}\right)=\operatorname{span}\{1\} .
$$

The Fredholm alternative in Theorem 54.8 implies that there exists a (non unique solution) $u \in D(\Delta)$ to $\Delta u=f$ for precisely those $f \in L^{2}(\Omega)$ such that $f \perp$ 1, i.e. $\int_{\Omega} f 1 d m=0$.

Remark 54.11. Suppose $\mathcal{E}$ is an elliptic Dirichlet form and $L=L_{\mathcal{E}}$ is the associated operator on $L^{2}(\Omega)$. If $\mathcal{E}$ has the property that the only solution to $\mathcal{E}(u, u)=0$ is $u=0$, then the equation $L u=f$ always has a unique solution for any $f \in L^{2}(\Omega)$.

Example 54.12. Let $A_{i j}=A_{j i}, A_{i}$ and $A_{0}$ be in $C^{\infty}(\bar{\Omega})$ with $A_{0} \geq 0$ and $\left(A_{i j}\right) \geq \epsilon I$ for some $\epsilon>0$. For $u, v \in H_{0}^{1}(\Omega)$ let

$$
\begin{equation*}
\mathcal{E}(u, v)=\int_{\Omega}\left(\sum A_{i j} \partial_{i} u \partial_{j} v+A_{0} u v\right) d m \tag{54.14}
\end{equation*}
$$

and $L=L_{\mathcal{E}}$, then

$$
\begin{equation*}
L=-\partial_{j} A_{i j} \partial_{i} u+A_{0} u \tag{54.15}
\end{equation*}
$$

with $D(L):=H^{2}(\Omega) \cap H_{0}^{1}(\Omega)$. If $u \in \operatorname{Nul}(L)$, then $0=(L u, u)=\mathcal{E}(u, u)=0$ implies $\partial_{i} u=0$ a.e. and hence $u$ is constant on each connected component of $\Omega$. Since $u \in H_{0}^{1}(\Omega),\left.u\right|_{\partial \Omega}=0$ from which we learn that $u \equiv 0$. Therefore $L u=f$ has a (unique) solution for all $f \in L^{2}(\Omega)$.

Example 54.13. Keeping the same notation as Example 54.12, except now we view $\mathcal{E}$ as a Dirichlet form on $H^{1}(\Omega)$. Now $L=L_{\mathcal{E}}$ is the operator given in Eq. (54.15) but now

$$
D(L)=\left\{u \in H^{2}(\Omega): B u=0 \text { on } \partial \Omega\right\}
$$

where $B u=n_{j} A_{i j} \partial_{i} u$. Again if $u \in \operatorname{Nul}(L)$ it follows that $u$ is constant on each connected component of $\Omega$. If we further assume that $A_{0}>0$ at some point in each connected component of $\Omega$, we could then conclude from

$$
0=\mathcal{E}(u, u)=\int_{\Omega} A_{0} u^{2} d m
$$

that $u=0$. So again $\operatorname{Nul}(L)=\{0\}$ and $\operatorname{Ran}(L)=L^{2}(\Omega)$.

### 54.3 Interior Regularity Revisited

Theorem 54.14 (Jazzed up interior regularity). Let $L$ be a second order elliptic differential operator on $\Omega$. If $u \in L_{\text {Loc }}^{2}(\Omega)$ such that $L u \in H_{L o c}^{k}(\Omega)$ then $u \in H_{L o c}^{k+2}(\Omega)$ and for any open precompact open sets $\Omega_{0}$ and $\Omega_{1}$ contained in $\Omega$ such that $\bar{\Omega}_{0} \subset \Omega_{1} \subset \bar{\Omega}_{1} \subset \Omega$ there is a constant $C<\infty$ independent of $u$ such that

$$
\|u\|_{H^{k+2}\left(\Omega_{0}\right)} \leq C\left(\|L u\|_{H^{k}\left(\Omega_{1}\right)}+\|u\|_{L^{2}\left(\Omega_{1}\right)}\right) .
$$

Proof. When $k>0$ the theorem follows from Theorem 52.5. So it suffices to consider the case, $k=0$, i.e. $u \in L_{\text {Loc }}^{2}(\Omega)$ such that $L u \in L_{L o c}^{2}(\Omega)$. To finish the proof, again because of Theorem 52.5 , it suffices to show $u \in H_{l o c}^{1}(\Omega)$. By replacing $\Omega$ by a precompact open subset of $\Omega$ which contains $\Omega_{1}$ we may further assume that $u \in L^{2}(\Omega)$ and $L u \in L^{2}(\Omega)$. Further, by replacing $L$ by $L+C$ for some constant $C>0$, the Lax-Milgram method implies we may assume $L: H_{0}^{1}(\Omega) \rightarrow H^{-1}(\Omega)$ is an isomorphism of Banach spaces. We will now finish the proof by showing $u \in H_{l o c}^{1}(\Omega)$ under the above assumptions.

If $\chi \in C_{c}^{\infty}(\Omega)$, then by Lemma 50.7 [ $L, M_{\chi}$ ] is a first order operator so,

$$
L(\chi u)=\chi L u+\left[L, M_{\chi}\right] u=: f_{\chi} \in L^{2}(\Omega)+H^{-1}(\Omega)=H^{-1}(\Omega)
$$

Let $u_{0}=L^{-1} f_{\chi} \in H_{0}^{1}(\Omega), \psi \in C_{c}^{\infty}(\Omega)$ such that $\psi=1$ on a neighborhood of $\operatorname{supp}(\chi)$ and

$$
v:=\psi\left(\chi u-u_{0}\right)=\chi u-\psi u_{0} \in L_{c}^{2}(\Omega)
$$

Then, because $\operatorname{supp}\left(f_{\chi}\right) \subset \operatorname{supp}(\chi)$, we have $f_{\chi}=\psi f_{\chi}$ and

$$
\begin{aligned}
L v & =L\left[\chi u-\psi u_{0}\right]=f_{\chi}-\psi L u_{0}-\left[L, M_{\psi}\right] u_{0}=f_{\chi}-\psi f_{\chi}-\left[L, M_{\psi}\right] u_{0} \\
& =-\left[L, M_{\psi}\right] u_{0}=: g \in L_{c}^{2}(\Omega)
\end{aligned}
$$

Let $D\left(L_{D}\right):=H^{2}(\Omega) \cap H_{0}^{1}(\Omega)$ and $L_{D} u=L u$ for all $u \in D\left(L_{D}\right)$ so that $L_{D}$ is $L$ with Dirichlet boundary conditions on $\Omega$. I now claim that $v \in D\left(L_{D}\right) \subset H_{0}^{1}(\Omega)$. To prove this suppose

$$
\xi \in D\left(L_{D}^{\dagger}\right)=\left\{\xi \in H_{0}^{1}(\Omega): L_{D}^{\dagger} \xi \in L^{2}(\Omega)\right\}=H_{0}^{1}(\Omega) \cap H^{2}(\Omega)
$$

and let $\xi_{m}:=\eta_{m} * \xi$ where $\eta_{m}$ is an approximate $\delta$ - sequence so that $\xi_{m} \rightarrow$ $\xi$ in $H_{l o c}^{2}(\Omega)$. Choose $\phi \in C_{c}^{\infty}(\Omega)$ such that $\phi=1$ on a neighborhood of $\operatorname{supp}(v) \supset \operatorname{supp}(g)$, then

$$
\begin{aligned}
(g, \xi) & =\lim _{m \rightarrow \infty}\left(g, \xi_{m}\right)=\lim _{m \rightarrow \infty}\left(L v, \phi \xi_{m}\right)=\lim _{m \rightarrow \infty}\left(v, L^{\dagger}\left(\phi \xi_{m}\right)\right) \\
& =\lim _{m \rightarrow \infty}\left(v, L^{\dagger} \xi_{m}\right)=\left(v, L^{\dagger} \xi\right)=\left(v, L_{D}^{\dagger} \xi\right)
\end{aligned}
$$

Since this holds for all $\xi \in D\left(L_{D}^{\dagger}\right)$ we see that
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$$
v \in D\left(\left(L_{D}^{\dagger}\right)^{*}\right)=D\left(L_{D}\right)=H_{0}^{1}(\Omega) \cap H^{2}(\Omega) \subset H^{1}(\Omega)
$$

where the first equality is a consequence of Theorem 54.2 which states $L_{D}=$ $\left(L_{D}^{\dagger}\right)^{*}$. Therefore, $\chi u=\psi u_{0}+v \in H^{1}(\Omega)$ and since $\chi \in C_{c}^{\infty}(\Omega)$ was arbitrary we learn that $u \in H_{l o c}^{1}(\Omega)$.

### 54.4 Classical Dirichlet Problem

Let $\bar{\Omega}$ be a $C^{\infty}$ - manifold with boundary such that $\Omega=\bar{\Omega}^{o}$ and let $L=\Delta$ with Dirichlet boundary conditions, so $D(\Delta):=H^{1}(\Omega) \cap H^{2}(\Omega)$.

Theorem 54.15. To each $f \in C(\partial \Omega)$, there exists a unique solution $u \in$ $C^{\infty}(\Omega) \cap C(\bar{\Omega})$ to the equation

$$
\Delta u=0 \text { with } u=f \text { on } \partial \Omega
$$

Proof. Choose $f_{n} \in C^{\infty}(\bar{\Omega})$ such that $\lim _{n \rightarrow \infty}\left\|\left.f_{n}\right|_{\partial \Omega}-f\right\|_{L^{\infty}(\partial \Omega)}=0$. We will now show that there exists $u_{n} \in C^{\infty}(\bar{\Omega})$ such that

$$
\begin{equation*}
\Delta u_{n}=0 \text { with } u_{n}=f_{n} \text { on } \partial \Omega \tag{54.16}
\end{equation*}
$$

To prove this let us write the desired solution as $u_{n}=v_{n}+f_{n}$ in which case $v_{n}=0$ on $\partial \Omega$ and $0=\Delta u_{n}=\Delta v_{n}+\Delta f_{n}$. Hence if $v_{n}$ solves $\Delta v_{n}=-\Delta f_{n}$ on $\Omega$ with $v_{n}=0$ on $\partial \Omega$ then $u_{n}=v_{n}+f_{n}$ solves the Dirichlet problem in Eq. (54.16).

By the maximum principle,

$$
\left\|u_{n}-u_{m}\right\|_{L^{\infty}(\bar{\Omega})} \leq\left\|f_{n}-f_{m}\right\|_{L^{\infty}(\partial \Omega)} \rightarrow 0 \text { as } m, n \rightarrow \infty
$$

and so $\left\{u_{n}\right\}_{n=1}^{\infty} \subset C(\bar{\Omega})$ is uniformly convergent sequence. Let $u:=$ $\lim _{n \rightarrow \infty} u_{n} \in C(\bar{\Omega})$. The proof will be completed by showing $u \in C^{\infty}(\Omega)$ and $\Delta u=0$. This can be done in one stroke by showing $u$ satisfies the mean value property. This it the case since each function $u_{n}$ satisfies the mean value property and this property is preserved under uniform limits.

Remark 54.16. Theorem 54.15 is more generally valid in the case $\Delta$ is replaced by an elliptic operator of the form $L=-\sum_{i j} \frac{1}{\rho} \partial^{i}\left(\rho a_{i j} \partial^{j}\right)$ with $\rho \in C^{\infty}(\bar{\Omega},(0, \infty))$ and $a_{i j} \in C^{\infty}(\bar{\Omega})$ such that $\sum_{i, j=1}^{n} a_{i j}(x) \xi_{i} \xi_{j} \geq \epsilon|\xi|^{2}$ for all $x \in \bar{\Omega}, \xi \in \mathbb{R}^{d}$. Then again for all $f \in C(\partial \Omega)$ there exist a solution $u \in C^{\infty}(\Omega) \cap C(\bar{\Omega})$ such that

$$
L u=0 \text { with } u=f \text { on } \partial \Omega .
$$

The proof is the same as that of Theorem 54.15 except the last step needs to be changed as follows. As above, we construct solution to $L u_{n}=0$ with
$u_{n}=f_{n}$ on $\partial \Omega$ and we then still have $u_{n} \rightarrow u \in C(\bar{\Omega})$ via the maximum principle. To finish the proof, because of Theorem 54.14 , it suffices to show $L u=0$ in the sense of distributions. This is valid becasue if $\phi \in \mathcal{D}(\Omega)$ then

$$
0=\lim _{n \rightarrow \infty}\left\langle L u_{n}, \phi\right\rangle=\lim _{n \rightarrow \infty}\left\langle u_{n}, L^{\dagger} \phi\right\rangle=\left\langle u, L^{\dagger} \phi\right\rangle=\langle L u, \phi\rangle,
$$

i.e. $L u=0$ in the sense of distribution.

### 54.5 Some Non-Compact Considerations

In this section we will make use of the results from Section 53.3. Let $\rho \in$ $C^{\infty}\left(\mathbb{R}^{d},(0, \infty)\right), A_{i j} \in B C^{\infty}\left(\mathbb{R}^{d}\right)$ such that $\sum A_{i j} \xi_{i} \xi_{j}>0$ for all $|\xi| \neq 0$ and define

$$
Q(u, v)=\int_{\mathbb{R}^{d}} \sum A_{i j} \partial^{i} u \partial^{j} v \rho d x \text { for } u, v \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)
$$

Then as we have seen $Q$ has a closed extension $\hat{Q}$ and unique self adjoint operator $L$ on $L^{2}(\rho d x)$ such that $u \in D(L)$ iff $v \rightarrow Q(u, v)$ is $L^{2}(\rho d x)$ continuous on $\mathcal{D}(\hat{Q})$ and in which case $Q(u, v)=(L u, v)_{L^{2}(\rho d x)}$. Standard integration by parts shows

$$
u \in C_{c}^{2}\left(\mathbb{R}^{d}\right) \subset \mathcal{D}(L) \text { and } L u=-\frac{1}{\rho} \sum \partial_{j}\left(\rho A_{i j} \partial_{i} u\right)
$$

Proposition 54.17. Let $\bar{Q}(u, u):=\int_{\mathbb{R}^{d}} \sum A_{i j} \partial^{i} u \cdot \partial^{j} u \rho d x$, then

$$
\mathcal{D}(\hat{Q}) \subset\left\{u \in L^{2}(\rho d x) \cap H_{l o c}^{1}\left(\mathbb{R}^{d}\right): \bar{Q}(u, u)<\infty\right\}
$$

Proof. By definition of the closure of $Q, C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ is dense in $\left(\mathcal{D}(\hat{Q}), \hat{Q}_{1}\right)$. Since for all $\Omega \subset \subset \mathbb{R}^{d}$ there exists an $\epsilon=\epsilon(\Omega)$ such that $\rho \sum A_{i j} \xi_{i} \xi_{j} \geq \epsilon|\xi|^{2}$ on $\Omega$, we learn

$$
\begin{equation*}
Q(u, u) \geq \epsilon\|u\|_{H^{1}(\Omega)}^{2} \text { for all } u \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right) \tag{54.17}
\end{equation*}
$$

Therefore if $u \in \mathcal{D}(\hat{Q})$ and $u_{n} \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ such that $\hat{Q}_{1}\left(u-u_{n}\right) \rightarrow 0$ as $n \rightarrow \infty$ then $\left\|u_{n}-u_{m}\right\|_{H^{1}(\Omega)} \rightarrow 0$, i.e. $u=\lim _{n \rightarrow \infty} u_{n}$ in $H^{1}(\Omega)$. Hence a simple limiting argument shows Eq. (54.17) holds for all $u \in \mathcal{D}(\hat{Q})$ :

$$
Q_{1}(u, u) \geq \epsilon\|u\|_{H^{1}(\Omega)}^{2} \text { for all } u \in \mathcal{D}(\hat{Q})
$$

This shows $\mathcal{D}(\hat{Q}) \subset H_{\text {Loc }}^{1}\left(\mathbb{R}^{d}\right)$. Moreover

$$
\begin{array}{rlc}
\hat{Q}\left(u_{n}, u_{n}\right) & \geq \int_{\Omega} A_{i j} \partial_{i} u_{n} \partial_{j} u_{n} \rho d x \\
\downarrow & n & \rightarrow \infty \\
\hat{Q}(u, u) & \geq & \int_{\Omega} A_{i j} \partial_{i} u \partial_{j} u \rho d x
\end{array}
$$

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Since $\Omega \subset \subset \mathbb{R}^{d}$ is arbitrary this implies that

$$
\hat{Q}(u, u) \geq \int_{\Omega} A_{i j} \partial_{i} u \partial_{i} u \rho d x
$$

- 

Proposition 54.18. Suppose $u \in D\left(L^{k}\right)$ then $u \in H_{L o c}^{2 k}\left(\mathbb{R}^{d}\right)$ and for all $\Omega \subset \subset \Omega_{1} \subset \subset \mathbb{R}^{d}$ there exist $C=C_{k}(\Omega)$ such that

$$
\begin{equation*}
\|u\|_{H^{2 k}(\Omega)} \leq C\left(\left\|L^{k} u\right\|_{L^{2}(\rho d x)}+\|u\|_{L^{2}(\rho d x)}\right) \tag{54.18}
\end{equation*}
$$

Proof. Suppose $u \in D(L)$ and $L u=f$. Then for all $\phi \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$, $(L u, \phi)=(f, \phi)_{L^{2}(\rho)}$. Therefore

$$
\int \sum A_{i j} \partial_{i} u \partial_{j} u \rho d x=\int f \phi \rho d x
$$

so $-\partial_{j}\left(\rho A_{i j} \partial_{i} u\right)=\rho f$ in the sense if distributions and hence

$$
-\sum A_{i j} \partial_{j} \partial_{i} u+\text { L.O.T. }=f
$$

in the distributional sense. Since $D(L) \subset D(\hat{Q}) \subset H_{L o c}^{1}\left(\mathbb{R}^{d}\right)$, by local elliptic regularity it follows that $D(L) \subset H_{L o c}^{2}\left(\mathbb{R}^{d}\right)$ and for all $\Omega_{1} \supset \Omega_{0}$

$$
\|u\|_{H^{2}\left(\Omega_{0}\right)} \leq C\left(\|L u\|_{L^{2}\left(\Omega_{1}\right)}+\|u\|_{L^{2}\left(\Omega_{1}\right)}\right) .
$$

Now suppose $u \in D\left(L^{2}\right)$, then $u \in D(L) \subset H_{L o c}^{2}\left(\mathbb{R}^{d}\right)$ and $L u \in D(L) \subset$ $H_{L o c}^{2}\left(\mathbb{R}^{d}\right)$ implies $u \in H_{L o c}^{4}\left(\mathbb{R}^{d}\right)$ and

$$
\begin{aligned}
\|u\|_{H^{4}(\Omega)} & \leq C\left(\|L u\|_{H^{2}\left(\Omega_{1}\right)}+\|L u\|_{L^{2}\left(\Omega_{1}\right)}\right) \\
& \leq C\left(L^{2} u\left\|_{L^{2}\left(\Omega_{2}\right)}+\right\| u\left\|_{L^{2}\left(\Omega_{2}\right)}+\right\| u \|_{H^{2}\left(\Omega_{1}\right)}\right) \\
& \leq C\left(\left\|L^{2} u\right\|_{L^{2}\left(\Omega_{2}\right)}+\|L u\|_{L^{2}\left(\Omega_{2}\right)}+\|u\|_{L^{2}\left(\Omega_{2}\right)}\right) .
\end{aligned}
$$

If $u \in D\left(L^{3}\right)$ then $u \in H_{L o c}^{4}\left(\mathbb{R}^{d}\right)$ and $L u \in H_{L o c}^{4}\left(\mathbb{R}^{d}\right)$ implies $u \in H_{L o c}^{6}\left(\mathbb{R}^{d}\right)$ and

$$
\begin{aligned}
\|u\|_{H^{6}\left(\Omega_{0}\right)} & \leq C\left(\|L u\|_{H^{4}\left(\Omega_{1}\right)}+\|u\|_{L^{2}\left(\Omega_{1}\right)}\right) \\
& \leq C\left(\left\|L^{3} u\right\|_{L^{2}\left(\Omega_{2}\right)}+\left\|L^{2} u\right\|_{L^{2}\left(\Omega_{2}\right)}+\|L u\|_{L^{2}\left(\Omega_{2}\right)}+\|u\|_{L^{2}}\right) .
\end{aligned}
$$

$u \in D\left(L^{k}\right)$ implies $u \in H_{L o c}^{2 k}\left(\mathbb{R}^{d}\right)$ and

$$
\begin{aligned}
\|u\|_{H^{2}\left(\Omega_{0}\right)} & \leq C \sum_{j=0}^{k}\left\|L^{j} u\right\|_{L^{2}(\tilde{\Omega})} \leq C \sum_{j=0}^{k}\left\|L^{j} u\right\|_{L^{2}(\rho)} \\
& \leq C\left(\left\|L^{k} u\right\|_{L^{2}(\rho d s)}+\|u\|_{L^{2}(\rho d x)}\right)
\end{aligned}
$$

### 54.5.1 Heat Equation

Let $u(t)=e^{-t L} u_{0}$ where $u_{0} \in L^{2}(\rho)$. Then $u(t) \in D\left(L^{k}\right)$ for all $k$ when $t>0$ and hence $u(t) \in H_{\text {Loc }}^{2 k}\left(\mathbb{R}^{d}\right)$ for all $k$. But this implies for each $t>0$ that $u(t)$ has a continuous in fact $C^{\infty}$ - version because $H^{2 k}(\Omega) \hookrightarrow C^{2 k-2 / d}(\bar{\Omega})$ for $k>\frac{1}{d}$. Moreover

$$
\left\|L^{k}\left[\frac{u(t+h)-u(t)}{h}-L u(t)\right]\right\|_{L^{2}(\Omega)} \rightarrow 0 \text { as } h \rightarrow 0
$$

for all $k=0,1,2, \ldots$ and therefore,

$$
\left\|\frac{u(t+h)-u(t)}{h}-L u(t)\right\|_{C^{2 k-2 / d}(\bar{\Omega})} \rightarrow 0 \text { as } h \rightarrow 0
$$

when $k>\frac{1}{d}$. This shows $t \rightarrow u(t)$ is differentiable and in $C^{2 k-\frac{2}{n}(\bar{\Omega})}$ for all $k>\frac{1}{d}$. Thus we conclude that $u(t, x)$ is in $C^{1, \infty}\left((0, \infty) \times \mathbb{R}^{d}\right)$ and $\frac{\partial u}{\partial t}(t, x)=$ $L u(t, x)$, i.e. $u$ is a classical solution to the heat equation.

### 54.5.2 Wave Equation

Now consider the generalized solution to the wave equation

$$
u(t)=\cos (\sqrt{L} t) f+\frac{\sin (\sqrt{L} t)}{\sqrt{L}} g
$$

where $f, g \in L^{2}(\rho)$. If $f, g \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$, then $f, g \in D\left(L^{k}\right)$ for all $k$ and hence $u(t) \in D\left(L^{k}\right)$ for all $k$. It now follows that $u(t)$ is $C^{\infty}$-differentiable in $t$ relative to the norm $\|f\|_{k}:=\|f\|_{L^{2}(\rho d x)}+\left\|L^{k} f\right\|_{L^{2}(\rho d x)}$ for all $k \in \mathbb{N}$. So by the above ideas $u(t, x) \in C^{\infty}\left(\mathbb{R} \times \mathbb{R}^{d}\right)$ and

$$
\begin{aligned}
\ddot{u}(t, x)+L u(t, x) & =0 \text { with } \\
u(0, x) & =f_{0}(x) \text { and } \\
\dot{u}(0, x) & =g_{0}(x) .
\end{aligned}
$$

## Spectral Considerations

For this section, let $\Omega$ be a bounded open subset of $\mathbb{R}^{d}$ such that $\bar{\Omega}^{o}=\Omega$ and $\bar{\Omega}$ is $C^{\infty}$ - manifold with boundary. Also let $\mathcal{E}$ be a symmetric Dirichlet form with domain being either $X=H^{1}(\Omega)$ or $X=H_{0}^{1}(\Omega)$ and let $L:=L_{\mathcal{E}}$ be the corresponding self adjoint operator.

Theorem 55.1. There exist $\left\{\lambda_{i}\right\}_{i=1}^{\infty} \subset \mathbb{R}$ such that $\lambda_{1} \leq \lambda_{2} \leq \lambda_{3} \leq \cdots \rightarrow \infty$ and $\left\{\phi_{n}\right\} \subset D(L) \subset L^{2}(\Omega)$ such that $\left\{\phi_{n}\right\}$ is an orthonormal basis for $L^{2}(\Omega)$ and $L \phi_{n}=\lambda_{n} \phi_{n}$ for all $n$.

Proof. Choose $C>0$ such that $(L+C): D(L) \rightarrow L^{2}(\Omega)$ is invertible and let $T:=(L+C)^{-1}$ which is a compact operator (see Theorem 54.7) when viewed as an operator from $L^{2}(\Omega)$ to $L^{2}(\Omega)$. Since $L=L^{*}$,

$$
((L+C) u, v)=(u,(L+C) v) \text { for all } u, v \in D(L)
$$

and using this equation with $u, v$ being replaced by $T u, T v$ respectively shows $(u, T v)=(T u, v)$ for all $u, v \in L^{2}$. Moreover if $u \in L^{2}(\Omega)$ and $v=T u \in D(L)$,

$$
(T u, u)=(T(L+C) v,(L+C) v)=((L+C) v, v) \geq 0
$$

and so we have shown $T=T^{*}$ and $T>0$. By the spectral Theorem 16.17 for self-adjoint compact operators, there exist $\left\{\mu_{n}\right\}_{n=1}^{\infty} \subset \mathbb{R}_{+}$and an orthonormal basis $\left\{\phi_{n}\right\}_{n=1}^{\infty}$ of $L^{2}(\Omega)$ such that $T \phi_{n}=\mu_{n} \phi_{n}$ and $\mu_{1} \geq \mu_{2} \geq \mu_{3} \geq \ldots$ and $\lim \mu_{i}=0$. Since $T \phi_{n}=\mu_{n} \phi_{n}$ iff $(L+C)^{-1} \phi_{n}=\mu_{n} \phi_{n} \in D(L)$ iff $\mu_{n}(L+C) \phi_{n}=\phi_{n}$ iff

$$
L \phi_{n}=\frac{1}{\mu_{n}}\left(1-C \mu_{n}\right) \phi_{n}=\lambda_{n} \phi_{n}
$$

where $\lambda_{n}:=\left(\frac{1}{\mu_{n}}-C\right) \uparrow \infty$ as $n \rightarrow \infty$.
Corollary 55.2. Let $L$ be as above, then

$$
D(L)=\left\{u \in L^{2}(\Omega): \sum_{n=1}^{\infty} \lambda_{n}^{2}\left(u, \phi_{n}\right)^{2}<\infty\right\}
$$

Moreover $L u=\sum_{n=1}^{\infty} \lambda_{n}\left(u, \phi_{n}\right) \phi_{n}$ for all $u \in D(L)$, i.e. $L$ is unitarily equivalent to the operator $\Lambda: \ell^{2} \rightarrow \ell^{2}$ defined by $(\Lambda x)_{n}=\lambda_{n} x_{n}$ for all $n \in \mathbb{N}$.

Proof. Suppose $u \in D(L)$, then

$$
L u=\sum\left(L u, \phi_{n}\right) \phi_{n}=\sum\left(u, L \phi_{n}\right) \phi_{n}=-\sum \lambda_{n}\left(u, \phi_{n}\right) \phi_{n}
$$

with the above sums being $L^{2}$ convergent and hence

$$
\sum \lambda_{n}^{2}\left(u, \phi_{n}\right)^{2}=\|L u\|_{L^{2}}^{2}<\infty .
$$

Conversely if $\sum \lambda_{n}^{2}\left(u, \phi_{n}\right)^{2}<\infty$, let

$$
u_{N}:=\sum_{n=1}^{N}\left(u, \phi_{n}\right) \phi_{n} \in D(L) .
$$

Then $u_{N} \rightarrow u$ in $L^{2}(\Omega)$ and

$$
L u_{N}=\sum_{1}^{N}\left(u, \phi_{n}\right) \lambda_{n} \phi_{n} \rightarrow \sum_{1}^{\infty}\left(u, \phi_{n}\right) \lambda_{n} \phi_{n} \text { in } L^{2}(\Omega) .
$$

Since $L$ is a closed operator, $u \in D(L)$ and $L u=\sum_{1}^{\infty}\left(u, \phi_{n}\right) \lambda_{n} \phi_{n}$.

### 55.1 Growth of Eigenvalues I

Example 55.3. Let $\Omega=(0, \ell)$.

1. Suppose $L=-\frac{d^{2}}{d x^{2}}$ with $D(L)=H^{2}(\Omega) \cap H_{0}^{1}(\Omega)$, i.e. we impose Dirichlet boundary conditions. Because $L=L^{*}$ and $L \geq 0$ (in fact $L \geq \epsilon I$ for some $\epsilon>0$ by the by the Poincaré Lemma in Theorem 49.31) if $L u=\lambda u$ then $\lambda>0$. Let $\lambda=\omega^{2}>0$, then the general solution to $L u=\omega^{2} u$ is given by

$$
u(x)=A \cos (\omega x)+B \sin (\omega x)
$$

where $A, B \in \mathbb{C}$. Because we want $u(0)=0=u(\ell)$ we must require $A=0$ and $\ell \omega=n \pi$. Hence we have $\lambda_{n}=\frac{n^{2} \pi^{2}}{\ell^{2}}$ and $u_{n}(x)=\sqrt{\frac{2}{\ell}} \sin \left(\frac{n \pi}{\ell} x\right)$ for $n \in \mathbb{N}$ is an orthonormal basis of eigenvectors for $L$.

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2. The reader is invited to show that if $L=\frac{-d^{2}}{d x^{2}}$ with Neuwmann boundary conditions then $u_{0}(x):=\frac{1}{\sqrt{\ell}}$ and $u_{n}(x)=\sqrt{\frac{2}{\ell}} \cos \left(\frac{n \pi}{\ell} x\right)$ for $n \in \mathbb{N}$ forms an orthonormal basis of eigenfunctions of $L$ with eigenvalues given by $\lambda_{n}=\frac{n^{2} \pi^{2}}{\ell^{2}}$ for $n \in \mathbb{N}_{0}$.
3. Suppose that $L=-\Delta$ on $\Omega^{d}$ with Dirichlet boundary conditions and for $m \in \mathbb{N}^{d}$ let

$$
U_{m}(x)=u_{m_{1}}\left(x_{1}\right) \ldots u_{m_{d}}\left(x_{d}\right)
$$

where each $u_{i}$ is given as in Item 1. Then $\left\{U_{m}: m \in \mathbb{N}^{d}\right\}$ is an orthonormal basis of eigenfunctions of $L$ with eigenvalues given

$$
\lambda_{m}=\frac{\pi^{2}}{\ell^{2}} \sum_{i=1}^{d} m_{i}^{2}=\frac{\pi^{2}}{\ell^{2}}|m|_{\mathbb{R}^{d}}^{2} \text { for all } m \in \mathbb{N}^{d}
$$

Remark 55.4. Keeping the notation of item 3. of Example 55.3, for $\lambda>0$ let

$$
E_{\lambda}=\operatorname{span}\left\{\phi_{m}: \lambda_{m} \leq \lambda\right\}
$$

Then

$$
\operatorname{dim}\left(E_{\lambda}\right)=\#\left\{m \in \mathbb{N}^{d}: \lambda_{m} \leq \lambda\right\}=\#\left\{m \in \mathbb{N}^{d}:|m|_{\mathbb{R}^{d}}^{2} \leq \frac{\lambda \ell^{2}}{\pi^{2}}\right\}
$$

from which it follows that

$$
\operatorname{dim}\left(E_{\lambda}\right) \asymp m_{d}\left(B\left(0, \sqrt{\frac{\lambda \ell^{2}}{\pi^{2}}}\right)\right)=\omega_{d}\left(\frac{\lambda \ell^{2}}{\pi^{2}}\right)^{d / 2}=C m_{d}\left(\Omega^{\ell}\right) \lambda^{d / 2}
$$

Lemma 55.5. Let $\Omega, \mathcal{E}$ and $L$ be as described at the beginning of this section. Given $k \in \mathbb{N}$, there exists $C=C_{k}<\infty$ such that

$$
\begin{equation*}
\|u\|_{H^{2 k}(\Omega)} \leq C\left(\|u\|_{L^{2}(\Omega)}+\left\|L^{k} u\right\|_{L^{2}(\Omega)}\right) \text { for all } u \in D\left(L^{k}\right) \tag{55.1}
\end{equation*}
$$

Proof. We first claim that for $u \in D\left(L^{k}\right)$.

$$
\begin{equation*}
\|u\|_{H^{2 k}(\Omega)} \leq C_{k}\left(\|u\|_{L^{2}(\Omega)}+\|L u\|_{L^{2}(\Omega)}+\cdots+\left\|L^{k} u\right\|_{L^{2}(\Omega)}\right) \tag{55.2}
\end{equation*}
$$

We prove Eq. (55.2) by induction. When $k=0$, Eq. (55.2) is trivial. Consider $u \in D\left(L^{(k+1)}\right) \subset D\left(L^{k}\right)$. By elliptic regularity (Theorem 52.15) and then using the induction hypothesis,

$$
\|u\|_{H^{2(k+1)}(\Omega)} \leq C\left(\|u\|_{L^{2}(\Omega)}+\|L u\|_{H^{k}(\Omega)}\right) \leq C_{k+1}\left(\sum_{j=0}^{k+1}\left\|L^{j} u\right\|_{L^{2}(\Omega)}\right)
$$

This proves Eq. (55.2). Because $\sup \left\{\frac{\lambda^{2 m}}{\left(1+\lambda^{2 p}\right)}: \lambda \geq 0\right\}$ for any $p>m$, it follows by the spectral theorem that

$$
\left\|L^{m} u\right\|_{L^{2}(\Omega)}^{2} \leq C\left(\|u\|_{L^{2}(\Omega)}^{2}+\left\|L^{p} u\right\|_{L^{2}(\Omega)}^{2}\right) \text { for all } p>m
$$

Combining this fact with Eq. (55.2) implies Eq. (55.1).

Theorem 55.6. Continue the notation in Lemma 55.5 and let $\left\{\phi_{n}\right\}_{n=1}^{\infty}$ be the orthonormal basis described in Theorem 55.1 and for $\lambda \in \mathbb{R}$ let

$$
E_{\lambda}:=\operatorname{span}\left\{\phi_{n}: \lambda_{n} \leq \lambda\right\}
$$

If $k$ is the smallest integer such that $k>d / 4$, there exist $C<\infty$ such that

$$
\begin{equation*}
\operatorname{dim}\left(E_{\lambda}\right) \leq C\left(1+\lambda^{2 k}\right) \tag{55.3}
\end{equation*}
$$

for all $\lambda \geq \inf \sigma(L)$.
Proof. By the Sobolev embedding Theorem 49.18, $H^{2 k}(\Omega) \hookrightarrow C^{2 k-\frac{d}{2}}(\Omega) \subset$ $C(\bar{\Omega})$. Combining this with Lemma 55.5 implies $D\left(L^{k}\right) \hookrightarrow C(\bar{\Omega})$ and for $u \in D\left(L^{k}\right)$,

$$
\begin{equation*}
\|u\|_{C^{0}(\Omega)} \leq\|u\|_{C^{2 k-\frac{d}{2}}(\Omega)} \leq C\|u\|_{H^{2 k}(\Omega)} \leq C\left(\|u\|_{L^{2}(\Omega)}+\left\|L^{k} u\right\|_{L^{2}(\Omega)}\right) \tag{55.4}
\end{equation*}
$$

Let $\lambda \geq \inf \sigma(L)$ and $u \in E_{\lambda} \subset D\left(L^{k}\right)$. Since

$$
\begin{gather*}
u=\sum_{n: \lambda_{n} \leq \lambda}\left(u, \phi_{n}\right) \phi_{n} \text { and } L^{k} u=\sum_{\lambda_{n} \leq \lambda} \lambda_{n}^{k}\left(u, \phi_{n}\right) \phi_{n} \\
\left\|L^{k} u\right\|_{L^{2}(\Omega)}^{2}=\sum_{\lambda_{n} \leq \lambda}\left|\lambda_{n}\right|^{2 k}\left|\left(u, \phi_{n}\right)\right|^{2} \leq|\lambda|^{2 k}\|u\|_{L^{2}(\Omega)}^{2} \tag{55.5}
\end{gather*}
$$

Combining Eqs. (55.4) and (55.5) implies

$$
\begin{equation*}
\|u\|_{C^{0}} \leq C\left(1+\lambda^{k}\right)\|u\|_{L^{2}(\Omega)}=C\left(1+\lambda^{k}\right) \sqrt{\sum_{\lambda_{n} \leq \lambda}\left|\left(u, \phi_{n}\right)\right|^{2}} \tag{55.6}
\end{equation*}
$$

Let $N=\operatorname{dim}\left(E_{\lambda}\right), y \in \Omega$ and take $u(x):=\sum_{n=1}^{N} \phi_{n}(y) \phi_{n}(x)$ in Eq. (55.6) to find

$$
\sum_{n=1}^{N}\left|\phi_{n}(y)\right|^{2} \leq \sup _{x \in \Omega}\left|\sum_{n=1}^{N} \phi_{n}(y) \phi_{n}(x)\right| \leq C\left(1+\lambda^{k}\right) \sqrt{\sum\left|\phi_{n}(y)\right|^{2}}
$$

from which it follows that

$$
\sum_{n=1}^{N}\left|\phi_{n}(y)\right|^{2} \leq C^{2}\left(1+\lambda^{k}\right)^{2}
$$

Integrating this estimate over $y \in \Omega$ then shows

$$
\operatorname{dim}\left(E_{\lambda}\right)=\sum_{1}^{N} 1=\sum_{1}^{N} \int_{\Omega}\left|\phi_{n}(y)\right|^{2} d y \leq C^{2}\left(1+\lambda^{k}\right)^{2}|\Omega|
$$

which implies Eq. (55.3).

Corollary 55.7. Let $k$ be the smallest integer larger than $d / 4$. Then there exists $\epsilon>0$ such that $\lambda_{n} \geq \epsilon n^{1 / 2 k}$ for $n$ sufficiently large. Noting that $k \sim d / 4$ this says roughly that $\lambda_{n} \sim n^{d / 2}$, which is the correct result.

Proof. Since ${ }^{1}$

$$
n \leq \operatorname{dim}\left(E_{\lambda_{n}}\right) \leq c\left(1+\lambda_{n}^{2 k}\right)
$$

$\frac{n}{c}-1 \leq \lambda_{n}^{2 k}$ or $\lambda_{n} \geq\left(\frac{n}{c}-1\right)^{\frac{1}{2 k}}$. $\square$

[^7] $\operatorname{dim} E_{\lambda_{n}}$.

Heat Kernel Properties

## Construction of Heat Kernels by Spectral Methods

A couple of references for this and later sections are Davies [3, 4] and L. Saloff-Coste [12].

For this section, again let $\Omega$ be a bounded open subset of $\mathbb{R}^{d}$ such that $\bar{\Omega}^{o}=\Omega$ and $\bar{\Omega}$ is $C^{\infty}$ - manifold with boundary. Also let $\mathcal{E}$ be a symmetric Dirichlet form with domain being either $X=H^{1}(\Omega)$ or $X=H_{0}^{1}(\Omega)$ and let $L:=L_{\mathcal{E}}$ be the corresponding self adjoint operator. Let $\left\{\phi_{n}\right\}_{n=1}^{\infty}$ be the orthonormal basis of eigenvectors of $L$ as described in Theorem 55.1 and $\left\{\lambda_{n}\right\}_{n=1}^{\infty}$ denote the corresponding eigenvalues, i.e. $L \phi_{n}=\lambda_{n} \phi_{n}$.

As we have seen abstractly before,

$$
\begin{aligned}
u(t) & =\cos (\sqrt{L} t) f+\frac{\sin (\sqrt{L} t)}{\sqrt{L}} g \\
& :=\sum_{n=1}^{\infty}\left\{\cos \left(\sqrt{\lambda_{n}} t\right)\left(f, \phi_{n}\right)+\frac{\sin \left(\sqrt{\lambda_{n}} t\right)}{\sqrt{\lambda_{n}}}\left(g, \phi_{n}\right)\right\} \phi_{n}
\end{aligned}
$$

solves the wave equation

$$
\frac{\partial^{2} u}{\partial t^{2}}+L u=0 \text { with } u(0, x)=f(x) \text { and } \dot{u}(0, x)=g(x)
$$

and

$$
u(t)=e^{-t L} u_{0}:=\sum_{n=1}^{\infty} e^{-t \lambda_{n}}\left(u_{0}, \phi_{n}\right) \phi_{n}
$$

solves the heat equation,

$$
\begin{equation*}
\frac{\partial u}{\partial t}=-L u \text { with } u(0, x)=u_{0}(x) . \tag{56.1}
\end{equation*}
$$

Here we will concentrate on some of the properties of the solutions to the heat equation (56.1). Let us begin by writing out $u(t, x)$ more explicitly as
$u(t, x)=\sum_{n=1}^{\infty} e^{-t \lambda_{n}} \int_{\Omega} u_{0}(y) \phi_{n}(y) d y \phi_{n}(x)=\lim _{N \rightarrow \infty} \int_{\Omega} u_{0}(y) \sum_{n=1}^{N} e^{-t \lambda_{n}} \phi_{n}(x) \phi_{n}(y) d$

Theorem 56.1. Let $p_{t}(x, y)$ denote the heat kernel associated to $L$ defined by

$$
\begin{equation*}
p_{t}(x, y):=\sum_{1}^{\infty} e^{-t \lambda_{n}} \phi_{n}(x) \phi_{n}(y) \tag{56.3}
\end{equation*}
$$

Then

1. the sum in Eq. (56.3) is uniformly convergent for all $t>0$
2. $(t, x, y) \rightarrow p_{t}(x, y) \in C^{\infty}\left(\mathbb{R}^{+} \times \bar{\Omega} \times \bar{\Omega}\right)$
3. $u(t, x)=\int_{\Omega} p_{t}(x, y) u_{0}(y) d y$ solves Eq. (56.1).

Proof. Let

$$
p_{t}^{N}(x, y)=\sum_{n=1}^{N} e^{-t \lambda_{n}} \phi_{n}(x) \phi_{n}(y),
$$

then $(t, x, y) \rightarrow p_{t}^{N}(x, y) \in C^{\infty}(\mathbb{R} \times \bar{\Omega} \times \bar{\Omega})$. Since $L^{k} \phi_{n}(x)=\lambda_{n}^{k} \phi_{n}$, by the Elliptic regularity Theorem 52.15 ,

$$
\left\|\phi_{n}\right\|_{H^{2 k}(\Omega)} \leq C\left(\left\|\phi_{n}\right\|_{L^{2}(\Omega)}+\left\|L^{k} \phi_{n}\right\|_{L^{2}(\Omega)}\right) \leq C\left(1+\lambda_{n}^{k}\right) .
$$

Taking $k>d / 4$, the Sobolev embedding Theorem 49.18 implies

$$
\left\|\phi_{n}\right\|_{C^{0}(\bar{\Omega})} \leq\left\|\phi_{n}\right\|_{C^{2 k-d / 2}(\bar{\Omega})} \leq C\left\|\phi_{n}\right\|_{H^{2 k}(\Omega)} \leq C\left(1+\lambda_{n}^{k}\right)
$$

Therefore $\sup _{x, y \in \Omega}\left|\phi_{n}(x) \phi_{n}(y)\right| \leq C^{2}\left(1+\lambda_{n}^{k}\right)^{2}$ while by Corollary 55.7, $\lambda_{n} \geq n^{2 / d}$ and therefore while $\sum_{n=1}^{\infty} e^{-t \lambda_{n}}\left(1+\lambda_{n}^{k}\right)^{2}<\infty$. More generally if $|\alpha|=2 m$

$$
\left\|\partial^{\alpha} \phi_{n}\right\|_{C^{0}} \leq\left\|\partial^{\alpha} \phi_{n}\right\|_{H^{2 k}} \leq\left\|\phi_{n}\right\|_{H^{2(k+m)}} \leq C\left(1+\lambda_{n}^{k+m}\right)
$$

and hence

$$
\left\|\phi_{n} \otimes \phi_{n}\right\|_{C^{2 m}(\bar{\Omega} \times \Omega)} \leq C^{2}\left(1+\lambda_{n}^{k+m}\right)^{2}
$$

from which it follows that

$$
\begin{aligned}
\sum_{n=1}^{\infty} \sup _{t \geq \epsilon} e^{-t \lambda_{n}}\left\|\phi_{n} \otimes \phi_{n}\right\|_{C^{2 m}(\bar{\Omega} \times \Omega)} & =\sum_{n=1}^{\infty} e^{-\epsilon \lambda_{n}}\left\|\phi_{n} \otimes \phi_{n}\right\|_{C^{2 m}(\bar{\Omega} \times \Omega)} \\
& \leq C^{2} \sum_{n=1}^{\infty} e^{-\epsilon \lambda_{n}}\left(1+\lambda_{n}^{k+m}\right)^{2}<\infty
\end{aligned}
$$

So $p_{t}^{N}(x, y)$ and all of its derivatives converge uniformly in $t \geq \epsilon$ and $x, y \in \bar{\Omega}$ as $N \rightarrow \infty$. Therefore $p_{t}(x, y):=\lim _{N \rightarrow \infty} p_{t}^{N}(x, y)$ exists and $(t, x, y) \rightarrow p_{t}(x, y)$

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is $C^{\infty}(\bar{\Omega})$ for $t>0$ and $x, y \in \bar{\Omega}$. It is now easy to justify passing the limit under the integral sign in Equation (56.2) to find $u(t, x)=\int_{\Omega} p_{t}(x, y) u_{0}(y) d y$.

Remark 56.2. $p_{t}(x, y)$ solves the following problem $\frac{\partial p_{t}}{\partial t}=-L_{x} p_{t}, p_{t}(\cdot, y)$ satisfies the boundary conditions and $\lim _{t \downarrow 0} p_{t}(\cdot, y)=\delta_{y}$.

Definition 56.3. A bounded operator $T: L^{2}(\Omega) \rightarrow L^{2}(\Omega)$ is positivity preserving if for every $f \in L^{2}(\Omega)$ with $f \geq 0$ a.e. on $\Omega$ has the property that $T f \geq 0$ a.e. on $\Omega$.

Proposition 56.4 (Positivity of heat kernel's). Suppose $D(L)=H_{0}^{1}(\Omega) \cap$ $H^{2}(\Omega)$, i.e. L has Dirichlet boundary conditions, then the operator $e^{-t L}$ is positivity preserving for all $t>0$ and the associated heat kernel $p_{t}(x, y)$ is non-negative for all $t \in(0, \infty)$ and $x, y \in \Omega$.

