# UNIVERSITY OF CALIFORNIA, SAN DIEGO Hypoelliptic heat kernel inequalities on Lie groups 

A dissertation submitted in partial satisfaction of the requirements for the degree

Doctor of Philosophy
in

Mathematics
by

Tai Alexis Melcher

Committee in charge:
Professor Bruce Driver, Chair
Professor Peter Ebenfelt
Professor Patrick Fitzsimmons
Professor George Fuller
Professor Costas Pozrikidis

## Copyright

Tai Alexis Melcher, 2004 All rights reserved.

The dissertation of Tai Alexis Melcher is approved, and it is acceptable in quality and form for publication on microfilm:
$\qquad$
$\qquad$
$\qquad$

Chair

University of California, San Diego

2004

To my parents,
Michael and Sue May.

## TABLE OF CONTENTS

Signature Page ..... iii
Dedication ..... iv
Table of Contents ..... V
Acknowledgements ..... vii
Vita and Publications ..... viii
Abstract of the Dissertation ..... ix
1 Introduction ..... 1
1.1 Background ..... 1
1.2 Statement of results ..... 6
2 A special case: the Heisenberg group ..... 12
2.1 Realization of the Heisenberg Lie group ..... 12
2.2 Differential operators on $G$ ..... 13
2.3 Scaling on $G$ ..... 16
2.4 The constant $K_{p}>1$ ..... 18
2.5 Review of calculus on Wiener space ..... 21
2.6 Heisenberg rolling map ..... 24
2.7 Lifted vector fields and their $L^{2}$-adjoints ..... 29
2.8 An $L^{p}$-type gradient estimate ( $p>1$ ) and a Poincaré inequality ..... 31
2.9 The $L^{1}$-type gradient estimate ..... 35
2.9.1 Method fails for the $p=1$ case ..... 35
2.9.2 Function classes for which the $p=1$ inequality holds ..... 38
2.9.3 A Green's function inequality ..... 43
3 Wiener calculus over $G$ ..... 46
3.1 The main theorems ..... 47
3.1.1 The rolling map on $G$ ..... 47
3.1.2 Covariance matrix of $\xi_{t}$ ..... 51
3.1.3 Lifted vector fields and their $L^{2}$-adjoints ..... 52
3.2 Results needed in the proof of Theorem 3.2 ..... 55
3.2.1 Matrix group computations ..... 56
3.2.2 Lie group computations ..... 70
4 Lie group inequalities ..... 87
4.1 An $L^{p}$-type gradient estimate $(p>1)$ ..... 88
4.2 A Poincaré inequality ..... 91
4.2.1 Stratified nilpotent Lie groups ..... 92
4.2.2 Nilpotent Lie groups ..... 95
5 Appendix ..... 98
5.1 Proof of Theorem 1.2||. ..... 98
5.2 Proof of Theorem 3.5 || ..... 100
5.3 Spectral decomposition of the Heisenberg heat kernel ..... 105
5.4 Covariance matrix of $\xi_{t}$ on the Heisenberg group ..... 108
5.5 Heisenberg lifted vector fields ..... 112
Bibliography ..... 116

## ACKNOWLEDGEMENTS

I thank my advisor Bruce Driver for patience, kindness, and the incredible gift of his knowledge. He has taught me how to be a mathematician. I thank Mark Lynch for suggesting that I become one in the first place. And I thank the members of my thesis committee, Peter Ebenfelt, Patrick Fitzsimmons, George Fuller, and Costas Pozrikidis.

The love and support of my parents, Michael and Sue May; of my sister and brothers, Li, Glenn, Robert, Danny, and Philip; and of their families, have carried me this far and will surely take me over many other obstacles. I hope I have made them proud. I thank Christian Gromoll for his strength and commitment and love, even when so far away.

I give my thanks to Imre Tuba, Ben Raphael, Alison Raphael, Arif Dowla, Rino Sanchez, Graham Hazel, Tim McMurry, Josh Griffin, Jason Bell, Brian Kurkoski, Joseph Soriaga, Jon Callihan, and Karen Warner - the friends who have been there for my graduate school years. And to Katy Stowe and Will and Angela Duke - the friends who have been there for everything.

To the people I have named here and to the people I have not - the usual suspects for Sunday night dinner, the weekend ultimate players, Doug Trenton and the staff at RT's Longboard, and many others - I could not have done this without you.

November 3, 1977 Born, Portsmouth, Virginia.
B. S., Tufts University.

2001
M. A., University of California, San Diego.

2004
Ph. D., University of California, San Diego.

## PUBLICATIONS

Hypoelliptic heat kernel inequalities on the Heisenberg group (with B. K. Driver), 2003. Submitted.

## ABSTRACT OF THE DISSERTATION

# Hypoelliptic heat kernel inequalities on Lie groups 

by

Tai Alexis Melcher<br>Doctor of Philosophy in Mathematics<br>University of California San Diego, 2004<br>Professor Bruce Driver, Chair

We study the existence of gradient estimates for second order hypoelliptic heat kernels on manifolds. It is now standard that such inequalities, in the elliptic case, are equivalent to a lower bound on the Ricci tensor of the Riemannian metric. For hypoelliptic operators, the associated "Ricci curvature" takes on the value $-\infty$ at points of degeneration of the semi-Riemannian metric associated to the operator. For this reason, many of the standard proofs for the elliptic theory no longer work in the hypoelliptic setting.