Proof. Since $e^{-t(L+C)}=e^{-t C} e^{-t L}$ is positivity preserving iff $e^{-t L}$ is positivity preserving, we may assume

$$
L=-\sum a_{i j} \partial_{i} \partial_{j}+\sum a_{i} \partial_{i}+a
$$

with $a>0$. Let $f \in C^{\infty}(\bar{\Omega},(0, \infty))$ and $u(t, x):=e^{-t L} f(x)$, in which case $u$ solves

$$
\begin{gathered}
\frac{\partial u}{\partial t}=-L u \text { with } u(0, x)=f(x) \geq 0 \text { for }(t, x) \in[0, T] \times \Omega \\
\text { with } u(t, x)=0 \text { for } x \in \partial \Omega
\end{gathered}
$$

If there exist $\left(t_{0}, x_{0}\right) \in(0, T] \times \Omega$ such that

$$
u\left(t_{0}, x_{0}\right)=\min \{u(t, x): 0 \leq t \leq T, x \in \bar{\Omega}\}<0
$$

then $\frac{\partial u}{\partial t}\left(t_{0}, x_{0}\right) \leq 0, \partial_{i} u\left(t_{0}, x_{0}\right)=0$ for all $i$ and by ellipticity,

$$
a_{i j}\left(x_{0}\right) \partial_{i} \partial_{j} u\left(t_{0}, x_{0}\right) \geq 0 .
$$

Therefore at $\left(t_{0}, x_{0}\right)$,

$$
\begin{aligned}
0 & =\frac{\partial u}{\partial t}+L u=\frac{\partial u}{\partial t}-a_{i j} \partial_{i} \partial_{j} u+a_{i} \partial_{i} u+a u \\
& =\frac{\partial u}{\partial t}-a_{i j} \partial_{i} \partial_{j} u+a u=(\leq 0)-(\geq 0)+(<0)<0
\end{aligned}
$$

which is a contradiction. Hence we have shown

$$
\begin{equation*}
0 \leq u(t, x)=\int_{\Omega} p_{t}(x, y) f(y) d y \text { for all }(t, x) \in[0, T] \times \bar{\Omega} \tag{56.4}
\end{equation*}
$$

By a simple limiting argument, Eq. (56.4) also holds for all non-negative bounded measurable functions $f$ on $\Omega$. Indeed, let $f_{n}:=\left(f 1_{\Omega}+n^{-1}\right) *$ $\eta_{n}$ where $\eta_{n} \in C_{c}^{\infty}\left(\mathbb{R}^{d},(0, \infty)\right)$ is a spherically symmetric approximate $\delta$ sequence. Then $f_{n} \in C^{\infty}(\bar{\Omega},(0, \infty))$ and hence

$$
0 \leq \int_{\Omega} p_{t}(x, y) f_{n}(y) d y \rightarrow \int_{\Omega} p_{t}(x, y) f(y) d y .
$$

From this equation it follows that $p_{t}(x, y) \geq 0$ and that $e^{-t L}$ is positivity preserving.
Lemma 56.5. Suppose $f_{n} \geq 0$ on $\Omega$ and $f_{n} \rightarrow \delta_{x}$ as $n \rightarrow \infty$, then $\lim _{n \rightarrow \infty} \int_{\Omega} f_{n}^{2}(x) d x=\infty$.

Proof. For sake of contradiction assume $\lim _{n \rightarrow \infty} \int f_{n}^{2}(x) d x \neq \infty$. By passing to a subsequence if necessary we may then assmume $M=: \sup _{n} \int_{\Omega} f_{n}^{2}(x) d x<$ $\infty$ and that $f_{n}$ converges weakly to some $g \in L^{2}(\Omega)$. In which case we would have

$$
\phi(x)=\lim _{n \rightarrow \infty} \int_{\Omega} f_{n}(y) \phi(y) d y=\int_{\Omega} g(y) \phi(y) d y \text { for all } \phi \in C_{c}^{\infty}(\Omega)
$$

But this would imply that $g=\delta_{x}$ which is incomensurate with $g$ being an $L^{2}(\Omega)$ function and we have reached the desired contradiction.

Theorem 56.6. Let $p_{t}$ be a Dirichlet heat kernel, then $\lim _{t \downarrow 0} p_{t}(x, x)=\infty$ and $p_{t}(x, y)>0$ for all $x, y \in \Omega$ and $t>0$.

Proof. We have seen $p_{t}(x, \cdot) \rightarrow \delta_{x}$ for all $x \in \Omega$. Therefore by lemma

$$
\lim _{t \downarrow 0} \int_{\Omega} p_{t}(x, y)^{2} d y=\infty
$$

Now

$$
\int_{\Omega} p_{t}(x, y)^{2} d y=\int_{\Omega} p_{t}(x, y) p_{t}(y, x) d y=p_{2 t}(x, x)
$$

Therefore $\lim _{t \downarrow 0} p_{2 t}(x, x)=\infty$ for all $x \in \Omega$.
Sketch of the rest: Choose a compact set $K$ in $\Omega$, then by continuity there exists $t_{0}>0$ and $\epsilon>0$ such that $p_{t}(x, y) \geq 1$ for all $x \in K,|x-y|<$ $\epsilon$ and $0<t \leq t_{0}$. (Note $\tanh p_{t}(x, x)$ is continuous on $[0,1] \times \Omega$, where $\tanh p_{0}(x, x)=1$.) Now if $x \in K$ and $y \in K$ and $|x-y|<2 \epsilon$, we have for $t<t_{0}$ that

$$
\begin{aligned}
p_{t}(x, y) & =\int_{\Omega} p_{t / 2}(x, z) p_{t / 2}(z, y) d z \geq \int_{B(x, \epsilon) \cap B(y, \epsilon) \cap \Omega} p_{t / 2}(x, z) p_{t / 2}(z, y) d z \\
& \geq m(B(x, \epsilon) \cap B(y, \epsilon) \cap \Omega)>0
\end{aligned}
$$

Working inductively, we may use this same idea to prove that $p_{t}(x, y)>0$ for all $t<t_{0}$ and $x, y \in K$. Moreover the same semi-group argument allows one to show $p_{t}(x, y)>0$ for all $t>0$ as well.

## Second Proof using the strong maximum principle.

Now for $y$ fixed $0<t \leq T, f(t, x)=p_{t}(x, y)$ solves $\frac{\partial f}{\partial t}=$ $L_{x} f(t, x), f(\epsilon, x)=p_{\epsilon}(x, y)$. With out loss of generality, assume

$$
L=-\sum-\partial_{i} a_{i j} \partial_{j}+a \text { with } a \geq 0
$$

For if $a$ is not greater than 0 , replace $L$ by $L+\lambda$ and observe that $e^{-t(L+\lambda)}=$ $e^{-t \lambda} e^{t L}$. Therefore $p_{t}^{\lambda}(x, y)=e^{-t \lambda} p_{t}(x, y)$ so $p_{t}(x, y)>0$ iff $p_{t}^{\lambda}(x, y) \geq 0$. Then $\frac{\partial p_{t}}{\partial t}=L_{x} p_{t}(x, y)$ for all $T \geq t \geq \epsilon$ and $x \in \Omega$. By the strong maximum principle of Theorem 12 on page 340 if there exist $\left(x_{0}, t_{0}\right) \in \Omega \times(\epsilon, T]$ such that $p_{t_{0}}\left(x_{0}, y\right)=0$ then $x \rightarrow p_{t}(x, y)$ is a constant on $\left(0 t_{0}\right) \times \Omega$ which is false because the constant would have to be 0 , but $\int p_{t_{0}}(x, y) d x>0$ for $T$ small.

### 56.1 Positivity of Dirichlet Heat Kernel by Beurling Deny Methods

Assumption 6 Suppose $L=-\partial_{i} a_{i j} \partial_{j}+a$ where $a_{i j}=a_{j i} \geq \epsilon I$ and $a \in$ $C^{\infty}(\bar{\Omega})$ and $D(L)=H_{0}^{1}(\Omega) \cap H^{2}(\Omega)$.
Theorem 56.7. Let $\lambda_{0}=\max (-a)$ i.e. $-\lambda_{0}=\min (a)$. Then for all $\lambda \geq \lambda_{0}$, $L_{\lambda}:=L+\lambda I: D(L) \rightarrow L^{2}(\Omega)$ invertible and if $f \in L^{2}(\Omega), f \geq 0$ a.e then $L_{\lambda}^{-1} f \geq 0$ a.e., i.e. $L_{\lambda}^{-1}$ is positivity preserving.

Proof. If $\lambda \geq \lambda_{0}$ then $\lambda+a \geq 0$ and hence if $L_{\lambda} u=0$ then

$$
\left(L_{\lambda} u, u\right)=\int_{\Omega}\left(a_{i j} \partial_{i} u \partial_{j} u+(a+\lambda) u^{2}\right) d x=0 .
$$

This implies $\nabla u=0$ a.e. and so $u$ is constant and hence $u=0$ because $u \in H_{0}^{1}$. Therefore $\operatorname{Nul}\left(L_{\lambda}\right)=\{0\}$ and so $L_{\lambda}$ is invertible by the Fredholm alternative. Now suppose $f \in C^{\infty}(\bar{\Omega})$ such that $f>0$, then $u=L_{\lambda}^{-1} f \in C^{\infty}(\bar{\Omega})$ with $L_{\lambda} u=f>0$ and $u=0$ on $\partial \Omega$. We may now use the maximum principle idea in Theorem 45.16 to conclude that $u \geq 0$. Indeed if there exists $x_{0} \in \Omega$ such that $u\left(x_{0}\right)=\min u<0$ at $x_{0}$, then

$$
\begin{aligned}
0<f\left(x_{0}\right)=\left(L_{\lambda} u\right)\left(x_{0}\right)=- & \underbrace{\left(a_{i j} \partial_{i} \partial_{j} u\right)\left(x_{0}\right)}_{\geq 0}+\underbrace{\left(\partial_{i} a_{i j}\right) \partial_{j} u\left(x_{0}\right)}_{=0} \\
& +\underbrace{(a+\lambda) u\left(x_{0}\right)}_{\leq 0} \leq 0
\end{aligned}
$$

which is a contradiction. Thus we have shown $u=L_{\lambda}^{-1} f \geq 0$ if $f \in$ $C^{\infty}(\bar{\Omega},(0, \infty))$. Given $f \in L^{2}(\Omega)$ such that $f \geq 0$ a.e. on $\bar{\Omega}$, choose $f_{n} \in C^{\infty}(\bar{\Omega})$ such that $f_{n} \geq \frac{1}{n}$ and $f_{n} \rightarrow f$ in $L^{2}(\Omega)$ and $f_{n} \rightarrow f$ a.e. on $\bar{\Omega}$. (For example, take $f_{n}=\eta_{\delta_{n}} *\left(f 1_{\bar{\Omega}}+\frac{1}{n}\right)$ say.) Then $u_{n}=L_{\lambda}^{-1} f_{n} \geq 0$ for all $n$ and $u_{n} \rightarrow u=L_{\lambda}^{-1} f$ in $H^{2}(\Omega)$. By passing to a subsequence if necessary we may assume that $u_{n} \rightarrow u$ a.e. from which it follows that $u \geq 0$ a.e. on $\Omega$.

Theorem 56.8. Keeping $L$ as above, $e^{-t L}: L^{2}(\Omega) \rightarrow L^{2}(\Omega)$ is positivity preserving for all $t \geq 0$.

Proof. By the spectral theorem and the fact that $\left(1+\frac{t \lambda}{n}\right)^{-1 n} \rightarrow e^{-t \lambda}$ boundedly for and $\lambda \geq 0$,

$$
e^{-t L} f=\lim _{n \rightarrow \infty}\left(1+\frac{t L}{n}\right)^{-n} f=\lim _{n \rightarrow \infty}\left(\frac{t}{n}\right)^{-n}\left[\left(\frac{n}{t}+L\right)^{-1}\right]^{n} f
$$

Now $\left(\frac{n}{t}+L\right)^{-1}$ is positivity preserving operator on $L^{2}(\Omega)$. Where $\frac{n}{t} \geq \lambda_{0}$ and hence so is the $n$-fold product. Thus if

$$
u_{n}:=\left(\frac{t}{n}\right)^{-n}\left[\left(\frac{n}{t}+L\right)^{-1}\right]^{n} f
$$

then $u_{n} \geq 0$ a.e. and $u_{n} \rightarrow e^{-t L} f$ in $L^{2}(\Omega)$ implies $e^{-t L} f \geq 0$ a.e.
Theorem 56.9. $p_{t}(x, y) \geq 0$ for all $x, y \in \Omega$.
Proof. $f \in L^{2}(\Omega)$ with $f \geq 0$ a.e. on $\Omega, e^{-t L} f \in C^{\infty}(\bar{\Omega})$ and $e^{-t L} f \geq 0$ a.e. by above. Thus $e^{-t L} f \geq 0$ everywhere. Now

$$
\int_{\Omega} p_{t}(x, y) f(y) d y=\left(e^{-t L} f\right)(x) \geq 0
$$

for all $f \geq 0$. Since $p_{t}(x, y)$ is smooth this implies $p_{t}(x, y) \geq 0$ for all $y \in \bar{\Omega}$ and since $x \in \bar{\Omega}$ was arbitrary we learn $p_{t}(x, y) \geq 0$ for all $x, y \in \bar{\Omega}$.
???? BRUCE for $f \in C_{c}^{\infty}(\Omega) \subset \cap D\left(L^{n}\right)=C^{\infty}(L)$ we have $\|\left(e^{-t L} f\right)-$ $f \|_{H^{k}} \rightarrow 0$ as $t \downarrow 0$ for all $k$. By Sobolev embedding this implies that $\left(e^{-t L} f\right) \rightarrow$ $f(x)$ as $t \downarrow 0$ for all $x \in \Omega$ i.e. $\int_{\Omega} p_{t}(x, y) f(y) d y \rightarrow f(x)$.

## Nash Type Inequalities and Their Consequences

Corollary 57.1. Suppose $d>2$, then there is a constant $C_{d}<\infty$ such that

$$
\begin{equation*}
\|u\|_{2}^{2+4 / d} \leq C_{d}\|\nabla u\|_{2}^{2}\|u\|_{1}^{4 / d} \tag{57.1}
\end{equation*}
$$

for all $u \in C_{c}^{1}\left(\mathbb{R}^{d}\right)$.
Proof. By Corollary 49.15, $\|u\|_{2^{*}} \leq C\|\nabla u\|_{2}$ where $2^{*}=\frac{2 d}{d-2}$ and by interpolation

$$
\|u\|_{2} \leq\|u\|_{p}^{\theta}\|u\|_{q}^{1-\theta}
$$

where $\frac{\theta}{p}+\frac{1-\theta}{q}=\frac{1}{2}$. Taking $p=2^{*}$ and $q=1$ implies $\frac{\theta}{2^{*}}+1-\theta=\frac{1}{2}$, i.e. $\theta\left(\frac{1}{2^{*}}-1\right)=-\frac{1}{2}$ and hence

$$
\begin{aligned}
\theta & =\frac{\frac{1}{2}}{1-\frac{1}{2^{*}}}=\frac{2^{*}}{2\left(2^{*}-1\right)}=\frac{d}{(d-2)} \cdot \frac{1}{\frac{2 d}{d-2}-1} \\
& =\frac{d}{d-2} \frac{d-2}{d+2}=\frac{d}{d+2}
\end{aligned}
$$

and $1-\theta=\frac{2}{d+2}$. Hence

$$
\|u\|_{2} \leq\|u\|_{2^{*}}^{\frac{d}{d+2}}\|u\|_{1}^{\frac{2}{d+2}} \leq C^{\frac{d}{d+2}}\|\nabla u\|_{2}^{\frac{d}{d=2}}\|u\|_{1}^{\frac{2}{d+2}}
$$

and therefore

$$
\|u\|_{2}^{\frac{d+2}{d}} \leq C\|\nabla u\|_{2}\|u\|_{1}^{\frac{2}{d}}
$$

and squaring this equation then gives the estimate in Eq. (57.1)
Proposition 57.2 (Nash). Corollary 57.1 holds for all d.
Proof. Since the Fourier transform is unitary, for any $R>0$ and $\|\hat{u}\|_{\infty} \leq$ $\|u\|_{L^{1}}$,

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$$
\begin{aligned}
\|u\|_{2}^{2} & =\int_{\mathbb{R}^{d}}|\hat{u}(\xi)|^{2} d \xi=\int_{|\xi| \leq R}|\hat{u}|^{2} d \xi+\int_{|\xi|>R}|\hat{u}|^{2} d \xi \\
& \leq \sigma\left(S^{d-1}\right) R^{d}\|u\|_{L^{1}}^{2}+\frac{1}{R^{2}} \int_{|\xi|>R}|\xi|^{2}|\hat{u}|^{2} d \xi \\
& \leq \sigma\left(S^{d-1}\right) R^{d}\|u\|_{L^{1}}^{2}+\frac{1}{R^{2}}\|D u\|_{2}^{2}=f(R)
\end{aligned}
$$

where $f(R)=a R^{d}+\frac{b}{R^{2}}$ and $a=\sigma\left(S^{d-1}\right)\|u\|_{L^{1}}^{2}$ and $b=\|D u\|_{2}^{2}$. To minimize $f$, we set $f^{\prime}(R)=0$ to find $d a R^{d-1}-2 b R^{-3}=0$, i.e. $R^{d+2}=\frac{2 b}{d a}$ and hence $R=\left(\frac{2 b}{d a}\right)^{\frac{1}{d+2}}$. With this value of $R$, we find

$$
\begin{aligned}
f(R) & =\frac{1}{R^{2}}\left(b+a R^{d+2}\right)=\left(\frac{d a}{2 b}\right)^{\frac{2}{d+2}}\left(b+\frac{2 b}{d}\right) \\
& =\left(\frac{d+2}{d}\right) b\left(\frac{d a}{2 b}\right)^{\frac{2}{d+2}}=C_{d} a^{\frac{2}{d+2}} b^{1-\frac{2}{d+2}}=C_{d} a^{\frac{2}{d+2}} b^{\frac{d}{d+2}}
\end{aligned}
$$

which gives the estimate

$$
\|u\|_{2}^{2} \leq C_{d}\|D u\|_{2}^{\frac{2 d}{d+2}}\|u\|_{1}^{\frac{4}{d+2}}
$$

which is equivalent to

$$
\|u\|_{2}^{2+4 / d}=\|u\|_{2}^{2\left(\frac{d+2}{d}\right)} \leq C_{d}\|D u\|_{2}^{2}\|u\|_{1}^{4 / d} .
$$

Proposition 57.3. Suppose $A(x)=\left\{a_{i j}(x)\right\}_{i, j=1}^{d}$ such that there exists $\epsilon>0$ and $M<\infty$ such that $\epsilon I \leq A(x) \leq M I$ for all $x \in \mathbb{R}^{d}$. Define $\mathcal{D}(\mathcal{E})=$ $W^{1,2}\left(\mathbb{R}^{d}\right)$ and

$$
\mathcal{E}(u, v)=\sum_{i j} \int_{\mathbb{R}^{d}} a_{i j} D_{i} u D_{i} v d x
$$

Then $\mathcal{E}$ is a closed symmetric quadratic form. Moreover $C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ is a core for $\mathcal{E}$.

Proof. Clearly $\|u\|_{W^{1,2}} \leq \frac{1}{\epsilon}\left(\|u\|_{L^{2}(\Omega)}^{2}+\mathcal{E}(u, u)\right)$ and

$$
\mathcal{E}(u, u) \leq M\|D u\|^{2} \leq M\|u\|_{W^{1,2}} .
$$

and hence

$$
\|\cdot\|_{2}^{2}+\mathcal{E}(\cdot, \cdot) \asymp\|\cdot\|_{W^{1,2}}
$$

so that $\left(\mathcal{D}(\mathcal{E}), \sqrt{\|\cdot\|_{2}^{2}+\mathcal{E}(\cdot, \cdot)}\right)$ is complete.
$1120 \quad 57$ Nash Type Inequalities and Their Consequences and hence for all $f \in C_{c}^{\infty}(\Omega)$,

$$
\left\langle\phi_{\epsilon}(u) D u, f\right\rangle=\left\langle D\left[\psi_{\epsilon}(u)\right], f\right\rangle=\left\langle\psi_{\epsilon}(u),-D f\right\rangle=0(\epsilon) \rightarrow 0
$$

Combining this with the observation that

$$
\phi_{\epsilon}(u) D u \xrightarrow{L^{p}} 1_{\{u=0\}} D u \text { as } \epsilon \downarrow 0
$$

implies

$$
\int_{\Omega} 1_{\{u=0\}} D u \cdot f d m=0 \text { for all } f \in C_{c}^{\infty}(\Omega)
$$

which proves $1_{\{u=0\}} D u=0$ a.e.
Exercise 57.7. Let $u \in W^{1, p}$. Show

1. If $\phi \in C^{1}(\mathbb{R}), \phi(0)=0$ and $\left|\phi^{\prime}\right| \leq M<\infty$, then $\phi(u) \in W^{1, p}$ and $D \phi(u)=\phi^{\prime}(u) D u$ a.e.
2. $|u| \in W^{1, p}$ and $D|u|=\operatorname{sgn}(u) D u$.
3. Eq. (57.2) holds.

Solution 57.8. Let $u_{n} \in C^{\infty}\left(\mathbb{R}^{d}\right) \cap W^{1, p}$ such that $u_{n} \rightarrow u$ in $W^{1, p}$. By passing to a subsequence if necessary we may further assuem that $u_{n}(x) \rightarrow$ $u(x)$ for a.e. $x \in \mathbb{R}^{d}$. Since $|\phi(u)| \leq M|u|$

$$
\left|\phi(u)-\phi\left(u_{n}\right)\right| \leq M\left|u-u_{n}\right|
$$

it follows that $|\phi(u)|,\left|\phi\left(u_{n}\right)\right| \in L^{p}$ and $\phi\left(u_{n}\right) \rightarrow \phi(u)$ in $L^{p}$. Since $\phi^{\prime}\left(u_{n}\right)$ is bounded, $\phi^{\prime}\left(u_{n}\right) \rightarrow \phi^{\prime}(u)$ a.e. and $\partial_{i} u_{n} \rightarrow \partial_{i} u$ in $L^{p}$, it follows that

$$
\begin{aligned}
& \left\|\partial_{i} \phi\left(u_{n}\right)-\partial_{i} \phi(u)\right\|_{p} \\
& \quad=\left\|\phi^{\prime}\left(u_{n}\right) \partial_{i} u_{n}-\phi^{\prime}(u) \partial_{i} u\right\|_{p} \\
& \quad \leq\left\|\phi^{\prime}\left(u_{n}\right)\left[\partial_{i} u_{n}-\partial_{i} u\right]\right\|_{p}+\left\|\left[\phi^{\prime}(u)-\phi^{\prime}\left(u_{n}\right)\right] \partial_{i} u\right\|_{p} \\
& \quad \leq M\left\|\partial_{i} u_{n}-\partial_{i} u\right\|_{p}+\left\|\left[\phi^{\prime}(u)-\phi^{\prime}\left(u_{n}\right)\right] \partial_{i} u\right\|_{p} \rightarrow 0
\end{aligned}
$$

where the second term is handled by the dominated convergence theorem. Therefore $\phi(u) \in W^{1, p}$ and $\partial_{i} \phi(u)=\phi^{\prime}(u) \partial_{i} u$.

Let $\phi_{\epsilon}(x):=\sqrt{x^{2}+\epsilon^{2}}$, then $\phi_{\epsilon}^{\prime}(x)=\frac{x}{\sqrt{x^{2}+\epsilon^{2}}},\left|\phi_{\epsilon}^{\prime}(x)\right| \leq 1$ for all $x$ and

$$
\lim _{\epsilon \downarrow 0} \phi_{\epsilon}^{\prime}(x)=\left\{\begin{array}{cc}
0 & \text { if } x=0 \\
\operatorname{sgn}(x) & \text { if } x \neq 0
\end{array}\right.
$$

From part 1., $\phi_{\epsilon}(u) \in W^{1, p}$ and $\partial_{i} \phi_{\epsilon}(u)=\phi_{\epsilon}^{\prime}(u) \partial_{i} u$. Since $\phi_{\epsilon}(u) \rightarrow|u|$ in $L^{p}$ and

$$
\partial_{i} \phi_{\epsilon}(u)=\phi_{\epsilon}^{\prime}(u) \partial_{i} u \rightarrow 1_{u \neq 0} \operatorname{sgn}(u) \partial_{i} u=\operatorname{sgn}(u) \partial_{i} u \text { a.e. }
$$

where the last equality is a consequence of Proposition 57.6. Hence we see that $|u| \in W^{1, p}$ and $D|u|=\operatorname{sgn}(u) D u$.

Remark: (BRUCE) I think using the absolute continuity of $u$ along lines could be used to simplify and generalize the above exercise to the case where $\phi \in A C(\mathbb{R})$ with $\left|\phi^{\prime}(x)\right| \leq M<\infty$ for $m$ - a.e. $x \in \mathbb{R}$.

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Remark 57.11. From Eq. (57.3),

$$
\begin{aligned}
\left\|e^{2 t L} f\right\|_{L^{\infty}}^{2} & =\left\|e^{t L} e^{t L} f\right\|_{L^{\infty}}^{2} \leq \frac{C}{t^{d / 2}}\left\|e^{t L} f\right\|_{L^{2}}^{2}=\frac{C}{t^{d / 2}}\left|\left(e^{2 t L} f, f\right)\right| \\
& \leq \frac{C^{2}}{t^{d / 2}}\left\|e^{2 t L} f\right\|_{L^{\infty}}\|f\|_{L^{1}}
\end{aligned}
$$

and hence Eq. (57.3) implies the inequality,

$$
\begin{equation*}
\left\|e^{2 t L} f\right\|_{L^{\infty}} \leq \frac{C^{2}}{t^{d / 2}}\|f\|_{L^{1}} \tag{57.5}
\end{equation*}
$$

Lemma 57.12 (Duality Lemma). Let $T$ be a linear operator on $\cup_{p \in[1, \infty]} L^{p}$ such that $T\left(L^{1} \cap L^{\infty}\right)=\left(L^{1} \cap L^{\infty}\right)$ and $T: L^{2} \rightarrow L^{2}$ is self adjoint. If

$$
\|T f\|_{p} \leq C\|f\|_{q} \text { for all } f \in L^{1} \cap L^{\infty}
$$

then

$$
\|T f\|_{q} \leq C\|f\|_{p^{\prime}} \text { for all } f \in L^{1} \cap L^{\infty}
$$

where $\frac{1}{q^{\prime}}+\frac{1}{q}=1$ and $\frac{1}{p}+\frac{1}{p^{\prime}}=1$.
Proof.

$$
\begin{aligned}
\|T f\|_{q^{\prime}} & =\sup _{\|g\|_{q}=1}|(T f, g)|=\sup _{\|g\|_{q}=1}|(f, T g)| \leq \sup _{\|g\|_{q}=1}\|f\|_{p^{\prime}}\|T g\|_{p} \\
& \leq C\|f\|_{p^{\prime}} \sup _{\|g\|_{q}=1}\|g\|_{q}=C\|f\|_{p^{\prime}} .
\end{aligned}
$$

Proposition 57.13 (Converse of Theorem 57.10. ). If

$$
\begin{equation*}
\left\|e^{t L} f\right\|_{\infty} \leq C t^{-d / 4}\|f\|_{2} \text { for all } f \in L^{2} \tag{57.6}
\end{equation*}
$$

then

$$
\|f\|_{2}^{2+4 / d} \leq \frac{4 C}{d} \mathcal{E}(f, f)\|f\|_{1}
$$

Proof. By the duality, Lemma 57.12, (57.6) implies

$$
\left\|e^{t L} f\right\|_{2} \leq C t^{-d / 4}\|f\|_{1} \text { for all } f \in L^{2}
$$

and therefore

$$
C^{2} t^{-d / 2}\|f\|_{1}^{2} \geq\left\|e^{t L} f\right\|_{2}^{2}=\left(e^{2 t L} f, f\right)=(f, f)+\int_{0}^{t} 2\left(L e^{2 \tau L} f, f\right) d \tau
$$

Since $x e^{2 \tau x} \geq x$ for $x \leq 0$, it follows from the spectral theorem that $L e^{2 \tau L} \geq L$ for $L \leq 0$. Using this in the above equation gives the estimate

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$$
\begin{aligned}
C^{2} t^{-d / 2}\|f\|_{1}^{2} & \geq\|f\|_{2}^{2}+2 \int_{0}^{t}(L f, f) d \tau=\|f\|_{2}^{2}+2 t(L f, f) \\
& \geq\|f\|_{2}^{2}-2 t \mathcal{E}(f, f)
\end{aligned}
$$

Optimizing this inequality over $t>0$ by taking $t=\mathcal{E}(f, f)^{-2 / d+2}\|f\|_{1}^{4 / d+2}$ implies
$\|f\|_{2}^{2} \leq C^{2}\left(\mathcal{E}(f, f)^{-2 / d+2}\|f\|_{1}^{4 / d+2}\right)^{-d / 2}\|f\|_{1}^{2}+2 \mathcal{E}(f, f)^{-2 / d+2}\|f\|_{1}^{4 / d+2} \mathcal{E}(f, f)$

$$
\|f\|_{2}^{2+4 / d} \leq C \mathcal{E}(f, f)\|f\|_{1}^{4 / d}
$$

## T. Coulhon Lecture Notes

Notes from coulhon.tex
Theorem 58.1. Let $(X, \mu)$ be a measure space and $f$ be a positive measurable function. Then for $1 \leq p<\infty$,

$$
\|f\|_{p}^{p}=p \int_{0}^{\infty} \mu(f>t) t^{p-1} d t
$$

Proof. We have

$$
\begin{aligned}
\|f\|_{p}^{p} & =\int_{X} f^{p} d \mu=\int_{X}\left(\int_{0}^{f^{p}} p t^{p-1} d t\right) d \mu=\int_{X \times \mathbb{R}_{+}} 1_{f>t} p t^{p-1} d \mu d t \\
& =\int_{0}^{\infty} \mu(f>t) p t^{p-1} d t
\end{aligned}
$$

In these notes we are going to work in either one of the two following settings.

### 58.1 Weighted Riemannian Manifolds

Here we assume that $M$ is a non-compact, connected Riemannian manifold with Riemannian metric $g$ which is also equipped with a smooth measure $\mu$. We let $|\nabla f|^{2}=g(\nabla f, \nabla f)$ and

$$
\mathcal{E}(f, f):=\int_{M}|\nabla f|^{2} d \mu=\left(\Delta_{\mu} f, f\right) \text { on } L^{2}(\mu) .
$$

Here $(f, g):=\int_{M} f \bar{g} d \mu$. We have the following general important facts. The heat kernel is the smooth integral kernel for the heat operator, $e^{-t \Delta_{\mu}}$.

Definition 58.2. Given a smooth hypersurface $A \subset M$ let $|A|$ denote the surface measure, i.e.

$$
|A|=\int_{A}|\mu(N,-)|
$$

where $N$ is a normal vector to $A$.
Theorem $58.3\left(\left\|\nabla 1_{\Omega}\right\|_{1}=|\partial \Omega|_{s}\right)$. Let $\Omega \subset \subset M$ be a precompact domain with smooth boundary and let $f_{\epsilon}(x)=h_{\epsilon}\left(d_{\Omega^{c}}(x)\right)$ where $h_{\epsilon}(x)=\left(\frac{x}{\epsilon}\right) \wedge 1$. Here $d_{\Omega^{c}}(x)$ denotes the Riemannian distance of $x$ to $\Omega^{c}$. Then

$$
\lim _{\epsilon \downarrow 0}\left\|\nabla f_{\epsilon}\right\|_{1}=|\partial \Omega|_{s}
$$

which we write heuristically as

$$
\left\|\nabla 1_{\Omega}\right\|_{1}=|\partial \Omega|_{s} .
$$

Proof. We have

$$
\nabla f_{\epsilon}(x)=h_{\epsilon}^{\prime}\left(d_{\Omega^{c}}(x)\right) \nabla d_{\Omega^{c}}(x) \text { for a.e. } x
$$

and hence

$$
\left|\nabla f_{\epsilon}(x)\right|=\frac{1}{\epsilon} 1_{d_{\Omega^{c}}(x)<\epsilon} \text { for a.e. } x
$$

and therefore,

$$
\lim _{\epsilon \downarrow 0}\left\|\nabla f_{\epsilon}\right\|_{1}=\lim _{\epsilon \downarrow 0} \frac{1}{\epsilon} \mu\left(d_{\Omega^{c}}<\epsilon\right)=|\partial \Omega|
$$

Theorem 58.4 (Coarea Formula). Let $(M, g)$ be a Riemannian manifold, $f: M \rightarrow[0, \infty)$ be a reasonable function and $\mu$ be a smooth volume form on M. Then

$$
\|\nabla f\|_{L^{1}(\mu)}=\int_{0}^{\infty}|\partial\{f>t\}| d t
$$

Proof. See $[5,1,2]$ for a complete Rigorous proof. We will only give the idea here. Locally choose coordinates $x=\left(x^{1}, \ldots, x^{n}\right)$ on $M$ such that $x^{1}=f$. This is possible in neighborhood of points where $d f$ is non-zero. Then $\frac{\partial}{\partial x^{i}}$ is tangential to the level surface $\{f=t\}$ for $i=2, \ldots, n$ and therefore

$$
\begin{aligned}
\mu\left(\frac{\partial}{\partial x^{1}}, \frac{\partial}{\partial x^{2}}, \ldots, \frac{\partial}{\partial x^{n}}\right) & =\mu\left(\left(\frac{\partial}{\partial x^{1}}, \frac{\nabla f}{|\nabla f|}\right) \frac{\nabla f}{|\nabla f|}, \frac{\partial}{\partial x^{2}}, \ldots, \frac{\partial}{\partial x^{n}}\right) \\
& =\left(\frac{\partial}{\partial x^{1}}, \frac{\nabla f}{|\nabla f|}\right) \mu\left(N, \frac{\partial}{\partial x^{2}}, \ldots, \frac{\partial}{\partial x^{n}}\right) .
\end{aligned}
$$

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$$
\begin{aligned}
\left(\frac{\partial}{\partial x^{1}}, \frac{\nabla f}{|\nabla f|}\right) & =|\nabla f|^{-1}\left(\frac{\partial}{\partial x^{1}}, \nabla f\right) \\
& =|\nabla f|^{-1} \frac{\partial}{\partial x^{1}} f=|\nabla f|^{-1} \frac{\partial}{\partial x^{1}} x^{1}=|\nabla f|^{-1}
\end{aligned}
$$

and therefore,

$$
\left|\mu\left(N, \frac{\partial}{\partial x^{2}}, \ldots, \frac{\partial}{\partial x^{n}}\right)\right|=|\nabla f|\left|\mu\left(\frac{\partial}{\partial x^{1}}, \frac{\partial}{\partial x^{2}}, \ldots, \frac{\partial}{\partial x^{n}}\right)\right| .
$$

Integrating this equation with respect to $d x^{1} \ldots d x^{n}$ then gives

$$
\int\left|\mu\left(N, \frac{\partial}{\partial x^{2}}, \ldots, \frac{\partial}{\partial x^{n}}\right)\right| d x^{1} \ldots d x^{n}=\|\nabla f\|_{L^{1}(\mu)}
$$

on one hand, on the other

$$
\begin{aligned}
& \int\left|\mu\left(N, \frac{\partial}{\partial x^{2}}, \ldots, \frac{\partial}{\partial x^{n}}\right)\right| d x^{1} \ldots d x^{n} \\
&=\int_{\mathbb{R}}\left[\int_{\mathbb{R}^{n-1}}\left|\mu\left(N, \frac{\partial}{\partial x^{2}}, \ldots, \frac{\partial}{\partial x^{n}}\right)\right| d x^{2} \ldots d x^{n}\right] d x^{1} \\
& \quad=\int_{\mathbb{R}}\left[\int_{\left\{f=x^{1}\right\}}|\mu(N,-)|\right] d x^{1}=\int_{0}^{\infty}|\partial\{f>t\}| d t .
\end{aligned}
$$

### 58.2 Graph Setting

Let $\Gamma=(V, E)$ be a non-oriented graph with vertices $V$ and edges $E$. We assume the graph is connected and locally finite, in fact I think he assumes the graph is finitely ramified, i.e. there is a bound $K<\infty$ on the number of edges that are attached to any vertex. Let $d$ denote the graph distance and

$$
\mu:=\left\{\mu_{y x}=\mu_{x y} \geq 0: x, y \in V \quad \ni x y \in E\right\}
$$

be a measure on $E$. Extend $\mu$ to all pairs of points $x y$ with $x, y \in V$ by setting $\mu_{x y}=0$ if $x y \notin E$. Using this notation we let

$$
\begin{aligned}
\mu_{x} & :=\sum_{y \in V} \mu_{x y}=\sum_{y \in V: x y \in E} \mu_{x y} \\
p(x, y) & :=\frac{\mu_{x y}}{\mu_{x}}(\text { the Markov Kernel }), \\
\mu(\Omega) & :=\sum_{x \in \Omega} \mu_{x} \text { when } \Omega \subset \subset V, \\
\partial \Omega & :=\{x y \in E: x \in \Omega \text { and } y \in \partial \Omega\}, \\
|\partial \Omega|_{s} & :=\sum_{x \in \Omega, y \notin \Omega} \mu_{x y} \text { and } \\
|\nabla f(x)|^{p} & :=\sum_{y}|f(x)-f(y)|^{p} p(x, y) .
\end{aligned}
$$

In this setting $\operatorname{Pf}(x):=\sum_{y} p(x, y) f(y)$ corresponds to $e^{-t \Delta_{\mu}}$ and $P$ is a selfadjoint operator on $L^{2}(V, \mu)$. Also recall that $\Delta_{\mu}$ corresponds to $1-P$. As in the manifold case we still have

$$
\begin{align*}
\left\|\nabla 1_{\Omega}\right\|_{1} & =2|\partial \Omega|_{s} \text { and }  \tag{58.1}\\
\int_{0}^{\infty}|\partial\{f>t\}| d t & =\frac{1}{2}\|\nabla f\|_{L^{1}(\mu)} \tag{58.2}
\end{align*}
$$

as we will now verify. Using the definitions,

$$
\begin{aligned}
\left\|\nabla 1_{\Omega}\right\|_{1} & =\sum_{x} \mu_{x} \sum_{y}\left|1_{\Omega}(x)-1_{\Omega}(y)\right| p(x, y)=\sum_{x, y}\left|1_{\Omega}(x)-1_{\Omega}(y)\right| \mu_{x y} \\
& =\sum_{x \in \Omega, y \notin \Omega} \mu_{x y}+\sum_{x \notin \Omega, y \in \Omega} \mu_{x y}=2|\partial \Omega|_{s}
\end{aligned}
$$

verifying Eq. (58.1). For the graph co-area formula (58.2):

$$
\|\nabla f\|_{L^{1}(\mu)}=\sum_{x} \mu_{x} \sum_{y}|f(x)-f(y)| p(x, y)=\sum_{x, y}|f(x)-f(y)| \mu_{x y} .
$$

while
$|\partial\{f>t\}|=\sum_{x \in\{f>t\}, y \notin\{f>t\}} \mu_{x y}=\sum_{f(x)>t \text { and } f(y) \leq t} \mu_{x y}=\sum 1_{f(y) \leq t<f(x)} \mu_{x y}$
so that

$$
\begin{aligned}
\int_{0}^{\infty}|\partial\{f>t\}| d t & =\int_{0}^{\infty} \sum 1_{f(y) \leq t<f(x)} \mu_{x y} d t=\sum \int_{0}^{\infty} 1_{f(y) \leq t<f(x)} \mu_{x y} d t \\
& =\sum 1_{f(y)<f(x)}(f(x)-f(y)) \mu_{x y} \\
& =\frac{1}{2} \sum|f(x)-f(y)| \mu_{x y}=\frac{1}{2}\|\nabla f\|_{L^{1}(\mu)}
\end{aligned}
$$

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### 58.3 Basic Inequalities

We begin with a couple of very simple Sobolev inequalities. Suppose that $f \in C_{c}^{\infty}(\mathbb{R})$, then

$$
f(x)=\frac{1}{2}\left(\int_{-\infty}^{x} f^{\prime}(t) d t-\int_{x}^{\infty} f^{\prime}(t) d t\right)
$$

from which it follows that

$$
|f(x)| \leq \frac{1}{2} \int_{\mathbb{R}}\left|f^{\prime}(t)\right| d t
$$

By the Mean value inequality we have the oscillation inequality,

$$
|f(y)-f(x)| \leq|y-x|\left\|f^{\prime}\right\|_{\infty}
$$

The first inequality is dimension dependent and probes the global structure of $\mathbb{R}$ while the second inequality works in arbitrary generality and hence does not probe the global structure of $\mathbb{R}$ in any way. Let us now list a number of inequalities which are true in $\mathbb{R}^{n}$ for all $f \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$.

$$
\begin{align*}
\|f\|_{\frac{2 n}{n-2}} & \leq C\|\nabla f\|_{2} \quad \text { (Sobolev Inequality) }  \tag{58.3}\\
\|f\|_{2}^{1+2 / n} & \leq C\|f\|_{1}^{2 / n}\|\nabla f\|_{2} \quad \text { (Nash Inequality) }  \tag{58.4}\\
\|f\|_{2(1+2 / n)}^{(1+2 / n)} & \leq C\|f\|_{2}^{2 / n} \cdot\|\nabla f\|_{2} \quad \text { (Moser Inequality) } \tag{58.5}
\end{align*}
$$

$|f(x)-f(y)| \leq C_{p}|x-y|^{1-n / p}\|\nabla f\|_{p}$ for $p>n$ (oscillation inequality).

The last inequality is valid for all $f \in C^{\infty}\left(\mathbb{R}^{n}\right)$ as well.
Let $W(f)=\|\nabla f\|_{p}$ and for a positive $f$ let

$$
f_{k}:=\left(f-2^{k}\right)+\wedge 2^{k}=\left\{\begin{array}{ccc}
2^{k} & \text { if } f \geq 2^{k+1} \\
f-2^{k} & \text { if } 2^{k} \leq f \leq 2^{k+1} \\
0 & \text { if } & f \leq 2^{k}
\end{array}\right.
$$

Then

$$
\nabla f_{k}=\nabla f 1_{2^{k} \leq f \leq 2^{k+1}}
$$

and therefore

$$
\|\nabla f\|_{p}^{p}=\int_{M}|\nabla f|^{p} d \mu=\int_{M} \sum_{k}|\nabla f|^{p} 1_{2^{k} \leq f \leq 2^{k+1}} d \mu=\sum_{k} \int_{M}\left|\nabla f_{k}\right|^{p} d \mu
$$

This shows that

$$
\|\nabla f\|_{p}=\left(\sum_{k}\left\|\nabla f_{k}\right\|_{p}\right)^{1 / p}
$$

This is a key truncation property for positive Lipschitz functions $f$. As an application we have the following theorem.

Theorem 58.5. The Nash and Sobolev inequalities are equivalent.
Proof. Sobolev $\Longrightarrow$ Nash. Recall the Hölder interpolation inequality,

$$
\|f\|_{p_{\theta}} \leq\|f\|_{p_{0}}^{1-\theta}\|f\|_{p_{1}}^{\theta}
$$

where

$$
\frac{1}{p_{\theta}}=\frac{1-\theta}{p_{0}}+\frac{\theta}{p_{1}} .
$$

Taking $p_{\theta}=2, p_{0}=1$ and $p_{1}=\frac{2 n}{n-2}$ we solve for $\theta$ to find

$$
\frac{1}{2}=\frac{1-\theta}{1}+\theta \frac{n-2}{2 n}=1+\theta\left(\frac{n-2}{2 n}-1\right)=1-\theta\left(\frac{n+2}{2 n}\right)
$$

and hence

$$
\theta=\frac{n}{n+2} \text { and } 1-\theta=\frac{2}{n+2}
$$

and therefore,

$$
\|f\|_{2} \leq\|f\|_{1}^{\frac{2}{n+2}}\|f\|_{\frac{2 n}{n-2}}^{\frac{n}{n+2}}
$$

This inequality along with the Sobolev inequality (58.3) shows,

$$
\|f\|_{2}^{1+2 / n}=\|f\|_{2}^{\frac{n+2}{n}} \leq\|f\|_{1}^{\frac{2}{n}}\|f\|_{\frac{2 n}{n-2}} \leq C\|f\|_{1}^{\frac{2}{n}}\|\nabla f\|_{2}
$$

which is the Nash inequality (58.4) with the same constant.
Nash $\Longrightarrow$ Sobolev. Conversely suppose the Nash inequality (58.4) is valid. Applying Nash to $f_{k}$ we find

$$
\begin{aligned}
\left(\int f_{k}^{2}\right)^{1+2 / n} & =\left(\left\|f_{k}\right\|_{2}^{1+2 / n}\right)^{2} \\
& \leq C^{2}\left\|f_{k}\right\|_{1}^{4 / n}\left\|\nabla f_{k}\right\|_{2}^{2}=C^{2}\left\|f_{k}\right\|_{1}^{4 / n} \int_{B_{k}}|\nabla f|^{2} d \mu
\end{aligned}
$$

where $B_{k}:=\left\{2^{k} \leq f<2^{k+1}\right\}$. Combining this inequality with the following two elementary inequalities;

$$
\begin{aligned}
& \int f_{k}^{2} \geq \int_{\left\{f \geq 2^{k+1}\right\}} f_{k}^{2}=2^{2 k} \mu\left(\left\{f \geq 2^{k+1}\right\}\right) \text { and } \\
& \left\|f_{k}\right\|_{1} \leq 2^{k} \mu\left(f \geq 2^{k}\right)
\end{aligned}
$$

gives

$$
\left(2^{2 k} \mu\left(\left\{f \geq 2^{k+1}\right\}\right)\right)^{1+2 / n} \leq C^{2}\left(2^{k} \mu\left(f \geq 2^{k}\right)\right)^{4 / n} \int_{B_{k}}|\nabla f|^{2} d \mu
$$

Let $q=\frac{2 n}{n-2}$ be the exponent in the Sobolev inequality and $\nu=\frac{n}{n+2}<1$, $a_{k}:=2^{q k} \mu\left(f \geq 2^{k}\right)$ and $b_{k}:=\int_{B_{k}}|\nabla f|^{2} d \mu$. Then the above inequality says,

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$$
a_{k+1} \leq C^{\prime} b_{k}^{\nu} a_{k}^{2(1-\nu)}
$$

and summing this equation of $k \in \mathbb{Z}$ then gives

$$
\sum_{k} a_{k}=\sum_{k} a_{k+1} \leq C^{\prime} \sum_{k} b_{k}^{\nu} a_{k}^{2(1-\nu)} \leq C^{\prime}\left(\sum_{k} b_{k}\right)^{\nu}\left(\sum_{k} a_{k}^{2}\right)^{(1-\nu)}
$$

wherein we have used Holder's inequalities with the conjugate indices $1 / \nu$ and $1 /(1-\nu)$ in the last inequality. Since we are using counting measure, we also have

$$
\|a\|_{2}^{2} \leq \sum_{k}\|a\|_{1} a_{k} \leq\|a\|_{1}^{2}
$$

which combined with the previous inequality gives

$$
\begin{aligned}
\sum_{k} a_{k} \leq C^{\prime} & \left(\sum_{k} b_{k}\right)^{\nu}\left(\sum_{k} a_{k}\right)^{2(1-\nu)} \\
& =C^{\prime}\left(\int|\nabla f|^{2} d \mu\right)^{\nu}\left(\sum_{k} a_{k}\right)^{2(1-\nu)}
\end{aligned}
$$

which then shows

$$
\left(\sum_{k} a_{k}\right)^{2 \nu-1} \leq C^{\prime}\left(\int|\nabla f|^{2} d \mu\right)^{\nu}
$$

and hence

$$
\sum_{k} a_{k} \leq C^{\prime}\left(\int|\nabla f|^{2} d \mu\right)^{\frac{\nu}{2 \nu-1}}
$$

We also have

$$
\|f\|_{q}^{q}=\sum_{k} \int_{2^{k}<f \leq 2^{k+1}} f^{q} d \mu \leq \sum_{k} 2^{q(k+1)} \mu\left(f \geq 2^{k}\right)=2^{q} \sum_{k} a_{k}
$$

and it then follows that

$$
\|f\|_{q}^{q} \leq 2^{q} C^{\prime}\left(\int|\nabla f|^{2} d \mu\right)^{\frac{\nu}{2 \nu-1}}
$$

which proves the Sobolev inequality (58.3).

### 58.4 A Scale of Inequalities

In this section let $\phi: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$be an increasing function, for example $\phi(t)=$ $c t^{1 / n}$ and $\phi(t)=c \log (t)$ for $t \geq 2$.

Definition $58.6\left(S_{\phi}^{p}\right)$. Given $p \in[1, \infty]$ and $\phi$ as above, we say $S_{\phi}^{p}$ holds provided

$$
\begin{equation*}
\|f\|_{p} \leq \phi(|\Omega|)\|\nabla f\|_{p} \text { for all } \Omega \ni|\Omega|<\infty \text { and } f \in \operatorname{Lip}_{o}(\Omega) \tag{58.7}
\end{equation*}
$$

where $\operatorname{Lip}_{o}(\Omega)$ denotes those functions $f$ on $M$ or $V$ such that $f$ is Lipschitz and $\operatorname{supp}(f) \subset \bar{\Omega}$ and $f=0$ on $\partial \Omega$.

Proposition 58.7. Suppose $1 \leq p<q<\infty$ and $S_{\phi}^{p}$ holds then $S_{\frac{q}{p} \phi}^{q}$ holds.
Proof. Apply $S_{\phi}^{p}$ to the function $f^{q / p}$ to find

$$
\|f\|_{q}^{q / p}=\left\|f^{q / p}\right\|_{p} \leq \phi(|\Omega|)\left\|\nabla f^{q / p}\right\|_{p}=\phi(|\Omega|) \frac{q}{p}\left\|f^{(q / p-1)} \nabla f\right\|_{p}
$$

where we should check that the approximate chain rule holds in the graph case here. Now apply Hölder's inequality with

$$
\frac{1}{p}=\frac{1}{q}+\frac{1}{p q /(q-p)}
$$

to the last expression to find

$$
\begin{aligned}
\left\|f^{(q / p-1)} \nabla f\right\|_{p} & \leq\|\nabla f\|_{q}\left\|f^{(q / p-1)}\right\|_{p q /(q-p)}=\|\nabla f\|_{q}\left(\int\left(f^{(q-p) / p}\right)^{p q /(q-p)}\right)^{\frac{q-p}{p q}} \\
& =\|\nabla f\|_{q}\left(\int f^{q}\right)^{\frac{q-p}{p q}}=\|\nabla f\|_{q}\|f\|_{q}^{\frac{q-p}{p}}
\end{aligned}
$$

This gives

$$
\|f\|_{q}^{q / p} \leq \phi(|\Omega|) \frac{q}{p}\|\nabla f\|_{q}\|f\|_{q^{p}}^{\frac{q-p}{p}}
$$

that is to say

$$
\|f\|_{q} \leq \phi(|\Omega|) \frac{q}{p}\|\nabla f\|_{q}
$$

We now give some equivalent inequalities to $S_{\phi}^{p}$ in the following theorem.
Theorem 58.8. We have

$$
\begin{aligned}
S_{\phi}^{\infty} & \Longleftrightarrow \phi(|B(x, r)|) \geq r \Longleftrightarrow|B(x, r)| \geq \phi^{-1}(r) \\
S_{\phi}^{1} & \left.\Longleftrightarrow|\Omega| \leq \phi(|\Omega|)|\partial \Omega| \quad \text { i.e. } \frac{|\partial \Omega|}{|\Omega|} \geq \frac{1}{\phi(|\Omega|)}\right) \text { for all reasonable } \Omega \\
S_{\phi}^{2} & \Longleftrightarrow \text { to the } \phi-\text { Nash inequality }
\end{aligned}
$$

(up to constants) where the $\phi$ - Nash inequality is the inequality,

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The $\phi-$ Nash inequality is clearly equivalent to

$$
\theta\left(\|f\|_{2}^{2}\right):=\frac{\|f\|_{2}^{2}}{\phi^{2}\left(\frac{1}{\|f\|_{2}^{2}}\right)} \leq\|\nabla f\|_{2}^{2} \text { for all } f \in C_{c}^{\infty}(M) \ni\|f\|_{1}=1
$$

where $\theta(x):=x / \phi^{2}(1 / x)$.
Proof. $\left(S_{\phi}^{\infty}\right)$ Let $f(x):=\left(r-d\left(x_{0}, x\right)\right)_{+}=h_{r}\left(d\left(x_{0}, x\right)\right)$ where $h_{r}(t)=$ $\max ((r-t), 0)$. Then


Fig. 58.1. Plot of $h_{r}$ when $r=3$.

$$
|\nabla f(x)|=\left|h_{r}^{\prime}\left(d\left(x_{0}, x\right)\right)\right|\left|\nabla_{x} d\left(x_{0}, x\right)\right|=1_{d\left(x_{0}, x\right) \leq r}
$$

and $\|f\|_{\infty}=r$. So putting this function into $\left(S_{\phi}^{\infty}\right)$ then implies

$$
r=\|f\|_{\infty} \leq \phi\left(\left|B\left(x_{0}, r\right)\right|\right)\|\nabla f\|_{\infty}=\phi\left(\left|B\left(x_{0}, r\right)\right|\right) .
$$

For the converse, if $\operatorname{supp}(f) \subset \Omega$, then by the mean value theorem,

$$
\begin{equation*}
\|f\|_{\infty} \leq i n(\Omega)\|\nabla f\|_{\infty} \tag{58.9}
\end{equation*}
$$

where $\operatorname{in}(\Omega):=$ (in radius of $\Omega$ ) is the radius of the largest ball contained in $\Omega$. To prove this last equation, let $x_{0} \in \Omega$ and $y \in \partial \Omega$, then by the mean value theorem and the definition of $\operatorname{in}(\Omega)$,

$$
\left|f\left(x_{0}\right)\right|=\left|f(y)-f\left(x_{0}\right)\right| \leq d\left(x_{0}, y\right)\|\nabla f\|_{\infty} \leq i n(\Omega)\|\nabla f\|_{\infty}
$$

This proves Eq. (58.9). Now suppose $\phi(|B(x, r)|) \geq r$ holds for all $x$ and $r$, and let $B(x, r) \subset \Omega$. Then since $\phi$ is increasing,

$$
r \leq \phi(|B(x, r)|) \leq \phi(|\Omega|)
$$

and taking sups over all $B(x, r) \subset \Omega$ we learn that $\operatorname{in}(\Omega) \leq \phi(|\Omega|)$. Using this inequality in the estimate in Eq. (58.9) shows

$$
\|f\|_{\infty} \leq \phi(|\Omega|)\|\nabla f\|_{\infty},
$$

as desired.
$\left(S_{\phi}^{1}\right)$ Applying $\left(S_{\phi}^{1}\right)$ to the function $1_{\Omega}$ shows

$$
|\Omega|=\left\|1_{\Omega}\right\|_{1} \leq \phi(|\Omega|)\left\|\nabla 1_{\Omega}\right\|_{1}=\phi(|\Omega|)|\partial \Omega|_{1}
$$

as desired. For the converse, we have by the co-area formula, the above inequality and the fact that $\phi$ is increasing for any $f$ with $\operatorname{supp}(f) \subset \Omega$ and positive that

$$
\begin{aligned}
\|\nabla f\|_{L^{1}(\mu)} & =\int_{0}^{\infty}|\partial\{f>t\}| d t \geq \int_{0}^{\infty} \frac{|f>t|}{\phi(|f>t|)} d t \\
& \geq \frac{1}{\phi(|\Omega|)} \int_{0}^{\infty} \mu(f>t) d t=\frac{1}{\phi(|\Omega|)}\|f\|_{1} .
\end{aligned}
$$

$\left(S_{\phi}^{2}\right)$ Clearly $\left(S_{\phi}^{2}\right)$, i.e.

$$
\|f\|_{2} \leq \phi(|\Omega|)\|\nabla f\|_{2} \text { for all } \operatorname{supp}(f) \subset \Omega
$$

is equivalent to

$$
\lambda_{1}(\Omega)=\sup _{f: \operatorname{supp}(f) \subset \Omega} \frac{(A f, f)}{\|f\|_{2}^{2}}=\sup _{f: \operatorname{supp}(f) \subset \Omega} \frac{\|\nabla f\|_{2}^{2}}{\|f\|_{2}^{2}} \geq \frac{1}{\phi^{2}(|\Omega|)}
$$

Now suppose that $\phi$ - Nash of Eq. (58.8) holds. Since

$$
\|f\|_{1}^{2}=\left(f, 1_{\Omega}\right)^{2} \leq\|f\|_{2}^{2}\left\|1_{\Omega}\right\|_{2}^{2}=|\Omega|\|f\|_{2}^{2}
$$

so that $|\Omega| \geq \frac{\|f\|_{1}^{2}}{\|f\|_{2}^{2}}$. Since $\phi$ is increasing $\phi(|\Omega|) \geq \phi\left(\frac{\|f\|_{1}^{2}}{\|f\|_{2}^{2}}\right)$, so that $\phi-$ Nash implies

$$
\|f\|_{2} \leq \phi\left(\frac{\|f\|_{1}^{2}}{\|f\|_{2}^{2}}\right)\|\nabla f\|_{2} \leq \phi(|\Omega|)\|\nabla f\|_{2}
$$

which is $\left(S_{\phi}^{2}\right)$.
Conversely suppose $\left(S_{\phi}^{2}\right)$ holds and let $f \in C_{c}^{\infty}(M,[0, \infty))$ and $t>0(t$ to be chosen later). Then using $f<2(f-t)$ on $f>2 t$ we find

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$$
\begin{equation*}
\theta\left(\|f\|_{2}^{2}\right) \leq \operatorname{Re}(A f, f) \text { for all } f \in \mathcal{D}(A) \ni\|f\|_{1} \leq M \tag{58.10}
\end{equation*}
$$

Then $T_{t}$ is ultracontractive, i.e. $\left\|T_{t}\right\|_{1 \rightarrow \infty}<\infty$ for all $t>0$, and moreover we have

$$
\left\|T_{t}\right\|_{1 \rightarrow \infty} \leq m(t) \text { for all } t>0
$$

where $m$ satisfies

$$
t=\int_{m(t)}^{\infty} \frac{d x}{\theta(x)}
$$

Remark 58.11. This type of result appears implicitly in Nash 1958. Also see Carlen, Kusuoka and Stroock 1986 and Tmisoki in 1990.

Proof. Let $f \in L^{1}$ with $\|f\|_{1}=1$ and so by assumption $\left\|T_{t} f\right\|_{1} \leq M$ for all $t \geq 0$. Letting $I(t):=\left\|T_{t} f\right\|_{2}^{2}$ we have using Eq. (58.10) to find

$$
I^{\prime}(t)=-2 \operatorname{Re}\left(A T_{t} f, T_{t} f\right) \leq-2 \theta\left(\left\|T_{t} f\right\|_{2}^{2}\right)=-2 \theta(I(t))
$$

Thus $-\frac{I^{\prime}(t)}{\theta(I(t))} \geq 2$ and upon integration gives

$$
\begin{aligned}
\int_{I(T)}^{\infty} \frac{d x}{\theta(x))} & \geq \int_{I(T)}^{I(0)} \frac{d x}{\theta(x))}=-\int_{0}^{T} \frac{I^{\prime}(t)}{\theta(I(t))} d t \\
& \geq \int_{0}^{T} 2 d t=2 T=\int_{m(2 T)}^{\infty} \frac{d x}{\theta(x)}
\end{aligned}
$$

and therefore we have $I(T) \leq m(2 T)$ for all $T$. From this we conclude that

$$
\left\|T_{t} f\right\|_{2}^{2}=I(t) \leq m(2 t)\|f\|_{1}^{2}
$$

showing $\left\|T_{t}\right\|_{1 \rightarrow 2}^{2} \leq m(2 t)$.
We will now apply this same result to $T_{t}^{*}$ using the following comments:

1. $A^{*}=-\left.\frac{d}{d t}\right|_{0} T_{t}^{*}$ and $\operatorname{Re}\left(A^{*} f, f\right)=\operatorname{Re}(f, A f)=\operatorname{Re}(A f, f)$ we have

$$
\theta\left(\|f\|_{2}^{2}\right) \leq \operatorname{Re}(A f, f)=\operatorname{Re}\left(A^{*} f, f\right) \text { for all } f \in \mathcal{D}\left(A^{*}\right) \ni\|f\|_{1} \leq M
$$

(Actually I am little worried about domain issues here but I do not pause to worry about them now.)
2. We also have

$$
\begin{aligned}
\left\|T_{t}^{*}\right\|_{1 \rightarrow 1} & =\sup _{\|f\|_{1}=1}\left\|T_{t}^{*} f\right\|_{1}=\sup _{\|f\|_{1}=1\|g\|_{\infty}=1} \sup _{t}\left|\left(T_{t}^{*} f, g\right)\right| \\
& =\sup _{\|g\|_{\infty}=1\|f\|_{1}=1}\left|\left(f, T_{t} g\right)\right|=\sup _{\|g\|_{\infty}=1}\left\|T_{t} g\right\|_{\infty} \leq M .
\end{aligned}
$$

Using these comments we have $\left\|T_{t}^{*}\right\|_{1 \rightarrow 2}^{2} \leq m(2 t)$ and hence by duality again,

$$
\begin{aligned}
\left\|T_{t}\right\|_{2 \rightarrow \infty} & =\sup _{\|f\|_{2}=1}\left\|T_{t} f\right\|_{\infty}=\sup _{\|f\|_{2}=1} \sup _{\|g\|_{1}=1}\left|\left(T_{t} f, g\right)\right| \\
& =\sup _{\|g\|_{1}=1} \sup _{\|f\|_{2}=1}\left|\left(f, T_{t}^{*} g\right)\right| \\
& =\sup _{\|g\|_{1}=1}\left\|T_{t}^{*} g\right\|_{2}=\left\|T_{t}^{*}\right\|_{1 \rightarrow 2} \leq \sqrt{m(2 t)} .
\end{aligned}
$$

$$
\left\|T_{t}\right\|_{1 \rightarrow \infty}=\left\|T_{t / 2} T_{t .2}\right\|_{1 \rightarrow \infty} \leq\left\|T_{t / 2}\right\|_{2 \rightarrow \infty}\left\|T_{t}\right\|_{1 \rightarrow 2}
$$

$$
\leq \sqrt{m(t)} \sqrt{m(t)}=m(t)
$$

These notes are on the construction and the asymptotic expansion for heat kernels on vector bundles over compact manifolds using Levi's method. The construction described here follows closely the presentation given in Berline, Getzler, and Vergne, "Heat Kernels and Dirac Operators."

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Proof. Direct computations show that $\left(\partial_{t}-\frac{1}{2} \Delta_{x}\right) p(t, x, y)=0$, see

## Heat Equation on $\mathbb{R}^{n}$

Let $\Delta=\sum_{i=1}^{n} \partial^{2} / \partial x_{i}^{2}$ be the usual Laplacian on $\mathbb{R}^{n}$ and consider the Heat Equation

$$
\begin{equation*}
\left(\partial_{t}-\frac{1}{2} \Delta\right) u=0 \text { with } u(0, x)=f(x) \tag{59.1}
\end{equation*}
$$

where $f$ is a given function on $\mathbb{R}^{n}$. By Fourier transforming the equation in the $x$ - variables one finds that (59.1) implies that

$$
\begin{equation*}
\left(\partial_{t}+\frac{1}{2}|\xi|^{2}\right) \hat{u}(t, \xi)=0 \text { with } \hat{u}(0, \xi)=\hat{f}(\xi) . \tag{59.2}
\end{equation*}
$$

and hence that $\hat{u}(t, \xi)=e^{-t|\xi|^{2} / 2} \hat{f}(\xi)$. Inverting the Fourier transform then shows that

$$
u(t, x)=\mathcal{F}^{-1}\left(e^{-t|\xi|^{2} / 2} \hat{f}(\xi)\right)(x)=\left(\mathcal{F}^{-1}\left(e^{-t|\xi|^{2} / 2}\right) * f\right)(x)
$$

Now by well known Gaussian integral formulas one shows that

$$
\mathcal{F}^{-1}\left(e^{-t|\xi|^{2} / 2}\right)(x)=(2 \pi)^{-n} \int_{\mathbb{R}^{n}} e^{-t|\xi|^{2} / 2} e^{i \xi \cdot x} d \xi=(2 \pi t)^{-n / 2} e^{-|x|^{2} / 2 t}
$$

Let us summarize the above computations in the following Theorem.
Theorem 59.1. Let

$$
\begin{equation*}
p(t, x, y):=(2 \pi t)^{-n / 2} e^{-|x-y|^{2} / 2 t} \tag{59.3}
\end{equation*}
$$

be the heat kernel on $\mathbb{R}^{n}$. Then

$$
\begin{equation*}
\left(\partial_{t}-\frac{1}{2} \Delta_{x}\right) p(t, x, y)=0 \text { and } \lim _{t \downarrow 0} p(t, x, y)=\delta_{x}(y), \tag{59.4}
\end{equation*}
$$

where $\delta_{x}$ is the $\delta$ - function at $x$ in $\mathbb{R}^{n}$. More precisely, if $f$ is a continuous bounded (can be relaxed considerably) function on $\mathbb{R}^{n}$, then $u(t, x)=$ $\int_{\mathbb{R}^{n}} p(t, x, y) f(y) d y$ is a solution to Eq. (59.1) where $u(0, x):=\lim _{t \downarrow 0} u(t, x)$.

Proposition 63.1 and Remark 63.2 below. The main issue is to prove that $\lim _{t \downarrow 0} p(t, x, y)=\delta_{x}(y)$ or equivalently that $\lim _{t \downarrow 0} \int_{\mathbb{R}^{n}} p(t, x, y) f(y) d y=f(x)$. To show this let $p_{t}(v):=(2 \pi t)^{-n / 2} e^{-|v|^{2} / 2 t}$ and notice

$$
\begin{align*}
\left|\int_{\mathbb{R}^{n}} p(t, x, y) f(y) d y-f(x)\right| & \leq \int_{\mathbb{R}^{n}} p(t, x, y)|f(y)-f(x)| d y \\
& =\int_{\mathbb{R}^{n}} p_{t}(v)|f(x+v)-f(x)| d y \tag{59.5}
\end{align*}
$$

Now for a bounded function $g$ on $\mathbb{R}^{n}$ we have that

$$
\begin{align*}
\int_{\mathbb{R}^{n}}|g(v)| p_{t}(v) d v & =\int_{B(\delta)}|g(v)| p_{t}(v) d v+\int_{B(\delta)^{c}}|g(v)| p_{t}(v) d v \\
& \leq \sup _{v \in B(\delta)}|g(v)|+\|g\|_{0} \int_{B(\delta)^{c}} p_{t}(v) d v \\
& \leq \sup _{v \in B(\delta)}|g(v)|+C\|g\|_{0} e^{-\delta^{2} / 4 t}, \tag{59.6}
\end{align*}
$$

where $\|g\|_{0}$ denotes the supremum norm of $g$. Applying this estimate to Eq. (59.5) implies

$$
\left|\int_{\mathbb{R}^{n}} p(t, x, y) f(y) d y-f(x)\right| \leq \sup _{v \in B(\delta)}|f(x+v)-f(x)|+C\|f\|_{0} e^{-\delta^{2} / 4 t} .
$$

Therefore if $K$ is a compact subset of $\mathbb{R}^{n}$, then

$$
\begin{aligned}
& \overline{\lim }_{t \downarrow 0} \sup _{x \in K}\left|\int_{\mathbb{R}^{n}} p(t, x, y) f(y) d y-f(x)\right| \\
& \quad \leq \sup _{v \in B(\delta)} \sup _{x \in K}|f(x+v)-f(x)| \rightarrow 0 \text { as } \delta \rightarrow 0
\end{aligned}
$$

by uniform continuity. This shows that $\lim _{t \downarrow 0} u(t, x)=f(x)$ uniformly on compact subsets of $\mathbb{R}^{n}$.

Notation 59.2 We will write $\left(e^{t \Delta / 2} f\right)(x)$ for $\int_{\mathbb{R}^{n}} p(t, x, y) f(y) d y$.

## An Abstract Version of E. Levi's Argument

The idea for the construction of the heat kernel for more general heat equations will be based on a method due to E. Levi. Let us illustrate the method with the following finite dimensional analogue. Suppose that $L$ is a linear operator on a finite dimensional vector space $V$ and let $P_{t}:=e^{t L}$, i.e. $P_{t}$ is the unique solution to the ordinary differential equation

$$
\begin{equation*}
\frac{d}{d t} P_{t}=L P_{t} \text { with } P_{0}=I . \tag{60.1}
\end{equation*}
$$

In this finite dimensional setting it is very easy to solve Eq. (60.1), namely one may take

$$
P_{t}=\sum_{k=0}^{\infty} \frac{t^{k}}{k!} L^{k} .
$$

Such a series solution will in general not converge when $L$ is an unbounded operator on an infinite dimensional space as are differential operators. On the other hand for the heat equation we can find quite good parametrix (approximate solution) to Eq. (60.1). Let us model this by a map $t \in \mathbb{R}_{+} \rightarrow K_{t} \in$ $\operatorname{End}(V)$ such that $K_{0}=I$ and

$$
\begin{equation*}
\frac{d}{d t} K_{t}-L K_{t}=-R_{t} \tag{60.2}
\end{equation*}
$$

where $\left\|R_{t}\right\|=O\left(t^{\alpha}\right)$ for some $\alpha>-1$. Using du Hamell's principle (or variation of parameters if you like) we see that $K_{t}$ is given by

$$
\begin{equation*}
K_{t}=P_{t}-\int_{0}^{t} P_{t-s} R_{s} d s=P_{t}-(Q P)_{t} \tag{60.3}
\end{equation*}
$$

where

$$
\begin{equation*}
(Q f)_{t}:=\int_{0}^{t} f_{t-s} R_{s} d s=\int_{0}^{t} f_{s} R_{t-s} d s \tag{60.4}
\end{equation*}
$$

We may rewrite Eq. (60.3) as $K=(I-Q) P$ and hence we should have that

$$
\begin{equation*}
P=(I-Q)^{-1} K=\sum_{m=0}^{\infty} Q^{m} K \tag{60.5}
\end{equation*}
$$

Now a simple change of variables shows that $(Q f)_{t-s}=\int_{s}^{t} f_{t-r} R_{r-s} d r$ and by induction one shows that

$$
\begin{equation*}
\left(Q^{m} f\right)_{t}=\int_{t \Delta_{m}} f_{t-s_{m}} R_{s_{m}-s_{m-1}} \ldots R_{s_{2}-s_{1}} R_{s_{1}} d \mathbf{s} \tag{60.6}
\end{equation*}
$$

where

$$
\begin{equation*}
t \Delta_{m}:=\left\{\mathbf{s}=\left(s_{1}, s_{2}, \ldots, s_{m}\right): 0 \leq s_{1} \leq s_{2} \cdots \leq s_{m} \leq t\right\} \tag{60.7}
\end{equation*}
$$

and $d \mathbf{s}=d s_{1} d s_{2} \ldots d s_{m}$. Alternatively one also shows that

$$
\begin{equation*}
\left(Q^{m} K\right)_{t}=\int_{t \Delta_{m}} K_{s_{1}} R_{s_{2}-s_{1}} R_{s_{3}-s_{2}} \ldots R_{s_{m}-s_{m-1}} R_{t-s_{m}} d \mathbf{s} \tag{60.8}
\end{equation*}
$$

Equation (60.6) implies that

$$
\begin{equation*}
\left(Q^{m} f\right)_{t}=\int_{0}^{t} f_{t-s}\left(Q^{m-1} R\right)_{s} d s \tag{60.9}
\end{equation*}
$$

Using this result, we may write Eq. (60.5) as

$$
\begin{equation*}
P_{t}=K_{t}+\int_{0}^{t} K_{t-s} V_{s} d s=K_{t}+\int_{0}^{t} K_{s} V_{t-s} d s \tag{60.10}
\end{equation*}
$$

where

$$
\begin{equation*}
V_{t}=\sum_{m=0}^{\infty}\left(Q^{m} R\right)_{t}=R_{t}+\sum_{m=1}^{\infty} \int_{t \Delta_{m}} R_{t-s_{m}} R_{s_{m}-s_{m-1}} \ldots R_{s_{2}-s_{1}} R_{s_{1}} d \mathbf{s} \tag{60.11}
\end{equation*}
$$

Let us summarize these results in the following proposition.
Proposition 60.1. Let $\alpha \geq 0, K, P, R=L K-\dot{K}, Q$ and $V$ be as above. Then the series in Eq. (60.5) and Eq. (60.11) are convergent and Eq. (60.10) holds, where $P_{t}=e^{t L}$ is the unique solution to Eq. (60.1). Moreover,

$$
\begin{equation*}
\|P-K\|_{t} \leq \frac{C}{\alpha+1} e^{C t^{\alpha+1}}\|K\|_{t} t^{\alpha+1}=O\left(t^{1+\alpha}\right) \tag{60.12}
\end{equation*}
$$

where $\|f\|_{t}:=\max _{0 \leq s \leq t}\left|f_{s}\right|$.
Remark 60.2. In the finite dimensional case or where $L$ is a bounded operator, we may take $K=I$ in the previous proposition. Then $R_{s}=L$ is constant independent of $s$ and

$$
V_{t}=\sum_{m=0}^{\infty} \frac{t^{m}}{m!} L^{m+1}
$$

which used in Eq. (60.10) gives the standard formula:

$$
P_{t}=I+\int_{0}^{t} \sum_{m=0}^{\infty} \frac{(t-s)^{m}}{m!} L^{m+1} d s=I+\sum_{m=0}^{\infty} \frac{t^{m+1}}{(m+1)!} L^{m+1}=e^{t L}
$$

Proof. From Eq. (60.6),

$$
\begin{aligned}
\left|\left(Q^{m} R\right)_{t}\right| & =\left|\int_{t \Delta_{m}} R_{t-s_{m}} R_{s_{m}-s_{m-1}} \ldots R_{s_{2}-s_{1}} R_{s_{1}} d \mathbf{s}\right| \\
& \leq(C t)^{(m+1) \alpha} \int_{t \Delta_{m}} d \mathbf{s}=\left(C t^{\alpha}\right)^{m+1} \frac{t^{m}}{m!}
\end{aligned}
$$

Therefore the series in Eq. (60.11) is absolutely convergent and

$$
\left|V_{t}\right| \leq C t^{\alpha} \sum_{m=0}^{\infty} \frac{1}{m!}\left(C t^{\alpha+1}\right)^{m}=C e^{C t^{\alpha+1}} t^{\alpha}
$$

Using this bound on $V$ and the uniform boundedness of $K_{t}$,

$$
\begin{equation*}
\int_{0}^{t}\left|K_{s} V_{t-s}\right| d s \leq C e^{C t^{\alpha+1}}\|K\|_{t} \int_{0}^{t}(t-s)^{\alpha} d s=\frac{C}{\alpha+1} e^{C t^{\alpha+1}}\|K\|_{t} t^{\alpha+1} \tag{60.13}
\end{equation*}
$$

and hence $P_{t}$ defined in Eq. (60.10) is well defined and is continuous in $t$. Moreover, (60.13) implies Eq. (60.12) once we shows that $P_{t}=e^{t L}$. This is checked as follows,

$$
\begin{aligned}
\frac{d}{d t} \int_{0}^{t} K_{t-s} V_{s} d s & =V_{t}+\int_{0}^{t} \dot{K}_{t-s} V_{s} d s=V_{t}+\int_{0}^{t}\left(L K_{t-s}-R_{t-s}\right) V_{s} d s \\
& =V_{t}+L \int_{0}^{t} K_{t-s} V_{s} d s-(Q V)_{t}=L \int_{0}^{t} K_{t-s} V_{s} d s+R_{t}
\end{aligned}
$$

Thus we have,

$$
\begin{aligned}
\frac{d}{d t} P_{t} & =\dot{K}_{t}+L \int_{0}^{t} K_{t-s} V_{s} d s+R_{t} \\
& =L K_{t}+L \int_{0}^{t} K_{t-s} V_{s} d s=L P_{t}
\end{aligned}
$$

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## Statement of the Main Results

Let $M$ be a compact Riemannian Manifold of dimension $n$ and $\Delta$ denote the Laplacian on $C^{\infty}(M)$. We again wish to solve the heat equation (59.1). It is natural to define a kernel $\rho(t, x, y)$ in analogy with the formula for $p(t, x, y)$ in Eq. (59.3), namely let

$$
\begin{equation*}
\rho(t, x, y):=(2 \pi t)^{-n / 2} e^{-d^{2}(x, y) / 2 t} \tag{61.1}
\end{equation*}
$$

where $d(x, y)$ is the Riemannian distance between two point $x, y \in M$. We may then define the operator $T_{t}$ on $C^{\infty}(M)$ by

$$
\begin{equation*}
T_{t} f(x)=\int_{M} \rho(t, x, y) f(y) d \lambda(y) \tag{61.2}
\end{equation*}
$$

where $\lambda$ is the volume measure on $M$. Although, $\lim _{t \downarrow 0} T_{t} f=f$, it is not the case that $u(t, x)=T_{t} f(x)$ is a solution to the heat equation on $M$. This is because $\rho$ does not satisfy the Heat equation. Nevertheless, $\rho$ is an approximate solution as will be seen Proposition 63.1 below. Moreover, $\rho$ will play a crucial role in constructing the true heat kernel $p(t, x, y)$ for $M$. Let us now summarize the main theorems to be proved.

### 61.1 The General Setup: the Heat Eq. for a Vector Bundle

Let $\pi: E \rightarrow M$ be a Vector bundle with connection $\nabla^{E}$. We will usually denote the covariant derivatives on $T M$ and $E$ all by $\nabla$. For a section $S \in$ $\Gamma(E)$, let $\square S:=\operatorname{tr}\left(\nabla^{T^{*} M \otimes E} \nabla^{E} S\right)$ be the rough or Bochner Laplacian on $E$ and let

$$
L:=\frac{1}{2} \square+\mathcal{R},
$$

where $\mathcal{R}$ is a section of $\operatorname{End}(E)$. We are interested in solving

$$
\begin{equation*}
\left(\partial_{t}-L\right) u=0 \text { with } u(0, x)=f(x) \tag{61.3}
\end{equation*}
$$

where $u(t, \cdot)$ and $f(\cdot)$ are section of $\Gamma(E)$.

### 61.2 The Jacobian ( $J$ - function)

Definition 61.1. Let $D: T M \rightarrow \mathbb{R}$ be defined so that for each $y \in M$, $\exp _{y}^{*} \lambda=D \lambda_{y}$ where $\lambda=\lambda_{M}$ is the volume form on $M$ and $\lambda_{y}$ is the volume form on $T_{y} M$. More explicitly, if $\left\{e_{i}\right\}_{i=1}^{n}$ is an oriented orthonormal basis for $T_{y} M$, then and $v \in T_{y} M$, then

$$
\begin{align*}
D(v) & =\lambda\left(\exp _{*}\left(e_{1}\right)_{v}, \exp _{*}\left(e_{2}\right)_{v} \ldots, \exp _{*}\left(e_{n}\right)_{v}\right) \\
& =\sqrt{\operatorname{det}\left\{\left(\exp _{y *}\left(e_{i}\right)_{v}, \exp _{y *}\left(e_{j}\right)_{v}\right)\right\}_{i, j=1}^{n}}, \tag{61.4}
\end{align*}
$$

where $\exp _{*} w_{v}=\left.\frac{d}{d t}\right|_{0} \exp (v+t w)$. Further define $J(x, y)=D\left(\exp _{y}^{-1}(x)\right)$.
Notice that $J(\cdot, y)$ satisfies

$$
\exp _{y}^{*}\left(\frac{1}{J(\cdot, y)} \lambda\right)=\lambda_{y}
$$

Alternatively we have that $J(x, y)=\operatorname{det}\left(\exp _{y *}(\cdot)_{v}\right)$ where $v=\exp _{y}^{-1}(x)$. To be more explicit, let $\left\{e_{i}\right\}_{i=1}^{n}$ be an orthonormal basis for $T_{y} M$, then

$$
J(x, y)=\sqrt{\operatorname{det}\left\{\left(\exp _{y *}\left(e_{i}\right)_{v}, \exp _{y *}\left(e_{j}\right)_{v}\right)\right\}_{i, j=1}^{n}}
$$

Remark 61.2 (Symmetry of $J$ ). It is interesting to notice that $J$ is a symmetric function. We will not need this fact below so the proof may be skipped. We will also be able to deduce the symmetry of $J$ using the asymptotic expansion of the heat kernel along with the symmetry of the heat kernel.

Proof. Let $w_{y} \in T_{y} M$, then

$$
w_{y} d^{2}(x, \cdot)=2\left(V(x, y), w_{y}\right)=-2\left(\exp _{y}^{-1}(x), w_{y}\right)
$$

where $V(x, y):=\left.\frac{d}{d t}\right|_{t=1} \exp \left(t \exp _{x}^{-1}(y)\right)=-\exp _{y}^{-1}(x)$. Thus if $u_{x} \in T_{x} M$, then

$$
u_{x} w_{y} d^{2}(x, y)=-2\left(\left(\exp _{y}^{-1}\right)_{*} u_{x}, w_{y}\right)
$$

Now

$$
J(x, y)=\operatorname{det}\left[\left(\exp _{y}^{-1}\right)_{*}\right]=\operatorname{det}\left[\left\{-\frac{1}{2} u_{i} w_{j} d^{2}(x, y)\right\}_{i, j=1}^{n}\right]
$$

where $\left\{u_{i}\right\}$ and $\left\{w_{j}\right\}$ is an orthonormal basis of $T_{x} M$ and $T_{y} M$ respectively. From this last formula it is clear from the fact that $d(x, y)=d(y, x)$ that $J(x, y)=J(y, x)$.
Lemma 61.3 (Expansion of $D$ ). The function $J$ is symmetric, $J(x, y)=$ $J(y, x)$. Moreover

$$
D(v)=1-\frac{1}{6}(\operatorname{Ric} v, v)+O\left(v^{3}\right)
$$

and hence

$$
J(x, y)=\left(\operatorname{Ric} \exp _{y}^{-1}(x), \exp _{y}^{-1}(x)\right)+O\left(d^{3}(x, y)\right)
$$

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The proof of this result will be given in the Appendix below since the result is not really needed for our purposes.

### 61.3 The Approximate Heat Kernels

Theorem 61.4 (Approximate Heat Kernel). For $x, y \in M$, let

$$
\gamma_{y, x}(t):=\exp _{y}\left(t \exp _{y}^{-1}(x)\right)
$$

so that $\gamma_{x, y}$ is the geodesic connecting $y$ to $x$. Also let $/ /_{t}\left(\gamma_{x, y}\right)$ denote parallel translation along $\gamma_{x, y}$. (Read $\gamma_{x, y}$ as $\gamma_{x \leftarrow y .}$.) Define, for $(x, y)$ near the diagonal $\Delta \subset M \times M$ and $k=0,1,2 \ldots, u_{k}(x, y): E_{y} \rightarrow E_{x}$ inductively by
$u_{k+1}(x, y)=u_{0}(x, y) \int_{0}^{1} s^{k} u_{0}\left(\gamma_{x, y}(s), y\right)^{-1}\left(L_{x} u_{k}\right)\left(\gamma_{x, y}(s), y\right) d s$ for $k=0,1,2 \ldots$ and

$$
u_{0}(x, y)=\frac{1}{\sqrt{J(x, y)}} / / 1\left(\gamma_{x, y}\right)
$$

For $q=0,1,2 \ldots$, let

$$
\begin{equation*}
\Sigma_{q}(t, x, y):=\rho(t, x, y) \sum_{k=0}^{q} t^{k} u_{k}(x, y) \tag{61.7}
\end{equation*}
$$

where

$$
\begin{equation*}
\rho(t, x, y):=(2 \pi t)^{-n / 2} e^{-d^{2}(x, y) / 2 t} . \tag{61.8}
\end{equation*}
$$

Then

$$
\begin{equation*}
\left(\partial_{t}-L_{x}\right) \Sigma_{q}(t, x, y)=-t^{q} \rho(t, x, y) L_{x} u_{q}(x, y) \tag{61.9}
\end{equation*}
$$

Definition 61.5 (Cut off function). Let $\epsilon>0$ be less than the injectivity radius of $M$ and choose $\Psi \in C_{c}^{\infty}\left(-\epsilon^{2}, \epsilon^{2}\right)$ such that $0 \leq \Psi \leq 1$ and $\Psi$ is 1 in a neighborhood of 0 . Set $\psi(x, y)=\Psi\left(d^{2}(x, y)\right)$, a cutoff function which is one in a neighborhood of the diagonal and such that $\psi(x, y)=0$ if $d(x, y) \geq \epsilon$.

Corollary 61.6 (Approximate Heat Kernel). Let $k_{q}(t, x, y):=\psi(x, y) \Sigma_{q}(t, x$ Define

$$
r_{q}(t, x, y):=-\left(\partial_{t}-L_{x}\right) k_{q}(t, x, y)
$$

then

$$
\begin{equation*}
\left\|\partial_{t}^{k} r_{q}(t, \cdot, \cdot)\right\|_{l} \leq C t^{q-n / 2-l / 2-k} \tag{61.10}
\end{equation*}
$$

where $\|r\|_{l}$ denotes the supremum norm of $r$ and all of its derivatives up to order $l$.

### 61.4 The Heat Kernel and its Asymptotic Expansion

Theorem 61.7 (Existence of Heat Kernels). There exists a heat kernel $p(t, x, y): E_{y} \rightarrow E_{x}$ for $L$, i.e. $p$ is a $C^{1}$-function in $t$ and $C^{2}$ in $(x, y)$ such that

$$
\begin{equation*}
\left(\partial_{t}-L_{x}\right) p(t, x, y)=0 \text { with } \lim _{t \downarrow 0} p(t, x, y)=\delta_{x}(y) i d \tag{61.11}
\end{equation*}
$$

Remark 61.8. The explicit formula for $p(t, x, y)$ is derived by a formal application of equations (60.10) and (60.11) above. The results are:

$$
\begin{equation*}
p(t, x, y)=k_{q}(t, x, y)+\int_{0}^{t} \int_{M} k_{q}(t-s, x, z) v_{q}(s, z, y) d \lambda(z) d s \tag{61.12}
\end{equation*}
$$

where

$$
\begin{equation*}
v_{q}(s, x, y)=\sum_{m=1}^{\infty} r^{m}(s, x, y) \tag{61.13}
\end{equation*}
$$

and the kernels $r^{m}$ are defined inductively by where $r^{1}=r$ and for $m \geq 2$

$$
\begin{equation*}
r^{m}(s, x, y)=\int_{0}^{s} \int_{M} r(s-r, x, z) r^{m-1}(r, z, y) d \lambda(z) d r \tag{61.14}
\end{equation*}
$$

Corollary 61.9 (Uniqueness of Heat Kernels). The heat equation (61.3) has a unique solution. Moreover, there is exactly one solution to (61.11).

Proof. Let $u(t, x):=\int_{M} p(t, x, y) f(y) d \lambda(y)$, then $u$ solves the heat equation (61.3). We will prove uniqueness of $u$ using the existence of the adjoint problem. In order to carry this out we will need to know that $L^{t}=\frac{1}{2} \square_{E^{*}}+\mathcal{R}^{t}$ : $\Gamma\left(E^{*}\right) \rightarrow \Gamma\left(E^{*}\right)$ is the formal transpose of $L$ in the sense that

$$
\begin{equation*}
\int_{M}\langle L f, g\rangle d \lambda=\int_{M}\left\langle f, L^{t} g\right\rangle d \lambda \tag{61.15}
\end{equation*}
$$

for all sections $f \in \Gamma(E)$ and $g \in \Gamma\left(E^{*}\right)$. Here $\square_{E^{*}}$ is the rough Laplacian on $E^{*}$. Indeed, let $X$ be the vector field on $M$ such that $(X, v)=\left\langle\nabla_{v} f, g\right\rangle$ for all $v \in T_{m} M$. Then

$$
\left(\nabla_{v} X, \cdot\right)=\nabla_{v}\langle\nabla \cdot f, g\rangle=\left\langle\nabla_{v}(\nabla \cdot f), g\right\rangle+\left\langle\nabla \cdot f, \nabla_{v} g\right\rangle
$$

so that in particular

$$
\left(\nabla_{v} X, v\right)=\left\langle\nabla_{v \otimes v}^{2} f, g\right\rangle+\left\langle\nabla_{v} f, \nabla_{v} g\right\rangle
$$

Let $v=e_{i}$, where $\left\{e_{i}\right\}$ an orthonormal frame, and sum this equation on $i$ to find that

$$
\begin{equation*}
\div(X)=\langle\square f, g\rangle+\langle\nabla f, \nabla g\rangle . \tag{61.16}
\end{equation*}
$$

Using the Riemannian metric on $M$ to identify one forms with vector fields, we may write this equality as:
61.4 The Heat Kernel and its Asymptotic Expansion 1153

$$
\div\langle\nabla \cdot f, g\rangle=\langle\square f, g\rangle+\langle\nabla f, \nabla g\rangle .
$$

Similarly one shows that

$$
\div\langle f, \nabla \cdot g\rangle=\langle f, \square g\rangle+\langle\nabla f, \nabla g\rangle,
$$

which subtracted from the previous equation gives "Green's identity,"

$$
\begin{equation*}
\langle f, \square g\rangle-\langle\square f, g\rangle=\div(\langle f, \nabla \cdot g\rangle-\langle\nabla \cdot f, g\rangle) \tag{61.17}
\end{equation*}
$$

Integrating equations (61.16) and (61.17) over $M$, we find, using the divergence theorem, that

$$
\int_{M}\langle\square f, g\rangle d \lambda=\int_{M}\langle f, \square g\rangle d \lambda=-\int_{M}\langle\nabla f, \nabla g\rangle d \lambda
$$

Thus

$$
\begin{aligned}
\int_{M}\langle L f, g\rangle d \lambda & =\frac{1}{2} \int_{M}\langle\square f, g\rangle d \lambda+\int_{M}\langle\mathcal{R} f, g\rangle d \lambda \\
& =\frac{1}{2} \int_{M}\langle f, \square g\rangle d \lambda+\int_{M}\left\langle f, \mathcal{R}^{t} g\right\rangle d \lambda=\int_{M}\left\langle f, L^{t} g\right\rangle d \lambda
\end{aligned}
$$

proving Eq. (61.15).
Suppose that $u$ is a solution to Eq. (61.3) with $u(0, x)=0$ for all $x$. By applying Theorem 61.7, we can construct a heat kernel $q_{t}$ for $L^{t}$. Given $g \in \Gamma\left(E^{*}\right)$, let

$$
v(t, x):=\int_{M} q_{T-t}(x, y) g(y) d \lambda(y)
$$

for $t<T$. Now consider

$$
\begin{aligned}
& \frac{d}{d t} \int_{M}\langle u(t, x), v(t, x)\rangle d \lambda(x) \\
& \quad=\int_{M}\langle L u(t, x), v(t, x)\rangle d \lambda(x)-\int_{M}\left\langle u(t, x), L^{t} v(t, x)\right\rangle d \lambda(x)=0
\end{aligned}
$$

and therefore, $\int_{M}\langle u(t, x), v(t, x)\rangle d \lambda(x)$ is constant in $t$. Considering this expression in the limit that $t$ tends to 0 and $T$ implies that

$$
0=\int_{M}\langle 0, v(0, x)\rangle d \lambda(x)=\int_{M}\langle u(T, x), g(x)\rangle d \lambda(x)
$$

Since $g$ is arbitrary, this implies that $u(T, x)=0$ for all $x$. Hence the solution to equation (61.3) is unique. It is now easy to use this result to show that $p(t, x, y)$ must be unique as well.
Theorem 61.10 (Assymptotics of the Heat Kernel). Let $p(t, x, y)$ be the heat kernel described by Eq. (61.11), then $p$ is smooth in $(t, x, y)$ for $t>0$. Moreover if $K^{q}$ is as in Corollary 61.6, then

$$
\begin{equation*}
\left\|\partial_{t}^{k}\left(p(t, \cdot, \cdot)-K^{q}(t, \cdot, \cdot)\right)\right\|_{l}=O\left(t^{q-n / 2-l / 2-k}\right) \tag{61.18}
\end{equation*}
$$

provided that $q>n / 2+l / 2+k$.
and from this it follows that $\sum_{m=1}^{\infty}\left\|\partial_{s}^{k} r^{m}(s, \cdot, \cdot)\right\|_{l}=O\left(s^{(q-n / 2)-l / 2-k}\right)$.

## Proof of Theorems 61.7 and 61.10

In this section we will give the proof of Theorems 61.7 and 61.10 assuming Theorem 61.4 and Corollary 61.6.

### 62.1 Proof of Theorem 61.7

Let $l$ and $k$ be given and fix $q>n / 2+l / 2+k .$. Let $k(t, x, y)=k_{q}(t, x, y)$ and $r(t, x, y)=r_{q}(t, x, y)$ as in Corollary 61.6. Let

$$
K_{t} f(x):=\int_{M} k(t, x, y) f(y) d \lambda(y) \text { and } R_{t} f(x):=\int_{M} r(t, x, y) f(y) d \lambda(y)
$$

Following the strategy described in Section 60, we will let $v(t, x, y)$ be the kernel of the operator $\sum_{m=0}^{\infty} Q^{m} R$, where $Q$ is as in Eq. (60.4). That is

$$
\begin{equation*}
v(s, x, y)=\sum_{m=1}^{\infty} r^{m}(s, x, y) \tag{62.1}
\end{equation*}
$$

where $r^{1}=r$ and for $m \geq 2$

$$
r^{m}(s, x, y)
$$

$$
\begin{align*}
& =\int_{s \Delta_{m-1}} \int_{M^{m-1}} r\left(s-s_{m-1}, x, y_{m-1}\right) r\left(s_{m-1}-s_{m-2}, y_{m-1}, y_{m-2}\right) \ldots r\left(s_{1}\right. \\
& =\int_{0}^{s} \int_{M} r(s-r, x, z) r^{m-1}(r, z, y) d \lambda(z) d r \tag{62.2}
\end{align*}
$$

and $d \mathbf{y}=d \lambda\left(y_{1}\right) \ldots d \lambda\left(y_{m}\right)$. The kernel $r^{m}$ is easy to estimate using (61.10) to find that

$$
\begin{align*}
\left\|\partial_{s}^{k} r^{m}(s, \cdot, \cdot)\right\|_{l} & \leq\left(C s^{q-n / 2}\right)^{m} s^{-l / 2-k} \operatorname{Vol}(M)^{m-1} s^{m-1} /(m-1)! \\
& =C^{m} s^{m(q-n / 2)-l / 2-k} \operatorname{Vol}(M)^{m-1} s^{m-1} /(m-1)! \tag{62.3}
\end{align*}
$$

## Therefore, $v$ is well defined with $\partial_{s}^{k} v(s, x, y)$ in $\Gamma_{l}$.

Proof. Let $p(t, x, y)$ be the kernel of the operator $P_{t}$ defined in Eq. (60.10) i.e.

$$
\begin{align*}
p(t, x, y) & =k(t, x, y)+\int_{0}^{t} \int_{M} k(t-s, x, z) v(s, z, y) d \lambda(z) d s  \tag{62.4}\\
& =k(t, x, y)+\int_{0}^{t} \int_{M} k(s, x, z) v(t-s, z, y) d \lambda(z) d s \tag{62.5}
\end{align*}
$$

Using Eq. (62.5), we find (since $\left.\partial_{t} v(0, z, y)=0\right)$ that

$$
\begin{align*}
\partial_{t} p(t, x, y) & =\partial_{t} k(t, x, y)+\int_{0}^{t} \int_{M} k(s, x, z) \partial_{t} v(t-s, z, y) d \lambda(z) d s \\
& =L_{x} k(t, x, y)-r(t, x, y) \\
& +\int_{0}^{t} \int_{M} k(s, x, z) \partial_{t} v(t-s, z, y) d \lambda(z) d s \tag{62.6}
\end{align*}
$$

More generally,

$$
\partial_{t}^{i} p(t, x, y)=\partial_{t}^{i} k(t, x, y)+\int_{0}^{t} \int_{M} k(s, x, z) \partial_{t}^{i} v(t-s, z, y) d \lambda(z) d s
$$

from which it follows that $\partial_{t}^{i} p$ is continuous in $\Gamma_{l}$ for all $i \leq k$. Furthermore,

$$
\begin{align*}
\left\|\partial_{t}^{i} p(t, \cdot, \cdot)-\partial_{t}^{i} k(t, \cdot, \cdot)\right\|_{l} & \leq C \int_{0}^{t}\left\|\partial_{t}^{i} v(t-s, \cdot, \cdot)\right\|_{l} d s \\
& \leq C \int_{0}^{t}(t-s)^{q-n / 2-l / 2-i} d s \\
& =O\left(t^{q-n / 2-l / 2-i+1}\right) \tag{62.7}
\end{align*}
$$

To finish the proof of Theorem 61.7, we need only verify Eq. (61.11). The assertion that $\lim _{t \downarrow 0} p(t, x, y)=\delta_{x}(y) i d$ follows from the previous estimate and the analogous property of $k(t, x, y)$. Fubini's theorem and integration by parts shows that

$$
\begin{aligned}
& \int_{\epsilon}^{t} \int_{M} k(s, x, z) \partial_{t} v(t-s, z, y) d \lambda(z) d s \\
&=-\left.\int_{M} k(s, x, z) v(t-s, z, y) d \lambda(z)\right|_{s=\epsilon} ^{s=t} \\
&+\int_{\epsilon}^{t} \int_{M} \partial_{s} k(s, x, z) v(t-s, z, y) d \lambda(z) d s \\
&=\int_{M} k(\epsilon, x, z) v(t-\epsilon, z, y) d \lambda(z) \\
&+\int_{\epsilon}^{t} \int_{M}\left(L_{x} k(s, x, z)-r(s, x, z)\right) v(t-s, z, y) d \lambda(z) d s \\
&=\int_{M} k(\epsilon, x, z) v(t-\epsilon, z, y) d \lambda(z) \\
&+L_{x} \int_{\epsilon}^{t} \int_{M} k(s, x, z) v(t-s, z, y) d \lambda(z) d s \\
&-\int_{\epsilon}^{t} \int_{M} r(s, x, z) v(t-s, z, y) d \lambda(z) d s
\end{aligned}
$$

Making use of the fact that $K$ is uniformly bounded on $\Gamma_{l}$ and that the strong$\lim _{\epsilon \downarrow 0} K_{\epsilon}=I$, we may pass to the limit, $\epsilon \rightarrow 0$, in this last equality to find that

$$
\begin{align*}
\int_{0}^{t} \int_{M} k(s, x, z) & \partial_{t} v(t-s, z, y) d \lambda(z) d s \\
& =v(t, x, y)+L_{x} \int_{0}^{t} \int_{M} k(s, x, z) v(t-s, z, y) d \lambda(z) d s \\
& -\int_{0}^{t} \int_{M} r(s, x, z) v(t-s, z, y) d \lambda(z) d s \\
& =r(t, x, y)+L_{x} \int_{0}^{t} \int_{M} k(s, x, z) v(t-s, z, y) d \lambda(z) d s \tag{62.8}
\end{align*}
$$

wherein the last equality we have made use of equations (62.1) and (62.2) to conclude that

$$
v(t, x, y)-\int_{0}^{t} \int_{M} r(s, x, z) v(t-s, z, y) d \lambda(z) d s=r(t, x, y)
$$

Combining (62.6) and (62.8) implies that $\left(\partial_{t}-L_{x}\right) p(t, x, y)=0$.