This thesis gives recent results for hypoelliptic operators. Using Malliavin calculus methods, we transfer the problem to one of determining certain infinite dimensional estimates. We study the case where the underlying manifold is a Lie group and the hypoelliptic operators are invariant under left translations. In particular, we are able to show that " $L^{p}$-type" gradient estimates hold for $p \in(1, \infty)$, and the $p=2$ gradient estimate implies a Poincaré estimate in this context. The case $p=1$ (which would imply a logarithmic Sobolev inequality) is still under investigation; however, in the special case of the Heisenberg Lie group, we are able to determine several large classes of functions for which the inequality holds.

## Chapter 1

## Introduction

### 1.1 Background

Let $M^{d}$ be a manifold of dimension $d$. A second order differential operator $L$ is subelliptic if, given a chart $x$ on $M$, one may write

$$
L=a^{i j}(x) \partial_{x^{i}} \partial_{x^{j}}+b^{i}(x) \partial_{x^{i}},
$$

where $\partial_{x^{i}}=\frac{\partial}{\partial x^{i}}$ and $\left(a^{i j}\right)$ is a real symmetric non-negative matrix. We observe here the summation convention of summing over repeated upper and lower indices. In the nondegenerate case, that is, $\left(a^{i j}\right)>0, L$ is said to be elliptic.

On a Riemannian manifold $M$ with metric $g$, a standard example of a subelliptic operator is the Laplace-Beltrami operator $\Delta_{g}$. Given a chart $x$ on $M, \Delta_{g}$ may be written as

$$
L=\Delta_{g}=\frac{1}{\sqrt{g}} \partial_{x^{i}}\left(\sqrt{g} g^{i j} \partial_{x^{j}}\right),
$$

where $g=g_{i j} d x^{i} d x^{j}=g\left(\partial_{x^{i}}, \partial_{x^{j}}\right) d x^{i} d x^{j}, \sqrt{g}=\sqrt{\operatorname{det}\left(g_{i j}\right)}$, and $\left(g^{i j}\right)$ is the inverse matrix to the matrix $\left(g_{i j}\right)$. It is well known that $\Delta_{g}$ is in fact an elliptic operator. On a general manifold $M$, given a set of smooth vector fields $\left\{X_{1}, \ldots, X_{k}\right\}$ on $M$, the operator

$$
L=\sum_{i=1}^{k} X_{i}^{2}
$$

is also subelliptic, which may be seen by writing the vector fields in local coordinates, $X_{i}=X_{i}^{j} \partial_{x^{j}}$.

Notation 1.1. Let $C_{c}^{\infty}(M)$ denote the set of smooth functions on $M$ with compact support. When $M=\mathbb{R}^{n}$, let $C_{p}^{\infty}\left(\mathbb{R}^{n}\right)$ denote those functions $f \in C^{\infty}\left(\mathbb{R}^{n}\right)$ such that $f$ and all of its partial derivatives have at most polynomial growth.

For any subelliptic operator $L$, we may define

$$
\Gamma_{1}(f, h):=\frac{1}{2}[L(f h)-f L h-h L f],
$$

for $f, h \in C^{\infty}(M)$. By iterating this construction, we also have

$$
\Gamma_{2}(f, h):=\frac{1}{2}\left[L \Gamma_{1}(f, h)-\Gamma_{1}(L f, h)-\Gamma_{1}(f, L h)\right] .
$$

If $M$ is a Riemannian manifold and $\nabla:=\left(X_{1}, \ldots, X_{k}\right)$, for $\left\{X_{i}\right\}_{i=1}^{k}$ a collection of smooth vector fields on $M$, then

$$
\Gamma_{1}(f, h)=(\nabla f, \nabla h) .
$$

In particular, along the diagonal $f=h$, we write

$$
\begin{equation*}
\Gamma_{1}(f):=\Gamma_{1}(f, f)=\frac{1}{2} L f^{2}-f L f=|\nabla f|^{2} . \tag{1.1}
\end{equation*}
$$

Similarly,

$$
\Gamma_{2}(f):=\Gamma_{2}(f, f)=\frac{1}{2} L \Gamma_{1}(f)-\Gamma_{1}(f, L f)=\frac{1}{2} L|\nabla f|^{2}-(\nabla f, \nabla L f) .
$$

When $L$ is the Laplace-Beltrami operator $\Delta_{g}$,

$$
\begin{equation*}
\Gamma_{2}(f)=\left\|\nabla^{2} f\right\|^{2}+(\operatorname{Ric} \nabla f, \nabla f), \tag{1.2}
\end{equation*}
$$

where $\nabla^{2}$ is the Hessian tensor, Ric is the Ricci tensor, and $\|\cdot\|$ is the Hilbert-Schmidt norm in the Riemannian metric.