### 62.2 Proof of Theorem 61.10

Because $q$ in the above proof was arbitrary, we may construct a kernel $p(t, x, y)$ as in the previous section which is $C^{N}$ for any $N$ we desire. By the uniqueness

## Properties of $\rho$

For the time being let $y \in M$ be a fixed point and $r(x):=d(x, y)$. Also let $v(x):=\exp _{y}^{-1}(x), \gamma_{x}(t):=\exp (t v(x))$ and $V$ be the "radial" vector field,

$$
\begin{equation*}
V(x)=V_{y}(x):=\left.\frac{d}{d t}\right|_{1} \exp (t v(x))=/ / 1_{1}\left(\gamma_{x}\right) v(x) \tag{63.1}
\end{equation*}
$$

where $/ /_{t}\left(\gamma_{x}\right)$ is used to denote parallel translation along $\gamma_{x}$ up to time $t$. Notice that $V$ is a smooth vector field on a neighborhood of $y$. To simplify notation we will write $\rho(t, x)$ for $\rho(t, x, y)$, i.e.

$$
\begin{equation*}
\rho(t, x)=(2 \pi t)^{-n / 2} e^{-r^{2}(x) / 2 t} \tag{63.2}
\end{equation*}
$$

The main proposition of this section is as follows.
Proposition 63.1. Fix $y \in M$, let $J(x)=J(x, y)$ (see Definition 61.1 above) and $\rho(t, x)$ be as in Eq. (63.2), then

$$
\begin{align*}
\left(\partial_{t}-\frac{1}{2} \Delta\right) \rho & =\frac{1}{2 t} r \partial \ln J / \partial r \rho=\frac{1}{2 t}(V \ln J) \rho \\
& =\frac{1}{2 t}(\nabla \cdot V-n) \rho \tag{63.3}
\end{align*}
$$

Remark 63.2. If $M=\mathbb{R}^{n}$ with the standard metric, then $V(x)=x, \nabla \cdot V=n$ (and $J \equiv 1$ ) so that

$$
p(t, x)=(2 \pi t)^{-n / 2} e^{-x^{2} / 2 t}
$$

is an exact solution to the heat equation as is seen from Eq. (63.3). Moreover, the constants have been chosen such that $\int_{\mathbb{R}^{n}} p(t, x) d x=1$ for all $t>0$. From this fact and the fact that $(2 \pi t)^{-n / 2} e^{-x^{2} / 2 t}$ has most of its mass within a radius of size order $\sqrt{t}$, it follows that $\lim _{t \downarrow 0} p(t, x)=\delta(x)$. Similar statements hold for $\rho(t, x)$ given in Eq. (63.2).

In order to prove the Proposition we will need to introduce some more notation which will allow us to compute the Laplacian on radial functions $f(r)$, see Lemma 65.3 below.

Notation 63.3 (Geodesic Polar Coordinates) Let $y \in M$ be fixed, $r(x):=d(x, y)$ and $\theta(x):=\exp _{y}^{-1}(x) / r(x)$. So that $(r, \theta): M \rightarrow \mathbb{R}_{+} \times S_{y}$ where $S=S_{y}$ is the unit sphere in $T_{y} M$. We also write $\partial f / \partial r=g_{r}(r, \theta)$ when $f=g(r, \theta)$. Alternatively, where

$$
\begin{align*}
\partial f / \partial r & =\frac{1}{r} V f=\left.\frac{1}{r} \frac{d}{d t}\right|_{1} f\left(\exp _{y}\left(t \exp _{y}^{-1}(\cdot)\right)\right)  \tag{63.4}\\
& =\left.\frac{d}{d t}\right|_{0} f\left(\exp _{y}((r+t) \theta)\right. \tag{63.5}
\end{align*}
$$

and $V$ is the vector field given in Eq. (63.2) above.
Notice that with this notation $\exp _{y}^{-1}(x)=r(x) \theta(x)$ and $J(x, y)=$ $D(r(x) \theta(x), y)$.

### 63.0.1 Proof of Proposition 63.1

We begin with the logaritheoremic derivatives of $\rho$,

$$
\partial_{t} \ln \rho(t, x)=-\frac{n}{2 t}+\frac{r^{2}}{2 t^{2}}
$$

and

$$
\nabla \ln \rho(t, x)=-\frac{\nabla r^{2}}{2 t}=-\frac{1}{t} V
$$

Therefore,

$$
\begin{aligned}
\Delta \rho & =\nabla \cdot(\rho \nabla \ln \rho(t, x))=-\frac{1}{t} \nabla \cdot(\rho V) \\
& =\frac{1}{t^{2}}|V|^{2} \rho-\frac{1}{t} \rho \nabla \cdot V=\rho\left(\frac{r^{2}}{t^{2}}-\frac{1}{2 t} \Delta r^{2}\right)
\end{aligned}
$$

and hence

$$
\begin{aligned}
\left(\partial_{t}-\frac{1}{2} \Delta\right) \rho & =\rho\left(-\frac{n}{2 t}+\frac{r^{2}}{2 t^{2}}-\frac{r^{2}}{2 t^{2}}+\frac{1}{4 t} \Delta r^{2}\right) \\
& =\frac{1}{2 t}\left(\frac{1}{2} \Delta r^{2}-n\right) \rho \\
& =\frac{1}{2 t}(\nabla \cdot V-n) \rho=\frac{1}{2 t}(r \partial \ln J / \partial r) \rho .
\end{aligned}
$$

### 63.0.2 On the Operator Associated to the Kernel $\rho$

We now modify the definition of $T_{t}$ in Eq. (61.2) by inserting the cutoff function $\psi$ as in Definition 61.5, that is let

$$
\begin{equation*}
T_{t} f(x):=\int_{M} \psi(x, y) \rho(t, x, y) f(y) d \lambda(y) . \tag{63.6}
\end{equation*}
$$

We will end this section with some basic properties of $T_{t}$.
Theorem 63.4. Let $T_{t} f$ be as in Eq. (63.6). Then for $t>0, T_{t}: C(M) \rightarrow$ $C^{\infty}(M)$, for each $l$, there is a constant $C_{l}$ such that $\left\|T_{t} f\right\|_{l} \leq C_{l}\|f\|_{l}$ for all $0<t \leq 1$ and $f \in C^{l}(M)$ and moreover $\lim _{t \downarrow 0}\left\|T_{t} f-f\right\|_{l}=0$. Here, $\|f\|_{l}$ denotes the sup-norm of $f$ and all of its derivatives up to order $l$.

Proof. First off, since $\psi(x, y) \rho(t, x, y)$ is a smooth function in $(x, y)$, it is clear that $T_{t} f(x)$ is smooth. To prove the remaining two assertions, let us make the change of variables, $y=\exp _{x}(v)$ in the definition of $T_{t} f$. This gives,

$$
\begin{aligned}
T_{t} f(x) & =\int_{B_{x}(\epsilon)} \psi\left(x, \exp _{x}(v)\right) \rho\left(t, x, \exp _{x}(v)\right) f\left(\exp _{x}(v)\right) D(v) d v \\
& =\int_{T_{x} M} \Psi\left(|v|^{2}\right)(2 \pi t)^{-n / 2} e^{-|v|^{2} / 2 t} f\left(\exp _{x}(v)\right) D(v) d v
\end{aligned}
$$

where $B_{x}(\epsilon)$ be the ball of radius $\epsilon$ centered at $0_{x} \in T_{x} M$. Now let $u(x)$ be a local orthonormal frame on $M$, so that $u(x): \mathbb{R}^{n} \rightarrow T_{x} M$ is a smoothly varying orthogonal isomorphism for $x$ in some neighborhood of $M$. We now make the change of variables $v \rightarrow u(x) v$ and $v \rightarrow \sqrt{t} u(x) v$ with $v \in \mathbb{R}^{n}$ in the above displayed equation to find,

$$
\begin{align*}
T_{t} f(x) & =\int_{\mathbb{R}^{n}} \Psi\left(|v|^{2}\right) f\left(\exp _{x}(u(x) v)\right) D(u(x) v) p_{t}(v) d v  \tag{63.7}\\
& =\int_{\mathbb{R}^{n}} \Psi\left(t|v|^{2}\right) f\left(\exp _{x}(u(x) \sqrt{t} v)\right) D(u(x) \sqrt{t} v) p_{1}(v) d v \tag{63.8}
\end{align*}
$$

where $p_{t}(v):=(2 \pi t)^{-n / 2} e^{-|v|^{2} / 2 t}$
Suppose that $L$ is a $l$ 'th order differential operator on $M$, then from Eq. (63.7) we find that

$$
\left(L T_{t} f\right)(x)=\int_{\mathbb{R}^{n}} \Psi\left(|v|^{2}\right) L_{x}\left[f\left(\exp _{x}(u(x) v)\right) D(u(x) v)\right] p_{t}(v) d v
$$

from which we see that

$$
\left|\left(L T_{t} f\right)(x)\right| \leq C_{l}(L)\|f\|_{l} \int_{\mathbb{R}^{n}} p_{t}(v) d v=C_{l}(L)\|f\|_{l}
$$

This shows that $\left\|T_{t} f\right\|_{l} \leq C_{l}\|f\|_{l}$.

Using the product and the chain rule,

$$
L_{x}\left[f\left(\exp _{x}(u(x) v)\right) D(u(x) v, x)\right]=\sum_{k} a_{k}(x, v)\left(\mathcal{L}_{k} f\right)\left(\exp _{x}(u(x) v)\right)
$$

where $a_{k}(x, v)$ are smooth functions of $(x, v)$ with $|v| \leq \epsilon$ and $\mathcal{L}_{k}$ are differential operators of degree at most $l$. Noting that

$$
L f(x)=\sum_{k} a_{k}(x, 0)\left(\mathcal{L}_{k} f\right)(x)
$$

we find that

$$
\begin{aligned}
\mid\left(L T_{t} f\right)(x) & -L f(x) \mid \\
& \leq \sum_{k} \int_{\mathbb{R}^{n}}\left|\begin{array}{c}
a_{k}(x, v)\left(\mathcal{L}_{k} f\right)\left(\exp _{x}(u(x) v)\right) \\
-a_{k}(x, 0)\left(\mathcal{L}_{k} f\right)(x)
\end{array}\right| \Psi\left(|v|^{2}\right) p_{t}(v) d v .
\end{aligned}
$$

Applying the estimate in Eq. (59.6) to the previous equation implies that

$$
\begin{aligned}
& \left\|L T_{t} f-L f\right\|_{0} \\
& \quad \leq C\|f\|_{l} e^{-\delta^{2} / 4 t}+\sum_{k} \sup _{v \in B(\delta)} \sup _{x}\left|\begin{array}{c}
a_{k}(x, v)\left(\mathcal{L}_{k} f\right)\left(\exp _{x}(u(x) v)\right) \\
-a_{k}(x, 0)\left(\mathcal{L}_{k} f\right)(x)
\end{array}\right|
\end{aligned}
$$

and therefore

$$
\varlimsup_{t \downarrow 0}\left\|L T_{t} f-L f\right\|_{0} \leq \sum_{k} \sup _{v \in B(\delta)} \sup _{x}\left|\begin{array}{c}
a_{k}(x, v)\left(\mathcal{L}_{k} f\right)\left(\exp _{x}(u(x) v)\right) \\
-a_{k}(x, 0)\left(\mathcal{L}_{k} f\right)(x)
\end{array}\right|
$$

which tends to zero as $\delta \rightarrow 0$ by uniform continuity. From this we conclude that $\lim _{t \downarrow 0}\left\|T_{t} f-f\right\|_{l}=0$.

To conclude this section we wish to consider $\lim _{t \downarrow 0}\left(\partial_{t}-\frac{1}{2} \Delta\right) T_{t}$.
Theorem 63.5. Let $T_{t}$ be as above and $S$ be the scalar curvature on $M$. Then $\partial_{t} T_{t}=\left(\frac{1}{2} \Delta-\frac{1}{6} S\right) T_{t}+O(\sqrt{t})$. So if we used $T_{t}$ for $K$ in the construction in Proposition 60.1, we would construct $e^{t(\Delta / 2-S / 6)}$ rather than $e^{t \Delta / 2}$.

Proof. We will start by computing,

$$
\begin{align*}
\left(\partial_{t}\right. & \left.\frac{1}{2} \Delta\right) T_{t} f(x) \\
& =\int_{M}\left(\partial_{t}-\frac{1}{2} \Delta_{x}\right) \psi(x, y) \rho(t, x, y) f(y) d \lambda(y) \\
& =\int_{M} \psi(x, y)\left(\partial_{t}-\frac{1}{2} \Delta_{x}\right) \rho(t, x, y) f(y) d \lambda(y)+O\left(t^{\infty}\right) \\
& =\int_{M} \psi(x, y) \frac{1}{2 t}\left(V_{y}(x) \ln J(\cdot, y)\right) \rho_{t}(x, y) f(y) d \lambda(y)+O\left(t^{\infty}\right) \tag{63.9}
\end{align*}
$$

$$
\begin{aligned}
\left(V_{y}(x) \ln J(\cdot, y)\right) & =\partial_{v} \ln D(v)=\partial_{v} \ln \left(1-\frac{1}{6}(\operatorname{Ric} v, v)+O\left(v^{3}\right)\right) \\
& =\partial_{v}\left(-\frac{1}{6}(\operatorname{Ric} v, v)+O\left(v^{3}\right)\right) \\
& =-\frac{1}{3}\left(\operatorname{Ric} \exp _{y}^{-1}(x), \exp _{y}^{-1}(x)\right)+O\left(d^{3}(x, y)\right)
\end{aligned}
$$

Using the symmetry of $J$ or by direct means one may conclude that
$\left(\operatorname{Ric} \exp _{y}^{-1}(x), \exp _{y}^{-1}(x)\right)=\left(\right.$ Ric $\left.\exp _{x}^{-1}(y), \exp _{x}^{-1}(y)\right)+O\left(d^{3}(x, y)\right)$ (63.10)
so that

$$
\left(V_{y}(x) \ln J(\cdot, y)\right)=-\frac{1}{3}\left(\operatorname{Ric} \exp _{x}^{-1}(y), \exp _{x}^{-1}(y)\right)+O\left(d^{3}(x, y)\right)
$$

To check Eq. (63.10) directly, let $\gamma_{v}(t)=\exp (t v)$ and notice that

$$
\frac{d}{d t}(\operatorname{Ric} \dot{\gamma}(t), \dot{\gamma}(t))=\left(\left(\nabla_{\dot{\gamma}(t)} \operatorname{Ric} \dot{\gamma}(t)\right), \dot{\gamma}(t)\right)=O\left(v^{3}\right)
$$

Integrating this expression implies that $(\operatorname{Ric} \dot{\gamma}(1), \dot{\gamma}(1))=(\operatorname{Ric} v, v)+O\left(v^{3}\right)$. Taking $v=\exp _{x}^{-1}(y)$ implies Eq. (63.10).

Using this result in (63.9) and making the change of variables $y=\exp _{x} v$ as above we find that

$$
\begin{aligned}
\left(\partial_{t}-\right. & \left.\frac{1}{2} \Delta\right) T_{t} f(x) \\
=- & \frac{1}{3} \frac{1}{2 t} \int_{T_{x} M} \Psi\left(|v|^{2}\right) \frac{e^{-|v|^{2} / 2 t}}{(2 \pi t)^{n / 2}}\left\{(\text { Ric } v, v)+O\left(v^{3}\right)\right\} f\left(\exp _{x}(v)\right) D(v) d v \\
& +O\left(t^{\infty}\right) \\
=- & \frac{1}{6 t}\left\{S(x) f(x) t+O\left(t^{3 / 2}\right)\right\}=-\frac{1}{6} S(x) f(x)+O\left(t^{1 / 2}\right)
\end{aligned}
$$

Therefore, $\partial_{t} T_{t}=\left(\frac{1}{2} \Delta-\frac{1}{6} S\right) T_{t}+O(\sqrt{t})$.

Similarly we have the same estimate for $\left|W_{y} e^{-d^{2} / 2 t}\right|$. Let us now consider higher order spatial derivatives, for example
$U_{y} W_{x} e^{-d^{2} / 2 t}=-U_{y}\left(\frac{(V, W)}{t} e^{-d^{2} / 2 t}\right)=-\left(\frac{U_{y}(V, W)}{t} e^{-d^{2} / 2 t}+\frac{(V, W)}{t} U_{y} e^{-d^{2} /:}\right.$
from which we find that

$$
\left|U_{y} W_{x} e^{-d^{2} / 2 t}\right| \leq \frac{C}{t} e^{-d^{2} / 2 t}+C \frac{d^{2}}{t^{2}} e^{-d^{2} / 2 t} \leq \frac{C}{t}
$$

Continuing in this way we learn that

$$
\left\|e^{-d^{2} / 2 t}\right\|_{l} \leq C t^{-l / 2}
$$

Let us consider the $t$-derivatives of $\rho(t, x, y)$,

$$
\left|\partial_{t} \rho\right|=\left|-\rho\left(\frac{n}{2 t}+\frac{d^{2}}{2 t^{2}}\right)\right| \leq C t^{-n / 2} t^{-1}
$$

Similarly,

$$
\left|\partial_{t}^{2} \rho\right|=\left|\rho\left(\frac{n}{2 t^{2}}+\frac{d^{2}}{t^{3}}\right)+\rho\left(\frac{n}{2 t}+\frac{d^{2}}{2 t^{2}}\right)^{2}\right| \leq C t^{-n / 2} t^{-2}
$$

Continuing this way, one learns that $\left|\partial_{t}^{k} \rho\right| \leq C C t^{-(n / 2+k)}$. Putting this all together gives Eq. (61.10).

### 64.2 Proof of Theorem 61.4

Proposition 64.1. Let $y \in M$ be fixed, $r(x)=d(x, y), J(x)=J(x, y), V$ and $\rho(t, x)$ be as above. Suppose that $g(t, x): E_{y} \rightarrow E_{x}$ is a time dependent section of $\operatorname{hom}\left(E_{y} \rightarrow E\right)$ and $u(t, x)=\rho(t, x) g(t, x)$. Then

$$
\begin{align*}
\left(\partial_{t}-L\right) u & =\rho\left(\partial_{t}-L+\frac{1}{t}\left(\nabla_{V}+\frac{1}{2} r \partial \ln J / \partial r\right)\right) g \\
& =\rho\left(\partial_{t}-L+\frac{1}{t} S\right) g \tag{64.1}
\end{align*}
$$

where

$$
\begin{equation*}
S=\nabla_{V}+\frac{1}{2} r \partial \ln J / \partial r \tag{64.2}
\end{equation*}
$$

Proof. First let us recall that

$$
\begin{aligned}
\square(\rho g) & =\operatorname{tr} \nabla^{2}(\rho g)=\operatorname{tr}\left(\nabla^{2} \rho g+2 \nabla \rho \otimes \nabla g+\rho \nabla^{2} g\right) \\
& =\Delta \rho g+2 \nabla_{\nabla \rho} g+\rho \square g \\
& =\Delta \rho g+2 \rho \nabla_{\nabla \ln \rho} g+\rho \square g
\end{aligned}
$$

and that

$$
\nabla \ln \rho=\nabla\left(-\frac{1}{2 t} r^{2}\right)=-\frac{V}{t}
$$

Hence

$$
\begin{aligned}
\left(\partial_{t}-L\right) u & =\left(\partial_{t}-\frac{1}{2} \square\right)(\rho g)-\rho \mathcal{R} g \\
& =\left(\partial_{t}-\frac{1}{2} \Delta\right) \rho g-\rho \nabla_{\nabla \ln \rho} g+\rho\left(\partial_{t}-L\right) g \\
& =\frac{1}{2 t}(r \partial \ln J / \partial r) \rho g+\frac{1}{t} \rho \nabla_{V} g+\rho\left(\partial_{t}-L\right) g \\
& =\rho\left(\partial_{t}-L+\frac{1}{t}\left(\nabla_{V}+\frac{1}{2} r \partial \ln J / \partial r\right)\right) g
\end{aligned}
$$

Now let

$$
\begin{equation*}
g_{q}(t, x)=\sum_{k=0}^{q} t^{k} u_{k}(x) \text { and } \Sigma_{q}(t, x)=\rho(t, x) g_{q}(t, x) \tag{64.3}
\end{equation*}
$$

where $u_{k}(x): E_{y} \rightarrow E_{x}$ are to be determined. Then

$$
\begin{aligned}
\left(\partial_{t}+\frac{1}{t} S-L\right) g & =\sum_{k=0}^{q}\left\{t^{k-1}\left(k u_{k}+S u_{k}\right)-t^{k} L u_{k}\right\} \\
& =\frac{1}{t} S u_{0}+\sum_{k=0}^{q-1} t^{k}\left((k+1) u_{k+1}+S u_{k+1}-L u_{k}\right)-t^{q} L u_{q}
\end{aligned}
$$

Thus if we choose $u_{0}$ such that

$$
\begin{equation*}
S u_{0}(x)=\left(\nabla_{V}+\frac{1}{2} V \ln J\right) u_{0}(x)=0 \tag{64.4}
\end{equation*}
$$

and $u_{k}$ such that

$$
\begin{equation*}
(S+k+1) u_{k+1}-L u_{k}=0 \tag{64.5}
\end{equation*}
$$

then $\left(\partial_{t}+\frac{1}{t} S-L\right) g=-t^{q} L u_{q}$ or equivalently by Eq. (64.1),

$$
\begin{equation*}
\left(\partial_{t}-L\right) k_{q}=\left(\partial_{t}-L\right)(\rho g)=-t^{q} \rho L u_{q} \tag{64.6}
\end{equation*}
$$

Let us begin by solving (64.4) for $u_{0}$. For $x, y \in M$, let $\gamma(t)=\gamma_{x, y}(t):=$ $\exp _{y}\left(t \exp _{y}^{-1}(x)\right)$ so that $\gamma_{x, y}$ is the geodesic connecting $y$ to $x$. Notice that $V(\gamma(t))=t \dot{\gamma}(t)$ and therefore $\nabla_{V}=t \nabla_{\dot{\gamma}(t)}$. Therefore, the equation $S u_{0}=0$ is implies that

$$
t \frac{\nabla}{d t} u_{0}\left(\gamma_{x, y}(t)\right)+\frac{1}{2} t\left[\frac{\partial}{\partial t} \ln J\left(\gamma_{x, y}(t)\right)\right] u_{0}\left(\gamma_{x, y}(t)\right)=0
$$

or equivalently that

$$
\begin{equation*}
\left(\frac{d}{d t}+\frac{1}{2} \frac{d}{d t} \ln J\left(\gamma_{x, y}(t)\right)\right)\left\{/ /_{t}\left(\gamma_{x, y}\right)^{-1} u_{0}\left(\gamma_{x, y}(t)\right)\right\}=0 \tag{64.7}
\end{equation*}
$$

We may solve this last equation to find that

$$
/ / t\left(\gamma_{x, y}\right)^{-1} u_{0}\left(\gamma_{x, y}(t)\right)=\frac{1}{\sqrt{J\left(\gamma_{x, y}(t)\right)}} u_{0}(y)
$$

and hence that

$$
u_{0}(x)=/ /{ }_{1}\left(\gamma_{x, y}\right) \frac{1}{\sqrt{J(x, y)}} u_{0}(y)
$$

Since we are going to want $u$ to be a fundamental solution, it is natural to require the $u_{0}(y)=I d_{E_{y}}$. This gives a first order parametrix,

$$
k_{0}(t, x, y):=\frac{1}{\sqrt{J(x, y)}} \rho(t, x, y) \tau(x, y): E_{y} \rightarrow E_{x}
$$

where

$$
\tau(x, y):=/ /_{1}\left(\gamma_{x, y}\right) \text { and } \rho(t, x, y):=\rho(t, x)=(2 \pi t)^{-n / 2} e^{-d^{2}(x, y) / 2 t}
$$

This kernel satisfies,

$$
\left(\partial_{t}-L_{x}\right) k_{0}(t, x, y)=-\rho(t, x, y) L_{x}\left(\frac{1}{\sqrt{J(x, y)}} \tau(x, y)\right)
$$

Proposition 64.2. Let $y \in M$ be fixed and set, for $x$ near $y$,

$$
\begin{equation*}
u_{0}(x, y)=/ /_{1}\left(\gamma_{x, y}\right) \frac{1}{\sqrt{J(x, y)}} \tag{64.8}
\end{equation*}
$$

Then $u_{0}(x, y)$ is smooth for $(x, y)$ near the diagonal in $M \times M$ and $S_{x} u_{0}(x, y)=0$.

Proof. Because of smooth dependence of differential equations on initial conditions and parameters, it follows that $u_{0}(x, y)$ is smooth for $(x, y)$ near the diagonal in $M \times M$. To simplify notation, let $u_{0}(x):=u_{0}(x, y), \tau(x)=\tau(x, y)$, and $J(x)=J(x, y)$. We must verify that $S u_{0}=0$. This is seen as follows:

$$
\begin{aligned}
\nabla_{V} u_{0}(x) & =\nabla_{V}\left(\frac{1}{\sqrt{J(x)}} \tau(x)\right) \\
& =-\frac{1}{2} \frac{1}{\sqrt{J(x)}}(V \ln J)(x) \tau(x)+\frac{1}{\sqrt{J(x)}} \nabla_{V} \tau(x) \\
& =\frac{1}{\sqrt{J(x)}}\left\{\nabla_{V} \tau(x)-\frac{1}{2}(V \ln J)(x) \tau(x)\right\} \\
& =-\frac{1}{2}(V \ln J)(x) u_{0}(x)
\end{aligned}
$$

117064 Proof of Theorem 61.4 and Corollary 61.6
since

$$
\nabla_{V} \tau(x)=\left.r(x) \frac{\nabla}{d t}\right|_{1} \tau\left(\gamma_{x, y}(t)\right)=\left.r(x) \frac{\nabla}{d t}\right|_{1} / /_{t}\left(\gamma_{x, y}\right)=0
$$

Hence

$$
S u_{0}(x)=\nabla_{V} u_{0}(x)+\frac{1}{2}(V \ln J)(x) u_{0}(x)=0 .
$$

We now consider solving Eq. (64.5). Fixing $x$ and $y$ and letting $\gamma(t):=$ $\gamma_{x, y}(t)$, Eq. (64.5) may be written as

$$
\left(t \frac{\nabla}{d t}+\frac{1}{2} t \frac{d}{d t} \ln J(\gamma(t))+k+1\right) u_{k+1}(\gamma(t))-L u_{k}(\gamma(t))=0
$$

or equivalently that
$\left(\frac{d}{d t}+\frac{1}{2} \frac{d}{d t} \ln J(\gamma(t))+\frac{k+1}{t}\right) / / t(\gamma)^{-1} u_{k+1}(\gamma(t))-\frac{1}{t} / /{ }_{t}(\gamma)^{-1} L u_{k}(\gamma(t))=0$.
Letting $f$ be a solution to

$$
\left(\frac{d}{d t}-\frac{1}{2} \frac{d}{d t} \ln J(\gamma(t))-\frac{k+1}{t}\right) f(t)=0
$$

it follows that Eq. (64.9) may be written as

$$
\begin{equation*}
\frac{d}{d t}\left[f(t) / /_{t}(\gamma)^{-1} u_{k+1}(\gamma(t))\right]-\frac{f(t)}{t} / /_{t}(\gamma)^{-1} L u_{k}(\gamma(t))=0 \tag{64.10}
\end{equation*}
$$

We now let $f$ be given by

$$
\begin{aligned}
f(t) & =\exp \left(\int\left[\frac{1}{2} \frac{d}{d t} \ln J(\gamma(t))+\frac{k+1}{t}\right] d t\right) \\
& =\exp \left(\frac{1}{2} \ln J(\gamma(t))+(k+1) \ln t\right)=\sqrt{J(\gamma(t))} t^{(k+1)}
\end{aligned}
$$

Integrating (64.10) over $[0, t]$ implies that

$$
f(t) / /_{t}(\gamma)^{-1} u_{k+1}(\gamma(t))=\int_{0}^{t} \frac{f(\tau)}{\tau} / /_{\tau}(\gamma)^{-1} L u_{k}(\gamma(\tau)) d \tau
$$

Evaluating this equation at $t=1$ and solving for $u_{k+1}(x)$ gives:

$$
\begin{aligned}
u_{k+1}(x) & =\frac{1}{\sqrt{J(x, y)}} \tau(x, y) \int_{0}^{1} \frac{\sqrt{J(\gamma(t))} s^{(k+1)}}{s} / / s(\gamma)^{-1} L u_{k}(\gamma(s)) d s \\
& =\frac{1}{\sqrt{J(x, y)}} \tau(x, y) \int_{0}^{1} s^{k} \sqrt{J\left(\gamma_{x, y}(s), y\right)} /{ }_{s}\left(\gamma_{x, y}\right)^{-1} L u_{k}\left(\gamma_{x, y}(s)\right) d s
\end{aligned}
$$

Theorem 64.3. Let $u_{0}(x, y)$ be given as in Equation (64.8), and define the smooth sections $u_{k}(x, y)$ inductively by

$$
\begin{equation*}
u_{k+1}(x, y)=u_{0}(x, y) \int_{0}^{1} s^{k} u_{0}\left(\gamma_{x, y}(s), y\right)^{-1} L_{x} u_{k}\left(\gamma_{x, y}(s), y\right) d s \tag{64.12}
\end{equation*}
$$

for $k=0,1,2, \ldots$. Then $u_{k}$ solves Eq.(64.5).
Proof. Let us begin by noting that Eq. (64.11) and (64.12) are the same equation because

$$
u_{0}\left(\gamma_{x, y}(s), y\right)=\frac{1}{\sqrt{J\left(\gamma_{x, y}(s), y\right)}} / /_{s}\left(\gamma_{x, y}\right)
$$

Let $x, y \in M$ be fixed and set $\gamma(t)=\gamma_{x, y}(t)$. Since $\gamma_{y, \gamma(t)}(s)=\gamma(t s)$ and $/ / s\left(\gamma_{y, \gamma(t)}\right)=/ / t s(\gamma)$,it follows that

$$
\begin{aligned}
u_{k+1}(\gamma(t), y) & =u_{0}(\gamma(t), y) \int_{0}^{1} s^{k} 1 \sqrt{J(\gamma(t s), y)} / / t s(\gamma)^{-1} L_{x} u_{k}(\gamma(t s), y) d s \\
& =t^{-(k+1)} u_{0}(\gamma(t), y) \int_{0}^{t} r^{k} \sqrt{J(\gamma(r), y)} / / r(\gamma)^{-1} L_{x} u_{k}(\gamma(r), y) d r
\end{aligned}
$$

From this equation we learn that

$$
\begin{aligned}
\left(\nabla_{V} u_{0}(x)\right. & \left.+\frac{1}{2}(V \ln J)(x)+k+1\right) u_{k+1}(x, y) \\
& =\left.\left(t \frac{\nabla}{d t}+\frac{1}{2} t \frac{d}{d t} \ln J(\gamma(t))+k+1\right)\right|_{t=1} u_{k+1}(\gamma(t)) \\
& =\left.u_{0}(x, y) \frac{d}{d t}\right|_{1} \int_{0}^{t} r^{k} \sqrt{J(\gamma(r), y)} / /_{r}(\gamma)^{-1} L_{x} u_{k}(\gamma(r), y) d r \\
& =u_{0}(x, y) \sqrt{J(\gamma(r), y)} / /{ }_{1}(\gamma)^{-1} L_{x} u_{k}(x, y)=L_{x} u_{k}(x, y)
\end{aligned}
$$

wherein the second equality we have used the product rule and the fact that

$$
\left(t \frac{\nabla}{d t}+\frac{1}{2} t \frac{d}{d t} \ln J(\gamma(t))+k+1\right) t^{-(k+1)} u_{0}(\gamma(t), y)=0
$$

which is verified using Eq. (64.7).

## Appendix: Gauss' Lemma \& Polar Coordinates

## Lemma 65.1 (Gauss' Lemma). Let $y \in M, v, w \in T_{y} M$, then

$$
\left(\exp _{* y} v_{v}, \exp _{* y} w_{v}\right)_{\exp (v)}=(v, w)_{y}
$$

Proof. Let $\Sigma(t, s):=\exp (t(v+s w))$, then

$$
\begin{aligned}
\frac{d}{d t}\left(\dot{\Sigma}(t, 0), \Sigma^{\prime}(t, 0)\right) & =\left(\frac{\nabla}{d t} \dot{\Sigma}(t, 0), \Sigma^{\prime}(t,)\right)+\left(\dot{\Sigma}(t, 0), \frac{\nabla}{d t} \Sigma^{\prime}(t, 0)\right) \\
& =\left(\dot{\Sigma}(t, 0),\left.\frac{\nabla}{d s}\right|_{0} \dot{\Sigma}(t, s)\right)=\left.\frac{1}{2} \frac{d}{d s}\right|_{0}(\dot{\Sigma}(t, s), \dot{\Sigma}(t, s)) \\
& =\left.\frac{1}{2} \frac{d}{d s}\right|_{0}|v+s w|^{2}=(v, w)
\end{aligned}
$$

Combining this equation with the observation that $\left.\left(\dot{\Sigma}(t, 0), \Sigma^{\prime}(t, 0)\right)\right|_{t=0}=0$ implies that

$$
\left(\dot{\Sigma}(1,0), \Sigma^{\prime}(1,0)\right)=(v, w)_{0} .
$$

Corollary 65.2. Suppose that $y \in M$ and choose $\delta>0$ such that $\exp _{y}$ is a diffeomorphism on $B\left(0_{y}, \delta\right)$. Then $d(x, y)=\left|\exp _{y}^{-1}(x)\right|$ for all $x \in V:=$ $\exp _{y}\left(B\left(0_{y}, \delta\right)\right)$.

Proof. Let $\sigma(t)$ be a curve in $M$ such that $\sigma(0)=y$ and $\sigma(1)=x$. Suppose for the moment that $\sigma(t)$ is contained in $V$ and write $\sigma(t)=\exp _{y}(c(t))$. Set $u=c(1) /|c(1)|$ and decompose $c(t)=(c(t), u) u+d(t)$ where $(d(t), u)=$ 0 .Then

$$
\begin{aligned}
|\dot{\sigma}(t)|^{2} & =\left|\exp _{y *} \dot{c}(t)_{c(t)}\right|^{2}=\left|\exp _{y *}\left((\dot{c}(t), u) u_{c(t)}+\dot{d}(t)_{c(t)}\right)\right|^{2} \\
& =\left|\exp _{y *}(\dot{c}(t), u) u_{c(t)}\right|^{2}+\left|\exp _{y *}\left(\dot{d}(t)_{c(t)}\right)\right|^{2} \\
& =|(\dot{c}(t), u)|^{2}+\left|\exp _{y *}\left(\dot{d}(t)_{c(t)}\right)\right|^{2} \geq|(\dot{c}(t), u)|^{2}
\end{aligned}
$$

From this we learn that

$$
\operatorname{Length}(\sigma)=\int_{0}^{1}|\dot{\sigma}(t)| d t \geq \int_{0}^{1}|(\dot{c}(t), u)| d t \geq \int_{0}^{1}(\dot{c}(t), u) d t=|c(1)|
$$

That is $\operatorname{Length}(\sigma) \geq\left|\exp _{y}^{-1}(x)\right|$. It is easily to use the same argument to show that if $\sigma$ leaves the open set $V$ then $\operatorname{Length}(\sigma) \geq \delta>\left|\exp _{y}^{-1}(x)\right|$ and hence Length $(\sigma) \geq\left|\exp _{y}^{-1}(x)\right|$ for all path $\sigma$ such that $\sigma(0)=y$ and $\sigma(1)=x$. Moreover we have equality if $\sigma(t)$ is the geodesic joining $y$ to $x$. This shows that

$$
d(x, y)=\inf _{\sigma} \operatorname{Length}(\sigma)=\left|\exp _{y}^{-1}(x)\right|
$$

For more on geodesic coordinates, see Appendix 67.

### 65.1 The Laplacian of Radial Functions

Lemma 65.3. Let $r(x):=d(x, y)$ and $J(x)=J(x, y)$ as in Definition 61.1. Then

$$
\begin{equation*}
\Delta f(r)=\frac{\partial\left(J r^{n-1} f^{\prime}(r)\right) / \partial r}{J r^{n-1}}=f^{\prime \prime}(r)+\left(\frac{n-1}{r}+\frac{\partial \ln J}{\partial r}\right) f^{\prime}(r) \tag{65.1}
\end{equation*}
$$

We also have that

$$
\begin{equation*}
\Delta f(r)=f^{\prime \prime}(r)+\frac{\nabla \cdot V-1}{r} f^{\prime}(r) \tag{65.2}
\end{equation*}
$$

and that

$$
\Delta r^{2}=2 \nabla \cdot V=2\left(n+r \frac{\partial \ln J}{\partial r}\right)
$$

Proof. We will give two proofs of this result. For the first proof recall that if $\left\{z^{i}\right\}$ is a chart on $M$, then

$$
\begin{equation*}
\Delta F=\frac{1}{\sqrt{g}} \frac{\partial}{\partial z^{i}}\left(\sqrt{g} g^{i j} \frac{\partial}{\partial z^{j}} F\right) \tag{65.3}
\end{equation*}
$$

where $d s^{2}=g_{i j} d z^{i} d z^{j}, g^{i j}$ is the inverse of $\left(g_{i j}\right)$ and $\sqrt{g}=\sqrt{\operatorname{det}\left(g_{i j}\right)}$ We now choose the coordinate system $z$ to be $z^{n}:=r$, and $z^{i}:=\alpha^{i} \circ \theta$, where $\left\{\alpha^{i}\right\}_{i=1}^{n-1}$ is a chart on $S=S_{y} \subset T_{y} M$. We now need to compute $g_{i j}$ in this case. Let us begin by noting that $\partial / \partial z^{i}=\exp _{y *}\left(\left.r \frac{\partial}{\partial \alpha^{2}}\right|_{\theta}\right)$ for $i=1,2, \ldots, n-1$ and $\partial / \partial z^{n}=\partial / \partial r=\exp _{y *}\left(\left.\frac{d}{d t}\right|_{0}(r+t) \theta\right)$. By Gauss' lemma, it follows that $d s^{2}=d r^{2}+h_{i j}(r, \theta) d \theta^{i} d \theta^{j}$, i.e. $g=\left[\begin{array}{ll}h & 0 \\ 0 & 1\end{array}\right]$. Therefore $g^{-1}=\left[\begin{array}{cc}h^{-1} & 0 \\ 0 & 1\end{array}\right]$ and $\sqrt{g}=\sqrt{h}$. So if $F=f(r)=f\left(z^{n}\right)$, it follows from Eq. (65.3) that

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$$
\begin{align*}
\Delta f(r) & =\frac{1}{\sqrt{g}} \frac{\partial}{\partial z^{i}}\left(\sqrt{g} g^{i n} \frac{\partial}{\partial z^{n}} F\right)=\frac{1}{\sqrt{h}} \frac{\partial}{\partial z^{n}}\left(\sqrt{h} \frac{\partial}{\partial z^{n}} F\right) \\
& =f^{\prime \prime}(r)+\frac{\partial \ln (\sqrt{h})}{\partial r} f^{\prime}(r) . \tag{65.4}
\end{align*}
$$

So to finish the proof we need to describe $\sqrt{g}=\sqrt{h}$ in terms of $J$. In order to do this, let us notice that $D$ in Eq. (61.4) may also be expressed as

$$
\begin{align*}
D(v) & =\frac{\lambda\left(\exp _{*}\left(w_{1}\right)_{v}, \exp _{*}\left(w_{2}\right)_{v} \ldots, \exp _{*}\left(w_{n}\right)_{v}\right)}{\lambda_{y}\left(w_{1}, w_{2}, \ldots, w_{n}\right)} \\
& =\sqrt{\frac{\operatorname{det}\left\{\left(\exp _{y *}\left(w_{i}\right)_{v}, \exp _{y *}\left(w_{j}\right)_{v}\right)\right\}_{i, j=1}^{n}}{\operatorname{det}\left\{\left(w_{i}, w_{j}\right)\right\}_{i, j=1}^{n}}} \tag{65.5}
\end{align*}
$$

where now $\left\{w_{i}\right\}$ is any oriented basis for $T_{y} M$. From this expression it follows that

$$
\sqrt{g}=D(r \theta) \sqrt{\operatorname{det}\left\{\left(\left.r \frac{\partial}{\partial \alpha^{i}}\right|_{\theta},\left.r \frac{\partial}{\partial \alpha^{j}}\right|_{\theta}\right)\right\}_{i, j=1}^{n-1}}=J(\cdot, y) r^{n-1} K(\theta)
$$

where

$$
K(\theta):=\sqrt{\operatorname{det}\left\{\left(\left.\frac{\partial}{\partial \alpha^{i}}\right|_{\theta},\left.\frac{\partial}{\partial \alpha^{j}}\right|_{\theta}\right)\right\}_{i, j=1}^{n-1}}
$$

Using these expression in Eq. (65.4) along with the observation that $\partial K(\theta) / \partial r=0$ proves Eq. (65.1).
(Second more direct Proof.) Let $(\rho, \omega)$ denote a generic point in $\mathbb{R}_{+} \times S$ and $d \omega$ denote the volume form on $S=S_{y} \subset T_{y} M$. Then

$$
\begin{aligned}
\int_{M} f(r, \theta) d \lambda & =\int_{T_{y} M} \tilde{f} D d \lambda_{y} \\
& =\int_{\mathbb{R}_{+} \times S} f(\rho, \omega) D(\rho \omega) \rho^{n-1} d \rho d \omega
\end{aligned}
$$

where $\tilde{f}(\rho \omega):=f(\rho, \omega)$. Therefore,

$$
\begin{aligned}
\int_{M} & \Delta f(r) g(r, \theta) d \lambda=-\int_{M}(\nabla f(r), \nabla g(r, \theta)) d \lambda \\
& =-\int_{M} f^{\prime}(r) \partial g(r, \theta) / \partial r d \lambda \\
& =-\int_{\mathbb{R}_{+} \times S} f^{\prime}(\rho) \partial g(\rho, \omega) / \partial \rho D(\rho \omega) \rho^{n-1} d \rho d \omega \\
& =\int_{\mathbb{R}_{+} \times S} \frac{\partial\left(f^{\prime}(\rho) D(\rho \omega) \rho^{n-1}\right) / \partial \rho}{D(\rho \omega) \rho^{n-1}} g(\rho, \omega) D(\rho \omega) \rho^{n-1} d \rho d \omega \\
& =\int_{M} \frac{\partial\left(f^{\prime}(r) J r^{n-1}\right) / \partial r}{J r^{n-1}} g(r, \theta) d \lambda
\end{aligned}
$$

which proves Eq. (65.1)
To prove Eq. (65.2), we compute more directly:

$$
\begin{aligned}
\Delta f(r) & =\nabla \cdot \nabla f(r)=\nabla \cdot\left(f^{\prime}(r) \partial / \partial r\right)=\nabla \cdot\left(\frac{f^{\prime}(r)}{r} V\right) \\
& =\left(\frac{\partial}{\partial r} \frac{f^{\prime}(r)}{r}\right)\left(\frac{\partial}{\partial r}, V\right)+\frac{f^{\prime}(r)}{r} \nabla \cdot V \\
& =f^{\prime \prime}(r)-\frac{f^{\prime}(r)}{r}+\frac{f^{\prime}(r)}{r} \nabla \cdot V
\end{aligned}
$$

which proves Eq. (65.2).
In particular we have that $\Delta r^{2}=2 \nabla \cdot V$ and

$$
\Delta r^{2}=2+\left(\frac{n-1}{r}+\frac{\partial \ln J}{\partial r}\right) 2 r=2\left(n+r \frac{\partial \ln J}{\partial r}\right)
$$

Lemma 65.4. Let $r(x):=d(x, y)$ as above, then

$$
\begin{gathered}
V r^{2}(x)=2 r^{2}(x) \\
\nabla r^{2}=2 V
\end{gathered}
$$

and

$$
\nabla \cdot V(y)=n
$$

Proof. We first claim that $|V|=r$. Moreover,

$$
V r^{2}(x)=\left.\frac{d}{d t}\right|_{1} r^{2}\left(\gamma_{x}(t)\right)=\left.\frac{d}{d t}\right|_{1} t^{2} r^{2}(x)=2 r^{2}(x)
$$

That is to say $V r^{2}=2 r^{2}$. Moreover, if $w \in T_{x} M$ is perpendicular to $V(x)$, then $w r^{2}=0$ so that $\nabla r^{2}$ is proportional to $V$. One way to argue this is that $V(x)$ points in the direction of maximum increase of $r^{2}$ by the triangle inequality. Hence

$$
\nabla r^{2}=\frac{\left(\nabla r^{2}, V\right)}{(V, V)} V=\frac{2 r^{2}}{r^{2}} V=2 V
$$

If we do not like this explanation, then use Gauss's lemma I guess. Come back to this point.

Now we wish to compute $\Delta r^{2}=2 \nabla \cdot V$. We would like to at least do this at $x=y$. To this end, let us work out $\nabla_{w} V$ for $w \in T_{y} M$. Setting $\sigma(t)=\exp (t w)$, we find that

$$
\nabla_{w} V=\left.\frac{d}{d t}\right|_{0} / /_{t}^{-1}(\sigma) V(\sigma(t))=\left.\frac{d}{d t}\right|_{0} / /_{t}^{-1}(\sigma) / / t(\sigma) t w=w
$$

Therefore $\nabla \cdot V(y)=\sum_{i=1}^{n}\left(\nabla_{e_{i}} V, e_{i}\right)=n$ and hence $\Delta r^{2}(y)=2 n$.

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## The Dirac Equation a la Roe's Book

In this section, we consider the Dirac equation

$$
\begin{equation*}
\partial_{t} S=i D S \text { with } S_{t=0}=S_{0} \text { given } . \tag{66.1}
\end{equation*}
$$

Here $D=\gamma_{e_{i}} \nabla_{e_{i}}$ is the Dirac operator on some spinor bundle over $M$. The most interesting statement made by Roe about the Dirac equation is it's finite speed of propagation property. Given a compact region $\Omega \subset M$ with smooth boundary and a solution to Eq. (66.1), let

$$
E_{\Omega}(t):=\int_{\Omega}\left|S_{t}(x)\right|^{2} d \lambda(x)
$$

be the energy of $S_{t}$ in the region $\Omega$. Let us begin by computing the derivative of $E_{\Omega}$.

Lemma 66.1. Let $\Omega \subset M$ be a compact region with smooth boundary, and $S$ be a solution to Eq. (66.1). Then

$$
\begin{equation*}
\frac{d}{d t} E_{\Omega}(t)=i \int_{\Omega} \nabla \cdot X_{t} d \lambda=i \int_{\partial \Omega}\left(X_{t}, N\right) d \sigma, \tag{66.2}
\end{equation*}
$$

where $X_{t}$ is the smooth vector field on $M$ such that $\left(X_{t}, Y\right)=\left(\gamma_{Y} S_{t}, S_{t}\right)$ for all vector fields $Y$ on $M$ and $N$ is the outward pointing normal to $\Omega$ and $\sigma$ is surface measure on $\Omega_{t}$.

Proof. Differentiating under the integral sign implies that

$$
\begin{align*}
\frac{d}{d t} E_{\Omega}(t) & =\int_{\Omega}\left\{\left(\dot{S}_{t}, S_{t}\right)+\left(S_{t}, \dot{S}_{t}\right)\right\} d \lambda \\
& =\int_{\Omega}\{(i D S, S)+(S, i D S)\} d \lambda \\
& =i \int_{\Omega}\{(D S, S)-(S, D S)\} d \lambda  \tag{66.3}\\
& =-\int_{\Omega} \operatorname{Im}(D S, S) d \lambda \tag{66.4}
\end{align*}
$$

We also have,

$$
\begin{aligned}
\left(\nabla_{*} X, \cdot\right) & =\nabla_{*}(X, \cdot)=\nabla_{*}(\gamma \cdot S, S) \\
& =\left(\left(\nabla_{*} \gamma\right) \cdot S, S\right)+\left(\gamma \cdot \nabla_{*} S, S\right)+\left(\gamma \cdot S, \nabla_{*} S\right) \\
& =\left(\gamma \cdot \nabla_{*} S, S\right)-\left(S, \gamma \cdot \nabla_{*} S\right),
\end{aligned}
$$

where we have made use of the fact that $\nabla \gamma=0$ and $\gamma_{X}^{*}=-\gamma_{X}$. Taking $*=e_{i}$ and $\cdot=e_{i}$ and summing on $i$ in the above equation implies that

$$
\begin{equation*}
\nabla \cdot X=(D S, S)-(S, D S) \tag{66.5}
\end{equation*}
$$

Combining Eq. (66.3) and (66.5) along with the divergence (Stoke's) theorem proves Eq. (66.2).

Corollary 66.2. The total energy $E(t)=E_{M}(t)$ remains constant and solutions to Eq. (66.1) are unique if they exist.

We now want to examine how $E_{\Omega}(t)$ depends on $\Omega$.
Lemma 66.3. Suppose that $\Omega \subset M$ is as above and $\phi_{t}: \Omega \rightarrow M$ is a one parameter family of smooth injective local diffeomorphisms depending smoothly on $t$ and let $\Omega_{t}:=\phi_{t}(\Omega)$. Also define a vector field $Y_{t}$ on $\Omega_{t}$ by $\dot{\phi}_{t}=Y_{t} \circ \phi_{t}$, i.e. $Y_{t}:=\phi_{t} \circ \phi_{t}^{-1}$. If $f: M \rightarrow \mathbb{R}$ is a smooth function, then

$$
\begin{equation*}
\frac{d}{d t} \int_{\Omega_{t}} f d \lambda=\int_{\Omega_{t}} \nabla \cdot\left(f Y_{t}\right) d \lambda=\int_{\partial \Omega_{t}} f\left(Y_{t}, N\right) d \sigma \tag{66.6}
\end{equation*}
$$

where again $N$ is the outward pointing normal to $\partial \Omega_{t}$ and $\sigma$ is surface measure on $\Omega_{t}$.

Proof. Since

$$
\begin{aligned}
\int_{\Omega_{t}} f d \lambda & =\int_{\phi_{t}(\Omega)} f \lambda=\int_{\Omega} \phi_{t}^{*}(f \lambda)=\int_{\Omega} f \circ \phi_{t} \phi_{t}^{*} \lambda \\
& =\int_{\Omega} f \circ \phi_{t} \phi_{t}^{*} \lambda
\end{aligned}
$$

$\frac{d}{d t} f \circ \phi_{t}=Y_{t} f \circ \phi_{t}$ and

$$
\frac{d}{d t} \phi_{t}^{*} \lambda=\phi_{t}^{*}\left(\left(d i_{Y_{t}}+i_{Y_{t}} d\right) \lambda\right)=\phi_{t}^{*}\left(\nabla \cdot Y_{t} \lambda\right)
$$

it follows that

$$
\begin{aligned}
\frac{d}{d t} \int_{\Omega_{t}} f d \lambda & =\int_{\Omega} Y_{t} f \circ \phi_{t} \phi_{t}^{*} \lambda+\int_{\Omega} f \circ \phi_{t} \phi_{t}^{*}\left(\nabla \cdot Y_{t} \lambda\right) \\
& =\int_{\Omega_{t}} Y_{t} f \lambda+\int_{\Omega_{t}} f \nabla \cdot Y_{t} \lambda \\
& =\int_{\Omega_{t}} \nabla \cdot\left(f Y_{t}\right) \lambda
\end{aligned}
$$

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from which Eq. (66.6) follows.
Fix a point $m \in M$ and let $B(m, R)$ be the geodesic ball centered at $m \in M$ with radius $R$. If we believe that the speed of propagation of the Dirac equation is 1 , then we should have

$$
\begin{equation*}
e(t):=\int_{B(m, R-t)}\left|S_{t}\right|^{2} d \lambda \tag{66.7}
\end{equation*}
$$

is non-increasing as $t$ increases to $R$. The reason is that, we are shrinking the ball at a rate equal to the speed of propagation, so no energy which was in the wave at time $T$ outside the ball $B(m, R-T)$ can enter the region $B(m, R-t)$ for $t \leq T$. We will now verify that $e(t)$ is non-increasing.

Proposition 66.4. For $R$ smaller than the injectivity radius of $M$, the function $e(t)$ in Eq. (66.7) is non-increasing as increases to $R$.

Proof. Let $\phi_{t}: T_{m} M \rightarrow M$ be given by $\phi_{t}(v)=\exp (t v)$ and $Y_{t}$ be the locally defined vector field on $M$ such that $\dot{\phi}_{t}(v)=Y_{t} \circ \phi_{t}(v)$ for all $v$ small. Since $\phi_{t}(D)=B(m, t)$, where $D$ is the unit disc in $T_{m} M$, we have that $\left(Y_{t}, N\right)=1$, by Gauss's lemma. So By Lemmas 66.1 and 66.3 ,

$$
\begin{aligned}
\frac{d}{d t} e(t) & =i \int_{\partial B(m, R-t)}\left(X_{t}, N\right) d \sigma-\int_{\partial B(m, R-t)}\left|S_{t}\right|^{2}\left(Y_{R-t}, N\right) d \sigma \\
& =i \int_{\partial B(m, R-t)}\left(X_{t}, N\right) d \sigma-\int_{\partial B(m, R-t)}\left|S_{t}\right|^{2} d \sigma
\end{aligned}
$$

Now

$$
\left|\left(X_{t}, N\right)\right|=\left|\left(\gamma_{N} S, S\right)\right| \leq\left|\gamma_{N} S\right||S| \leq|S|^{2}
$$

since $N$ is a unit vector and $\gamma_{N}$ is an isometry. (Recall that $\gamma_{N}$ is skew adjoint and $\gamma_{N}^{2}=-I$.). This shows that

$$
\begin{aligned}
\frac{d}{d t} e(t) & =-\operatorname{Im} \int_{\partial B(m, R-t)}\left(X_{t}, N\right) d \sigma-\int_{\partial B(m, R-t)}\left|S_{t}\right|^{2} d \sigma \\
& \leq \int_{\partial B(m, R-t)}|S|^{2} d \sigma-\int_{\partial B(m, R-t)}\left|S_{t}\right|^{2} d \sigma=0
\end{aligned}
$$

Corollary 66.5. Suppose that the support of $S_{0}$ is contained in $\Omega$. Then the support of $S_{t}$ is contained in

$$
\Omega_{t}:=\{x \in M: d(x, m)<t \text { for all } m \in \Omega\}
$$

Proof. By repeating the argument and using the semi-group property of $e^{i t D}$, we may and do assume that $t$ is positive and less than the injectivity radius of $M$. Let $x \notin \bar{\Omega}_{t}$, so that there exists $R>t$ such that $B(x, R) \cap \Omega=\emptyset$. By the previous proposition, $e(\tau):=\int_{B(x, R-\tau)}\left|S_{\tau}\right|^{2} d \lambda$ is decreasing and hence

$$
\int_{B(x, R-t)}\left|S_{t}\right|^{2} d \lambda=e(t) \leq e(0)=\int_{B(x, R)}\left|S_{0}\right|^{2} d \lambda=0
$$

This shows that $S_{t} \equiv 0$ on $B(x, R-t)$ and in particular at $x$.

### 66.1 Kernel Construction

Lemma 66.6. Suppose that $U \subset \mathbb{R}^{N}$ is a open set and $A: L^{2}(E) \rightarrow C^{r+1}(U)$ is bounded linear map. For each $x \in U$, let $T_{x} \in L^{2}(E)$ such that $\left(T_{x}, S\right)=$ $A S(x)$ for all $S \in L^{2}(E)$. Then the map $x \in U \rightarrow T_{x} \in L^{2}(E)$ is $C^{r}-$ smooth. Moreover, we have estimates of the derivatives of $x \rightarrow T_{x}$ in terms of the operator norm $\|A\|_{o p}$ of $A$ as an operator from $A: L^{2}(E) \rightarrow C^{r+1}(U)$.
Remark 66.7. The above Lemma may has well been formulated with $L^{2}(E)$ replaced by an abstract Hilbert space $H$. The proof given below would still go through without any change.

Proof. First notice that

$$
\begin{aligned}
\left|\left(T_{x}-T_{y}, S\right)\right| & =|A S(x)-A S(y)| \leq|\nabla A S|_{U}|x-y| \\
& \leq|A S|_{C^{1}(U)}|x-y| \leq C|S|_{L^{2}(E)}|x-y|
\end{aligned}
$$

which shows that $\left|T_{x}-T_{y}\right|_{L^{2}(E)} \leq C|x-y|$, so the $T$. is continuous. Let us now consider the directional derivatives of $T_{x}$. For $x \in U$ and $v \in \mathbb{R}^{N}$, let $B S(x, v):=\partial_{v} A S(x)$. As above there exists $T_{x, v} \in L^{2}(E)$ such $\left(T_{x, v}, S\right)=$ $\partial_{v} A S(x)=B S(x, v)$ for all $S \in L^{2}(E)$ and moreover $(x, v) \rightarrow T_{x, v}$ is locally Lipschitz continuous and linear in $v$. Indeed,

$$
\begin{aligned}
\left|\left(T_{x, v}-T_{y, v}, S\right)\right| & =\left|\partial_{v} A S(x)-\partial_{v} A S(y)\right| \leq|A S|_{C^{2}(U)}|x-y \| v| \\
& \leq\|A\|_{o p}|S|_{L^{2}(E)}|x-y \| v|
\end{aligned}
$$

That is to say, $x \in U \rightarrow T_{x,} \in B\left(\mathbb{R}^{N}, L^{2}(E)\right)$ is a Lipschitz continuous map. Now let $x \in U$ and $v \in \mathbb{R}^{N}$, then

$$
\begin{aligned}
\left(T_{x+v}-T_{x}, S\right) & =A S(x+v)-A S(x)=\int_{0}^{1} \partial_{v} A S(x+t v) d t \\
& =\int_{0}^{1}\left(T_{x+t v, v}, S\right) d t
\end{aligned}
$$

which shows that

$$
\begin{aligned}
\left|T_{x+v}-T_{x}-T_{x, v}\right| & \leq \int_{0}^{1}\left|T_{x+t v, v}-T_{x, v}\right| d t \leq C|v| \int_{0}^{1}|t v| d t \\
& =\|A\|_{o p}|v|^{2} / 2
\end{aligned}
$$

This shows that $T_{x}$ is differentiable and that $T_{x}^{\prime} v=T_{x, v}$. We have already seen that $T_{x}^{\prime}$ is continuous. This shows that $T_{x}$ is $C^{1}$. We may continue this way inductively to finish the proof of the lemma.

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Proposition 66.8. Suppose that $A: L^{2}(E) \rightarrow C^{r+1}(E)$ is a bounded operator, then $A A^{*}$ has an integral kernel which is $C^{r}-$ smooth. Moreover the $C^{r}$ - norm of the kernel is bounded by the square of the operator norm for $A: L^{2}(E) \rightarrow C^{r+1}(E)$.

Proof. Let $x \in M$ the map $S \in L^{2}(E) \rightarrow A S(x) \in E_{x}$ is a bounded linear map and hence there is a unique element $T(x, \cdot) \in L^{2}\left(\operatorname{End}\left(E, E_{x}\right)\right)$ such that

$$
A S(x)=\int_{M} T(x, y) S(y) d y
$$

Notice that if $\xi \in \Gamma(E)$, then

$$
\begin{aligned}
(\xi(x), A S(x)) & =\int_{M}(\xi(x), T(x, y) S(y)) d y \\
& =\int_{M}\left(T^{*}(x, y) \xi(x), S(y)\right) d y
\end{aligned}
$$

which by the previous lemma shows that $x \rightarrow T^{*}(x, \cdot) \xi(x) \in L^{2}(E)$ is a $C^{r}-$ map with bounds determined by $\|A\|_{o p}$.
(Surely one can show that there is a version of $T(x, y)$ such that $(x, y) \rightarrow$ $T(x, y)$ is jointly measurable. We will avoid this issue here however.) Ignoring measurability issues, we know that

$$
A^{*} S(x)=\int_{M} T^{*}(y, x) S(y) d y
$$

so the

$$
\begin{aligned}
A A^{*} S(x) & =\int_{M \times M} T(x, y) T^{*}(z, y) S(z) d y d z \\
& =\int_{M} k(x, z) S(z) d z
\end{aligned}
$$

where

$$
k(x, z):=\int_{M} T(x, y) T^{*}(z, y) d y
$$

Even though the derivation of $k$ above was suspect because of measurability questions, the formula make perfect sense. Indeed suppose that $\xi$ and $\eta$ are in $\Gamma(E)$, then

$$
\begin{aligned}
(\eta(x), k(x, z) \xi(z))_{E_{x}} & =\int_{M}\left(\eta(x), T(x, y) T^{*}(z, y) \xi(z)\right)_{E_{x}} d y \\
& =\int_{M}\left(T^{*}(x, y) \eta(x), T^{*}(z, y) \xi(z)\right)_{E_{y}} d y \\
& =\left(T^{*}(x, \cdot) \eta(x), T^{*}(z, \cdot) \xi(z)\right)_{L^{2}(E)}
\end{aligned}
$$

Furthermore this shows that $(x, z) \rightarrow(\eta(x), k(x, z) \xi(z))_{E_{x}}$ is $C^{r}$ in $(x, y)$ with a $C^{r}$ norm which is controlled by the $C^{r}$ - norms of $\eta, \xi$, and $\|A\|_{o p}^{2}$. Since $\eta$ and $\xi$ are arbitrary, we find that $k(x, z)$ is $C^{r}$ as well and the $C^{r}-$ norm is bounded by a constant times $\|A\|_{o p}^{2}$.

So the only thing left to check is that

$$
A A^{*} S(x)=\int_{M} k(x, z) S(z) d \lambda(z)
$$

Letting $\xi, \eta \in \Gamma(E)$ as before, then

$$
\begin{aligned}
(\eta(x), k(x, z) \xi(z))_{E_{x}} & =\int_{M}\left(\eta(x), T(x, y) T^{*}(z, y) \xi(z)\right)_{E_{x}} d y \\
& =\left(\eta(x),\left(A T^{*}(z, \cdot) \xi(z)\right)(x)\right)_{E_{x}}
\end{aligned}
$$

so that

$$
\begin{aligned}
\int_{M}(\eta(x), k(x, z) \xi(z))_{E_{x}} d x & =\int_{M}\left(\eta(x),\left(A T^{*}(z, \cdot) \xi(z)\right)(x)\right)_{E_{x}} d x \\
& =\int_{M}\left(A^{*} \eta(x), T^{*}(z, x) \xi(z)\right)_{E_{x}} d x \\
& =\int_{M}\left(T(z, x) A^{*} \eta(x), \xi(z)\right)_{E_{x}} d x .
\end{aligned}
$$

Integrating this last expression over $z$ shows that

$$
\int_{M \times M}(\eta(x), k(x, z) \xi(z))_{E_{x}} d x d z=\left(A A^{*} \eta, \xi\right)_{L^{2}(E)}=\left(\eta, A A^{*} \xi\right)_{L^{2}(E)}
$$

Since $\eta$ is arbitrary we conclude that

$$
A A^{*} \xi(x)=\int_{M} k(x, z) \xi(z) d z
$$

Using the above results, one can show that $f(D)$ has a smooth kernel for any function $f: \mathbb{R} \rightarrow \mathbb{C}$ which has rapid decrease. To see this, by writing $f$ in its real and imaginary parts, we may assume that $f$ is real valued. Furthermore, by decomposing $f$ into its positive and negative parts we may assume that $f \geq 0$. Let $g=f^{1 / 2}$, a function with rapid decrease still, we see that $g(D)$ is a self-adjoint smoothing operator. Therefore $f(D)=g^{2}(D)$ has a smooth integral kernel. In this way we find that $e^{-t D^{2} / 2}$ has a smooth integral kernel. Let $k_{t}(x, y)=e^{-t D^{2} / 2}(x, y)$ denote the smooth kernel.

Proposition 66.9. The function $k_{t}(x, y) \rightarrow 0$ in $C^{\infty}$ as $t \rightarrow 0$ off of any neighborhood of the diagonal $x=y$ in $M \times M$.

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Proof. Let $\delta>0$ be given, and let $\phi$ and $\psi$ be smooth functions on $\mathbb{R}$ such that $\phi+\psi=1$, the support of $\phi$ is contained in $(-\delta, \delta)$ and $\psi$ is supported in $\{x:|x| \geq \delta / 2\}$. Also let $p_{t}(\lambda)=(2 \pi t)^{-1 / 2} e^{-\lambda^{2} / 2 t}$. Then

$$
e^{-t D^{2} / 2}=\int_{\mathbb{R}} p_{t}(\lambda) e^{i \lambda D} d \lambda=\int_{\mathbb{R}} p_{t}(\lambda) \phi(\lambda) e^{i \lambda D} d \lambda+\int_{\mathbb{R}} p_{t}(\lambda) \psi(\lambda) e^{i \lambda D} d \lambda .
$$

This can be written as $e^{-t D^{2} / 2}=h_{t}(D)+g_{t}(D)$, where $h_{t}(\xi):=$ $\int_{\mathbb{R}} p_{t}(\lambda) \phi(\lambda) e^{i \lambda \xi} d \lambda$ and $g_{t}(\xi)=\int_{\mathbb{R}} p_{t}(\lambda) \psi(\lambda) e^{i \lambda \xi} d \lambda$. Now we notice that

$$
\begin{aligned}
\left|\xi^{n} g_{t}(\xi)\right| & =\left|i^{-n} \int_{\mathbb{R}} p_{t}(\lambda) \psi(\lambda) \partial_{\lambda}^{n} e^{i \lambda \xi} d \lambda\right| \\
& =\left|i^{n} \int_{\mathbb{R}} e^{i \lambda \xi} \partial_{\lambda}^{n}\left(p_{t}(\lambda) \psi(\lambda)\right) d \lambda\right| \\
& \leq K_{n}(t)
\end{aligned}
$$

where $\lim _{t \downarrow 0} K_{n}(t)=0$ for each $n$. From this it follows for any $n$ that $D^{n} g_{t}(D): L^{2}(E) \rightarrow L^{2}(E)$ tends to zero in the operator norm as $t \rightarrow 0$. This fact, elliptic regularity, and the Sobolev embedding theorems implies that $g_{t}(D): L^{2}(E) \rightarrow C^{r}(E)$ tends to zero in operator norm for any $r \geq 0$. Using the previous proposition, this shows that the integral kernel of $g_{t}(D)$ goes to 0 in $C^{\infty}$. (Note, Roe proves some of this by appealing to the closed graph theorem for Frechet spaces.) Finally, $h_{t}(D)=\int_{\mathbb{R}} p_{t}(\lambda) \phi(\lambda) e^{i \lambda D} d \lambda$ is an operator which does not increase the support of a section by more than size $\delta$. This implies that the support of the integral kernel of $h_{t}(D)$ is contained $\{(x, y): d(x, y)<3 \delta\}$. Since $\delta$ is arbitrary, we are done.

### 66.2 Asymptotics by Sobolev Theory

Let me end this section by explaining how Roe shows that the formal asymptotic expansions of the heat kernel are close to the heat kernel.

Let $\xi \in E_{m}$ and let $w_{t}(x):=e^{t L}(x, m) \xi$, then

$$
\begin{equation*}
\left(\partial_{t}+D^{2}\right) w_{t}=0 \text { and } \lim _{t \downarrow 0} w_{t}=\xi \delta_{m} . \tag{66.8}
\end{equation*}
$$

Conversely, if $w_{t}$ solves Eq. (66.8), then for all smooth sections $S$ of $E$, $\partial_{t}\left(e^{-(T-t) D^{2}} S, w_{t}\right)=0$. Therefore

$$
\begin{aligned}
\left(S, w_{T}\right) & =\lim _{t \downarrow 0}\left(e^{-(T-t) D^{2}} S, w_{t}\right)=\left(\left(e^{-T D^{2}} S\right)(m), \xi\right) \\
& =\int\left(S(x),\left[e^{-T D^{2}}(m, x)\right]^{*} \xi\right) d x=\int\left(S(x), e^{-T D^{2}}(x, m) \xi\right) d x
\end{aligned}
$$

which shows that $w_{t}(x):=e^{-t D^{2}}(x, m) \xi$. since $S$ is arbitrary.

Now suppose that $S_{t}$ is an approximate fundamental solution at $\xi_{m} \in E_{m}$, so that

$$
\left(\partial_{t}+D^{2}\right) S_{t}=t^{k} r_{t} \text { and } \lim _{t \downarrow 0} S_{t}=\xi \delta_{m}
$$

where $r_{t}$ is a smooth section of $\Gamma(E)$. Let $\alpha_{t}$ denote the solution to

$$
\left(\partial_{t}+D^{2}\right) \alpha_{t}=t^{k} r_{t} \text { with } \alpha_{0}=0
$$

which can be written by du Hamel's principle as $\alpha_{t}=\int_{0}^{t} e^{-(t-\tau) D^{2}} \tau^{k} r_{\tau} d \tau$. Since $w_{t}:=S_{t}-\alpha_{t}$ satisfies

$$
\left(\partial_{t}+D^{2}\right) w_{t}=0 \text { with } \lim _{t \downarrow 0} w_{t}=\xi \delta_{m}
$$

we find the $w_{t}=e^{-t D^{2}}(m, \cdot)$. Therefore, for any $k>n / 2$,

$$
\begin{aligned}
\left|e^{-t D^{2}}(m, \cdot)-S_{t}\right| & =\left|\alpha_{t}\right|=\left|\int_{0}^{t} e^{-(t-\tau) D^{2}} \tau^{k} r_{\tau} d \tau\right| \\
& \leq C_{k}\left\|\int_{0}^{t} e^{-(t-\tau) D^{2}} \tau^{k} r_{\tau} d \tau\right\|_{L_{k}^{2}(E)} \\
& \leq C_{k}\left\|D^{k} \int_{0}^{t} e^{-(t-\tau) D^{2}} \tau^{k} r_{\tau} d \tau\right\|_{L^{2}(E)}+C_{k} \| \int_{0}^{t} e^{-(t-\tau) D^{2}} \tau^{k} r_{\tau} a
\end{aligned}
$$

wherein the second to last inequality we have used the Sobolev embedding theorem and in the last we use elliptic regularity. Since $e^{-t D^{2}}$ is a bounded operator, we find that

$$
\begin{aligned}
\left|e^{-t D^{2}}(m, \cdot)-S_{t}\right| & \leq C_{k}\left\|\int_{0}^{t} e^{-(t-\tau) D^{2}} \tau^{k} D^{k} r_{\tau} d \tau\right\|_{L^{2}(E)}+C_{k} \| \int_{0}^{t} e^{-(t-\tau) D^{2}} \tau^{k} r_{\tau} d \\
& \leq C_{k} t^{k+1} \sup _{0 \leq \tau \leq t}\left\{\left\|D^{k} r_{\tau}\right\|_{L^{2}(E)}+\left\|r_{\tau}\right\|_{L^{2}(E)}\right\}=K t^{k+1}
\end{aligned}
$$

## Appendix: VanVleck Determinant Properties

### 67.1 Proof of Lemma 61.3

The first step is to get a more explicit expression for $D$. To this end, fix $v \in T_{y} M$ and for any $w \in T_{y} M$, let

$$
\begin{equation*}
Y_{w}(t):=\left.\frac{d}{d s}\right|_{0} \exp (t(v+s w))=\exp _{*}(t w)_{t v} \tag{67.1}
\end{equation*}
$$

and $\gamma(t)=\exp (t v)$. Then $Y_{w}$ solves Jacobi's equation:

$$
\begin{aligned}
\frac{\nabla^{2}}{d t^{2}} Y_{w}(t) & =\left.\frac{\nabla^{2}}{d t^{2}} \frac{d}{d s}\right|_{0} \exp (t(v+s w)) \\
& =\left.\frac{\nabla}{d t} \frac{\nabla}{d s}\right|_{0} \frac{d}{d t} \exp (t(v+s w)) \quad \text { (No Torsion) } \\
& =\left[\frac{\nabla}{d t},\left.\frac{\nabla}{d s}\right|_{0}\right] \frac{d}{d t} \exp (t(v+s w)) \quad\left(\frac{\nabla}{d t} \frac{d}{d t} \exp (t(v+s w))=0\right) \\
& \left.=R\left(\dot{\gamma}(t), Y_{w}(t)\right) \dot{\gamma}(t) . \quad \text { (Definition of } R\right)
\end{aligned}
$$

$Y_{w}$ obeys the initial conditions

$$
Y_{w}(0)=\left.\frac{d}{d s}\right|_{0} \exp (0(v+s w))=0
$$

and

$$
\begin{aligned}
\left.\frac{\nabla}{d t}\right|_{0} Y_{w}(0) & =\left.\left.\frac{\nabla}{d s}\right|_{0} \frac{d}{d t}\right|_{0} \exp (t(v+s w)) \\
& =\left.\frac{\nabla}{d s}\right|_{0}(v+s w)=w
\end{aligned}
$$

Notice that $Y_{w}(t) \in T_{\gamma(t)} M$ for each $t$, so by using parallel translation $u(t):=$ $/{ }_{t}(\gamma)$ along $\gamma$ we may pull this back to $T_{y} M$. Set

$$
Z_{w}(t)=u^{-1}(t) Y_{w}(t)
$$

Using the fact that $\frac{\nabla}{d t} u(t)=0$, the previous equations imply that $Z_{w}$ satisfies:

$$
\begin{equation*}
\ddot{Z}_{w}(t)=A_{\gamma}(t) Z_{w}(t) \text { with } Z_{w}(0)=0 \text { and } \dot{Z}_{w}(0)=w \tag{67.2}
\end{equation*}
$$

where

$$
\begin{align*}
A_{\gamma}(t) w & =u^{-1}(t) R(\dot{\gamma}(t), u(t) w) \dot{\gamma}(t) \\
& =u^{-1}(t) R(u(t) v, u(t) w) u(t) v \tag{67.3}
\end{align*}
$$

Since $u(t)$ is orthogonal for all $t$, we may now compute $D(v)$ as

$$
\begin{equation*}
D(v)=\lambda_{y}\left(Z_{w_{1}}(1), \ldots, Z_{w_{n}}(1)\right) / \lambda_{y}\left(w_{1}, \ldots, w_{n}\right)=\operatorname{det} Z_{\gamma}(1) \tag{67.4}
\end{equation*}
$$

where $Z_{\gamma}$ is the matrix solution to the differential equation

$$
\begin{equation*}
\ddot{Z}_{\gamma}(t)=A_{\gamma}(t) Z_{\gamma}(t) \text { with } Z_{\gamma}(0)=0 \text { and } \dot{Z}_{\gamma}(0)=I \tag{67.5}
\end{equation*}
$$

By Taylor's theorem,

$$
\begin{equation*}
Z_{\gamma}(1)=I+\frac{1}{2} \ddot{Z}_{\gamma}(0)+\frac{1}{6} Z_{\gamma}^{(3)}(0)+\int_{0}^{1} Z_{\gamma}^{(4)}(t) d \mu(t) \tag{67.6}
\end{equation*}
$$

where $\mu$ is a positive measure such that $\mu([0,1])=1 / 4$ !. Now from the differential equation $\ddot{Z}_{\gamma}(0)=0$,

$$
\begin{aligned}
& Z_{\gamma}^{(3)}(t)=\dot{A}_{\gamma}(t) Z_{\gamma}(t)+A_{\gamma}(t) \dot{Z}_{\gamma}(t) \text { and } \\
& Z_{\gamma}^{(4)}(t)=\ddot{A}_{\gamma}(t) Z_{\gamma}(t)+2 \dot{A}_{\gamma}(t) \dot{Z}_{\gamma}(t)+A_{\gamma}(t) \ddot{Z}_{\gamma}(t)
\end{aligned}
$$

In particular $Z_{\gamma}^{(3)}(0)=A_{\gamma}(0)=R(v, \cdot) v$, and

$$
Z_{\gamma}^{(4)}(t)=\ddot{A}_{\gamma}(t) Z_{\gamma}(t)+2 \dot{A}_{\gamma}(t) \dot{Z}_{\gamma}(t)+A_{\gamma}(t) A_{\gamma}(t) Z_{w}(t)
$$

Now $A_{\gamma}(t)=O\left(v^{2}\right)$,

$$
\begin{gathered}
\dot{A}_{\gamma}(t)=u^{-1}(t)\left(\nabla_{u(t) v} R\right)(u(t) v, u(t) w) u(t) v=O\left(v^{3}\right), \\
\ddot{A}_{\gamma}(t)=u^{-1}(t)\left(\nabla_{u(t) v \otimes u(t) v}^{2} R\right)(u(t) v, u(t) w) u(t) v=O\left(v^{4}\right)
\end{gathered}
$$

and hence $Z_{\gamma}^{(4)}(t)=O\left(v^{3}\right)$. Using these estimates in Eq. (67.6) shows that

$$
\begin{equation*}
Z_{\gamma}(1)=I+\frac{1}{6} R(v, \cdot) v+O\left(v^{3}\right) . \tag{67.7}
\end{equation*}
$$

Taking the determinant of this equation shows that

$$
\begin{aligned}
D(v) & =1+\frac{1}{6} \operatorname{tr}(w \rightarrow R(v, w) v)+O\left(v^{3}\right) \\
& =1-\frac{1}{6}(\operatorname{Ric} v, v)+O\left(v^{3}\right) .
\end{aligned}
$$

Before finishing this section, let us write out Eq. (67.7) in detail.

Lemma 67.1. Let $v, w \in T_{y} M$, then

$$
\begin{equation*}
/ /_{1}^{-1}\left(\gamma_{v}\right) \exp _{*}\left(w_{v}\right)=w+\frac{1}{6} R(v, w) v+O\left(v^{3}\right) w \tag{67.8}
\end{equation*}
$$

In particular we have for $v, w, u \in T_{y} M$ that

$$
\begin{align*}
\left(\exp _{*}\left(w_{v}\right),\right. & \left.\exp _{*}\left(w_{v}\right)\right) \\
& =\left(w+\frac{1}{6} R(v, w) v+O\left(v^{3}\right) w, u+\frac{1}{6} R(v, u) v+O\left(v^{3}\right) u\right) \\
& =(w, u)+\frac{1}{6}(w, R(v, u) v)+\frac{1}{6}(R(v, w) v, u)+O\left(v^{3}\right)(w, u) \\
& =(w, u)-\frac{1}{3}(R(w, v) v, u)+O\left(v^{3}\right)(w, u) \tag{67.9}
\end{align*}
$$

### 67.2 Another Proof of Remark 61.2: The Symmetry of $J(x, y)$.

Recall that $\lambda$ denotes the Riemannian volume form on $M$ and

$$
J(x, y)=\left(\exp _{x}^{*} \lambda\right)_{v} / \lambda_{x}
$$

where $\exp _{x}(v)=y$. Also recall that $D(v):=\left(\exp _{x}^{*} \lambda\right)_{v} / \lambda_{x}$ where $x=\pi(v)$ and $\pi: T M \rightarrow M$ is the canonical projection map. The precise meaning of this equation is, given any basis $\left\{w_{i}\right\}_{i=1}^{n}$ for $T_{x} M$, then

$$
D(v)=\lambda\left(\exp _{*}\left(w_{1}\right)_{v}, \ldots, \exp _{*}\left(w_{n}\right)_{v}\right) / \lambda_{x}\left(w_{1}, \ldots, w_{n}\right)=\operatorname{det}\left(Z_{\gamma}(1)\right)
$$

where $w_{v}=\left.\frac{d}{d s}\right|_{0}(v+s w) \in T_{v} T_{x} M$ and $Z_{\gamma}$ is defined in Eq. (67.2) above. In particular,

$$
\exp _{*} w_{v}:=\left.\frac{d}{d s}\right|_{0} \exp (v+s w) .
$$

Notice that

$$
J(x, y)=D\left(\exp _{x}^{-1}(y)\right)
$$

Let $i: T M \rightarrow T M$ denote the involution given by $i(v)=-\dot{\gamma}(1)$, where $\gamma(t)=\exp (t v)$ is the geodesic determined by $v$. Alternatively we may describe $i(v)=-/ / 1(\gamma) v$. Now if $v=\exp _{x}^{-1}(y)$, i.e. $y=\exp _{x}(v)=\exp (v)$, then $\exp (i(v))=x$. That is to say, $i(v)=\exp _{y}^{-1}(x)$. Hence to show $J(x, y)=J(y, x)$ if and only if $D(v)=D(i(v))$. This is what is proved in A. L. Bess, "Manifolds all of whose Geodesics are Closed," see Lemma 6.12 on p. 156.

Now let us work out $D(i(v))$. Let $\sigma(t)=\gamma(1-t)=\exp (t i(v))$. Since $/ / t(\sigma)=/ /{ }_{1-t}(\gamma) / /_{1}(\gamma)^{-1}$, it follows after a short calculation that $A_{\sigma}(t)=$ $u(1) A_{\gamma}(1-t) u(1)^{-1}$. Let $W(t):=u(1)^{-1} Z_{\sigma}(1-t) u(1)$, then $W(1)=0$, $\dot{W}(1)=-I$ and

$$
\ddot{W}(t)=u(1)^{-1} A_{\sigma}(1-t) Z_{\sigma}(1-t) u(1)=W(t)=A_{\gamma}(t) W(t) .
$$

Notice that

$$
D(i(v))=\operatorname{det} Z_{\sigma}(1)=\operatorname{det}\left[u(1)^{-1} Z_{\sigma}(1) u(1)\right]=\operatorname{det} W(0) .
$$

So to finish the proof, we must show that $\operatorname{det} W(0)=\operatorname{det} Z(1)$. For this observe that $A_{\gamma}(t)$ is a symmetric operator (by symmetry properties of the curvature tensor) and hence

$$
\begin{aligned}
\frac{d}{d t}\left\{\dot{Z}^{*}(t) W(t)-Z^{*}(t) \dot{W}(t)\right\} & =\ddot{Z}^{*}(t) W(t)-Z^{*}(t) \ddot{W}(t) \\
& =Z^{*}(t) A_{\gamma}(t) W(t)-Z^{*}(t) A_{\gamma}(t) W(t)=0
\end{aligned}
$$

and hence

$$
\left.\left\{\dot{Z}^{*}(t) W(t)-Z^{*}(t) \dot{W}(t)\right\}\right|_{0} ^{1}=0 .
$$

This implies that

$$
W(0)=\dot{Z}^{*}(0) W(0)-Z^{*}(0) \dot{W}(0)=\dot{Z}^{*}(1) W(1)-Z^{*}(1) \dot{W}(1)=Z^{*}(1)
$$

Therefore $\operatorname{det} W(0)=\operatorname{det} Z^{*}(1)=\operatorname{det} Z(1)$ as desired.

### 67.3 Normal Coordinates

Notation 67.2 Suppose that $o \in M$ is given and let $x(m):=\exp _{o}^{-1}(m) \in$ $T_{o} M$ for $m$ in a neighborhood of $o$. The chart $x$ is called a geodesic normal coordinate system near o.

In geodesic coordinates, $t \rightarrow t x$ is a geodesic, therefore if $\Gamma$ is the Christofel symbols in the this coordinate system, we have

$$
0=\frac{\nabla}{d t} \frac{d}{d t}(t x)=\left(\frac{d}{d t}+\Gamma(t x)\langle x\rangle\right) \frac{d}{d t}(t x)=\Gamma(t x)\langle x\rangle x
$$

for all $x$ near 0 . Since $\nabla$ has zero Torsion we also have that

$$
\Gamma(z)\langle x\rangle y=\Gamma(z)\langle y\rangle x
$$

for all $x, y, z$. From the previous two equations it follows that

$$
0=\Gamma(0)\langle x\rangle y \text { for all } x, y,
$$

i.e. that $\Gamma(0)=0$ and that

$$
\partial_{x} \Gamma(0)\langle x\rangle x=0 .
$$

Let $B(x, y, z):=\partial_{x} \Gamma(0)\langle y\rangle z$, then we have shown that

$$
\begin{equation*}
B(x, y, z)=B(x, z, y) \text { and } B(x, x, x)=0 \text { for all } x, y, z \tag{67.10}
\end{equation*}
$$

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Thus

$$
0=\left.\frac{d}{d t}\right|_{0} B(x+t y, x+t y, x+t y)=B(y, x, x)+2 B(x, x, y)
$$

and therefore,

$$
\begin{aligned}
0 & =\left.\frac{d}{d t}\right|_{0}\{B(y, x+t z, x+t z)+2 B(x+t z, x+t z, y)\} \\
& =2 B(y, x, z)+2 B(z, x, y)+2 B(x, z, y) \\
& =2(B(y, z, x)+B(z, x, y)+B(x, y, z))
\end{aligned}
$$

wherin the last equality we use Eq. (67.10). Hence we have shown that

$$
\begin{equation*}
\partial_{x} \Gamma(0)\langle y\rangle z+\text { cyclic }=0 \tag{67.11}
\end{equation*}
$$

So at $x=0$ the curvature tensor is given by $R=d \Gamma$ and hence

$$
\begin{aligned}
R(x, y) x & =\partial_{x} \Gamma(0)\langle y\rangle x-\partial_{y} \Gamma(0)\langle x\rangle x \\
& =\partial_{x} \Gamma(0)\langle y\rangle x+\partial_{x} \Gamma(0)\langle y\rangle x+\partial_{x} \Gamma(0)\langle x\rangle y \\
& =3 \partial_{x} \Gamma(0)\langle y\rangle x=3 \partial_{x} \Gamma(0)\langle x\rangle y
\end{aligned}
$$

and hence

$$
\begin{aligned}
\partial_{x} \Gamma(0)\langle x\rangle y & =\frac{1}{3} R(x, y) x \text { or } \\
\Gamma(x)\langle x\rangle y & =\frac{1}{3} R(x, y) x+O\left(x^{2}\right)\langle x, y\rangle
\end{aligned}
$$

Therefore, if

$$
\begin{aligned}
\left.\partial_{z}^{2}(x, y)\right|_{z=0} & =\left.\partial_{z}(\Gamma\langle z\rangle x, y)\right|_{z=0}+\left.\partial_{z}(x, \Gamma\langle z\rangle y)\right|_{z=0} \\
& =\frac{1}{3}(R(z, x) z, y)+\frac{1}{3}(x, R(z, y) z) \\
& =\frac{2}{3}(R(z, x) z, y)
\end{aligned}
$$

and therefore by Taylor's theorem we learn that

$$
\begin{aligned}
\left(x_{z}, y_{z}\right) & =(x, y)+\frac{1}{2!} \frac{2}{3}(R(z, x) z, y)+O\left(z^{3}\right) \\
& =(x, y)-\frac{1}{3}(R(x, z) z, y)+O\left(z^{3}\right)
\end{aligned}
$$

and hence we have reproved Lemma 67.1
We now change notation a bit. Let $E \rightarrow M$ be a vector bundle with connection $\nabla$.

Notation 67.3 Let $x$ be a chart on $M$ such that $x(o)=0$ and let $\gamma_{v}(t)=$ $x^{-1}(t v)$ for all $v \in \mathbb{R}^{\operatorname{dim}(M)}=: V$. Let $u u: \mathcal{D}(x) \rightarrow G l(E)$ be the local orthonormal frame given by

$$
\begin{equation*}
u(m)=/ / 1\left(\gamma_{x(m)}\right) \text {, i.e. } u\left(x^{-1}(t v)\right)=/ / t\left(\gamma_{v}\right): E_{o} \rightarrow E_{m} \tag{67.12}
\end{equation*}
$$

for all $v \in V$ sufficiently small. Also let $\Gamma=\Gamma^{u}=u^{-1} \nabla u$ be the associated connection one form.

From Eq. (67.12) it follows that

$$
\frac{\nabla}{d t} u\left(x^{-1}(t v)\right)=\frac{\nabla}{d t} / / t\left(\gamma_{v}\right)=0
$$

and in particular at $t=0$ this shows that

$$
0=\left.v^{i} \nabla_{\partial_{i}} u\right|_{\gamma_{v}(t)}=v^{i} u\left(\gamma_{v}(t)\right) \Gamma\left(\left.\partial_{i}\right|_{\gamma_{v}(t)}\right)=u\left(\gamma_{v}(t)\right) \Gamma\left(\left.v^{i} \partial_{i}\right|_{\gamma_{v}(t)}\right) .
$$

That is to say

$$
\Gamma\left(\left.v^{i} \partial_{i}\right|_{\gamma_{v}(t)}\right)=0 \text { for all } v \in V
$$

In particular at $t=0$ we learn that $\Gamma\left(\left.\partial_{i}\right|_{o}\right)=0$ and

$$
0=\left.\frac{d}{d t}\right|_{0} \Gamma\left(\left.v^{i} \partial_{i}\right|_{\gamma_{v}(t)}\right)=\left.v^{i} v^{j} \partial_{j} \Gamma\left(\partial_{i}\right)\right|_{o} \text { for all } v \in V
$$

This shows that $\left.\partial_{j} \Gamma\left(\partial_{i}\right)\right|_{o}=-\left.\partial_{i} \Gamma\left(\partial_{j}\right)\right|_{o}$. Since

$$
\begin{aligned}
\left.R^{E}\left(\partial_{i}, \partial_{j}\right)\right|_{o} & =\left.u^{-1} \nabla_{\partial_{i} \wedge \partial_{j}}^{2} u\right|_{o} \\
& =\left.u^{-1}(o)\left\{\partial_{i} \Gamma\left(\partial_{j}\right)-\partial_{j} \Gamma\left(\partial_{i}\right)+\left[\Gamma\left(\partial_{i}\right), \Gamma\left(\partial_{j}\right)\right]\right\}\right|_{o} \\
& =-\left.2 \partial_{j} \Gamma\left(\partial_{i}\right)\right|_{o}
\end{aligned}
$$

from which it follows that

$$
\left.\partial_{j} \Gamma\left(\partial_{i}\right)\right|_{o}=-\left.\frac{1}{2} R^{E}\left(\partial_{i}, \partial_{j}\right)\right|_{o}
$$

From Taylor's theorem we find
$\left.\Gamma\left(\partial_{i}\right)\right|_{m}=-\left.\frac{1}{2} x^{j}(m) R^{E}\left(\partial_{i}, \partial_{j}\right)\right|_{o}+o\left(|x(m)|^{2}\right)=-\left.\frac{1}{2} R^{E}\left(\partial_{i}, x^{j}(m) \partial_{j}\right)\right|_{o}+o\left(|x(m)|^{2}\right)$
This result is summarized as follows.
Proposition 67.4. Keeping the notation as above and let $w \in \mathbb{R}^{\operatorname{dim}(M)}$, then

$$
\begin{equation*}
\Gamma\left(w^{i} \partial_{i}\right)=-\frac{1}{2} R^{E}\left(\left.w^{i} \partial_{i}\right|_{o},\left.x^{j} \partial_{j}\right|_{o}\right)+o\left(x^{2}\right)(w) \tag{67.13}
\end{equation*}
$$

near $o \in M$. In particular if $x(m)=\exp _{o}^{-1}(m)$ are normal coordinates on $M$, then

$$
\begin{aligned}
\left.w^{i} \partial_{i}\right|_{m} & =\left.\frac{d}{d t}\right|_{0} x^{-1}(x(m)+t w)=\left.\frac{d}{d t}\right|_{0} \exp (x(m)+t w) \\
& =\exp _{*}\left(w_{x(m)}\right)
\end{aligned}
$$

Therefore, Eq. (67.13) may be written as

$$
\Gamma\left(\exp _{*}\left(w_{x(m)}\right)\right)=-\frac{1}{2} R^{E}(w, x(m))+O\left(x^{2}(m)\right)(w)
$$

Proof. The quick proof of these results is as follows. We work in the local frame $u$. Write $/{ }_{t}(\sigma)=u(\sigma(t)) P_{t}(\sigma)$ and recall the formula

$$
\left.\frac{d}{d s}\right|_{0} P_{t}\left(\Sigma_{s}\right)=P_{t}\left(\Sigma_{s}\right) \int_{0}^{t} R_{/ / \tau\left(\Sigma_{s}\right)}^{E}\left(\dot{\Sigma}_{s}(\tau), \Sigma_{s}^{\prime}(\tau)\right) d \tau-\Gamma\left(\Sigma_{s}^{\prime}(t)\right) P_{t}\left(\Sigma_{s}\right)
$$

Apply this to $\Sigma_{s}(t)=x^{-1}(t(x(m)+s v))$ and use that fact that in the frame defined by $u, P_{t}\left(\Sigma_{s}\right)=i d$ so that

$$
0=\int_{0}^{t} R_{/ / \tau\left(\Sigma_{0}\right)}^{E}\left(\dot{\Sigma}_{0}(\tau), \Sigma_{0}^{\prime}(\tau)\right) d \tau-\Gamma\left(\Sigma_{0}^{\prime}(t)\right)
$$

Therefore, at $t=1$,

$$
\begin{aligned}
\Gamma\left(\left.v^{i} \partial_{i}\right|_{m}\right) & =\int_{0}^{1} R_{/ / \tau\left(\Sigma_{0}\right)}^{E}\left(\dot{\Sigma}_{0}(\tau), \Sigma_{0}^{\prime}(\tau)\right) d \tau \\
& =-\int_{0}^{1} R_{/ / \tau\left(\Sigma_{0}\right)}^{E}\left(\left.\tau v^{i} \partial_{i}\right|_{x^{-1}(\tau x(m))},\left.x^{j}(m) \partial_{j}\right|_{x^{-1}(\tau x(m))}\right) d \tau \\
& =-\int_{0}^{1} R_{/ / \tau\left(\Sigma_{0}\right)}^{E}\left(\left.v^{i} \partial_{i}\right|_{x^{-1}(\tau x(m))},\left.x^{j}(m) \partial_{j}\right|_{x^{-1}(\tau x(m))}\right) \tau d \tau \\
& =-\frac{1}{2} R^{E}\left(\left.v^{i} \partial_{i}\right|_{o},\left.x^{j}(m) \partial_{j}\right|_{o}\right)+O\left(x^{2}(m)\right) \\
& =-\frac{1}{2} R^{E}\left(\left.v^{i} \partial_{i}\right|_{o}, x(m)\right)+O\left(x^{2}(m)\right)(v)
\end{aligned}
$$

## Miscellaneous

### 68.1 Jazzed up version of Proposition 68.1

Proposition 68.1. Let $\alpha>-1, R, K, P, Q$ and $V$ be as above. Then the series in Eq. (60.5) and Eq. (60.11) is convergent and is equal to $P_{t}=e^{t L_{-}}$ the unique solution to Eq. (60.1). Moreover,

$$
\begin{equation*}
\|P-K\|_{t} \leq \kappa(t)\|K\|_{t} t^{\alpha+1}=O\left(t^{1+\alpha}\right) \tag{68.1}
\end{equation*}
$$

where $\|f\|_{t}:=\max _{0 \leq s \leq t}\left|f_{s}\right|$ and $\kappa(t)$ is an increasing function of $t$, see Eq. (68.4).