Over the last twenty years or more, a fairly complete and very beautiful theory has been developed applying to elliptic operators on Riemannian manifolds. This theory
relates properties of the solutions of elliptic and parabolic equations to properties of the Riemannian geometry. These geometric properties are determined by the principal symbol of the underlying elliptic operator. The following theorem is a typical example of the type of result we have in mind here; see for example [4, [5, 6]. In particular, (35] gives a nice survey of the progress made on the geometry of diffusion generators. We include the proof of this theorem in the appendix to the thesis.

Theorem 1.2. Suppose $(M, g)$ is a complete Riemannian manifold, and $\nabla$ and $\Delta$ are the gradient and Laplace-Beltrami operators acting on $C^{\infty}(M)$. We write $e^{t \Delta / 2}$ to denote the heat semigroup $e^{t \bar{\Delta} / 2}$, where $\bar{\Delta}$ is the self-adjoint extension of $\left.\Delta\right|_{C_{c}^{\infty}(M)}$ to $L^{2}(M, d V)$ with $d V$ volume measure on $M$. Let $|v|:=\sqrt{g(v, v)}$ for all $v \in T M$, Ric denote the Ricci curvature tensor, and $k$ be a constant. Then the following are equivalent:

1. $\operatorname{Ric}(\nabla f, \nabla f) \geq-2 k|\nabla f|^{2}$, for all $f \in C_{c}^{\infty}(M)$,
2. $\Gamma_{2}(f) \geq-2 k \Gamma_{1}(f)$, for all $f \in C_{c}^{\infty}(M)$,
3. $\left|\nabla e^{t \Delta / 2} f\right| \leq e^{k t} e^{t \Delta / 2}|\nabla f|$, for all $f \in C_{c}^{\infty}(M)$ and $t>0$,
4. $\left|\nabla e^{t \Delta / 2} f\right|^{2} \leq e^{2 k t} e^{t \Delta / 2}|\nabla f|^{2}$, for all $f \in C_{c}^{\infty}(M)$ and $t>0$,
5. there is a function $K(t)>0$ such that $K(0)=1, \dot{K}(0)$ exists, and

$$
\begin{equation*}
\left|\nabla e^{t \Delta / 2} f\right|^{2} \leq K(t) e^{t \Delta / 2}|\nabla f|^{2}, \tag{1.3}
\end{equation*}
$$

$$
\text { for all } f \in C_{c}^{\infty}(M) \text { and } t>0 .
$$

Remark 1.3. This theorem leaves open the possibility that a function $K(t)>0$ exists such that the inequality (1.3) is satisfied for all $f \in C_{c}^{\infty}(M)$ and $t>0$, but $K(0)>1$, or $\dot{K}(0)=\infty$ or does not exist, while at the same time there is no constant $k$ such that $\operatorname{Ric}(\nabla f, \nabla f) \geq-2 k|\nabla f|^{2}$, for all $f \in C_{c}^{\infty}(M)$.

Note that the estimate in (3) implies that

$$
\begin{equation*}
\left|\nabla e^{t \Delta / 2} f\right|^{p} \leq e^{p k t} e^{t \Delta / 2}|\nabla f|^{p}, \quad p \in[1, \infty), \tag{1.4}
\end{equation*}
$$

for all $f \in C_{c}^{\infty}(M)$ and $t>0$, by Hölder's inequality. Estimates like (11)-(5) are also equivalent to one parameter families of Poincaré and log Sobolev estimates for the heat kernel. The latter has implications for hypercontractivity of an associated semigroup; see Gross [23, 24].

As a simple illustration of this theorem, consider the manifold $M=\mathbb{R}^{d}$ with the usual vector fields $\partial_{x^{1}}, \ldots, \partial_{x^{d}}$. Let $\nabla$ and $\Delta$ be the standard gradient and Laplacian,

$$
\nabla=\left(\partial_{x^{1}}, \ldots, \partial_{x^{d}}\right) \text { and } \Delta=\partial_{x^{1}}^{2}+\ldots+\partial_{x^{d}}^{2} .
$$

On $\mathbb{R}^{d}$, the Ricci curvature is 0 , and the standard Laplacian is the Laplace-Beltrami operator. Also, $\Gamma_{2}(f)=\left\|\nabla^{2} f\right\|^{2} \geq 0$. In this case $e^{t \Delta / 2}$ is convolution by the well known probability density

$$
p_{t}(x):=\frac{1}{(2 \pi t)^{d / 2}} e^{-\frac{1}{2 t}|x|_{\mathbb{R}^{d}}^{2}},
$$

and

$$
\begin{equation*}
\nabla e^{t \Delta / 2} f=e^{t \Delta / 2} \nabla f \tag{1.5}
\end{equation*}
$$

for all $f \in C_{c}^{1}\left(\mathbb{R}^{d}\right)$, as follows from basic properties of convolutions; more abstractly, this follows from the commutativity of the Euclidean gradient and Laplacian. Equation 1.5) and an application of Hölder's inequality then imply that

$$
\left|\nabla e^{t \Delta / 2} f\right|^{p}=\left[e^{t \Delta / 2}|\nabla f|\right]^{p} \leq e^{t \Delta / 2}|\nabla f|^{p},
$$