Proof. By making the change of variables $r=s+u(t-s)$ we find that

$$
\begin{equation*}
\int_{s}^{t}(t-r)^{\alpha}(r-s)^{\beta} d r=C(\alpha, \beta)(t-s)^{\alpha+\beta+1} \tag{68.2}
\end{equation*}
$$

where

$$
\begin{equation*}
C(\alpha, \beta):=\int_{0}^{1} u^{\alpha}(1-u)^{\beta} d u=B(\alpha+1, \beta+1) \tag{68.4}
\end{equation*}
$$

and $B$ is the beta function. From Eq. (1.5.5) of Lebedev, "Special Functions and Their Applications", p.13,

$$
\begin{equation*}
C(\alpha, \beta)=B(\alpha+1, \beta+1)=\frac{\Gamma(\alpha+1) \Gamma(\beta+1)}{\Gamma(\alpha+\beta+2)} . \tag{68.3}
\end{equation*}
$$

(See below for a proof of Eq. (68.3).)
By repeated use of Eq. (68.2),

$$
\begin{aligned}
\left|\left(Q^{m} R\right)_{t}\right| & =\left|\int_{t \Delta_{m}} R_{t-s_{m}} R_{s_{m}-s_{m-1}} \ldots R_{s_{2}-s_{1}} R_{s_{1}} d \mathbf{s}\right| \\
& \leq C^{m} \int_{t \Delta_{m}}\left(t-s_{m}\right)^{\alpha}\left(s_{m}-s_{m-1}\right)^{\alpha} \ldots\left(s_{2}-s_{1}\right)^{\alpha} s_{1}^{a} d \mathbf{s} \\
& =C^{m} C(\alpha, \alpha) \int_{t \Delta_{m-1}}\left[\begin{array}{c}
\left(t-s_{m-1}\right)^{2 \alpha+1}\left(s_{m-1}-s_{m-2}\right)^{\alpha} \times \\
\ldots\left(s_{2}-s_{1}\right)^{\alpha} s_{1}^{a}
\end{array}\right] d \mathbf{s} \\
& =C^{m} C(\alpha, \alpha) C(\alpha, 2 \alpha+1) C(\alpha, 3 \alpha+2) \ldots \\
& \ldots C(\alpha,(m-1) \alpha+m-2) \int_{0}^{t}\left(t-s_{1}\right)^{m \alpha+m-1} d s_{1} \\
& =C^{m} C(\alpha, \alpha) C(\alpha, 2 \alpha+1) C(\alpha, 3 \alpha+2) \ldots \\
& \ldots C(\alpha,(m-1) \alpha+m-2) C(\alpha, m \alpha+m-1) t^{(m+1) \alpha+m}
\end{aligned}
$$

Now from Eq. (68.3) we find that

$$
\begin{aligned}
& C(\alpha, \alpha) C(\alpha, 2 \alpha+1) C(\alpha, 3 \alpha+2) \ldots C(\alpha, m \alpha+m-1) \\
&= \frac{\Gamma(\alpha+1) \Gamma(\alpha+1)}{\Gamma(2 \alpha+2)} \frac{\Gamma(\alpha+1) \Gamma(2 \alpha+2)}{\Gamma(3 \alpha+3)} \times \\
& \frac{\Gamma(\alpha+1) \Gamma(3 \alpha+3)}{\Gamma(4 \alpha+4)} \cdots \frac{\Gamma(\alpha+1) \Gamma(m \alpha+1)}{\Gamma((m+1) \alpha+m+1)} \\
&= \frac{\Gamma(\alpha+1)^{m}}{\Gamma((m+1) \alpha+m+1)} .
\end{aligned}
$$

Therefore,

$$
\left|\left(Q^{m} R\right)_{t}\right| \leq C^{m} \frac{\Gamma(\alpha+1)^{m}}{\Gamma((m+1) \alpha+m+1)} t^{\alpha} t^{m(\alpha+1)}
$$

and thus the series in Eq. (60.11) is absolutely convergent and $\left|V_{t}\right| \leq(a+$ 1) $\kappa(t) t^{\alpha}$ where

$$
\kappa(t)=(\alpha+1)^{-1} \sum_{m=0}^{\infty} C^{m} \frac{\Gamma(\alpha+1)^{m}}{\Gamma((m+1) \alpha+m+1)} t^{m(\alpha+1)}
$$

which is seen to be finite by Stirlings formula,

$$
\Gamma(m(\alpha+1)+1) \backsim(2 \pi)^{1 / 2} e^{-(m(\alpha+1)+1)}(m(\alpha+1)+1)^{(m(\alpha+1)+1 / 2)},
$$

see Eq. (1.4.12) of Lebedev.
Using the bound on $V$ and the uniform boundedness of $K_{t}$,

$$
\begin{equation*}
\int_{0}^{t}\left|K_{s} V_{t-s}\right| d s \leq \kappa(t)\|K\|_{t}(\alpha+1) \int_{0}^{t}(t-s)^{\alpha} d s=\kappa(t)\|K\|_{t} t^{\alpha+1} \tag{68.5}
\end{equation*}
$$

and hence $P_{t}$ defined in Eq. (60.10) is well defined and is continuous in $t$. Moreover, (68.5) implies Eq. (68.1) once we shows that $P_{t}=e^{t L}$. This is checked as follows,
68.1 Jazzed up version of Proposition 68.11193

$$
\begin{aligned}
\frac{d}{d t} \int_{0}^{t} K_{t-s} V_{s} d s & =V_{t}+\int_{0}^{t} \dot{K}_{t-s} V_{s} d s=V_{t}+\int_{0}^{t}\left(L K_{t-s}-R_{t-s}\right) V_{s} d s \\
& =V_{t}+L \int_{0}^{t} K_{t-s} V_{s} d s-(Q V)_{t}=L \int_{0}^{t} K_{t-s} V_{s} d s+R_{t}
\end{aligned}
$$

Thus we have,

$$
\begin{aligned}
\frac{d}{d t} P_{t} & =\dot{K}_{t}+L \int_{0}^{t} K_{t-s} V_{s} d s+R_{t} \\
& =L K_{t}+L \int_{0}^{t} K_{t-s} V_{s} d s=L P_{t}
\end{aligned}
$$

### 68.1.1 Proof of Eq. (68.3)

Let us recall that $\Gamma(x):=\int_{0}^{\infty} t^{x} e^{-t} d t / t$ and hence let $z=x+y$ and then $x=u z$ we derive,

$$
\begin{aligned}
\Gamma(\alpha+1) \Gamma(\beta+1) & =\int_{[0, \infty)^{2}} x^{\alpha} y^{\beta} e^{-(x+y)} d x d y \\
& =\int_{0}^{\infty} d x \int_{x}^{\infty} d z x^{\alpha}(z-x)^{\beta} e^{-z} \\
& =\int d x d z 1_{0<x<z<\infty} x^{\alpha}(z-x)^{\beta} e^{-z} \\
& =\int_{0}^{\infty} d z \int_{0}^{1} d u u^{\alpha}(1-u)^{\beta} z^{\alpha+\beta+1} e^{-z} \\
& =C(\alpha, \beta) \Gamma(\alpha+\beta+2),
\end{aligned}
$$

i.e.

$$
C(\alpha, \beta)=\frac{\Gamma(\alpha+1) \Gamma(\beta+1)}{\Gamma(\alpha+\beta+2)} .
$$

### 68.1.2 Old proof of Proposition 60.1

Proof. Taking norms of Eq. (60.8) shows that
$\left\|Q^{m} K\right\|_{t}=\|K\|_{t} C^{m} \int_{t \Delta_{m}}\left(s_{2}-s_{1}\right)^{\alpha}\left(s_{3}-s_{2}\right)^{\alpha} \ldots\left(s_{m}-s_{m-1}\right)^{\alpha}\left(t-s_{m}\right)^{\alpha} d \mathbf{s}$, where $d \mathbf{s}=d s_{1} d s_{2} \ldots d s_{m}$. To evaluate this last integral, we will make use Eq. (68.2) to find of Repeated use of Eq. (68.2) gives,

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$\int_{t \Delta_{m}}\left(s_{2}-s_{1}\right)^{\alpha}\left(s_{3}-s_{2}\right)^{\alpha} \ldots\left(s_{m}-s_{m-1}\right)^{\alpha}\left(t-s_{m}\right)^{\alpha} d \mathbf{s}$

$$
\begin{aligned}
& =C(\alpha, \alpha) \int_{t \Delta_{m-1}}\left[\begin{array}{c}
\left(s_{2}-s_{1}\right)^{\alpha}\left(s_{3}-s_{2}\right)^{\alpha} \times \\
\ldots\left(s_{m-1}-s_{m-2}\right)^{\alpha}\left(t-s_{m-1}\right)^{2 \alpha+1}
\end{array}\right] d \mathbf{s} \\
& =C(\alpha, \alpha) C(\alpha, 2 \alpha+1) \int_{t \Delta_{m-2}}\left[\begin{array}{c}
\left(s_{2}-s_{1}\right)^{\alpha}\left(s_{3}-s_{2}\right)^{\alpha} \times \\
= \\
=C(\alpha, \alpha) C(\alpha, 2 \alpha+1) C(\alpha, 3 \alpha+2) \times
\end{array}\right. \\
& \quad \ldots C(\alpha,(m-1) \alpha+m-2) \int_{0}^{t}\left(t-s_{1}\right)^{m \alpha+m-1} d s_{1} \\
& =C(\alpha, \alpha) C(\alpha, 2 \alpha+1) C(\alpha, 3 \alpha+2) \times \\
& \quad \ldots C(\alpha,(m-1) \alpha+m-2) \frac{t^{m \alpha+m}}{m \alpha+m}
\end{aligned}
$$

Now from Eq. (68.3) we find that

$$
\begin{gathered}
C(\alpha, \alpha) C(\alpha, 2 \alpha+1) C(\alpha, 3 \alpha+2) \ldots C(\alpha,(m-1) \alpha+m-2) \\
=\frac{\Gamma(\alpha+1)^{2}}{\Gamma(2 \alpha+2)} \frac{\Gamma(\alpha+1) \Gamma(2 \alpha+2)}{\Gamma(3 \alpha+3)} \times \\
\quad \ldots \frac{\Gamma(\alpha+1) \Gamma((m-1) \alpha+m-1)}{\Gamma(m \alpha+m)} \\
=\frac{\Gamma(\alpha+1)^{m}}{\Gamma(m \alpha+m)} .
\end{gathered}
$$

Combining these results gives the estimate,

$$
\begin{aligned}
\left\|Q^{m} K\right\|_{t} & \leq\|K\|_{t} \frac{(C \Gamma(\alpha+1))^{m}}{\Gamma(m \alpha+m)} \frac{t^{m \alpha+m}}{m \alpha+m} \\
& =\|K\|_{t} \frac{\left(C \Gamma(\alpha+1) t^{\alpha+1}\right)^{m}}{\Gamma(m(\alpha+1)+1)} \\
& \backsim\|K\|_{t}\left(C \Gamma(\alpha+1) t^{\alpha+1}\right)^{m} \frac{e^{m(\alpha+1)+1}}{(2 \pi)^{1 / 2}(m(\alpha+1)+1)^{(m(\alpha+1)+1 / 2)}} \\
& =\|K\|_{t} \frac{\left(C \Gamma(\alpha+1)(e t)^{\alpha+1}\right)^{m}}{(2 \pi)^{1 / 2} e(m(\alpha+1)+1)^{(m(\alpha+1)+1 / 2)}},
\end{aligned}
$$

where the second to last expression is a result of Stirlings formula, Eq. (1.4.12) of Lebdev.

From this estimate we learn that $\sum_{m=0}^{\infty} Q^{m} K$ is uniformly convergent on compact subsets of $[0, \infty)$ and that

$$
\begin{aligned}
\left\|\sum_{m=0}^{\infty} Q^{m} K-K\right\|_{t} & \leq \sum_{m=1}^{\infty}\left\|Q^{m} K\right\|_{t} \\
& \leq\|K\|_{t} \sum_{m=1}^{\infty} \frac{\left(C \Gamma(\alpha+1) t^{\alpha+1}\right)^{m}}{\Gamma(m(\alpha+1)+1)}=O\left(t^{1+\alpha}\right)
\end{aligned}
$$

So it only remains to prove that $P:=\sum_{m=0}^{\infty} Q^{m} K$ solves (60.1).
Now by the chain rule and the fundamental theorem of calculus, $\frac{d}{d t}(Q f)_{t}=$ $\int_{0}^{t} \dot{f}_{t-s} R_{s} d s+f_{0} R_{t}$ or equivalently

$$
\frac{d}{d t}(Q f)=Q \dot{f}+f_{0} R
$$

Applying this formula inductively using the fact that $\left(Q^{m} K\right)_{0}=0$ if $m \geq 1$ implies that

$$
\begin{aligned}
\frac{d}{d t}\left(Q^{m} K\right) & =Q^{m-1} \frac{d}{d t}(Q K)=Q^{m} \dot{K}+Q^{m-1}\left(K_{0} R\right) \\
& =Q^{m}(L K-R)+Q^{m-1} R \\
& =L Q^{m} K-Q^{m} R+Q^{m-1} R,
\end{aligned}
$$

wherein the last equality we have used the fact that $L$ commutes with $Q$. Setting $P_{t}^{N}:=\sum_{m=0}^{N} Q^{m} K$, we find using the previous equation that

$$
\frac{d}{d t} P_{t}^{N}=L P_{t}^{N}-\left(Q^{N} R\right)_{t}
$$

or equivalently that

$$
\begin{equation*}
P_{t}^{N}=I+\int_{0}^{t} L P_{s}^{N} d s-\int_{0}^{t}\left(Q^{N} R\right)_{s} d s \tag{68.6}
\end{equation*}
$$

Since,

$$
\left\|Q^{N} R\right\|_{t} \leq\|R\|_{t} \frac{\left(C \Gamma(\alpha+1) t^{\alpha+1}\right)^{N}}{\Gamma(N(\alpha+1)+1)} \rightarrow 0 \text { as } N \rightarrow \infty
$$

we may pass to the limit, $N \rightarrow \infty$, in Eq. (68.6) to conclude that

$$
P_{t}=I+\int_{0}^{t} L P_{s} d s
$$

This completes the proof.
For later purposes, let us rework the above derivative aspects of the proof. Let

$$
\begin{aligned}
R_{m}(s) & :=\int_{0 \leq s_{1} \leq s_{2} \leq \cdots \leq s_{m-1} \leq s} R_{s-s_{m-1}} R_{s_{m-1}-s_{m-2}} \ldots R_{s_{2}-s_{1}} R_{s_{1}} d \mathbf{s} \\
& =\left(Q^{m-1} R\right)_{s},
\end{aligned}
$$

then by Eq. (60.6)

$$
\left(Q^{m} K\right)_{t}=\int_{0}^{t} K_{t-s} R_{m}(s) d s
$$

Hence

$$
\begin{aligned}
\frac{d}{d t}\left(Q^{m} K\right)_{t} & =R_{m}(t)+\int_{0}^{t} \frac{d}{d t} K_{t-s} R_{m}(s) d s \\
& =R_{m}(t)+\int_{0}^{t}\left(L K_{t-s}-R_{t-s}\right) R_{m}(s) d s \\
& =R_{m}(t)+L\left(Q^{m} K\right)_{t}-R_{m+1}(t) \\
& =L\left(Q^{m} K\right)_{t}+\left(Q^{m-1} R\right)_{t}-\left(Q^{m} R\right)_{t}
\end{aligned}
$$

as before.
Taking norms of this equation implies that

$$
\begin{align*}
\left|\left(Q^{m} f\right)_{t}\right| & \leq C^{m} \int_{t \Delta_{m}}\left|f_{s_{1}}\right|\left(s_{2}-s_{1}\right)^{\alpha}\left(s_{3}-s_{2}\right)^{\alpha} \ldots\left(s_{m}-s_{m-1}\right)^{\alpha}\left(t-s_{m}\right)^{\alpha} d \mathbf{s} \\
& \leq C^{m} t^{\alpha m} \int_{t \Delta_{m}}\left|f_{s_{1}}\right| d \mathbf{s}=C^{m} t^{\alpha m} \int_{0}^{t}\left|f_{s_{1}}\right| \frac{\left(t-s_{1}\right)^{m-1}}{(m-1)!} d s_{1} \\
& \leq \frac{\left(C t^{1+\alpha}\right)^{m}}{m!} \int_{0}^{t}\left|f_{s}\right| d \mu_{m}(s) \leq \frac{\left(C t^{1+\alpha}\right)^{m}}{m!} \max _{0 \leq s \leq t}\left|f_{s}\right| \tag{68.7}
\end{align*}
$$

where $d \mu_{m}(s):=m(t-s)^{m-1} t^{-m} d s$, a probability measure on $[0, t]$. This shows that $\left\|Q^{m} f\right\|_{t} \leq\left(C t^{1+\alpha}\right)^{m}\|f\|_{t} / m$ !. From this estimate we learn that $\sum_{m=0}^{\infty} Q^{m} K$ is uniformly convergent on compact subsets of $[0, \infty)$ and that

$$
\begin{aligned}
&\left\|\sum_{m=0}^{\infty} Q^{m} K-K\right\|_{t} \leq \sum_{m=0}^{\infty}\left\|Q^{m} K\right\|_{t} \leq e^{C t^{1+\alpha}}-1=O\left(t^{1+\alpha)}\right) \\
& \frac{d}{d t} \int_{0}^{t} K_{t-s} V_{s} d s=\lim _{\delta \rightarrow 0} \frac{1}{\delta}\left(\int_{0}^{t+\delta} K_{t+\delta-s} V_{s} d s-\int_{0}^{t} K_{t-s} V_{s} d s\right) \\
&=\lim _{\delta \rightarrow 0} \frac{1}{\delta}\left(\int_{t}^{t+\delta} K_{t+\delta-s} V_{s} d s+\int_{0}^{t}\left(K_{t+\delta-s}-K_{t-s}\right) V_{s} d s\right) \\
&=\lim _{\delta \rightarrow 0}\left(\frac{1}{\delta} \int_{t}^{t+\delta} K_{t-s} V_{s} d s+\frac{1}{\delta} \int_{t}^{t+\delta}\left(K_{t+\delta-s}-K_{t-s}\right) V_{s} d s\right) \\
&+\frac{d}{d \delta} \int_{0}^{t} K_{t+\delta-s} V_{s} d s \\
&=V_{t}+\int_{0}^{t} \dot{K}_{t-s} V_{s} d s=V_{t}+\int_{0}^{t}\left(L K_{t-s}-R_{t-s}\right) V_{s} d s \\
&=V_{t}+L \int_{0}^{t} K_{t-s} V_{s} d s-(Q V)_{t}=L \int_{0}^{t} K_{t-s} V_{s} d s+R_{t}
\end{aligned}
$$

### 68.1.3 Old Stuff related to Theorem 61.7

## Proof. Let

$$
F(t, x, y)=\int_{0}^{t} \int_{M} k(t-s, x, z) v(s, z, y) d \lambda(z) d s
$$

so that $p(t, x, y)=k(t, x, y)+F(t, x, y)$. For $\delta>0$ set

$$
F_{\delta}(t, x, y)=\int_{0}^{t} \int_{M} k(t+\delta-s, x, z) v(s, z, y) d \lambda(z) d s
$$

then

$$
\begin{aligned}
\partial_{t} F_{\delta}(t, x, y) & =\int_{M} k(\delta, x, z) v(t, z, y) d \lambda(z) \\
& +\int_{0}^{t} \int_{M} \partial_{t} k(t+\delta-s, x, z) v(s, z, y) d \lambda(z) d s \\
& =\int_{M} k(\delta, x, z) v(t, z, y) d \lambda(z) \\
& +\int_{0}^{t} \int_{M} L_{x} k(t+\delta-s, x, z) v(s, z, y) d \lambda(z) d s \\
& -\int_{0}^{t} \int_{M} r(t+\delta-s, x, z) v(s, z, y) d \lambda(z) d s \\
& =\int_{M} k(\delta, x, z) v(t, z, y) d \lambda(z)+L_{x} F_{\delta}(t, x, y) \\
& -\int_{0}^{t} \int_{M} r(t+\delta-s, x, z) v(s, z, y) d \lambda(z) d s .
\end{aligned}
$$

We may let $\delta \rightarrow 0$ in this last expression using the fact that $K_{t}$ is uniformly bounded on $\Gamma_{l}$ to find that

$$
\lim _{\delta \downarrow 0} F_{\delta}(t, x, y)=F(t, x, y)
$$

and

$$
\begin{aligned}
\lim _{\delta \downarrow 0} \partial_{t} F_{\delta}(t, x, y) & =v(t, x, y)+L_{x} F(t, x, y) \\
& -\int_{0}^{t} \int_{M} r(t-s, x, z) v(s, z, y) d \lambda(z) d s
\end{aligned}
$$

with the limits being uniform in $t$. Also by equation (62.1) and (62.2),

$$
v(t, x, y)-\int_{0}^{t} \int_{M} r(t-s, x, z) v(s, z, y) d \lambda(z) d s=r(t, x, y)
$$

$$
\begin{aligned}
\partial_{t} p(t, x, y) & =\partial_{t} k(t, x, y)+\partial_{t} F(t, x, y) \\
& =L_{x} k(t, x, y)-r(t, x, y)+L_{x} F(t, x, y)+r(t, x, y) \\
& =L_{x}(k(t, x, y)+F(t, x, y))=L_{x} p(t, x, y) .
\end{aligned}
$$

For $s<t$, let

$$
b(t, s, x, y):=\int_{M} k(t-s, x, z) v(s, z, y) d \lambda(z)
$$

It is clear the $b$ is smooth in $t$ and $s$ and is $C^{l}$ in $(x, y)$, moreover

$$
\lim _{s \uparrow 0}\|b(t, s, x, y)-v(t, x, y)\|_{C^{l}(x, y)}=0
$$

Hence we set $b(t, t, x, y):=v(t, x, y)$ so that $b(t, s, \cdot, \cdot)$ is continuous for $s \in[0, t]$ in the space of $C^{l}$ sections $\Gamma_{l}$. Similarly, for $s<t$,

$$
\begin{aligned}
\partial_{t} b(t, s, x, y) & =\int_{M} \dot{k}(t-s, x, z) v(s, z, y) d \lambda(z) \\
& =\int_{M}\left(L_{x} k(t-s, x, z)-r(t-s, x, z)\right) v(s, z, y) d \lambda(z) \\
& =L_{x} b(t, s, x, y)-\int_{M} r(t-s, x, z) v(s, z, y) d \lambda(z) .
\end{aligned}
$$

From this last expression and our previous comments, $\lim _{s \uparrow 0} L_{x} b(t, s, x, y)=$ $L_{x} v(t, x, y)$ in $\Gamma_{l-2}$ and hence

$$
\lim _{s \uparrow 0} \partial_{t} b(t, s, x, y)=L_{x} v(t, x, y) \text { in } C^{l-2}
$$

More precisely, we will construct $p(t, x, y)$ as

$$
\begin{equation*}
p(t, x, y)=\sum_{m=0}^{\infty} \int_{t \Delta_{m}} \int_{M^{m}} k\left(t-s_{m}, x, y_{m}\right) r\left(s_{m}-s_{m-1}, y_{m}, y_{m-1}\right) \tag{68.8}
\end{equation*}
$$

$$
r\left(s_{m-1}-s_{m-2}, y_{m-1}, y_{m-2}\right) \ldots r\left(s_{1}, y_{1}, y\right) d \mathbf{s} d \mathbf{y}
$$

Consider,

$$
\partial_{t} p(t, x, y)=\partial_{t} k(t, x, y)+\int_{0}^{t} \int_{M} k(t-s, x, z) v(s, y, z) d s d \lambda(y)
$$

is a bounded operator and its derivatives in $s$ up to order $k$ is a convergent sum in converge.

## Remarks on Covariant Derivatives on Vector Bundles

Let $\pi: E \rightarrow M$ be a vector bundle with fiber $W$. A local frame $u$ on $E$ is a local section of the bundle $\operatorname{Aut}(V, E) \rightarrow M$, i.e. for $m \in \mathcal{D}(u)$ (the domain of $u) u(m): W \rightarrow E_{m}=\pi^{-1}(m)$ is a linear isomorphism of vector spaces. Notice that any local section $S$ of $E$ may be written as $S(m)=u(m) s(m)$, where $s \in C^{\infty}(M, W)$. Suppose that $\nabla$ is a covariant derivative on $E$, define for $v \in T_{m} M$, a linear transformation $\nabla_{v} u: W \rightarrow E_{m}$ by

$$
\left(\nabla_{v} u\right) w:=\nabla_{v}(u(\cdot) w) \text { for each } w \in W
$$

With this notation and the basic properties of $\nabla$, given $s \in C^{\infty}(M, W)$, we have that

$$
\nabla_{v}(u s)=\left(\nabla_{v} u\right) s(m)+u(m) \partial_{v} s,
$$

where $\partial_{v} s:=\left.\frac{d}{d t}\right|_{0} s(\sigma(t))$ provided that $\dot{\sigma}(0)=v$. In particular this shows that

$$
u(m)^{-1} \nabla_{v}(u s)=\partial_{v} s+A(v) s(m)
$$

where $A(v)=A^{u}(v):=u(m)^{-1}\left(\nabla_{v} u\right)$. So the local representation of $\nabla$ is $\nabla=d+A$, where $A$ is a one form with values in $\operatorname{End}(W)$.

Given a path $S(t) \in E$, let $\sigma(t)=\pi(S(t))$ and $s(t)=u(\sigma(t))^{-1} S(t)$. Then define

$$
\begin{equation*}
\nabla S(t) / d t:=u(\sigma(t))(\dot{s}(t)+A(\dot{\sigma}(t)) s(t)) \tag{69.1}
\end{equation*}
$$

i.e. the local version of $\frac{\nabla}{d t}=\frac{d}{d t}+A(\dot{\sigma}(t))$. Notice that if $S=u s$ is a local section of $E$, then

$$
\begin{aligned}
\nabla S(\sigma(t))) / d t & =u(\sigma(t))\left(\frac{d}{d t} s(\sigma(t))+A(\dot{\sigma}(t)) s(\sigma(t))\right) \\
& =u(\sigma(t))\left(\partial_{\dot{\sigma}(t)} s+A(\dot{\sigma}(t)) s(\sigma(t))\right) \\
& =\nabla_{\dot{\sigma}(t)} S
\end{aligned}
$$

This explains why $\nabla / d t$ is independent of the local frame $u$ used in Eq. (69.1), a property which follows by direct computation as well.

We say that a path $S(t) \in E$ is parallel provided that $\nabla S(t) / d t=0$ for all $t$. Given a curve $\sigma(t)$ in $M$ and a point $S_{0} \in E_{\sigma(0)}$, there is a unique path $S(t) \in E$ such that $\pi(S(t))=\sigma(t)$ and $\nabla S(t) / d t=0$. This path is constructed by solving (locally) the linear equation

$$
\dot{s}(t)+A(\dot{\sigma}(t)) s(t)=0 \text { with } s(0)=u(\sigma(0))^{-1} S_{0}
$$

and then setting $S(t)=u(\sigma(t)) s(t)$. It is easy to check that the map $S_{0} \in$ $E_{\sigma(0)} \rightarrow S(t) \in E_{\sigma(t)}$ is linear. In fact $S(t)=/ / t(\sigma) S_{0}$, where $/ / t(\sigma)=$ $u(\sigma(t)) g(t) u(\sigma(0))^{-1}$ and $g(t) \in \operatorname{End}(W)$ is the unique solution to the linear differential equation,

$$
\dot{g}(t)+A(\dot{\sigma}(t)) g(t)=0 \text { with } g(0)=i d \in \operatorname{End}(W) .
$$

We will call $/ / t(\sigma)$ parallel translation along $\sigma$. It is uniquely characterized as the solution to the differential equation

$$
\nabla / / t(\sigma) / d t=0 \text { with } / / 0(\sigma)=i d \in \operatorname{End}\left(E_{\sigma(0)}\right)
$$

(If $U(t) \in \operatorname{End}\left(E_{\sigma(0)}, E_{\sigma(t)}\right)$ for each $t$, then $\nabla U(t) / d t$ is by definition the linear transformation from $E_{\sigma(0)}$ to $E_{\sigma(t)}$ determined by $(\nabla U(t) / d t) \xi:=$ $\nabla(U(t) \xi) / d t$ for all $\xi \in E_{0}$.) We have the following properties of parallel translation which follow from the uniqueness theorem for ordinary differential equations and the chain rule for covariant derivatives. Namely if $S(t)$ is a smooth path in $E$ and $t=\tau(s)$, then

$$
\nabla S(\tau(s)) / d s=\tau^{\prime}(s) \nabla S(t) /\left.d t\right|_{t=\tau(s)}
$$

This property is easily verified from Eq. (69.1).
Proposition 69.1. Let $\sigma(t) \in M$ be a smooth curve for $t \in[0, T]$ and let $\tau:[0, S] \rightarrow[0, T]$ be a smooth function such that $\tau(0)=0$. Then $/ /_{s}(\sigma \circ \tau)=/ /_{\tau(s)}(\sigma)$, i.e. parallel translation does not depend on how the underlying curve is parametrized. Secondly, let $\tilde{\sigma}(t):=\sigma(T-t)$, then $/ /_{t}(\tilde{\sigma})=/ /{ }_{T-t}(\sigma) / /_{T}(\sigma)^{-1}$. In particular $/ /_{T}(\sigma)^{-1}=/ /_{T}(\tilde{\sigma})$.

Proof. We have that

$$
\nabla / /_{s}(\sigma \circ \tau) / d s=0 \text { with } / / 0(\sigma \circ \tau)=i d \in \operatorname{End}\left(E_{\sigma(0)}\right)
$$

and

$$
\begin{aligned}
\nabla / /_{\tau(s)}(\sigma) / d s & =\nabla / /_{t}(\sigma) /\left.d t\right|_{t=\tau(s)} \tau^{\prime}(s)=0 \text { with } \\
/ /\left._{\tau(s)}(\sigma)\right|_{s=0} & =i d \in \operatorname{End}\left(E_{\sigma(0)}\right)
\end{aligned}
$$

and hence by uniqueness of solutions to O.D.E.'s we must have that $/ / s(\sigma \circ$ $\tau)=/ /_{\tau(s)}(\sigma)$. Similarly,

69 Remarks on Covariant Derivatives on Vector Bundles

$$
\nabla / /_{t}(\tilde{\sigma}) / d t=0 \text { with } / /_{0}(\tilde{\sigma})=i d \in E_{\sigma(T)}
$$

and

$$
\begin{aligned}
\nabla / /_{T-t}(\sigma) / /_{T}(\sigma)^{-1} / d t & =-\nabla / /_{s}(\sigma) / /_{T}(\sigma)^{-1} /\left.d s\right|_{s=T-t}=0 \text { with } \\
\quad / /_{T-t}(\sigma) / /\left._{T}(\sigma)^{-1}\right|_{t=0} & =i d \in E_{\sigma(T)}
\end{aligned}
$$

Hence again by uniqueness of solutions to O.D.E.'s we must have that $/ /_{t}(\tilde{\sigma})=$ $/ /_{T-t}(\sigma) / /_{T}(\sigma)^{-1}$.

## Spin Bundle Stuff

Let $M^{n}$ be a Riemannian manifold, $V=\mathbb{R}^{n}, C l(V)$ be the Clifford algebra over $V$ such that $v^{2}=-(v, v) 1$. Let $\operatorname{Spin}(n) \subset C l(V)$ be the spin group, $\rho: \operatorname{Spin}(n) \rightarrow S O(n)$ be the spin representation and $W$ be a left $C l(V)$ module. The following compatibility condition is need below in the construction of Spinor bundles $S$ over $M$ such that $C l(T M)$ acts on $S$.

Assumption 7 For $h \in \operatorname{Spin}(n), v \in V$ and $w \in W, h(v w)=(\rho(h) v)(h w)$.
Now for the construction of spinor bundles. Let $\left\{U_{\alpha}\right\}_{\alpha \in A}$ be an open cover of $M$ such that there exists $u_{a}: U_{\alpha} \rightarrow \operatorname{Hom}\left(\mathbb{R}^{n}, T U_{\alpha}\right)$ which are isometries. For $m \in U_{\alpha} \cap U_{\beta}$ let $g_{\alpha \beta}(m):=u_{\alpha}(m)^{-1} u_{\beta}(m) \in O(n)$. Notice that the for $m \in U_{\alpha} \cap U_{\beta} \cap U_{\delta}$,
$g_{\alpha \beta}(m) g_{\beta \delta}(m)=u_{\alpha}(m)^{-1} u_{\beta}(m) u_{\beta}(m)^{-1} u_{\delta}(m)=u_{\alpha}(m)^{-1} u_{\delta}(m)=g_{\alpha \delta}(m)$.
We now assume that $M$ is orientable which means that we may choose $u_{\alpha}$ such that $g_{\alpha \beta}(m) \in S O(n)$. Now if $M$ is spin as well, we may choose $\tilde{g}_{\alpha \beta}(m) \in \operatorname{Spin}(n)$ such that

1. $g_{\alpha \beta}(m)=\rho\left(\tilde{g}_{\alpha \beta}(m)\right)$ for all $m \in U_{\alpha} \cap U_{\beta}$ and all $\alpha$ and $\beta$.
2. $\tilde{g}_{\alpha \beta}(m) \tilde{g}_{\beta \delta}(m)=\tilde{g}_{\alpha \delta}(m)$.

Given this data, it is now possible to build a spin bundle over $M$ as follows. For $m \in M$, let

$$
S_{m}:=\left\{(m, \alpha, w): w \in W \text { and } \alpha \in A \text { s.t. } m \in U_{\alpha}\right\} / \sim
$$

where $(m, \alpha, w) \backsim\left(m, \alpha^{\prime}, w^{\prime}\right)$ if and only if $w^{\prime}=\tilde{g}_{\alpha^{\prime} \alpha}(m) w$. Let $S:=\cup_{m \in M} S_{m}$ and $\pi: S \rightarrow M$ be the projection map which takes $S_{m}$ to $m$ for all $m \in M$. Given $\alpha \in A$, let $\tilde{u}_{\alpha}: U_{\alpha} \rightarrow S_{U_{\alpha}}:=\pi^{-1}\left(U_{\alpha}\right)=\cup_{m \in U_{\alpha}} S_{m}$ be given by $\tilde{u}_{\alpha}(m) w$ is the equivalence class containing $(m, \alpha, w)$. Notice that $\tilde{u}_{\alpha}(m)$ : $W \rightarrow S_{m}$ is a bijective map and that $\tilde{u}_{\alpha}(m)^{-1} \tilde{u}_{\beta}(m) w=\tilde{g}_{\alpha \beta}(m) w$. One may now easily check that we may make $S_{m}$ in a well defined way into a linear
space by defining $\tilde{u}_{\alpha}(m) w+c \tilde{u}_{\alpha}(m) w^{\prime}:=\tilde{u}_{\alpha}(m)\left(w+c w^{\prime}\right)$, i.e. by requiring each $\tilde{u}_{\alpha}(m): W \rightarrow S_{m}$ to be linear.

Let us now show that we can make $S_{m}$ into a $C l\left(T_{m} M\right)$ module. For $\eta \in T_{m} M$ and $\xi \in S_{m}$ choose $\alpha \in A$ such that $m \in U_{\alpha}$ and choose $v \in \mathbb{R}^{n}$ and $w \in W$ such that $\eta=u_{\alpha}(m) v$ and $\xi=\tilde{u}_{\alpha}(m) w$. We then define $\eta \xi:=$ $\tilde{u}_{\alpha}(m)(v w)$. To see this is well defined choose $\alpha^{\prime} \in A$ such that $m \in U_{\alpha^{\prime}}$ and choose $v^{\prime} \in \mathbb{R}^{n}$ and $w^{\prime} \in W$ such that $\eta^{\prime}=u_{\alpha^{\prime}}(m) v^{\prime}$ and $\xi^{\prime}=\tilde{u}_{\alpha^{\prime}}(m) w^{\prime}$. Then $w=\tilde{g}_{\alpha \alpha^{\prime}}(m) w^{\prime}$ and $v=g_{\alpha \alpha^{\prime}}(m) v^{\prime}=\rho\left(\tilde{g}_{\alpha \alpha^{\prime}}(m)\right) v^{\prime}$, and hence

$$
v w=\left(g_{\alpha \alpha^{\prime}}(m) v^{\prime}\right)\left(\tilde{g}_{\alpha \alpha^{\prime}}(m) w^{\prime}\right)=\left(\rho\left(\tilde{g}_{\alpha \alpha^{\prime}}(m)\right) v^{\prime}\right)\left(\tilde{g}_{\alpha \alpha^{\prime}}(m) w^{\prime}\right)=\tilde{g}_{\alpha \alpha^{\prime}}\left(v^{\prime} w^{\prime}\right)
$$

From this it follows that $\tilde{u}_{\alpha}(m)(v w)=\tilde{u}_{\alpha^{\prime}}(m)\left(v^{\prime} w^{\prime}\right)$ so that $\eta \xi$ is well defined independent of the choice of $\alpha \in A$. Since $\xi \in T_{m} M \xrightarrow{\phi} L_{\xi} \in \operatorname{End}\left(S_{m}\right)\left(L_{\xi}\right.$ denotes left multiplication by $\xi$ ) satisfies $\phi(\xi)^{2}=-(\xi, \xi)_{m} I$, it follows that the action of $T_{m} M$ on $S_{m}$ extends uniquely to an action of $C l\left(T_{m} M\right)$ on $S_{m}$.

Hence if $M$ is a spin manifold, we have produced a vector bundle $S \rightarrow M$ such that each fiber of $S_{m}$ of $S$ is a $C l\left(T_{m} M\right)$ Clifford module.

## The Case where $M=\mathbb{R}^{n}$

### 71.1 Formula involving $p$

Let $L:=\frac{1}{2} \Delta+B+c, L_{0}:=\frac{1}{2} \Delta+B$ and $\beta(x):=b(x) \cdot x$ where and $B=$ $b \cdot \nabla=\sum_{i=1}^{n} b_{i} \partial_{i}$ and with $b_{i}(x)$ and $c(x)$ in $\mathbb{R}^{N \times N}$. Let $g$ be an $\mathbb{R}^{N \times N}-$ valued function of $(t, x)$ with $t>0$ and $x \in \mathbb{R}^{n}$ and set

$$
u(t, x)=p(t, x) g(t, x)
$$

Then

$$
\begin{aligned}
\left(\partial_{t}-L\right) u & =\left(\partial_{t}-L\right) p g \\
& =\left(\partial_{t}-L_{0}\right) p \cdot g+p\left(\partial_{t}-L\right) g-\nabla p \cdot \nabla g \\
& =-B p \cdot g+p\left(\partial_{t}-L\right) g-\nabla p \cdot \nabla g \\
& =p\left\{-B \ln p+\partial_{t}-L-\nabla \ln p \cdot \nabla\right\} g .
\end{aligned}
$$

Now

$$
\nabla \ln p=\nabla\left(-x^{2} / 2 t\right)=-\frac{x}{t}
$$

so the above equation may be written as:

$$
\begin{aligned}
\left(\partial_{t}-L\right) u & =p\left(\frac{1}{t} b \cdot x+\partial_{t}-L+\frac{1}{t} x \cdot \nabla\right) g \\
& =p\left(\partial_{t}-L+\frac{1}{t} S\right) g
\end{aligned}
$$

where

$$
\begin{aligned}
S & =\partial_{x}+b \cdot x \\
& =\partial_{x}+\beta
\end{aligned}
$$

and

$$
\beta(x)=b(x) \cdot x
$$

### 71.2 Asymptotics of a perturbed Heat Eq. on $\mathbb{R}^{n}$

Let $L:=\frac{1}{2} \Delta+B+c, L_{0}:=\frac{1}{2} \Delta+B$ and $\beta(x):=b(x) \cdot x$ where and $B=$ $b \cdot \nabla=\sum_{i=1}^{n} b_{i} \partial_{i}$ and with $b_{i}(x)$ and $c(x)$ in $\mathbb{R}^{N \times N}$. As above let

$$
p(t, x)=(2 \pi t)^{-n / 2} e^{-x^{2} / 2 t}
$$

be the heat kernel with pole at 0 for $\mathbb{R}^{n}$.
Lemma 71.1. Let $u_{0} \in \mathbb{R}^{N \times N}$ be given, then there is a unique solution to the O.D.E

$$
\begin{equation*}
\dot{U}(t, x)=-\frac{1}{t} \beta(t x) U(t, x) \text { with } U(0, x)=u_{0} \in \mathbb{R}^{N \times N} \tag{71.1}
\end{equation*}
$$

Moreover, $U(t, x)$ is smooth in $(t, x)$ and

$$
\begin{equation*}
U(t, s x)=U(t s, x) \tag{71.2}
\end{equation*}
$$

for all $s, t \in \mathbb{R}$ and $x \in \mathbb{R}^{n}$.
Proof. Since $\beta(0)=0$,

$$
\frac{1}{t} \beta(t x)=\frac{1}{t} \int_{0}^{t} \partial_{x} \beta(\tau x) d \tau=\int_{0}^{1} \partial_{x} \beta(u t x) d u
$$

and hence the matrix function $(t, x) \rightarrow \frac{1}{t} \beta(t x)$ is smooth even for $t$ near 0 . By basic O.D.E. theory this shows that $U$ exists and is smooth. Since

$$
\begin{aligned}
\frac{d}{d t} U(t s, x) & =s \dot{U}(t s, x)=-s \frac{1}{s t} \beta(t s x) U(t s, x) \\
& =-\frac{1}{t} \beta(t s x) U(t s, x)
\end{aligned}
$$

it follows that $U(t s, x)$ satisfies the same O.D.E. as $U(t, s x)$. Hence by uniqueness of solutions to O.D.E.'s we find that Eq. (71.2) holds.

Remark 71.2. If $[\beta(x), \beta(y)]=0$ for all $x$ and $y$, then the solution to Eq. (71.1) is given by

$$
\begin{equation*}
U(1, x)=\exp \left(-\int_{0}^{1} \frac{1}{t} \beta(t x) d t .\right) \tag{71.3}
\end{equation*}
$$

The rest of this section is devoted to the proof of the following Theorem.
Theorem 71.3. Let $U$ be defined as in Lemma 71.1 and defined

$$
\begin{equation*}
u_{0}(x)=U(1, x) \tag{71.4}
\end{equation*}
$$

and $u_{k}$ inductively by
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$$
\begin{equation*}
u_{k+1}(x)=u_{0}(x) \int_{0}^{1} \tau^{k} u_{0}(\tau x)^{-1} L u_{k}(\tau x) d \tau \tag{71.5}
\end{equation*}
$$

for all $k=0,1,2 \ldots$ If

$$
\begin{equation*}
\Sigma_{q}(t, x):=p(t, x) \sum_{k=0}^{q} t^{k} u_{k}(x) \tag{71.6}
\end{equation*}
$$

then

$$
\begin{equation*}
\left(\partial_{t}-L\right) \Sigma_{q}(t, x)=-t^{q} p(t, x) L u_{q}(t, x) \tag{71.7}
\end{equation*}
$$

The proof of this theorem could be given by direct computation. However, we will take a longer route however and derive the formulas in the Theorem.

Let $g$ be an $\mathbb{R}^{N \times N}$ - valued function of $(t, x)$ with $t>0$ and $x \in \mathbb{R}^{n}$ and set

$$
u(t, x)=p(t, x) g(t, x)
$$

Then

$$
\begin{aligned}
\left(\partial_{t}-L\right) u & =\left(\partial_{t}-L\right) p g \\
& =\left(\partial_{t}-L_{0}\right) p \cdot g+p\left(\partial_{t}-L\right) g-\nabla p \cdot \nabla g \\
& =\left(p \frac{r}{2 t} \partial_{r} \ln J-B p\right) \cdot g+p\left(\partial_{t}-L\right) g-\nabla p \cdot \nabla g \\
& =p\left\{\frac{r}{2 t} \partial_{r} \ln J-B \ln p+\partial_{t}-L-\nabla \ln p \cdot \nabla\right\} g
\end{aligned}
$$

Now

$$
\nabla \ln p=\nabla\left(-x^{2} / 2 t\right)=-\frac{x}{t}
$$

so the above equation may be written as:

$$
\begin{align*}
\left(\partial_{t}-L\right) u & =p\left(b \cdot \frac{x}{t}+\partial_{t}-L+\frac{x}{t} \cdot \nabla\right) g \\
& =p\left(\partial_{t}-L+\frac{1}{t} S\right) g \tag{71.8}
\end{align*}
$$

where

$$
\begin{equation*}
S:=x \cdot \nabla+b(x) \cdot x=\partial_{x}+\beta(x) \tag{71.9}
\end{equation*}
$$

and $\beta(x):=b(x) \cdot x$ as above.
Now let

$$
\begin{equation*}
g(t, x)=g_{q}(t, x)=\sum_{k=0}^{q} t^{k} u_{k}(x) \tag{71.10}
\end{equation*}
$$

and consider

$$
\begin{aligned}
\left(\partial_{t}+\frac{1}{t} S-L\right) g & =\sum_{k=0}^{q}\left\{t^{k-1}\left(k u_{k}+S u_{k}\right)-t^{k} L u_{k}\right\} \\
& =\frac{1}{t} S u_{0}+\sum_{k=0}^{q-1} t^{k}\left((k+1) u_{k+1}+S u_{k+1}-L u_{k}\right)-t^{q} L u_{q}
\end{aligned}
$$

Thus if we choose $u_{0}$ such that

$$
\begin{equation*}
S u_{0}(x)=(x \cdot \nabla \cdot+\beta(x)) u_{0}(x)=0 \tag{71.11}
\end{equation*}
$$

and $u_{k}$ such that

$$
\begin{equation*}
(k+1) u_{k+1}+S u_{k+1}-L u_{k}=0 \tag{71.12}
\end{equation*}
$$

then $\left(\partial_{t}+\frac{1}{t} S-L\right) g=-t^{q} L u_{q}$ or equivalently by Eq. (71.8),

$$
\left(\partial_{t}-L\right) \Sigma_{q}=\left(\partial_{t}-L\right)(p g)=-t^{q} L u_{q}
$$

which then proves Theorem 71.3 show that $u_{k}$ defined in the theorem solve Equations (71.11) and (71.12).