for all $f \in C_{c}^{1}\left(\mathbb{R}^{d}\right)$, where $|\nabla f|:=\left\{\sum_{i=1}^{d}\left(\partial_{x^{i}} f\right)^{2}\right\}^{1 / 2}$.
This research program moves toward extending Theorem 1.2 to operators of the form

$$
\begin{equation*}
L=\sum_{i=1}^{k} X_{i}^{2} \tag{1.6}
\end{equation*}
$$

where $\left\{X_{i}\right\}_{i=1}^{k}$ is a set of smooth vector fields on $M$ satisfying

$$
\begin{equation*}
T_{m} M=\operatorname{span}(\{X(m): X \in \mathcal{L}\}), \quad \forall m \in M \tag{HC}
\end{equation*}
$$

where $\mathcal{L}$ is the Lie algebra of vector fields generated by the collection $\left\{X_{i}\right\}_{i=1}^{k}$. This assumption is the Hörmander bracket condition, and we say that $\left\{X_{i}\right\}_{i=1}^{k}$ is a Hörmander set. Under this assumption, by a celebrated theorem of Hörmander, the operator $L$ given in Equation (1.6) is hypoelliptic. Recall that a subelliptic operator $L$ is said to be hypoelliptic if $L u \in C^{\infty}(\Omega)$ implies that $u \in C^{\infty}(\Omega)$, for all distributions $u \in C^{\infty}(\Omega)^{\prime}$ on any open set $\Omega \subset_{o} M$.

However, $L$ need not be elliptic. The principle symbol of $L$ at $\xi \in T_{m}^{*} M$ is given by $\sigma_{L}(\xi)=\sum_{i=1}^{k}\left[\xi\left(X_{i}\right)\right]^{2}$. By definition, the operator $L$ is degenerate at points $m \in M$ where there exists $0 \neq \xi \in T_{m}^{*} M$ such that $\sigma_{L}(\xi)=0$. At points of degeneracy of $L$, the Ricci tensor is not well defined and should be interpreted to take the value $-\infty$ in some directions. Thus, there exist no lower bounds in the Ricci curvature in this case. $\Gamma_{1}$ and $\Gamma_{2}$ are well defined for any subelliptic operator $L$; however, it can be shown that inequalities analogous to that in item (2) of Theorem 1.2 also fail in the non-elliptic case. Hence it is not possible to directly generalize Theorem 1.2 in this setting. Nevertheless it is reasonable to ask if inequalities of the form (1.4) might still hold. More precisely, we let $\nabla=\left(X_{1}, \ldots, X_{k}\right)$ and address the following question: do functions $K_{p}(t)<\infty$ exist such that

$$
\left|\nabla e^{t L / 2} f\right|^{p} \leq K_{p}(t) e^{t L / 2}|\nabla f|^{p}, \quad p \in[1, \infty)
$$

for all $f \in C_{c}^{\infty}(M)$ and $t>0$ ?
Related results appear in Kusuoka and Stroock [34]. In particular, a special case of Theorem 2.18 of that paper implies that, for all $p \in(1, \infty)$, there exist finite constants $C_{p}$ such that

$$
\left|\nabla e^{t L / 2} f\right|^{p} \leq C_{p} t^{-p / 2} e^{t L / 2}|f|^{p}
$$

for all smooth, bounded functions $f$ with bounded derivatives of all orders and all $t>0$. A similar result may be found in Picard [42].

Additionally, Auscher, Coulhon, Duong, and Hofmann [3] show that on a noncompact connected Riemannian manifold satisfying certain spectral gap conditions, the

Laplace-Beltrami operator satisfies

$$
\left|\nabla e^{t \Delta} f\right|^{2} \leq C e^{-C^{\prime} t \Delta}|\nabla f|^{2}
$$

for some $C, C^{\prime}>0$, for all $f \in C_{c}^{\infty}(M)$ such that $f, \nabla f \in L^{2}(M, d V)$, and $t>0$. Their methods may be adapted for the sum of squares of vectors fields satisfying the Hörmander bracket condition (HC). In the case when the manifold is a Lie group, this is carried out in Alexopoulos [2]. The papers [3, 11] also include some potential nice applications of our own result.

### 1.2 Statement of results

Let $G$ be a $d$-dimensional Lie group with Lie algebra $\mathfrak{g}=\operatorname{Lie}(G)$ and identity element $e$. Suppose $\left\{X_{i}\right\}_{i=1}^{k} \subset \mathfrak{g}$ is a Lie generating set; that is, there exists some $m \in \mathbb{N}$ such that

$$
\begin{align*}
& \operatorname{span}\left\{X_{i},\left[X_{i_{1}}, X_{i_{2}}\right],\left[X_{i_{1}},\left[X_{i_{2}}, X_{i_{3}}\right]\right], \ldots,\left[X_{i_{1}},\left[\cdots,\left[X_{i_{m-1}}, X_{i_{m}}\right] \cdots\right]\right]:\right. \\
& \left.i, i_{r} \in\{1, \ldots, k\}, r \in\{1, \ldots, m\}\right\}=\mathfrak{g} . \tag{1.7}
\end{align*}
$$

Notation 1.4. Let $\Sigma=\Sigma_{0}:=\left\{X_{1}, \ldots, X_{k}\right\}$ and $\Sigma_{r}$ be defined inductively by

$$
\Sigma_{r}:=\left\{\left[X_{i}, V\right]: V \in \Sigma_{r-1}, i=1, \ldots, k\right\},
$$

for all $r \in \mathbb{N}$. Since $\left\{X_{i}\right\}_{i=1}^{k}$ is a Lie generating set, there exists a finite $m$ such that

$$
\operatorname{span}\left(\cup_{r=0}^{m} \Sigma_{r}\right)=\mathfrak{g} .
$$

Let $\mathfrak{g}_{0}:=\operatorname{span}\left(\Sigma_{0}\right)$, and let $\left\{Y_{j}\right\}_{j=1}^{d-k} \subset \cup_{r=1}^{m} \Sigma_{r}$ be a basis of $\mathfrak{g} / \mathfrak{g}_{0}$. Define an inner product $\langle\cdot, \cdot\rangle$ on $\mathfrak{g}$ by making $\left\{X_{i}\right\}_{i=1}^{k} \cup\left\{Y_{j}\right\}_{j=1}^{d-k}$ an orthonormal set. Note then that $\left\{X_{i}\right\}_{i=1}^{k}$ is an orthonormal basis of $\mathfrak{g}_{0}$. We may extend $\langle\cdot, \cdot\rangle$ to a right invariant metric on $G$ by defining $\langle\cdot, \cdot\rangle_{g}: T_{g} G \times T_{g} G \rightarrow \mathbb{R}$ as

$$
\langle v, w\rangle_{g}:=\left\langle R_{g^{-1} *} v, R_{g^{-1} *} w\right\rangle, \quad \text { for all } v, w \in T_{g} G .
$$

The $g$ subscript will be suppressed when there is no chance of confusion.

Notation 1.5. Let $L_{g}$ denote left translation by an element $g \in G$, and let $R_{g}$ denote right translation. Given an element $X \in \mathfrak{g}$, let $\tilde{X}$ denote the left invariant vector field on $G$ such that $\tilde{X}(e)=X$, where $e$ is the identity of $G$. Recall that $\tilde{X}$ being left invariant means that the vector field commutes with left translation in the following way:

$$
\tilde{X}\left(f \circ L_{g}\right)=(\tilde{X} f) \circ L_{g},
$$

for all $f \in C^{1}(G)$. Similarly, let $\hat{X}$ denote the right invariant vector field associated to $X$.

Definition 1.6. The left invariant gradient on $G$ is the operator on $C^{1}(G)$ defined by

$$
\nabla:=\left(\tilde{X}_{1}, \ldots, \tilde{X}_{k}\right) .
$$

The subLaplacian on $G$ is the second-order operator acting on $C^{2}(G)$

$$
L:=\sum_{i=1}^{k} \tilde{X}_{i}^{2} .
$$

Remark 1.7. Since $\left\{X_{i}\right\}_{i=1}^{k}$ is a Lie generating set, $\left\{\tilde{X}_{i}\right\}_{i=1}^{k}$ satisfies the Hörmander condition (HC) and Hörmander's theorem [26] implies that $L$ is a hypoelliptic operator.

Let $L^{2}(G)$ denote the space of twice integrable functions on $G$ with respect to right invariant Haar measure. Note that for any $X \in \mathfrak{g}, \tilde{X}$ is formally skew-adjoint,

$$
\begin{aligned}
\left(\tilde{X} f_{1}, f_{2}\right)_{L^{2}(G)} & =\int_{0}^{t} \tilde{X} f_{1}(g) f_{2}(g) d g=\left.\int_{0}^{t} \frac{d}{d \epsilon}\right|_{0} f_{1}\left(g e^{\epsilon X}\right) f_{2}(g) d g \\
& =\left.\frac{d}{d \epsilon}\right|_{0} \int_{0}^{t} f_{1}\left(g e^{\epsilon X}\right) f_{2}(g) d g \\
& =\left.\frac{d}{d \epsilon}\right|_{0} \int_{0}^{t} f_{1}(g) f_{2}\left(g e^{-\epsilon X}\right) d g \\
& =\left.\int_{0}^{t} \frac{d}{d \epsilon}\right|_{0} f_{1}(g) f_{2}\left(g e^{-\epsilon X}\right) d g=-\left(f_{1}, \tilde{X} f_{2}\right)_{L^{2}(G)},
\end{aligned}
$$

for all $f_{1}, f_{2} \in C_{c}^{\infty}(G)$. Thus,

$$
\left(f_{1}, L f_{2}\right)_{L^{2}(G)}=\left(L f_{1}, f_{2}\right)_{L^{2}(G)},
$$

for all $f_{1}, f_{2} \in C_{c}^{\infty}(G)$, and $L$ is a densely defined, symmetric operator on $L^{2}(G, d g)$. So we may associate to $L$ the symmetric bilinear form $\mathcal{E}^{0}\left(f_{1}, f_{2}\right):=\left(-L f_{1}, f_{2}\right)_{L^{2}(G)}$. Note that $\mathcal{E}^{0}$ is positive, since

$$
(-L f, f)_{L^{2}(G)}=\sum_{i=1}^{k}\left(\tilde{X}_{i} f, \tilde{X}_{i} f\right)_{L^{2}(G)}=\|\nabla f\|_{L^{2}(G)}^{2} \geq 0
$$

for all $f \in C_{c}^{\infty}(G)$. Thus, $\mathcal{E}^{0}$ is closable, and its minimal closure $\mathcal{E}$ is associated to a self-adjoint operator $\bar{L}$ which is an extension of $L$, called the Friedrichs extension of $L$. Now we may define the following.