Suppose that $u_{0}$ is a solution to Eq. (71.11). If $U(t, x):=u_{0}(t x)$, then

$$
\begin{aligned}
\dot{U}(t, x) & =x \cdot \nabla u_{0}(t x)=\frac{1}{t} t x \cdot \nabla u_{0}(t x) \\
& =\frac{1}{t} \partial_{t x} u_{0}(t x)=-\frac{1}{t} \beta(t x) u_{0}(t x) \\
& =-\frac{1}{t} \beta(t x) U(t, x)
\end{aligned}
$$

and hence $u_{0}(x)$ must be given by $U(1, x)$ as in Eq. (71.4). Conversely if $u_{0}$ is defined by Eq. (71.4) then from Lemma 71.1, $u_{0}(s x)=U(s, x)$ and hence

$$
\begin{aligned}
\partial_{x} u_{0}(x) & =\left.\frac{d}{d s}\right|_{1} u_{0}(s x)=\left.\frac{d}{d s}\right|_{1} U(s, x) \\
& =-\beta(x) \dot{U}(1, x)=-\beta(x) u_{0}(x) .
\end{aligned}
$$

Thus we have shown that $u_{0}$ solves Eq. (71.11).
We now turn our attention to solving Eq. (71.12). Assuming $u_{k+1}$ is a solution to Eq. (71.12), then $V_{k+1}(t, x):=u_{k+1}(t x)$ satisfies

$$
\begin{aligned}
t \dot{V}_{k+1}(t, x)(t) & =t x \cdot \nabla u_{k+1}(t x) \\
& =-\left(\frac{1}{t} \beta(t x)+k+1\right) V_{k+1}(t, x)+L u_{k}(t x)
\end{aligned}
$$

or equivalently

$$
\begin{equation*}
\dot{V}_{k+1}(t, x)=-\left(\frac{1}{t} \beta(t x)+\frac{k+1}{t}\right) V_{k+1}(t, x)+\frac{1}{t} L u_{k}(t x) \tag{71.13}
\end{equation*}
$$

This equation may be solved by introducing an integrating factor, i.e. let

$$
U_{k+1}(t, x):=t^{k+1} u_{0}(t x)^{-1} V_{k+1}(t, x)
$$

Then $U_{k+1}$ solves

$$
\begin{aligned}
\dot{U}_{k+1}(t, x) & =t^{k+1} u_{0}(t x)^{-1} \dot{V}_{k+1}(t, x)+(k+1) t^{k} u_{0}(t x)^{-1} V_{k+1}(t, x) \\
& +t^{k+1} u_{0}(t x)^{-1} \frac{\beta(t x)}{t} V_{k+1}(t, x) \\
& =t^{k} u_{0}(t x)^{-1} L u_{k}(t x)
\end{aligned}
$$

## Hence

$$
\begin{aligned}
U_{k+1}(t, x) & =U_{k+1}(0, x)+\int_{0}^{t} \tau^{k} u_{0}(\tau x)^{-1} L u_{k}(\tau x) d \tau \\
& =\int_{0}^{t} \tau^{k} u_{0}(\tau x)^{-1} L u_{k}(\tau x) d \tau
\end{aligned}
$$

Therefore if $u_{k+1}$ exists it must be given by Eq. (71.5).
Conversely if $u_{k+1}$ is defined by Eq. (71.5) then

$$
\begin{aligned}
u_{k+1}(s x) & =u_{0}(s x) \int_{0}^{1} \tau^{k} u_{0}(\tau s x)^{-1} L u_{k}(\tau s x) d \tau \\
& =u_{0}(s x) \int_{0}^{s} s^{-(k+1)} t^{k} u_{0}(t x)^{-1} L u_{k}(t x) d t
\end{aligned}
$$

and hence

$$
\begin{aligned}
\partial_{x} u_{k+1}(x) & =\left.\frac{d}{d s}\right|_{1} u_{k+1}(s x) \\
& =-\beta(x) u_{k+1}(x)+\left.u_{0}(x) \frac{d}{d s}\right|_{1} \int_{0}^{s} s^{-(k+1)} t^{k} u_{0}(t x)^{-1} L u_{k}(t x) d t \\
& =-\beta(x) u_{k+1}(x)+L u_{k}(x)-(k+1) u_{0}(x) \int_{0}^{1} t^{k} u_{0}(t x)^{-1} L u_{k}(t x) d t \\
& =-(\beta(x)+k+1)) u_{k+1}(x)+L u_{k}(x)
\end{aligned}
$$

which is shows that $u_{k+1}$ solves Eq. (71.12). This finishes the proof of Theorem 71.3 .
$1214 \quad 72$ Higher Order Elliptic Equations

## Higher Order Elliptic Equations

Definition 72.1. $H^{s}\left(\mathbb{R}^{d}, \mathbb{C}^{N}\right):=\left\{u \in \mathcal{S}^{\prime}\left(\mathbb{R}^{d}, \mathbb{C}^{N}\right):|\hat{u}| \in L^{2}\left(\left(1+|\xi|^{2}\right)^{s} d s\right)\right\}$.
Note For $s=0,1,2, \ldots$ this agrees with our previous Notations of Sobolev spaces.

Lemma 72.2. (A) $C_{C}^{\infty}\left(\mathbb{R}^{d}, \mathbb{C}^{N}\right)$ is dense in $H^{s}\left(\mathbb{R}^{d}, \mathbb{C}^{N}\right)$.
(B) If $s>0$ then $\|u\|_{-s}=\|u\|_{H^{-s}}=\sup _{\phi \in C_{C}^{\infty}} \frac{|<u, \phi>|}{\|\phi\|_{-s}}$.

Proof. (A) We first rate that $\mathcal{S}\left(\mathbb{R}^{d}, \mathbb{C}^{N}\right)$ is dense in $H^{s}\left(\mathbb{R}^{d}, \mathbb{C}^{N}\right)$ because $\mathcal{F} \mathcal{S}=\mathcal{S}$ and $\mathcal{S}$ is dense in $\left.L^{2}\left(1+|\xi|^{2}\right)^{s} d \xi\right)$ for all $s \in \mathbb{R}$. It is easily seen that $C_{c}^{\infty}$ is dense in $\mathcal{S}$ for all $s=0,1,2, \ldots$ relative to $H^{s}$-norm. But this enough to prove (A) since for all $s \in \mathbb{R}\|\cdot\|_{S} \leq\|\cdot\|_{k}$ for some $k \in \mathbb{N}$.
(B)

$$
\begin{aligned}
\|u\|_{-s}^{2} & =\|\hat{u}\|_{L^{2}\left(\left(1+|\xi|^{2}\right)^{-2} d s\right)}=\sup _{\phi \in \mathcal{S}} \frac{\left|(\hat{u}, \phi)_{-s}\right|}{\|\phi\|_{-s}} \\
& =\sup _{\phi \in \mathcal{S}} \frac{\left|\int \hat{u}(\xi) \phi(\xi)\left(1+|\xi|^{2}\right)^{-s} d \xi\right|}{\sqrt{S|\phi(\xi)|^{2}\left(1+|\xi|^{2}\right)^{-s} d \xi}} .
\end{aligned}
$$

Let $\hat{\psi}(\xi)=\hat{\varpi}(\xi)\left(a+|\xi|^{2}\right)^{-s}$, the arbitrary element of $\mathcal{S}$ still to find

$$
\begin{aligned}
\|u\|_{-s} & =\sup _{\psi \in \mathcal{S}} \frac{\left|\int \hat{u}(\xi) \hat{\psi}(\xi) d \xi\right|}{\sqrt{\int|\hat{\psi}(\xi)|^{2}\left(1+|\xi|^{2}\right)^{s} d \xi}} \\
& =\sup _{\psi \in \mathcal{S}} \frac{|<u, \psi>|}{\|\psi\|_{s}}=\sup _{\psi \in C_{c}^{\infty}} \frac{|<u, \psi>|}{\|\psi\|_{s}}
\end{aligned}
$$

Since $C_{c}^{\infty}$ is dense as well.
Definition 72.3. $H^{s}\left(\mathbb{R}^{d}, \mathbb{C}^{N}\right):=\left\{u \in \mathcal{S}^{\prime}\left(\mathbb{R}^{d}, \mathbb{C}^{N}\right):|\hat{u}| \in L^{2}\left(\left(1+|\xi|^{2}\right)^{s} d \xi\right)\right\}$.

Proposition 72.4. Suppose $L=L_{0}$ has constant coefficients then $\|u\|_{s} \leq$ $C\left(\|L u\|_{s-k}+\|u\|_{s-1}\right)$ for all $s \in \mathbb{R}$.

## Proof.

$$
\begin{aligned}
\widehat{L u}(s)^{"} & =" L \widehat{u(x)} e^{i x-\xi} d x=\int \sigma\left(L_{0}\right) \widehat{(x, \xi)} u(x) e^{i x-\xi} d x \\
& =\int \sigma(\xi) u(x) e^{i x \cdot \xi} d x=\sigma(\xi) \hat{u}(\xi)
\end{aligned}
$$

Therefore

$$
\begin{aligned}
\|\hat{u}(\xi)\| & =\left\|\sigma(\xi)^{-1} \sigma(\xi) \hat{u}(\xi)\right\| \leq\left\|\sigma(\xi)^{-1}\right\|\|\sigma(\xi) \hat{u}(\xi)\| \\
& \leq \frac{\epsilon}{|\xi|^{k}}\|\sigma(\xi) \hat{u}(\xi)\|
\end{aligned}
$$

Therefore

$$
\|\sigma(\xi) \hat{u}(\xi)\| \geq \frac{|\xi|^{k}}{C_{1}}\|\hat{u}(\xi)\|
$$

Notice there exist $C>0$ such that

$$
1 \leq C\left(|\xi|^{2 k}\left(\frac{1}{1+|x \imath|^{2}}\right)^{k}+\frac{1}{1+|\xi|^{2}}\right)
$$

Then

$$
\begin{aligned}
|\hat{u}|^{2}\left(1+|\xi|^{2}\right)^{2} & \leq C\left(|\xi|^{2 k}\left(1+|\xi|^{2}\right)^{s-k}|\hat{u}|^{2}+\left(1+|\xi|^{2}\right)^{s-1}|\hat{u}|^{2}\right) \\
& \leq \widetilde{C}\left(\left(1+|\xi|^{2}\right)^{s-k}|\sigma(\xi) \hat{u}|^{2}+\left(1+|\xi|^{2}\right)^{s-1}\left(\hat{u}(\xi)^{2}\right)\right. \\
& \leq \widetilde{C}\left(\left(1+|\xi|^{2}\right)^{s-k}|\widehat{L u}(\xi)|^{2}+\left(1+|\xi|^{2}\right)^{s-1}\left(\hat{u}(\xi)^{2}\right) .\right.
\end{aligned}
$$

Integrate this on $\xi$ to get the desired inequality.
Notation 72.5 Suppose $\Omega \subset \mathbb{R}^{d}$ is an open set and $a_{\alpha} \in C^{\infty}\left(\bar{\Omega}, \mathbb{C}^{N}\right)$ for some $N$ and $|\alpha| \leq k$. Set $L: C^{\infty}\left(\Omega, \mathbb{C}^{N}\right) \rightarrow C^{\infty}\left(\Omega, \mathbb{C}^{N}\right)$ to be the operator $L=\sum_{|\alpha| \leq k} a_{\alpha} \partial^{\alpha}$ and $L_{0}:=\sum_{|\alpha|=k} a_{\alpha} \partial^{\alpha}$

$$
\sigma\left(L_{0}\right)(\xi)=\sum_{|\alpha|=k} a_{\alpha}(i \xi)^{\alpha}=i^{|\alpha|} \sum_{|\alpha|=k} a_{\alpha} \xi^{\alpha}
$$

for $\xi \in \mathbb{R}^{d}$. Notice that $L_{0} e^{i \xi \cdot x}=\sigma\left(L_{0}\right)(x, \xi) e^{i \xi \cdot x}$.
Definition 72.6. L is elliptic at $x \in \Omega$ if $\sigma\left(L_{0}\right)(x, \xi)^{-1}$ exists for all $\xi \neq 0$ and $L$ is elliptic on $\bar{\Omega}$ if $\sigma\left(L_{0}\right)(x, \xi)^{-1}$ exist for all $\xi \neq 0, x \in \bar{\Omega}$.
Remark 72.7. If $L$ is elliptic on $\bar{\Omega}$ then there exist $\in>0$ and $C>0$ such that $\left\|\sigma\left(L_{0}\right)(x, \xi)\right\| \geq \epsilon|\xi|^{k}$ for all $x \in \bar{\Omega}, \xi \in \mathbb{R}^{d}$. Also $\left\|\sigma\left(L_{0}\right)(x, \xi)^{-1}\right\| \leq C|\xi|^{k}$ for all $x \in \bar{\Omega}, \xi \in \mathbb{R}^{d} \backslash\{0\}$.
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Indeed for $|\xi|=1,(x, \xi) \rightarrow\left\|\sigma\left(L_{0}\right)(x, \xi)^{-1}\right\|$ is a continuous function of $(x, \xi)$ on $\bar{\Omega} \times S^{d-1}$. Therefore it has a global maximum say $C$. Then for $|\xi|=t \neq 0$.

$$
\begin{gathered}
\left\|\sigma\left(L_{0}\right)\left(x, \frac{\xi}{t}\right)^{-1}\right\| \leq C \text { implies }\left\|\left[\left(\frac{1}{t}\right)^{k} \sigma\left(L_{0}, x, \xi\right)\right]^{-1}\right\| \leq C \text { or } \\
\left\|\sigma\left(L_{0}\right)(x, \xi)^{-1}\right\| \leq C t^{-k}=C|\xi|^{-k} .
\end{gathered}
$$

Given $\Omega$ as above let

$$
H_{0}^{s}\left(\Omega, \mathbb{C}^{N}\right):=\overline{C_{C}^{\infty}\left(\Omega, \mathbb{C}^{N}\right)} H^{s}\left(\mathbb{R}^{d}, \mathbb{C}^{N}\right)
$$

Theorem 72.8 (A priori Estimates). For all $s \in \mathbb{R}$ there exist $C>0$ such that $\forall u \in H_{s}^{0}(\Omega)$

$$
\begin{equation*}
\|u\|_{s} \leq C\left(\|L u\|_{s-k}+\|u\|_{s-1}\right) \tag{72.1}
\end{equation*}
$$

Reference Theorem (6.28) of Folland p. 210 chapter 6.
Proof. I will only prove the inequality (72.1) for $s=0$. (However, negative $s$ are needed to prove desired elliptic regularity results. This could also be done using the Theory of pseudo differential operators.) With out loss of generality we may assume $u \in C_{c}^{\infty}\left(\Omega, \mathbb{C}^{N}\right)$.

Step (1) The inequality holds if $L=L_{0}$ with constant coefficients by above proposition.

Step (2) Suppose now $L=L_{0}$ but does not have constant coefficients. Define $L_{x_{0}}=\sum_{|\alpha|=k} a_{\alpha}\left(x_{0}\right) \partial^{\alpha}$ for all $x_{0} \in \Omega$. Then there exist $C_{0}$ independent of $x_{0}$ such that

$$
\|u\|_{0} \leq C_{0}\left(\left\|L_{x_{0}} u\right\|_{-k}+\|u\|_{-1}\right)
$$

for all $x_{0} \in \Omega$. Suppose $\operatorname{supp}(u) \subset B\left(x_{0}, \delta\right)$ with $\delta$ small. Consider

$$
\left\|\left(L-L_{x_{0}}\right) u\right\|_{-k}=\sup _{\phi \in C_{C}^{\infty}} \frac{\left|\left(\left(L-L_{x_{0}}\right) u, \phi\right)\right|}{\|\phi\|_{s}}
$$

Now

$$
\begin{aligned}
\left(\left(L-L_{x_{0}}\right) u, \phi\right) & =(-1)^{k} \sum\left(u, \partial^{\alpha}\left(a_{\alpha}^{+}-a_{\alpha}\left(x_{0}\right)\right) \phi\right) \\
& =(-1)^{k}\left(u, \sum\left(a_{\alpha}^{+}-a_{\alpha}\left(x_{0}\right) \partial^{\alpha} \phi+s \phi\right) .\right.
\end{aligned}
$$

Where $s \phi=\sum_{|\alpha|=k}\left[\partial^{\alpha}, a_{\alpha}^{+}\right] \phi$ is a $k-1$ order differential operator. Therefore

$$
\begin{aligned}
&\left|\left(\left(L=L_{x_{0}}\right) u, \phi\right)\right| \leq\left|\left(s^{+} u, \phi\right)\right|+\sum_{\alpha}\left|\left(u,\left(a_{\alpha}^{+}-a_{\alpha}\left(x_{0}\right)\right) \partial^{\alpha} \phi\right)\right| \\
& \leq\left\|S^{+} u\right\|_{-k}\|\phi\|_{k} \\
&+\underbrace{\sum_{\alpha\left|x-x_{0}\right| \leq \delta} \sup \left|a_{\alpha}(x)-a_{\alpha}\left(x_{0}\right)\right|}_{C(\delta)}\|u\|_{0}\|\phi\|_{k}
\end{aligned}
$$

$$
\leq k\|u\|_{-1}\|\phi\|_{k}+C(\delta)\|u\|_{0}\|\phi\|_{k}
$$

Therefore

$$
\left\|\left(L-L_{x_{0}}\right) u\right\|_{-k} \leq C(\delta)\|u\|_{0}+K\|u\|_{-1}
$$

Choose $\delta$ small such that $C_{0} C(\delta) \leq \frac{1}{2}$ by unit continuous. Therefore

$$
\left\|\left(L-L_{x_{0}}\right) u\right\|_{-k} \leq \frac{1}{2}\|u\|_{0}+K\|u\|_{-1}
$$

Hence

$$
\begin{aligned}
\|u\|_{0} & \leq C_{0}\left(\left\|L_{x_{0}} u\right\|_{-k}+\|u\|_{-1}\right) \\
& \leq C_{0}\left(\left\|L-L_{x_{0}} u\right\|_{-k}+\|L u\|_{-k}+\|u\|_{-1}\right) \\
& \leq C_{0}\left(C(\delta)\|u\|_{0}+\|L u\|_{-k}+(1+K)\|u\|_{-1}\right) \\
& \leq \frac{1}{2}\|u\|_{0}+C_{0}(1+K)\left(\|L u\|_{-k}+\|u\|_{-1}\right)
\end{aligned}
$$

Therefore

$$
\|u\|_{0} \leq 2 C_{0}(1+K)\left(\|L u\|_{-k}+\|u\|_{-1}\right)
$$

provided supp $u \subset B\left(x_{0}, \delta\right)$ for some $x_{0} \in R$. Now cover $\bar{\Omega}$ by finite collection of balls with radius $\delta$ and choose a partition of unity subordinate to this cover. $\left\{x_{i}\right\}$. Therefore if $u \in C_{c}^{\infty}(\Omega)$, then $\sum \psi_{i} u=u$ and

$$
\begin{aligned}
\|u\|_{0} & \leq C\left(\|L u\|_{-k}+\|s u\|_{-k}+\|u\|_{-1}\right) \\
& \leq \widetilde{C}\left(\|L u\|_{-k}+\|u\|_{-1}\right)
\end{aligned}
$$

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$$
\begin{aligned}
e^{(t+h) L}-e^{t L} & =\sum_{n=0}^{\infty} \frac{1}{n!}\left\{(t+h)^{n}-t^{n}\right\} L^{n} \\
& =\sum_{n=1}^{\infty} \frac{n}{n!} c_{n}^{n-1}(h) h L^{n}=h L \sum_{n=0}^{\infty} \frac{1}{n!} c_{n+1}^{n}(h) L^{n}
\end{aligned}
$$

where $c_{n}(h)$ is some number between $t$ and $t+h$ for each $n$. Hence

$$
\frac{e^{(t+h) L}-e^{t L}}{h}-L e^{t L}=L \sum_{n=0}^{\infty} \frac{1}{n!}\left[c_{n+1}^{n}(h)-t^{n}\right] L^{n}
$$

and thus

$$
\left\|\frac{e^{(t+h) L}-e^{t L}}{h}-L e^{t L}\right\| \leq\|L\| \sum_{n=1}^{\infty} \frac{1}{n!}\left|c_{n+1}^{n}(h)-t^{n}\right|\|L\|^{n} \rightarrow 0
$$

as $n \rightarrow \infty$ by the dominated convergence theorem.
Before continuing, let us prove the basic group property of $e^{t L}$, namely:

$$
\begin{equation*}
e^{t L} e^{s L}=e^{(t+s) L} \tag{73.3}
\end{equation*}
$$

To prove this equation, notice that

$$
\frac{d}{d t} e^{-t L} e^{(t+s) L}=e^{-t L}(-L+L) e^{(t+s) L}=0
$$

Thus $e^{-t L} e^{(t+s) L}$ is independent of $t$ and hence

$$
\begin{equation*}
e^{-t L} e^{(t+s) L}=e^{s L} \tag{73.4}
\end{equation*}
$$

By choosing $s=0$ we find that $e^{-t L} e^{t L}=I$, and by replacing $t$ by $-t$ we can conclude that $e^{-t L}=\left(e^{t L}\right)^{-1}$. This last observation combined with (73.4) proves (73.3).

Alternate Proof of Eq. (73.3).

$$
\begin{aligned}
e^{(t+s) L} & =\sum_{n=0}^{\infty} \frac{(t+s)^{n}}{n!} L^{n}=\sum_{n=0}^{\infty} \sum_{k=0}^{n} \frac{t^{k} L^{k} s^{n-k} L^{n-k}}{k!(n-k)!} \\
& =\sum_{n=0}^{\infty} \sum_{\substack{k+\ell=n \\
b, \ell \geq 0}} \frac{t^{k} L^{k} s^{\ell} L^{\ell}}{k!\ell!}=e^{t L} e^{s L}
\end{aligned}
$$

where the above manipulations are justified since,

$$
\sum_{n=0}^{\infty} \sum_{\substack{k+\ell=n \\ b, \ell \geq 0}} \frac{|t|^{k}|s|^{\ell}\|L\|^{k}\|L\|^{\ell}}{k!l!}=e^{(|t|+|s|)\|L\|}<\infty
$$

Now clearly if

$$
v(t):=e^{t L} v_{0}:=\sum_{n=0}^{\infty} \frac{t^{n}}{n!} L^{n} v_{0}
$$

we have that $v \in C^{1}(\mathbb{R} \rightarrow X)$ and $\dot{v}=L v$. Therefore $v$ does solve Eq. (73.2). To see that this solution is unique, suppose that $v(t)$ is any solution to (73.2). Then

$$
\frac{d}{d t} e^{-t L} v(t)=e^{-t L}(-L+L) v(t)=0
$$

so that $e^{-t L} v(t)$ is constant and thus $e^{-t L} v(t)=v_{0}$. Therefore $v(t)=e^{t L} v_{0}$.

Theorem 73.4 (The Diagonal Case). Consider $t>0$ only now. Let $p \in$ $[1, \infty),(\Omega, \mathcal{F}, m)$ be a measure space and $a: \Omega \rightarrow \mathbb{R}$ be a measurable function such that $a$ is bounded above by a constant $C<\infty$. Define $D(L)=\{f \in$ $\left.L^{p}(m): a f \in L^{p}\right\}$, and for $f \in D(L)$ set $L f=a f=M_{a} f$. (In general $L$ is an unbounded operator.) Define $e^{t L}=M_{e}$ taf. (Note that $\left|e^{t a}\right| \leq e^{t C}<\infty$, so $e^{t L}$ is a bounded operator.) by D.C.T. So one has $e^{t L} \rightarrow I$ strongly as $t \downarrow 0$.

1. $e^{t L}$ is a strongly continuous semi-group of bounded operators.
2. If $f \in D(L)$ then $\frac{d}{d t} e^{t L} f=e^{t L} L f$ in $L^{p}(m)$.
3. For all $f \in L^{p}$ and $t>0, \frac{d}{d t} e^{t L} f=L e^{t L} f$ in $L^{p}(m)$.

Proof. By the dominated convergence theorem,

$$
\left\|e^{t L} f-f\right\|_{L^{p}}^{p}=\int_{\Omega}\left|\left(e^{t a}-1\right) f\right|^{p} d m \rightarrow 0 \text { as } t \downarrow 0
$$

which proves item 1 . For item 2 we see using the fundamental theorem of calculus that

$$
\begin{align*}
\left\|\left(\frac{e^{(t+h) L}-e^{t L}}{h}-e^{t L} L\right) f\right\|_{L^{p}} & =\left\|\frac{1}{h} \int_{0}^{h}\left(a e^{(t+\tau) a}-e^{t a} a\right) d \tau \cdot f\right\|_{L^{p}} \\
& =\left\|\frac{1}{h} \int_{0}^{h}\left(e^{(t+\tau) a}-e^{t a}\right) d \tau \cdot a f\right\|_{L^{p}} \tag{73.5}
\end{align*}
$$

Since $a f \in L^{p}$,

$$
\left|\frac{1}{h} \int_{0}^{h}\left(e^{(t+\tau) a}-e^{t a}\right) d \tau\right| \leq 2 e^{(t+|h|) C}
$$

and

$$
\frac{1}{h} \int_{0}^{h}\left(e^{(t+\tau) a}-e^{t a}\right) d \tau \rightarrow 0 \text { as } h \rightarrow 0
$$

the Dominated convergence shows that the last term in Eq. (73.5) tends to zero as $h \rightarrow 0$. The above computations also work at $t=0$ provided $h$ is restricted to be positive.

Item 3 follows by the same techniques as item 2 . We need only notice that by basic calculus if $t>0$ and $\tau \in(t / 2,3 t / 2)$ then

$$
\left|a e^{(t+\tau) a}\right| \leq \max \left\{(t / 2)^{-1}, C e^{3 C / 2}\right\}
$$

Example 73.5 (Nilpotent Operators). Let $L: X \rightarrow X$ be a nilpotent operator, i.e., $\forall v \in D(L)$ there exists $n=n(v)$ such that $L^{n} v=0$ Then

$$
\dot{v}(t)=L v(t) \text { with } v(0)=v \in D(L)
$$

has a solution (in $D(L)$ ) given by

$$
v(t)=e^{t L} v:=\sum_{n=0}^{\infty} \frac{t^{n}}{n!} L^{n} v
$$

A special case of the last example would be to take $X=\mathcal{D}(L)$ to be the space of polynomial functions on $\mathbb{R}^{d}$ and $L p=\Delta p$.

Example 73.6 (Eigenvector Case). Let $L: X \rightarrow X$ be a linear operator and suppose that $D(L)=\operatorname{span} X_{0}$, where $X_{0}$ is a subset of $X$ consisting of eigenvectors for $L$, i.e., $\forall v \in X_{0}$ there exists $\lambda(v) \in \mathbf{C}$ such the $L v=\lambda(v) v$. Then

$$
\dot{v}(t)=L v(t) \text { with } v(0)=v \in D(L)
$$

has a solution in $D(L)$ given by

$$
v(t)=e^{t L} v:=\sum_{n=0}^{\infty} \frac{t^{n}}{n!} L^{n} v
$$

More explicitly, if $v=\sum_{v=1}^{n} v_{i}$ with $v_{i} \in X_{0}$, then

$$
e^{t L} v=\sum_{i=1}^{n} e^{t \lambda\left(v_{i}\right)} v_{i}
$$

The following examples will be covered in more detail in the exercises.
Example 73.7 (Translation Semi-group). Let $X:=L^{2}\left(\mathbb{R}^{d}, d \lambda\right), w \in \mathbb{R}^{d}$ and

$$
\left(T_{w}(t) f\right)(x):=f(x+t)
$$

Then $T_{w}(t)$ is a strongly continuous contraction semi-group. In fact $T_{w}(t)$ is unitary for all $t \in \mathbb{R}$.

Example 73.8 (Rotation Semi-group). Suppose that $X:=L^{2}\left(\mathbb{R}^{d}, d \lambda\right)$ and $O: \mathbb{R} \rightarrow O(d)$ is a one parameter semi-group of orthogonal operators. Set $\left(T_{O}(t) f\right)(x):=f(O(t) x)$ for all $f \in X$ and $x \in \mathbb{R}^{d}$. Then $T_{O}$ is also a strongly continuous unitary semi-group.

### 73.2 General Theory of Contraction Semigroups

For this section, let $(X,\|\cdot\|)$ be a Banach space with norm $\|\cdot\|$. Also let $T:=\{T(t)\}_{t>0}$ be a collection of bounded operators on $X$.
Definition 73.9. Let $X$ and $T$ be as above.

1. $T$ is a semi-group if $T(t+s)=T(t) T(s)$ for all $s, t>0$.
2. A semi-group $T$ is strongly continuous if $\lim _{t \downarrow 0} T(t) v=v$ for all $v \in X$. By convention if $T$ is strongly continuous, set $T(0) \equiv I$-the identity operator on $X$.
3. A semi-group $T$ is a contraction semi-group if $\|T(t)\| \leq 1$ for $t>0$.

Definition 73.10. Suppose that $T$ is a contraction semi group. Set

$$
\mathcal{D}(L):=\left\{v \in X:\left.\frac{d}{d t}\right|_{0^{+}} T(t) v \text { exists in } X\right\}
$$

and for $v \in \mathcal{D}(L)$ set $L v:=\left.\frac{d}{d t}\right|_{0^{+}} T(t) v . L$ is called the infinitesimal generator of $T$.

Proposition 73.11. Let $T$ be a strongly continuous contraction semi-group, then

1. For all $v \in X, t \in[0, \infty) \rightarrow T(t) v \in X$ is continuous.
2. $D(L)$ is dense linear subspace of $X$.
3. Suppose that $v:[0, \infty) \rightarrow X$ is a continuous, then $w(t):=T(t) v(t)$ is also continuous on $[0, \infty)$.
Proof. By assumption $v(t):=T(t) v$ is continuous at $t=0$. For $t>0$ and $h>0$,

$$
\|v(t+h)-v(t)\|=\|T(t)(T(h)-I) v\| \leq\|v(h)-v\| \rightarrow 0 \text { as } h \downarrow 0 .
$$

Similarly if $h \in(0, t)$,

$$
\|v(t-h)-v(t)\|=\|T(t-h)(I-T(h)) v\| \leq\|v-v(h)\| \rightarrow 0 \text { as } h \downarrow 0 .
$$

This proves the first item.
Let $v \in X$ set $v_{s}:=\int_{0}^{s} T(\sigma) v d \sigma$, where, since $\sigma \rightarrow T(\sigma) v$ is continuous, the integral may be interpreted as $X$-valued Riemann integral. Note

$$
\left\|\frac{1}{s} v_{s}-v\right\|=\left\|\frac{1}{s} \int_{0}^{s}(T(\sigma) v-v) d \sigma\right\| \leq \frac{1}{|s|}\left|\int_{0}^{s}\|T(\sigma) v-v\| d \sigma\right| \rightarrow 0
$$

as $s \downarrow 0$, so that $\mathcal{D}:=\left\{v_{s}: s>0\right.$ and $\left.v \in X\right\}$ is dense in $X$. Moreover,

$$
\begin{aligned}
\left.\frac{d}{d t}\right|_{0+} T(t) v_{s} & =\left.\frac{d}{d t}\right|_{0+} \int_{0}^{s} T(t+\sigma) v d \sigma \\
& =\left.\frac{d}{d t}\right|_{0+} \int_{t}^{t+s} T(\tau) v d \tau=T(s) v-v .
\end{aligned}
$$

Therefore $v_{s} \in \mathcal{D}(L)$ and $L v_{s}=T(s) v-v$. In particular, $\mathcal{D} \subset \mathcal{D}(L)$ and hence $\mathcal{D}(L)$ is dense in $X$. It is easily checked that $\mathcal{D}(L)$ is a linear subspace of $X$

Finally if $v:[0, \infty) \rightarrow X$ is a continuous function and $w(t):=T(t) v(t)$, then for $t>0$ and $h \in(-t, \infty)$,

$$
w(t+h)-w(t)=(T(t+h)-T(t)) v(t)+T(t+h)(v(t+h)-v(t))
$$

The first term goes to zero as $h \rightarrow 0$ by item 1 and the second term goes to zero since $v$ is continuous and $\|T(t+h)\| \leq 1$. The above argument also works with $t=0$ and $h \geq 0$.

Definition 73.12 (Closed Operators). A linear operator $L$ on $X$ is said to be closed if $\Gamma(L):=\{(v, L v) \in X \times X: v \in \mathcal{D}(L)\}$ is closed in the Banach space $X \times X$. Equivalently, $L$ is closed iff for all sequences $\left\{v_{n}\right\}_{n=1}^{\infty} \subset \mathcal{D}(L)$ such that $\lim _{n \rightarrow \infty} v_{n}=: v$ exists and $\lim _{n \rightarrow \infty} L v_{n}=: w$ exists implies that $v \in \mathcal{D}(L)$ and $L v=w$.

Roughly speaking, a closed operator is the next best thing to a bounded (i.e. continuous) operator. Indeed, the definition states that $L$ is closed iff

$$
\begin{equation*}
\lim _{n \rightarrow \infty} L v_{n}=L \lim _{n \rightarrow \infty} v_{n} \tag{73.6}
\end{equation*}
$$

provided both limits in (73.6) exist. While $L$ is continuous iff Eq. (73.6) holds whenever $\lim _{n \rightarrow \infty} v_{n}$ exists: part of the assertion being that the limit on the left side of Eq. (73.6) should exist.
Proposition 73.13 ( $L$ is Closed). Let $L$ be the infinitesimal generator of a contraction semi-group, then $L$ is closed.

Proof. Suppose that $v_{n} \in D(L), v_{n} \rightarrow v$, and $L v_{n} \rightarrow w$ in $X$ as $n \rightarrow \infty$. Then, using the fundamental theorem of calculus,

$$
\begin{aligned}
\frac{T(t) v-v}{t}=\lim _{n \rightarrow \infty} \frac{T(t) v_{n}-v_{n}}{t} & =\lim _{n \rightarrow \infty} \frac{1}{t} \int_{0}^{t} T(\tau) L v_{n} d \tau \\
& \rightarrow \frac{1}{t} \int_{0}^{t} T(\tau) w d \tau
\end{aligned}
$$

Therefore $v \in D(L)$ and $L v=w$.

## Theorem 73.14 (Solution Operator).

1. For any $t>0$ and $v \in \mathcal{D}(L), T(t) v \in \mathcal{D}(L)$ and $\frac{d}{d t} T(t) v=L T(t) v$.
2. Moreover if $v \in \mathcal{D}(L)$, then $\frac{d}{d t} T(t) v=T(t) L v$.

Proof. $T(t): D(L) \rightarrow D(L)$ and $\frac{d}{d t} T(t) v=L T(t) v=T(t) L v$. Suppose that $v \in \mathcal{D}(L)$, then

$$
\frac{T(t+h)-T(t)}{h} v=\frac{(T(h)-I)}{h} T(t) v=\frac{T(t)(T(h)-I)}{h} v .
$$

$$
\begin{aligned}
\frac{w(t+h)-w(t)}{h} & =\frac{1}{h}[P(t+h) v(t+h)-P(t) v(t)] \\
& =\frac{(P(t+h)-P(t))}{h} v(t)+P(t+h) \frac{(v(t+h)-v(t))}{h} \\
& =\frac{(P(t+h)-P(t))}{h} v(t)+P(t+h)(\dot{v}(t)+\epsilon(h)) \\
& \rightarrow-P(t) L v(t)+P(t) \dot{v}(t) \text { as } h \rightarrow 0,
\end{aligned}
$$

wherein we have used $\|P(t+h) \epsilon(h)\| \leq\|\epsilon(h)\| \rightarrow 0$ as $h \rightarrow 0$. Hence we have shown that

$$
\dot{w}(t)=-P(t) L v(t)+P(t) \dot{v}(t)=-P(t) L v(t)+P(t) L v(t)=0
$$

Therefore $w(t)=T(t-t) v(t)$ is constant or $(0, t)$ and hence by continuity of $w(\tau)=w(0)$, i.e.

$$
v(\tau)=w(\tau)=w(0)=T(\tau) v(0)=T(\tau) v_{0}
$$

This proves uniqueness.
Corollary 73.17. Suppose that $T$ and $\hat{T}$ are two strongly continuous contraction semi-groups on a Banach space $X$ which have the same infinitesimal generators $L$. Then $T=\hat{T}$.

Proof. Let $v_{0} \in D(L)$ then $v(t)=T(t) v$ and $\hat{v}(t)=\hat{T}(t) \hat{v}$ both solve Eq. (73.8 with initial condition $v_{0}$. By Theorem $73.16, v=\hat{v}$ which implies that $T(t) v_{0}=\hat{T}(t) v_{0}$, i.e., $T=\hat{T}$.

Because of the last corollary the following notion is justified.
Notation 73.18 If $T$ is a strongly continuous contraction semi-group with infinitesimal generator $L$, we will write $T(t)$ as $e^{t L}$.

Remark 73.19. Since $T$ is a contraction, $L$ should be "negative." Thus, working informally,

$$
\int_{0}^{\infty} e^{-t \lambda} e^{t L} d t=\left.\frac{1}{L-\lambda} e^{t(L-\lambda)}\right|_{t=0} ^{\infty}=\frac{1}{\lambda-L}=(\lambda-L)^{-1} .
$$

Theorem 73.20. Suppose $T=e^{t L}$ is a strongly continuous contraction semigroup with infinitesimal generator $L$. For any $\lambda>0$ the integral

$$
\begin{equation*}
\int_{0}^{\infty} e^{-t \lambda} e^{t L} d t=: R_{\lambda} \tag{73.9}
\end{equation*}
$$

exists as a $B(X)$-valued improper Riemann integral. ${ }^{1}$ Moreover, $(\lambda-L)$ : $\mathcal{D}(L) \rightarrow X$ is an invertible operator, $(\lambda-L)^{-1}=R_{\lambda}$, and $\left\|R_{\lambda}\right\| \leq \lambda^{-1}$.

[^8]Proof. First notice that

$$
\int_{0}^{\infty} e^{-t \lambda}\left\|e^{t L}\right\| d t \leq \int_{0}^{\infty} e^{-t \lambda} d t=1 / \lambda
$$

Therefore the integral in Eq. (73.9) exists and the result, $R_{\lambda}$, satisfies $\left\|R_{\lambda}\right\| \leq$ $\lambda^{-1}$. So we now must show that $R_{\lambda}=(\lambda-L)^{-1}$.

Let $v \in X$ and $h>0$, then
$e^{h L} R_{\lambda} v=\int_{0}^{\infty} e^{-t \lambda} e^{(t+h) L} v d t=\int_{h}^{\infty} e^{-(t-h) \lambda} e^{t L} v d t=e^{h \lambda} \int_{h}^{\infty} e^{-t \lambda} e^{t L} v d t$.
Therefore

$$
\left.\frac{d}{d h}\right|_{0^{+}} e^{h L} R_{\lambda} v=-v+\int_{0}^{\infty} \lambda e^{-t \lambda} e^{t L} v d t=-v+\lambda R_{\lambda} v
$$

which shows that $R_{\lambda} v \in \mathcal{D}(L)$ and that $L R_{\lambda} v=-v+\lambda R_{\lambda} v$. So $(\lambda-L) R_{\lambda}=I$. Similarly,

$$
\begin{equation*}
R_{\lambda} e^{h L} v=e^{h \lambda} \int_{h}^{\infty} e^{-t \lambda} e^{t L} v d t \tag{73.11}
\end{equation*}
$$

and hence if $v \in \mathcal{D}(L)$, then

$$
R_{\lambda} L v=\left.\frac{d}{d h}\right|_{0^{+}} R_{\lambda} e^{h L} v=-v+\lambda R_{\lambda} v
$$

Hence $R_{\lambda}(\lambda-L)=I_{\mathcal{D}(L)}$.
Before continuing it will be useful to record some properties of the resolvent operators $R_{\lambda}:=(\lambda-L)^{-1}$. Again working formally for the moment, if $\lambda, \mu \in$ $(0, \infty)$, then we expect

$$
R_{\lambda}-R_{\mu}=\frac{1}{\lambda-L}-\frac{1}{\mu-L}=\frac{\mu-L-(\lambda-L)}{(\lambda-L)(\mu-L)}=(\mu-\lambda) R_{\lambda} R_{\mu}
$$

For each $\lambda>0$ define $L_{\lambda}:=\lambda L R_{\lambda}$. Working again formally we have that

$$
L_{\lambda}=\frac{\lambda L}{\lambda-L}=\frac{\lambda(L-\lambda+\lambda)}{\lambda-L}=-\lambda+\lambda^{2} R_{\lambda}
$$

and

$$
\frac{d}{d \lambda} L_{\lambda}=\frac{L(\lambda-L)-\lambda L}{(\lambda-L)^{2}}=-\frac{L^{2}}{(\lambda-L)^{2}}=-L R_{\lambda} L R_{\lambda} .
$$

These equations will be verified in the following lemma.
Lemma 73.21. Let $L: X \rightarrow X$ be an operator on $X$ such that for all $\lambda>0$, $\lambda-L$ is invertible with a bounded inverse $R_{\lambda}$ or $(\lambda-L)^{-1}$. Set $L_{\lambda}:=\lambda L R_{\lambda}$. Then

1. for $\lambda, \mu \in(0, \infty)$,

$$
\begin{equation*}
R_{\lambda}-R_{\mu}=(\mu-\lambda) R_{\lambda} R_{\mu} \tag{7.12}
\end{equation*}
$$

and in particular $R_{\lambda}$ and $R_{\mu}$ commute,
2. $L_{\lambda}=-\lambda+\lambda^{2} R_{\lambda}$, and
3. $\frac{d}{d \lambda} L_{\lambda}=-L R_{\lambda} L R_{\lambda}$.

Proof. Since $\lambda-L$ is invertible, $\lambda-L$ is injective. So in order to verify Eq. (73.12 it suffices to verify:

$$
\begin{equation*}
(\lambda-L)\left(R_{\lambda}-R_{\mu}\right)=(\mu-\lambda)(\lambda-L) R_{\lambda} R_{\mu} \tag{73.13}
\end{equation*}
$$

Now
$(\lambda-L)\left(R_{\lambda}-R_{\mu}\right)=I-(\lambda-\mu+\mu-L) R_{\mu}=I-(\lambda-\mu) R_{\mu}-I=-(\lambda-\mu) R_{\mu}$,
while

$$
(\mu-\lambda)(\lambda-L) R_{\lambda} R_{\mu}=(\mu-\lambda) R_{\mu}
$$

Clearly the last two equations show that Eq. (73.13) holds. The second item is easily verified since, $L_{\lambda}=\lambda(L-\lambda+\lambda) R_{\lambda}=-\lambda+\lambda^{2} R_{\lambda}$.

For the third item, first recall that $\lambda \rightarrow R_{\lambda}$ is continuous in the operator norm topology (in fact analytic). To see this let us first work informally,
$R_{\lambda+h}=\frac{1}{\lambda+h-L}=\frac{1}{(\lambda-L+h)}=\frac{1}{(\lambda-L)\left(I+h(\lambda-L)^{-1}\right)}=R_{\lambda}\left(I+h R_{\lambda}\right)^{-1}$.
To verify this last equation, first notice that for sufficiently small $h,\left\|h R_{\lambda}\right\|<$ 1 , so that $\sum_{n=0}^{\infty}\left\|-h R_{\lambda}\right\|^{n}<\infty$ and hence $\left(I+h R_{\lambda}\right)$ is invertible and

$$
\left(I+h R_{\lambda}\right)^{-1}=\sum_{n=0}^{\infty}\left(-h R_{\lambda}\right)^{n}
$$

To verify the ends of Eq. (73.14) are equal it suffices to verify that $R_{\lambda+h}(I+$ $\left.h R_{\lambda}\right)=R_{\lambda}$, i.e., $R_{\lambda+h}-R_{\lambda}=-h R_{\lambda+h} R_{\lambda}$. But this last equation follows directly form (73.12). Therefore, we have shown that for $h$ sufficiently close to zero,

$$
R_{\lambda+h}=R_{\lambda} \sum_{n=0}^{\infty}\left(-h R_{\lambda}\right)^{n}
$$

Differentiating this last equation at $h=0$ shows that

$$
\begin{equation*}
\frac{d}{d \lambda} R_{\lambda}=\left.\frac{d}{d h}\right|_{0} R_{\lambda+h}=-R_{\lambda}^{2} \tag{73.15}
\end{equation*}
$$

We now may easily compute:

$$
\frac{d}{d \lambda} L_{\lambda}=\frac{d}{d \lambda}\left(-\lambda+\lambda^{2} R_{\lambda}\right)=-I+2 \lambda R_{\lambda}-\lambda^{2} R_{\lambda}^{2}=-\left(\lambda R_{\lambda}-I\right)^{2} .
$$

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This finishes the proof since,

$$
\lambda R_{\lambda}-I=\lambda R_{\lambda}-(\lambda-L) R_{\lambda}=L R_{\lambda}
$$

We now show that $L_{\lambda}$ is a good approximation to $L$ when $\lambda \rightarrow \infty$.
Proposition 73.22. Let $L$ be an operator on $X$ such that for each $\lambda \in(0, \infty)$, $R_{\lambda}:=(\lambda-L)^{-1}$ exists as a bounded operator and $\left\|(\lambda-L)^{-1}\right\| \leq \lambda^{-1}$. Then $\lambda R_{\lambda} \rightarrow I$ strongly as $\lambda \rightarrow \infty$ and for all $v \in \mathcal{D}(L)$

$$
\begin{equation*}
\lim _{\lambda \rightarrow \infty} L_{\lambda} v=L v \tag{73.16}
\end{equation*}
$$

Proof. First notice, informally, that

$$
\lambda R_{\lambda}=\frac{\lambda}{\lambda-L}=\frac{\lambda-L+L}{\lambda-L}=I+R_{\lambda} L .
$$

So we expect that

$$
\begin{equation*}
\left.\lambda R_{\lambda}\right|_{\mathcal{D}(L)}=I+R_{\lambda} L . \tag{73.17}
\end{equation*}
$$

(This last equation is easily verified by applying $(\lambda-L)$ to both sides of the equation.) Hence, for $v \in \mathcal{D}(L)$,

$$
\lambda R_{\lambda} v=v+R_{\lambda} L v
$$

and $\left\|R_{\lambda} L v\right\| \leq\|L v\| / \lambda$. Thus $\lim _{\lambda \rightarrow \infty} \lambda R_{\lambda} v=v$ for all $v \in \mathcal{D}(L)$. Using the fact that $\left\|\lambda R_{\lambda}\right\| \leq 1$ and a standard $3 \epsilon$-argument, it follows that $\lambda R_{\lambda}$ converges strongly to $I$ as $\lambda \rightarrow \infty$. Finally, for $v \in \mathcal{D}(L)$,

$$
L_{\lambda} v=\lambda L R_{\lambda} v=\lambda R_{\lambda} L v \rightarrow L v \text { as } \lambda \rightarrow \infty
$$

See Dynkin,
Lemma 73.23. Suppose that $A$ and $B$ are commuting bounded operators on a Banach space, $X$, such that $\left\|e^{t A}\right\|$ and $\left\|e^{t B}\right\|$ are bounded by 1 for all $t>0$, then

$$
\begin{equation*}
\left\|\left(e^{t A}-e^{t B}\right) v\right\| \leq t\|(A-B) v\| \text { for all } v \in X \tag{73.18}
\end{equation*}
$$

Proof. The fundamental theorem of calculus implies that

$$
e^{-t A} e^{t B}-I=\int_{0}^{t} \frac{d}{d \tau} e^{-\tau A} e^{\tau B} d \tau=\int_{0}^{t} e^{-\tau A}(-A+B) e^{\tau B} d \tau
$$

and hence, by multiplying on the left by $e^{t A}$,

$$
\begin{aligned}
e^{t B}-e^{t A} & =\int_{0}^{t} e^{(t-\tau) A}(-A+B) e^{\tau B} d \tau \\
& =\int_{0}^{t} e^{(t-\tau) A} e^{\tau B}(-A+B) d \tau
\end{aligned}
$$

wherein the last line we have used the fact that $A$ and $B$ commute. Eq. (73.18) is an easy consequence of this equation.

Theorem 73.24 (Hile-Yosida). A closed linear operator $L$ or a Banach space $X$ generates a contraction semi-group iff for all $\lambda \in(0, \infty)$,

1. $(\lambda-L)^{-1}$ exists as a bounded operator and
2. $\left\|(\lambda-L)^{-1}\right\| \leq \lambda^{-1}$ for all $\lambda>0$.

Proof.

$$
T_{\lambda}(t):=e^{t L_{\lambda}}=e^{t L_{\lambda}}:=\sum_{n=0}^{\infty}\left(t L_{\lambda}\right)^{n} / n!
$$

The outline of the proof is: i) show that $T_{\lambda}(t)$ is a contraction for all $t>0$, ii) show for $t>0$ that $T_{\lambda}(t)$ converges strongly to an operator $T(t)$, iii) show $T(t)$ is a strongly continuous contraction semi-group, and iv) show that $L$ is the generator of $T$.

Step i) Using $L_{\lambda}=-\lambda+\lambda^{2} R_{\lambda}$, we find that $e^{t L_{\lambda}}=e^{-t \lambda} e^{t \lambda^{2} R_{\lambda}}$. Hence

$$
\left\|T_{\lambda}(t)\right\|=\left\|e^{t L_{\lambda}}\right\| \leq e^{-t \lambda} e^{t \lambda^{2}\left\|R_{\lambda}\right\|} \leq e^{-t \lambda} e^{t \lambda^{2} \lambda^{-1}}=1
$$

Step ii) Let $\alpha, \mu>0$ and $v \in \mathcal{D}(L)$, then by Lemma 73.23 and Proposition 73.22,

$$
\left\|\left(T_{\alpha}(t)-T_{\mu}(t)\right) v\right\| \leq t\left\|L_{\alpha} v-L_{\mu} v\right\| \rightarrow 0 \text { as } \alpha, \mu \rightarrow \infty
$$

This shows, for all $v \in \mathcal{D}(L)$, that $\lim _{\alpha \rightarrow \infty} T_{\alpha}(t) v$ exists uniformly for $t$ in compact subsets of $[0, \infty)$ For general $v \in X, w \in \mathcal{D}(L), \tau>0$, and $0 \leq t \leq \tau$, we have

$$
\begin{aligned}
\left\|\left(T_{\alpha}(t)-T_{\mu}(t)\right) v\right\| & \leq\left\|\left(T_{\alpha}(t)-T_{\mu}(t)\right) w\right\|+\left\|\left(T_{\alpha}(t)-T_{\mu}(t)\right)(v-w)\right\| \\
& \leq\left\|\left(T_{\alpha}(t)-T_{\mu}(t)\right) w\right\|+2\|v-w\| . \\
& \leq \tau\left\|L_{\alpha} w-L_{\mu} w\right\|+2\|v-w\| .
\end{aligned}
$$

Thus

$$
\limsup _{\alpha, \mu \rightarrow \infty} \sup _{t \in[0, \tau]}\left\|\left(T_{\alpha}(t)-T_{\mu}(t)\right) v\right\| \leq 2\|v-w\| \rightarrow 0 \text { as } w \rightarrow v
$$

Hence for each $v \in X, T(t) v:=\lim _{\alpha \rightarrow \infty} T_{\alpha}(t) v$ exists uniformly for $t$ in compact sets of $[0, \infty)$.

Step iii) It is now easily follows that $\|T(t)\| \leq 1$ and that $t \rightarrow T(t)$ is strongly continuous. Moreover,

$$
T(t+s) v=\lim _{\alpha \rightarrow \infty} T_{\alpha}(t+s) v=\lim _{\alpha \rightarrow \infty} T_{\alpha}(t) T_{\alpha}(s) v
$$

Letting $\epsilon(\alpha):=T_{\alpha}(s) v-T(s) v$, it follows that

$$
T(t+s) v=\lim _{\alpha \rightarrow \infty} T_{\alpha}(t)(T(s) v+\epsilon(\alpha))=T(t) T(s) v+\lim _{\alpha \rightarrow \infty} T_{\alpha}(t) \epsilon(\alpha)
$$

This shows that $T$ is also satisfies the semi-group property, since $\left\|T_{\alpha}(t) \epsilon(\alpha)\right\| \leq$ $\|\epsilon(\alpha)\|$ and $\lim _{\alpha \rightarrow \infty} \epsilon(\alpha)=0$.

$$
\begin{aligned}
\frac{P_{t} P_{\varphi} v-P_{\varphi} v}{t}=\int_{0}^{\infty} & \frac{\varphi(s-t)-\varphi(s)}{t} P_{s} v d s \\
& \rightarrow \int_{0}^{\infty}-\left.\varphi^{\prime}(s) P_{s} v d s \frac{d}{d t}\right|_{0} P_{t} P_{\varphi} v=-P_{\varphi^{\prime}} v
\end{aligned}
$$

so $P_{\varphi} v \in D(L)$ and $L P_{\varphi} v=-P_{\varphi^{\prime}} v$ Now $\varphi \rightarrow \delta_{0}$.
Proof of Theorem Key point is to prove uniqueness. Let $P_{t}^{*}=$ transpose semi group. Let $B=\left.\frac{d}{d t}\right|_{0+} P_{t}^{*}$ denote the adjoint generator.

Claim $B=L^{*}$.

$$
\begin{aligned}
\left.\varphi \in D(B) \Rightarrow \frac{d}{d t}\right|_{0}\left\langle P_{t}^{*} \varphi, v\right\rangle & =\langle B \varphi, v\rangle \\
\left.\frac{d}{d t}\right|_{0}\left\langle\varphi, P_{t} v\right\rangle & =\langle\varphi, L v\rangle \quad \forall v \in D(L t
\end{aligned}
$$

$\Rightarrow \varphi \in D\left(L^{*}\right)$ and $B \varphi=L \varphi$ so that $B \subseteq L^{*}$. Suppose $\varphi \in D\left(L^{*}\right) \Rightarrow$ $\exists \psi \ni L^{*} \varphi=\psi$. For example $\left\langle L^{*} \varphi, v\right\rangle=\langle\psi, v\rangle \quad \forall v \in D(L)$ for example $\langle\varphi, L v\rangle=\langle\psi, v\rangle$

$$
\frac{d}{d t}\left\langle\varphi, P_{t} v\right\rangle=\left\langle\psi, P_{t} v\right\rangle
$$

So $\left\langle\varphi, P_{t} v\right\rangle-\langle\varphi, v\rangle=\int_{0}^{t}\left\langle P_{t}^{*} \psi, v\right\rangle d t \quad \forall v \in D(L)$. So $P_{t}^{*} \varphi=\varphi+\int_{0}^{t} P_{t}^{*} \psi d t$. So $\left.\frac{d}{d t}\right|_{0+} P_{t}^{*} \varphi=\psi \Rightarrow \varphi \in D\left(L^{*}\right)$ and $L^{*} \varphi=\varphi$.

Uniqueness: Let $\dot{v}=L v, \quad v(0)=v_{0}$. Let $\varphi \in D\left(L^{*}\right), T>0$.

$$
\frac{d}{d t}\left\langle P_{T-t}^{*} \varphi, v(t)\right\rangle=\left\langle-L^{*} P_{T-t}^{*} \varphi, v(t)\right\rangle+\left\langle P_{T-t}^{*} \varphi, L v\right\rangle=0
$$

$\Rightarrow\left\langle P_{T-t}^{*} \varphi, v(T)\right\rangle=$ construction therefore $\left\langle P_{T}^{*} \varphi, v_{0}\right\rangle=\langle\varphi, v(T)\rangle \quad \forall \varphi \in$ $D\left(L^{*}\right)\left\langle\varphi^{\prime \prime}, P_{T} v_{0}\right\rangle$. Since $D\left(L^{*}\right)$ is dense $v(T)=P_{T} v_{0}$.

1234 A Multinomial Theorems and Calculus Results

## Multinomial Theorems and Calculus Results

Given a multi-index $\alpha \in \mathbb{Z}_{+}^{n}$, let $|\alpha|=\alpha_{1}+\cdots+\alpha_{n}, \alpha!:=\alpha_{1}!\cdots \alpha_{n}!$,

$$
x^{\alpha}:=\prod_{j=1}^{n} x_{j}^{\alpha_{j}} \text { and } \partial_{x}^{\alpha}=\left(\frac{\partial}{\partial x}\right)^{\alpha}:=\prod_{j=1}^{n}\left(\frac{\partial}{\partial x_{j}}\right)^{\alpha_{j}}
$$

We also write

$$
\partial_{v} f(x):=\left.\frac{d}{d t} f(x+t v)\right|_{t=0}
$$

## A. 1 Multinomial Theorems and Product Rules

For $a=\left(a_{1}, a_{2}, \ldots, a_{n}\right) \in \mathbb{C}^{n}, m \in \mathbb{N}$ and $\left(i_{1}, \ldots, i_{m}\right) \in\{1,2, \ldots, n\}^{m}$ let $\hat{\alpha}_{j}\left(i_{1}, \ldots, i_{m}\right)=\#\left\{k: i_{k}=j\right\}$. Then

$$
\left(\sum_{i=1}^{n} a_{i}\right)^{m}=\sum_{i_{1}, \ldots, i_{m}=1}^{n} a_{i_{1}} \ldots a_{i_{m}}=\sum_{|\alpha|=m} C(\alpha) a^{\alpha}
$$

where

$$
C(\alpha)=\#\left\{\left(i_{1}, \ldots, i_{m}\right): \hat{\alpha}_{j}\left(i_{1}, \ldots, i_{m}\right)=\alpha_{j} \text { for } j=1,2, \ldots, n\right\}
$$

I claim that $C(\alpha)=\frac{m!}{\alpha!}$. Indeed, one possibility for such a sequence $\left(a_{1}, \ldots, a_{i_{m}}\right)$ for a given $\alpha$ is gotten by choosing

$$
(\overbrace{a_{1}, \ldots, a_{1}}^{\alpha_{1}}, \overbrace{a_{2}, \ldots, a_{2}}^{\alpha_{2}}, \ldots, \overbrace{a_{n}, \ldots, a_{n}}^{\alpha_{n}}) .
$$

Now there are $m$ ! permutations of this list. However, only those permutations leading to a distinct list are to be counted. So for each of these $m$ ! permutations we must divide by the number of permutation which just rearrange the
groups of $a_{i}$ 's among themselves for each $i$. There are $\alpha!:=\alpha_{1}!\cdots \alpha_{n}!$ such permutations. Therefore, $C(\alpha)=m!/ \alpha!$ as advertised. So we have proved

$$
\begin{equation*}
\left(\sum_{i=1}^{n} a_{i}\right)^{m}=\sum_{|\alpha|=m} \frac{m!}{\alpha!} a^{\alpha} . \tag{A.1}
\end{equation*}
$$

Now suppose that $a, b \in \mathbb{R}^{n}$ and $\alpha$ is a multi-index, we have

$$
\begin{equation*}
(a+b)^{\alpha}=\sum_{\beta \leq \alpha} \frac{\alpha!}{\beta!(\alpha-\beta)!} a^{\beta} b^{\alpha-\beta}=\sum_{\beta+\delta=\alpha} \frac{\alpha!}{\beta!\delta!} a^{\beta} b^{\delta} \tag{A.2}
\end{equation*}
$$

Indeed, by the standard Binomial formula,

$$
\left(a_{i}+b_{i}\right)^{\alpha_{i}}=\sum_{\beta_{i} \leq \alpha_{i}} \frac{\alpha_{i}!}{\beta_{i}!\left(\alpha_{i}-\beta_{i}\right)!} a^{\beta_{i}} b^{\alpha_{i}-\beta_{i}}
$$

from which Eq. (A.2) follows. Eq. (A.2) generalizes in the obvious way to

$$
\begin{equation*}
\left(a_{1}+\cdots+a_{k}\right)^{\alpha}=\sum_{\beta_{1}+\cdots+\beta_{k}=\alpha} \frac{\alpha!}{\beta_{1}!\cdots \beta_{k}!} a_{1}^{\beta_{1}} \cdots a_{k}^{\beta_{k}} \tag{A.3}
\end{equation*}
$$

where $a_{1}, a_{2}, \ldots, a_{k} \in \mathbb{R}^{n}$ and $\alpha \in \mathbb{Z}_{+}^{n}$.
Now let us consider the product rule for derivatives. Let us begin with the one variable case (write $d^{n} f$ for $f^{(n)}=\frac{d^{n}}{d x^{n}} f$ ) where we will show by induction that

$$
\begin{equation*}
d^{n}(f g)=\sum_{k=0}^{n}\binom{n}{k} d^{k} f \cdot d^{n-k} g \tag{A.4}
\end{equation*}
$$

Indeed assuming Eq. (A.4) we find

$$
\begin{aligned}
d^{n+1}(f g) & =\sum_{k=0}^{n}\binom{n}{k} d^{k+1} f \cdot d^{n-k} g+\sum_{k=0}^{n}\binom{n}{k} d^{k} f \cdot d^{n-k+1} g \\
& =\sum_{k=1}^{n+1}\binom{n}{k-1} d^{k} f \cdot d^{n-k+1} g+\sum_{k=0}^{n}\binom{n}{k} d^{k} f \cdot d^{n-k+1} g \\
& =\sum_{k=1}^{n+1}\left[\binom{n}{k-1}+\binom{n}{k}\right] d^{k} f \cdot d^{n-k+1} g+d^{n+1} f \cdot g+f \cdot d^{n+1} g
\end{aligned}
$$

Since

$$
\begin{aligned}
\binom{n}{k-1}+\binom{n}{k} & =\frac{n!}{(n-k+1)!(k-1)!}+\frac{n!}{(n-k)!k!} \\
& =\frac{n!}{(k-1)!(n-k)!}\left[\frac{1}{(n-k+1)}+\frac{1}{k}\right] \\
& =\frac{n!}{(k-1)!(n-k)!} \frac{n+1}{(n-k+1) k}=\binom{n+1}{k}
\end{aligned}
$$

the result follows
Now consider the multi-variable case

$$
\begin{aligned}
\partial^{\alpha}(f g) & =\left(\prod_{i=1}^{n} \partial_{i}^{\alpha_{i}}\right)(f g)=\prod_{i=1}^{n}\left[\sum_{k_{i}=0}^{\alpha_{i}}\binom{\alpha_{i}}{k_{i}} \partial_{i}^{k_{i}} f \cdot \partial_{i}^{\alpha_{i}-k_{i}} g\right] \\
& =\sum_{k_{1}=0}^{\alpha_{1}} \cdots \sum_{k_{n}=0}^{\alpha_{n}} \prod_{i=1}^{n}\binom{\alpha_{i}}{k_{i}} \partial^{k} f \cdot \partial^{\alpha-k} g=\sum_{k \leq \alpha}\binom{\alpha}{k} \partial^{k} f \cdot \partial^{\alpha-k} g
\end{aligned}
$$

where $k=\left(k_{1}, k_{2}, \ldots, k_{n}\right)$ and

$$
\binom{\alpha}{k}:=\prod_{i=1}^{n}\binom{\alpha_{i}}{k_{i}}=\frac{\alpha!}{k!(\alpha-k)!} .
$$

So we have proved

$$
\begin{equation*}
\partial^{\alpha}(f g)=\sum_{\beta \leq \alpha}\binom{\alpha}{\beta} \partial^{\beta} f \cdot \partial^{\alpha-\beta} g \tag{A.5}
\end{equation*}
$$

## A. 2 Taylor's Theorem

Theorem A.1. Suppose $X \subset \mathbb{R}^{n}$ is an open set, $x:[0,1] \rightarrow X$ is a $C^{1}-$ path, and $f \in C^{N}(X, \mathbb{C})$. Let $v_{s}:=x(1)-x(s)$ and $v=v_{1}=x(1)-x(0)$, then

$$
\begin{equation*}
f(x(1))=\sum_{m=0}^{N-1} \frac{1}{m!}\left(\partial_{v}^{m} f\right)(x(0))+R_{N} \tag{A.6}
\end{equation*}
$$

where

$$
\begin{equation*}
R_{N}=\frac{1}{(N-1)!} \int_{0}^{1}\left(\partial_{\dot{x}(s)} \partial_{v_{s}}^{N-1} f\right)(x(s)) d s=\frac{1}{N!} \int_{0}^{1}\left(-\frac{d}{d s} \partial_{v_{s}}^{N} f\right)(x(s)) d s \tag{A.7}
\end{equation*}
$$ and $0!:=1$

Proof. By the fundamental theorem of calculus and the chain rule,

$$
f(x(t))=f(x(0))+\int_{0}^{t} \frac{d}{d s} f(x(s)) d s=f(x(0))+\int_{0}^{t}\left(\partial_{\dot{x}(s)} f\right)(x(s)) d s
$$

and in particular,

$$
f(x(1))=f(x(0))+\int_{0}^{1}\left(\partial_{\dot{x}(s)} f\right)(x(s)) d s
$$

This proves Eq. (A.6) when $N=1$. We will now complete the proof using induciton on $N$.

Applying Eq. (A.8) with $f$ replaced by $\frac{1}{(N-1)!}\left(\partial_{\dot{x}(s)} \partial_{v_{s}}^{N-1} f\right)$ gives

$$
\begin{aligned}
\frac{1}{(N-1)!}\left(\partial_{\dot{x}(s)} \partial_{v_{s}}^{N-1} f\right)(x(s))= & \frac{1}{(N-1)!}\left(\partial_{\dot{x}(s)} \partial_{v_{s}}^{N-1} f\right)(x(0)) \\
& +\frac{1}{(N-1)!} \int_{0}^{s}\left(\partial_{\dot{x}(s)} \partial_{v_{s}}^{N-1} \partial_{\dot{x}(t)} f\right)(x(t)) d t \\
= & -\frac{1}{N!}\left(\frac{d}{d s} \partial_{v_{s}}^{N} f\right)(x(0))-\frac{1}{N!} \int_{0}^{s}\left(\frac{d}{d s} \partial_{v_{s}}^{N} \partial_{\dot{x}(t)} f\right)
\end{aligned}
$$

wherein we have used the fact that mixed partial derivatives commute to show $\frac{d}{d s} \partial_{v_{s}}^{N} f=N \partial_{\dot{x}(s)} \partial_{v_{s}}^{N-1} f$. Integrating this equation on $s \in[0,1]$ shows, using the fundamental theorem of calculus,

$$
\begin{aligned}
R_{N} & =\frac{1}{N!}\left(\partial_{v}^{N} f\right)(x(0))-\frac{1}{N!} \int_{0 \leq t \leq s \leq 1}\left(\frac{d}{d s} \partial_{v_{s}}^{N} \partial_{\dot{x}(t)} f\right)(x(t)) d s d t \\
& =\frac{1}{N!}\left(\partial_{v}^{N} f\right)(x(0))+\frac{1}{(N+1)!} \int_{0 \leq t \leq 1}\left(\partial_{w_{t}}^{N} \partial_{\dot{x}(t)} f\right)(x(t)) d t \\
& =\frac{1}{N!}\left(\partial_{v}^{N} f\right)(x(0))+R_{N+1}
\end{aligned}
$$

which completes the inductive proof.
Remark A.2. Using Eq. (A.1) with $a_{i}$ replaced by $v_{i} \partial_{i}$ (although $\left\{v_{i} \partial_{i}\right\}_{i=1}^{n}$ are not complex numbers they are commuting symbols), we find

$$
\partial_{v}^{m} f=\left(\sum_{i=1}^{n} v_{i} \partial_{i}\right)^{m} f=\sum_{|\alpha|=m} \frac{m!}{\alpha!} v^{\alpha} \partial^{\alpha} .
$$

Using this fact we may write Eqs. (A.6) and (A.7) as

$$
f(x(1))=\sum_{|\alpha| \leq N-1} \frac{1}{\alpha!} v^{\alpha} \partial^{\alpha} f(x(0))+R_{N}
$$

and

$$
R_{N}=\sum_{|\alpha|=N} \frac{1}{\alpha!} \int_{0}^{1}\left(-\frac{d}{d s} v_{s}^{\alpha} \partial^{\alpha} f\right)(x(s)) d s
$$

Corollary A.3. Suppose $X \subset \mathbb{R}^{n}$ is an open set which contains $x(s)=(1-$ s) $x_{0}+s x_{1}$ for $0 \leq s \leq 1$ and $f \in C^{N}(X, \mathbb{C})$. Then

$$
\begin{align*}
f\left(x_{1}\right) & =\sum_{m=0}^{N-1} \frac{1}{m!}\left(\partial_{v}^{m} f\right)\left(x_{0}\right)+\frac{1}{N!} \int_{0}^{1}\left(\partial_{v}^{N} f\right)(x(s)) d \nu_{N}(s) \\
& =\sum_{|\alpha|<N} \frac{1}{\alpha!} \partial^{\alpha} f(x(0))\left(x_{1}-x_{0}\right)^{\alpha}+\sum_{\alpha:|\alpha|=N} \frac{1}{\alpha!}\left[\int_{0}^{1} \partial^{\alpha} f(x(s)) d \nu_{N}(s)\right]\left(x_{1}-\alpha\right. \tag{A.10}
\end{align*}
$$

where $v:=x_{1}-x_{0}$ and $d \nu_{N}$ is the probability measure on $[0,1]$ given by

$$
\begin{equation*}
d \nu_{N}(s):=N(1-s)^{N-1} d s \tag{A.11}
\end{equation*}
$$

If we let $x=x_{0}$ and $y=x_{1}-x_{0}\left(\right.$ so $\left.x+y=x_{1}\right) E q$. (A.10) may be written as

$$
\begin{equation*}
f(x+y)=\sum_{|\alpha|<N} \frac{\partial_{x}^{\alpha} f(x)}{\alpha!} y^{\alpha}+\sum_{\alpha:|\alpha|=N} \frac{1}{\alpha!}\left(\int_{0}^{1} \partial_{x}^{\alpha} f(x+s y) d \nu_{N}(s)\right) y^{\alpha} . \tag{A.12}
\end{equation*}
$$

Proof. This is a special case of Theorem A.1. Notice that

$$
v_{s}=x(1)-x(s)=(1-s)\left(x_{1}-x_{0}\right)=(1-s) v
$$

and hence
$R_{N}=\frac{1}{N!} \int_{0}^{1}\left(-\frac{d}{d s}(1-s)^{N} \partial_{v}^{N} f\right)(x(s)) d s=\frac{1}{N!} \int_{0}^{1}\left(\partial_{v}^{N} f\right)(x(s)) N(1-s)^{N-1} d s$

Example A.4. Let $X=(-1,1) \subset \mathbb{R}, \beta \in \mathbb{R}$ and $f(x)=(1-x)^{\beta}$. The reader should verify

$$
f^{(m)}(x)=(-1)^{m} \beta(\beta-1) \ldots(\beta-m+1)(1-x)^{\beta-m}
$$

and therefore by Taylor's theorem (Eq. (??) with $x=0$ and $y=x$ )

$$
\begin{equation*}
(1-x)^{\beta}=1+\sum_{m=1}^{N-1} \frac{1}{m!}(-1)^{m} \beta(\beta-1) \ldots(\beta-m+1) x^{m}+R_{N}(x) \tag{A.13}
\end{equation*}
$$

where

$$
\begin{aligned}
R_{N}(x) & =\frac{x^{N}}{N!} \int_{0}^{1}(-1)^{N} \beta(\beta-1) \ldots(\beta-N+1)(1-s x)^{\beta-N} d \nu_{N}(s) \\
& =\frac{x^{N}}{N!}(-1)^{N} \beta(\beta-1) \ldots(\beta-N+1) \int_{0}^{1} \frac{N(1-s)^{N-1}}{(1-s x)^{N-\beta}} d s .
\end{aligned}
$$

Now for $x \in(-1,1)$ and $N>\beta$,

$$
0 \leq \int_{0}^{1} \frac{N(1-s)^{N-1}}{(1-s x)^{N-\beta}} d s \leq \int_{0}^{1} \frac{N(1-s)^{N-1}}{(1-s)^{N-\beta}} d s=\int_{0}^{1} N(1-s)^{\beta-1} d s=\frac{N}{\beta}
$$

and therefore,

$$
\left|R_{N}(x)\right| \leq \frac{|x|^{N}}{(N-1)!}|(\beta-1) \ldots(\beta-N+1)|=: \rho_{N}
$$

Since

$$
\lim \sup _{N \rightarrow \infty} \frac{\rho_{N+1}}{\rho_{N}}=|x| \cdot \lim \sup _{N \rightarrow \infty} \frac{N-\beta}{N}=|x|<1
$$

and so by the Ratio test, $\left|R_{N}(x)\right| \leq \rho_{N} \rightarrow 0$ (exponentially fast) as $N \rightarrow \infty$. Therefore by passing to the limit in Eq. (A.13) we have proved

$$
\begin{equation*}
(1-x)^{\beta}=1+\sum_{m=1}^{\infty} \frac{(-1)^{m}}{m!} \beta(\beta-1) \ldots(\beta-m+1) x^{m} \tag{A.14}
\end{equation*}
$$

which is valid for $|x|<1$ and $\beta \in \mathbb{R}$. An important special cases is $\beta=-1$ in which case, Eq. (A.14) becomes $\frac{1}{1-x}=\sum_{m=0}^{\infty} x^{m}$, the standard geometric series formula. Another another useful special case is $\beta=1 / 2$ in which case Eq. (A.14) becomes

$$
\begin{align*}
\sqrt{1-x} & =1+\sum_{m=1}^{\infty} \frac{(-1)^{m}}{m!} \frac{1}{2}\left(\frac{1}{2}-1\right) \ldots\left(\frac{1}{2}-m+1\right) x^{m} \\
& =1-\sum_{m=1}^{\infty} \frac{(2 m-3)!!}{2^{m} m!} x^{m} \text { for all }|x|<1 \tag{A.15}
\end{align*}
$$

## Zorn's Lemma and the Hausdorff Maximal Principle

Definition B.1. A partial order $\leq$ on $X$ is a relation with following properties
(i) If $x \leq y$ and $y \leq z$ then $x \leq z$.
(ii) If $x \leq y$ and $y \leq x$ then $x=y$.
(iii) $\leq x$ for all $x \in X$.

Example B.2. Let $Y$ be a set and $X=\mathcal{P}(Y)$. There are two natural partial orders on $X$.

1. Ordered by inclusion, $A \leq B$ is $A \subset B$ and
2. Ordered by reverse inclusion, $A \leq B$ if $B \subset A$.

Definition B.3. Let $(X, \leq)$ be a partially ordered set we say $X$ is linearly a totally ordered if for all $x, y \in X$ either $x \leq y$ or $y \leq x$. The real numbers $\mathbb{R}$ with the usual order $\leq$ is a typical example.

Definition B.4. Let $(X, \leq)$ be a partial ordered set. We say $x \in X$ is a maximal element if for all $y \in X$ such that $y \geq x$ implies $y=x$, i.e. there is no element larger than $x$. An upper bound for a subset $E$ of $X$ is an element $x \in X$ such that $x \geq y$ for all $y \in E$.

Example B.5. Let

$$
X=\{a=\{1\} b=\{1,2\} c=\{3\} d=\{2,4\} e=\{2\}\}
$$

ordered by set inclusion. Then $b$ and $d$ are maximal elements despite that fact that $b \not \leq a$ and $a \not \leq b$. We also have

- If $E=\{a, e, c\}$, then $E$ has no upper bound.

Definition B.6. - If $E=\{a, e\}$, then $b$ is an upper bound.

- $E=\{e\}$, then $b$ and d are upper bounds.

Theorem B.7. The following are equivalent.

1. The axiom of choice: to each collection, $\left\{X_{\alpha}\right\}_{\alpha \in A}$, of non-empty sets there exists a"choice function," $x: A \rightarrow \coprod X_{\alpha}$ such that $x(\alpha) \in X_{\alpha}$ for all $\alpha \in A$, i.e. $\prod_{\alpha \in A} X_{\alpha} \neq \emptyset$.
2. The Hausdorff Maximal Principle: Every partially ordered set has a maximal (relative to the inclusion order) linearly ordered subset.
3. Zorn's Lemma: If $X$ is partially ordered set such that every linearly ordered subset of $X$ has an upper bound, then $X$ has a maximal element. ${ }^{1}$

Proof. $(2 \Rightarrow 3)$ Let $X$ be a partially ordered subset as in 3 and let $\mathcal{F}=$ $\{E \subset X: E$ is linearly ordered $\}$ which we equip with the inclusion partial ordering. By 2. there exist a maximal element $E \in \mathcal{F}$. By assumption, the linearly ordered set $E$ has an upper bound $x \in X$. The element $x$ is maximal, for if $y \in Y$ and $y \geq x$, then $E \cup\{y\}$ is still an linearly ordered set containing $E$. So by maximality of $E, E=E \cup\{y\}$, i.e. $y \in E$ and therefore $y \leq x$ showing which combined with $y \geq x$ implies that $y=x .{ }^{2}$
$(3 \Rightarrow 1)$ Let $\left\{X_{\alpha}\right\}_{\alpha \in A}$ be a collection of non-empty sets, we must show $\prod_{\alpha \in A} X_{\alpha}$ is not empty. Let $\mathcal{G}$ denote the collection of functions $g: D(g) \rightarrow$ $\coprod_{\alpha \in A} X_{\alpha}$ such that $D(g)$ is a subset of $A$, and for all $\alpha \in D(g), g(\alpha) \in X_{\alpha}$. Notice that $\mathcal{G}$ is not empty, for we may let $\alpha_{0} \in A$ and $x_{0} \in X_{\alpha}$ and then set $D(g)=\left\{\alpha_{0}\right\}$ and $g\left(\alpha_{0}\right)=x_{0}$ to construct an element of $\mathcal{G}$. We now put a partial order on $\mathcal{G}$ as follows. We say that $f \leq g$ for $f, g \in \mathcal{G}$ provided that $D(f) \subset D(g)$ and $f=\left.g\right|_{D(f)}$. If $\Phi \subset \mathcal{G}$ is a linearly ordered set, let $D(h)=\cup_{g \in \Phi} D(g)$ and for $\alpha \in D(g)$ let $h(\alpha)=g(\alpha)$. Then $h \in \mathcal{G}$ is an upper bound for $\Phi$. So by Zorn's Lemma there exists a maximal element $h \in \mathcal{G}$. To finish the proof we need only show that $D(h)=A$. If this were not the case, then let $\alpha_{0} \in A \backslash D(h)$ and $x_{0} \in X_{\alpha_{0}}$. We may now define $D(\tilde{h})=D(h) \cup\left\{\alpha_{0}\right\}$ and

$$
\tilde{h}(\alpha)=\left\{\begin{array}{cl}
h(\alpha) & \text { if } \alpha \in D(h) \\
x_{0} & \text { if } \alpha=\alpha_{0}
\end{array}\right.
$$

[^9]$$
\mathcal{F}_{0}=\cap\left\{\mathcal{F}^{\prime}: \mathcal{F}^{\prime} \text { is a tower }\right\}
$$

Then $h \leq \tilde{h}$ while $h \neq \tilde{h}$ violating the fact that $h$ was a maximal element.
$(1 \Rightarrow 2)$ Let $(X, \leq)$ be a partially ordered set. Let $\mathcal{F}$ be the collection of linearly ordered subsets of $X$ which we order by set inclusion. Given $x_{0} \in X$, $\left\{x_{0}\right\} \in \mathcal{F}$ is linearly ordered set so that $\mathcal{F} \neq \emptyset$.

Fix an element $P_{0} \in \mathcal{F}$. If $P_{0}$ is not maximal there exists $P_{1} \in \mathcal{F}$ such that $P_{0} \nsubseteq P_{1}$. In particular we may choose $x \notin P_{0}$ such that $P_{0} \cup\{x\} \in \mathcal{F}$. The idea now is to keep repeating this process of adding points $x \in X$ until we construct a maximal element $P$ of $\mathcal{F}$. We now have to take care of some details.

We may assume with out loss of generality that $\mathcal{F}=\{P \in \mathcal{F}: P$ is not maxima is a non-empty set. For $P \in \tilde{\mathcal{F}}$, let $P^{*}=\{x \in X: P \cup\{x\} \in \mathcal{F}\}$. As the above argument shows, $P^{*} \neq \emptyset$ for all $P \in \mathcal{F}$. Using the axiom of choice, there exists $f \in \prod_{P \in \tilde{\mathcal{F}}} P^{*}$. We now define $g: \mathcal{F} \rightarrow \mathcal{F}$ by

$$
g(P)=\left\{\begin{array}{cl}
P & \text { if } P \text { is maximal }  \tag{B.1}\\
P \cup\{f(x)\} & \text { if } P \text { is not maximal. }
\end{array}\right.
$$

The proof is completed by Lemma B. 8 below which shows that $g$ must have a fixed point $P \in \mathcal{F}$. This fixed point is maximal by construction of $g$.

Lemma B.8. The function $g: \mathcal{F} \rightarrow \mathcal{F}$ defined in Eq. (B.1) has a fixed point. ${ }^{3}$
Proof. The idea of the proof is as follows. Let $P_{0} \in \mathcal{F}$ be chosen arbitrarily. Notice that $\Phi=\left\{g^{(n)}\left(P_{0}\right)\right\}_{n=0}^{\infty} \subset \mathcal{F}$ is a linearly ordered set and it is therefore easily verified that $P_{1}=\bigcup_{n=0}^{\infty} g^{(n)}\left(P_{0}\right) \in \mathcal{F}$. Similarly we may repeat the process to construct $P_{2}=\bigcup_{n=0}^{\infty} g^{(n)}\left(P_{1}\right) \in \mathcal{F}$ and $P_{3}=\bigcup_{n=0}^{\infty} g^{(n)}\left(P_{2}\right) \in \mathcal{F}$, etc. etc. Then take $P_{\infty}=\cup_{n=0}^{\infty}{ }_{n}^{n=0} P_{n}$ and start again with $P_{0}{ }^{n=0}$ replaced by $P_{\infty}$. Then keep going this way until eventually the sets stop increasing in size, in which case we have found our fixed point. The problem with this strategy is that we may never win. (This is very reminiscent of constructing measurable sets and the way out is to use measure theoretic like arguments.)

Let us now start the formal proof. Again let $P_{0} \in \mathcal{F}$ and let $\mathcal{F}_{1}=\{P \in$ $\left.\mathcal{F}: P_{0} \subset P\right\}$. Notice that $\mathcal{F}_{1}$ has the following properties:

1. $P_{0} \in \mathcal{F}_{1}$.
2. If $\Phi \subset \mathcal{F}_{1}$ is a totally ordered (by set inclusion) subset then $\cup \Phi \in \mathcal{F}_{1}$.
3. If $P \in \mathcal{F}_{1}$ then $g(P) \in \mathcal{F}_{1}$.

Let us call a general subset $\mathcal{F}^{\prime} \subset \mathcal{F}$ satisfying these three conditions a tower and let

[^10]Standard arguments show that $\mathcal{F}_{0}$ is still a tower and clearly is the smallest tower containing $P_{0}$. (Morally speaking $\mathcal{F}_{0}$ consists of all of the sets we were trying to constructed in the "idea section" of the proof.)

We now claim that $\mathcal{F}_{0}$ is a linearly ordered subset of $\mathcal{F}$. To prove this let $\Gamma \subset \mathcal{F}_{0}$ be the linearly ordered set

$$
\Gamma=\left\{C \in \mathcal{F}_{0}: \text { for all } A \in \mathcal{F}_{0} \text { either } A \subset C \text { or } C \subset A\right\}
$$

Shortly we will show that $\Gamma \subset \mathcal{F}_{0}$ is a tower and hence that $\mathcal{F}_{0}=\Gamma$. That is to say $\mathcal{F}_{0}$ is linearly ordered. Assuming this for the moment let us finish the proof. Let $P \equiv \cup \mathcal{F}_{0}$ which is in $\mathcal{F}_{0}$ by property 2 and is clearly the largest element in $\mathcal{F}_{0}$. By 3 . it now follows that $P \subset g(P) \in \mathcal{F}_{0}$ and by maximality of $P$, we have $g(P)=P$, the desired fixed point. So to finish the proof, we must show that $\Gamma$ is a tower.

First off it is clear that $P_{0} \in \Gamma$ so in particular $\Gamma$ is not empty. For each $C \in \Gamma$ let

$$
\Phi_{C}:=\left\{A \in \mathcal{F}_{0}: \text { either } A \subset C \text { or } g(C) \subset A\right\}
$$

We will begin by showing that $\Phi_{C} \subset \mathcal{F}_{0}$ is a tower and therefore that $\Phi_{C}=\mathcal{F}_{0}$.

1. $P_{0} \in \Phi_{C}$ since $P_{0} \subset C$ for all $C \in \Gamma \subset \mathcal{F}_{0}$. 2. If $\Phi \subset \Phi_{C} \subset \mathcal{F}_{0}$ is totally ordered by set inclusion, then $A_{\Phi}:=\cup \Phi \in \mathcal{F}_{0}$. We must show $A_{\Phi} \in \Phi_{C}$, that is that $A_{\Phi} \subset C$ or $C \subset A_{\Phi}$. Now if $A \subset C$ for all $A \in \Phi$, then $A_{\Phi} \subset C$ and hence $A_{\Phi} \in \Phi_{C}$. On the other hand if there is some $A \in \Phi$ such that $g(C) \subset A$ then clearly $g(C) \subset A_{\Phi}$ and again $A_{\Phi} \in \Phi_{C}$.
2. Given $A \in \Phi_{C}$ we must show $g(A) \in \Phi_{C}$, i.e. that

$$
\begin{equation*}
g(A) \subset C \text { or } g(C) \subset g(A) \tag{B.2}
\end{equation*}
$$

There are three cases to consider: either $A \varsubsetneqq C, A=C$, or $g(C) \subset A$. In the case $A=C, g(C)=g(A) \subset g(A)$ and if $g(C) \subset A$ then $g(C) \subset A \subset g(A)$ and Eq. (B.2) holds in either of these cases. So assume that $A \varsubsetneqq C$. Since $C \in \Gamma$, either $g(A) \subset C$ (in which case we are done) or $C \subset g(A)$. Hence we may assume that

$$
A \varsubsetneqq C \subset g(A)
$$

Now if $C$ were a proper subset of $g(A)$ it would then follow that $g(A) \backslash A$ would consist of at least two points which contradicts the definition of $g$. Hence we must have $g(A)=C \subset C$ and again Eq. (B.2) holds, so $\Phi_{C}$ is a tower.

It is now easy to show $\Gamma$ is a tower. It is again clear that $P_{0} \in \Gamma$ and Property 2. may be checked for $\Gamma$ in the same way as it was done for $\Phi_{C}$ above. For Property 3., if $C \in \Gamma$ we may use $\Phi_{C}=\mathcal{F}_{0}$ to conclude for all $A \in \mathcal{F}_{0}$, either $A \subset C \subset g(C)$ or $g(C) \subset A$, i.e. $g(C) \in \Gamma$. Thus $\Gamma$ is a tower and we are done.

## References

1. Isaac Chavel, Riemannian geometry - a modern introduction, Cambridge University Press, Cambridge, 1993. MR 95j:53001
2. $\qquad$ , The Laplacian on Riemannian manifolds, Spectral theory and geometry (Edinburgh, 1998), Cambridge Univ. Press, Cambridge, 1999, pp. 30-75. MR 2001c:58029
3. E. B. Davies, Heat kernels and spectral theory, Cambridge Tracts in Mathematics, vol. 92, Cambridge University Press, Cambridge, 1990. MR 92a:35035
4. $\qquad$ , Heat kernel bounds for higher order elliptic operators, Journées "Équations aux Dérivées Partielles" (Saint-Jean-de-Monts, 1995), École Polytech., Palaiseau, 1995, pp. Exp. No. III, 11. MR 96i:35020
5. Enrico Giusti, Minimal surfaces and functions of bounded variation, Birkhäuser Verlag, Basel, 1984. MR 87a:58041
6. Ronald B. Guenther and John W. Lee, Partial differential equations of mathematical physics and integral equations, Dover Publications Inc., Mineola, NY, 1996, Corrected reprint of the 1988 original. MR 97e:35001
7. Serge Lang, Real and functional analysis, third ed., Graduate Texts in Mathematics, vol. 142, Springer-Verlag, New York, 1993. MR 94b:00005
8. Vladimir G. Maz'ja, Sobolev spaces, Springer Series in Soviet Mathematics, Springer-Verlag, Berlin, 1985, Translated from the Russian by T. O. Shaposhnikova. MR 87g:46056
9. A. Pazy, Semigroups of linear operators and applications to partial differential equations, Applied Mathematical Sciences, vol. 44, Springer-Verlag, New York, 1983. MR 85g:47061
10. Michael Reed and Barry Simon, Methods of modern mathematical physics. II. Fourier analysis, self-adjointness, Academic Press [Harcourt Brace Jovanovich Publishers], New York, 1975. MR 58 \#12429b
11. $\qquad$ , Methods of modern mathematical physics. I, second ed., Academic Press Inc. [Harcourt Brace Jovanovich Publishers], New York, 1980, Functional analysis. MR 85e:46002
12. Laurent Saloff-Coste, Aspects of Sobolev-type inequalities, London Mathematical Society Lecture Note Series, vol. 289, Cambridge University Press, Cambridge, 2002. MR 2003c:46048
13. Robert Schatten, Norm ideals of completely continuous operators, Second printing. Ergebnisse der Mathematik und ihrer Grenzgebiete, Band 27, SpringerVerlag, Berlin, 1970. MR 41 \#2449
14. B. Simon, Convergence in trace ideals, Proc. Amer. Math. Soc. 83 (1981), no. 1, 39-43. MR 82h:47042
15. Barry Simon, Trace ideals and their applications, London Mathematical Society Lecture Note Series, vol. 35, Cambridge University Press, Cambridge, 1979. MR 80k:47048
16. S. L. Sobolev and E. R. Dawson, Partial differential equations of mathematical physics, Translated from the third Russian edition by E. R. Dawson; English translation edited by T. A. A. Broadbent, Pergamon Press, Oxford, 1964. MR 31 \#2478
17. Elias M. Stein, Singular integrals and differentiability properties of functions, Princeton Mathematical Series, No. 30, Princeton University Press, Princeton, N.J., 1970. MR 44 \#7280
18. Frank W. Warner, Foundations of differentiable manifolds and Lie groups, Graduate Texts in Mathematics, vol. 94, Springer-Verlag, New York, 1983, Corrected reprint of the 1971 edition. MR 84k:58001

[^0]:    ${ }^{1}$ The argument shows that $v_{j}^{(n)}(0) \geq 0$ for all $n$. This is also easily seen directly by induction using Eq. (39.9) with $f$ replaced by $g$ and the fact that $\partial^{\alpha} g_{j}(0) \geq 0$ for all $\alpha$.

[^1]:    ${ }^{1}$ For example choose $u, v$ so that $L u=0=L v$ and $u(\alpha)=v^{\prime}(\alpha)=0$ and $u^{\prime}(\alpha)=$

[^2]:    ${ }^{2}$ Actually we can simply define $G^{*}$ to be a Green's function for $L^{*}$. It is not

[^3]:    $\overline{{ }^{4} \text { Basically, }}$ if $L^{k} u=g \in L^{2}(J)$ then $u=G^{k} g \in H_{0}^{2 k}(J)$.

[^4]:    ${ }^{1}$ We will do this again later using the Fourier transform.

[^5]:    ${ }^{1}$ Similarly since both sides of Eq. (45.10) are analytic functions in $t$, it follows that

[^6]:    $\left.\overline{{ }^{1} \text { A priori, }} L u \in H^{-1}(\Omega) \subset \mathcal{D}^{\prime}(\Omega)\right)$.

[^7]:    ${ }^{1}$ If $\lambda=\lambda_{n}$ has multiplicity larger than one, then $n<\operatorname{dim} E_{\lambda_{n}}$ otherwise $n=$

[^8]:    ${ }^{1}$ This may also be interpreted as a Bochner integral, since $T(t)$ is continuous and thus has separable range in $B(X)$.

[^9]:    ${ }^{1}$ If $X$ is a countable set we may prove Zorn's Lemma by induction. Let $\left\{x_{n}\right\}_{n=1}^{\infty}$ be an enumeration of $X$, and define $E_{n} \subset X$ inductively as follows. For $n=1$ let $E_{1}=\left\{x_{1}\right\}$, and if $E_{n}$ have been chosen, let $E_{n+1}=E_{n} \cup\left\{x_{n+1}\right\}$ if $x_{n+1}$ is an upper bound for $E_{n}$ otherwise let $E_{n+1}=E_{n}$. The set $E=\cup_{n=1}^{\infty} E_{n}$ is a linearly ordered (you check) subset of $X$ and hence by assumption $E$ has an upper bound, $x \in X$. I claim that his element is maximal, for if there exists $y=x_{m} \in X$ such that $y \geq x$, then $x_{m}$ would be an upper bound for $E_{m-1}$ and therefore $y=x_{m} \in E_{m} \subset E$. That is to say if $y \geq x$, then $y \in E$ and hence $y \leq x$, so $y=x$. (Hence we may view Zorn's lemma as a " jazzed" up version of induction.) ${ }^{2}$ Similalry one may show that $3 \Rightarrow 2$. Let $\mathcal{F}=\{E \subset X: E$ is linearly ordered $\}$ and order $\mathcal{F}$ by inclusion. If $\mathcal{M} \subset \mathcal{F}$ is linearly ordered, let $E=\cup \mathcal{M}=\bigcup_{A \in \mathcal{M}} A$.
    If $x, y \in E$ then $x \in A$ and $y \in B$ for some $A, B \subset \mathcal{M}$. Now $\mathcal{M}$ is linearly ordered by set inclusion so $A \subset B$ or $B \subset A$ i.e. $x, y \in A$ or $x, y \in B$. Sinse $A$ and $B$ are linearly order we must have either $x \leq y$ or $y \leq x$, that is to say $E$ is linearly ordered. Hence by 3 . there exists a maximal element $E \in \mathcal{F}$ which is the assertion in 2.

[^10]:    ${ }^{3}$ Here is an easy proof if the elements of $\mathcal{F}$ happened to all be finite sets and there existed a set $P \in \mathcal{F}$ with a maximal number of elements. In this case the condition that $P \subset g(P)$ would imply that $P=g(P)$, otherwise $g(P)$ would have more elements than $P$.