Definition 1.8. Let $P_{t}$ denote the heat semigroup $e^{t \bar{L} / 2}$, where $\bar{L}$ is the Friedrichs extension of $\left.L\right|_{C_{c}^{\infty}(G)}$ to $L^{2}(G, d g)$ with $d g$ right Haar measure on $G$. By the left invariance of $L$ and the satisfaction of the Hörmander condition, $P_{t}$ admits a left convolution kernel $p_{t}$ such that

$$
\begin{equation*}
P_{t} f(g)=f * p_{t}(g)=\int_{G} f(g h) p_{t}(h) d h \tag{1.8}
\end{equation*}
$$

for all $f \in C_{c}^{\infty}(G)$, where $g h$ is defined by the group operation of $G$ and $d h$ again denotes right Haar measure on $G$. We will call $p_{t}$ the heat kernel of $G$.
$P_{t}$ is a symmetric Markov semigroup. By Remark 1.7, $L$ is a hypoelliptic operator, and so $p_{t}$ is a smooth density on $G$. In the sequel, we will let $L$ denote its own Friedrichs extension. For the semigroup theory used here, see [13, 20, 41, [55].

Notation 1.9. Let $K_{p}(t)$ be the best function such that

$$
\begin{equation*}
\left|\nabla P_{t} f\right|^{p} \leq K_{p}(t) P_{t}|\nabla f|^{p}, \quad p \in[1, \infty) \tag{p}
\end{equation*}
$$

for all $f \in C_{c}^{\infty}(G)$ and $t>0$.
Note that in this context the inequality $\left(\overline{I_{p}}\right)$ could be equivalently written as

$$
\Gamma_{1}\left(P_{t} f\right)^{p / 2} \leq K_{p}(t) P_{t} \Gamma_{1}(f)^{p / 2}, \quad p \in[1, \infty)
$$

for all $f \in C_{c}^{\infty}(G)$ and $t>0$.

Theorem 1.10. Let $G$ be a Lie group. Then for all $p \in(1, \infty), K_{p}(t)<\infty$ for all $t>0$. In particular, if $G$ is a nilpotent Lie group, then there exists a constant $K_{p}<\infty$ such that $K_{p}(t)<K_{p}$ for all $t>0$.

This statement is verified in Chapter 2 in the context of the real 3 -dimensional Heisenberg Lie group, before being treated in full generality in Chapter 4. We approach this case first because the basic idea of the proof can be seen without the added geometric complications appearing in the more general formulation. So for now taking $G$ to be the Heisenberg group, one could attempt to determine finite constants satisfying ( $I_{p}$ ) on $G$ by proving bounds for the integral representation of $P_{t}$, Equation (1.8). In particular, there is an explicit integral formula for $p_{t}(g)$ on the Heisenberg group, given in Equation 2.5). However, our attempts at proving such bounds using this formula have been unsuccessful.

The Heisenberg heat kernel $p_{t}(g) d g$ may also be realized as the distribution in $t$ of the Cartan rolling map on $G$, the process $\xi$ satisfying Equation 2.20 , the Stratonovich differential equation

$$
\sum_{i=1}^{k} \tilde{X}_{i}\left(\xi_{t}\right) \circ d b_{t}^{i}, \text { with } \xi_{0}=e,
$$

where $b^{1}, \ldots, b^{k}$ are $k$ independent real-value Brownian motions. Thus, applying the heat semigroup $P_{t}$ to a function $f \in C_{c}^{\infty}(G)$ at the identity is equivalent to taking the expectation of $f\left(\xi_{t}\right)$; that is, $P_{t} f(e)=\mathbb{E}\left\{f\left(\xi_{t}\right)\right\}$. Using this representation of $P_{t}$, we may transform our finite dimensional problem to a problem on Wiener space. In a "lifting" procedure described in Section 2.7, we are able to construct vector fields $\mathbf{X}_{i}$ on Wiener space from the Heisenberg vector fields $\tilde{X}_{i}$ and the map $\xi$. Then we can apply Malliavin's probabilistic techniques on proving hypoellipticity to determine bounds of expressions like $\mathbb{E}\left\{\mathbf{X}_{i}\left[f\left(\xi_{t}\right)\right]\right\}$. Section 2.5 reviews some infinite dimensional calculus on Wiener space necessary for this argument.

Before proceeding with this proof, in Chapter 2 we show that the left invariance of the vector fields leaves the inequality $\left(\overline{I_{p}}\right)$ invariant under group translation, a result
which also holds in the general Lie group context. The Heisenberg group admits a family of dilations adapted to its group structure. These dilations allow a scaling argument which additionally proves that in this case the functions $K_{p}(t)=K_{p}$ for all $t>0$, for some constant $K_{p}$. We prove that $K_{p} \geq \sqrt{2}$ when $1 \leq p \leq 2$ and, in general, that $K_{p}>1$. Note that at $t=0$ the inequality is an empty statement and certainly holds for constant 1. Unlike the elliptic case where the inequality holds with a continuous coefficient (see Equation (1.4), there is now a jump discontinuity in $K_{p}(t)$ at $t=0$. This discontinuity at $t=0$ should be a feature which persists in the general hypoelliptic setting, although that is not verified in this thesis. Sections 2.6 to 2.8 give the proof of Theorem 1.10 in the Heisenberg group, following the outline described above. Finally, Section 2.9.1 shows that our method can not, without modification, be used to prove $K_{1}<\infty$. However, Section 2.9.2 gives hope, giving several large classes of functions on the Heisenberg group for which the estimate holds with $p=1$.

Using the notation for calculus on Wiener space introduced in Section 2.5, Chapter 3 begins the generalization of the argument made in the Heisenberg case to a general Lie group $G$. Section 3.1.1 introduces the rolling map on $G$, the solution $\xi$ to the Stratonovich equation (3.3) analogous to the Heisenberg equation (2.20). Section 3.2 provides the additional proofs of existence and differentiability on the path space over a group, not covered in the standard literature (which primarily addresses the case of Euclidean space or compact manifolds). In particular, a large part of this chapter is devoted to proving Theorem 3.2, which states that, for the solution $\xi_{t}$ to Equation (3.3) at fixed time $t>0, f\left(\xi_{t}\right)$ is smooth in some sense for reasonable functions $f$ and has derivatives in all $L^{p}$ with respect to Wiener measure. Section 3.1 .3 generalizes the procedure of "lifting" vector fields given in Section 2.7 and, together with Section 3.1.2, gives the results necessary for bounding expectations of the lifted vector fields.

Chapter 4 contains the proof of Theorem 1.10 in the most generality we have been able to prove it in. In Section 4.1, using the results from Chapter 3, we show that for a Lie group $G, K_{p}(t)<\infty$ for all $t>0$, although we are not able to estimate the
behavior of $K_{p}$ with respect to $t$. In a generalization of the Heisenberg scaling argument in Section 2.3. Section 4.2 addresses the special case of nilpotent and stratified groups. When $G$ is stratified, dilation arguments imply that the $K_{p}$ are independent of the $t$ parameter. When $G$ is nilpotent, we are able to prove that there is a constant $K_{p}$ such that $K_{p}(t)<K_{p}$ for all $t>0$, completing the proof of Theorem 1.10. This also allows us to prove the following Poincaré estimate for the heat kernel measure in this context.

Theorem 1.11. Suppose $G$ is a nilpotent Lie group with identity element e. Then

$$
P_{t} f^{2}(e)-\left(P_{t} f\right)^{2}(e) \leq K_{2} t P_{t}|\nabla f|^{2}(e),
$$

for all $f \in C_{c}^{\infty}(G)$ and $t>0$, where $K_{2}$ is the constant in Theorem 1.10 for $p=2$.
Note that this theorem gives a slight improvement in the elliptic case, where the estimate is known only to hold with coefficients of exponential growth.

## Chapter 2

## A special case: the Heisenberg group

Let $M=G$ be $\mathbb{R}^{3}$ equipped with the Heisenberg group multiplication given in Equation (2.2). In this setting, we will take $L=\tilde{X}^{2}+\tilde{Y}^{2}$, where $\tilde{X}$ and $\tilde{Y}$ are the vector fields

$$
\begin{equation*}
\tilde{X}:=\partial_{x}-\frac{1}{2} y \partial_{z} \text { and } \tilde{Y}:=\partial_{y}+\frac{1}{2} x \partial_{z} \tag{2.1}
\end{equation*}
$$

### 2.1 Realization of the Heisenberg Lie group

Recall that the real Heisenberg Lie algebra is $\mathfrak{g}=\operatorname{span}\{X, Y, Z\}$ where $Z=$ $[X, Y]$ and $Z$ is in the center of $\mathfrak{g}$. Thus, $\mathfrak{g}_{0}:=\operatorname{span}\{X, Y\}$ is a hypoelliptic subspace of $\mathfrak{g}$; that is, the Lie algebra generated by $\mathfrak{g}_{0}$ is $\mathfrak{g}$. The Heisenberg group $G$ is the simply connected real Lie group such that $\operatorname{Lie}(G)=\mathfrak{g}$. Letting $A=a X+b Y+c Z$ and $A^{\prime}=a^{\prime} X+b^{\prime} Y+c^{\prime} Z$, we have by the Baker-Campbell-Hausdorff formula that

$$
e^{A} e^{A^{\prime}}=e^{A+A^{\prime}+\frac{1}{2}\left[A, A^{\prime}\right]}
$$

Thus we may realize $G$ as $\mathbb{R}^{3}$ with the following group multiplication

$$
\begin{equation*}
(a, b, c)\left(a^{\prime}, b^{\prime}, c^{\prime}\right)=\left(a+a^{\prime}, b+b^{\prime}, c+c^{\prime}+\frac{1}{2}\left(a b^{\prime}-a^{\prime} b\right)\right) . \tag{2.2}
\end{equation*}
$$

### 2.2 Differential operators on $G$

Let $X=(1,0,0), Y=(0,1,0)$, and $Z=(0,0,1)$ at the identity $0 \in G$. We extend these to left invariant vector fields on $G$ in the standard way. For $g=(a, b, c) \in G$, let $L_{g}$ denote left translation by $g$, and compute as follows:

$$
\begin{aligned}
\tilde{X}(a, b, c)=L_{(a, b, c) *} X & =\left.\frac{d}{d t}\right|_{0}(a, b, c)(t, 0,0) \\
& =\left.\frac{d}{d t}\right|_{0}\left(a+t, b, c-\frac{1}{2} b t\right)=\left(1,0,-\frac{1}{2} b\right) .
\end{aligned}
$$

So if $(x, y, z)$ are the standard coordinates on $G=\mathbb{R}^{3}$, for $f \in C^{1}(G)$,

$$
(\tilde{X} f)(g)=\left.\frac{d}{d t}\right|_{0} f(g \cdot t X)=\frac{\partial f}{\partial x}(g)-\frac{1}{2} y \frac{\partial f}{\partial z}(g) .
$$

Performing similar computations for $Y$ and $Z$, we then have

$$
\begin{equation*}
\tilde{X}=\partial_{x}-\frac{1}{2} y \partial_{z}, \tilde{Y}=\partial_{y}+\frac{1}{2} x \partial_{z}, \text { and }[\tilde{X}, \tilde{Y}]=\tilde{Z}=\partial_{z}, \tag{2.3}
\end{equation*}
$$

where $\partial_{x}=\frac{\partial}{\partial x}, \partial_{y}=\frac{\partial}{\partial y}$, and $\partial_{z}=\frac{\partial}{\partial z}$; compare with Equation 2.1. Note then that $\{\tilde{X}, \tilde{Y}, \tilde{Z}\}$ forms a basis for the tangent space at every point of $G$. This combined with $[\tilde{X}, \tilde{Y}]=\tilde{Z}$ implies that $\{\tilde{X}, \tilde{Y}\}$ satisfies the Hörmander bracket condition HC). In the same manner one shows the right invariant vector fields associated to $X, Y$, and $Z$ are given by

$$
\begin{equation*}
\hat{X}=\partial_{x}+\frac{1}{2} y \partial_{z}, \hat{Y}=\partial_{y}-\frac{1}{2} x \partial_{z}, \text { and }[\hat{X}, \hat{Y}]=\hat{Z}=-\partial_{z} . \tag{2.4}
\end{equation*}
$$

Remark 2.1. The right invariant vector fields associated to $X$ and $Y$ may be expressed as the following linear combinations,

$$
\hat{X}=\tilde{X}+y \tilde{Z} \text { and } \hat{Y}=\tilde{Y}-x \tilde{Z}
$$

Definition 2.2. The left invariant gradient on $G=\mathbb{R}^{3}$ is the operator

$$
\nabla=(\tilde{X}, \tilde{Y})
$$

The subLaplacian is

$$
L=\tilde{X}^{2}+\tilde{Y}^{2}
$$

and we let $P_{t}=e^{t L / 2}$ be the heat semigroup associated to $L$, as given in Definition 1.8. Note that here

$$
L=\partial_{x}^{2}+\partial_{y}^{2}+\left(x \partial_{y}-y \partial_{x}\right) \partial_{z}+\frac{1}{4}\left(x^{2}+y^{2}\right) \partial_{z}^{2},
$$

and it is easy to see that $L$ is a positive, symmetric operator on $C_{c}^{\infty}(G)$. Finally, $p_{t}(g)=P_{t} \delta_{0}(g)=e^{t L / 2} \delta_{0}(g)$ denotes the fundamental solution associated to $L$, so that for $f \in C_{p}^{\infty}(G)$,

$$
P_{t} f(g)=p_{t} * f(g):=\int_{G} f(g h) p_{t}(h) d h
$$

where $d h$ denotes right Haar measure and $g h$ is computed relative to the Heisenberg group multiplication in Equation (2.2).

Remark 2.3. Since $\{X, Y\}$ is a Lie generating set, Hörmander's theorem [26] implies that $L$ is a hypoelliptic operator. Also Malliavin's techniques show $p_{t}$ is a smooth positive function on $\mathbb{R}^{3}$; see Section 2.6. In this simple setting, an explicit formula for $p_{t}(g)$ is

$$
\begin{equation*}
p_{t}(g)=\frac{1}{8 \pi^{2}} \int_{\mathbb{R}} \frac{w}{\sinh \left(\frac{w t}{2}\right)} \exp \left(-\frac{1}{4}|\vec{x}|^{2} w \operatorname{coth}\left(\frac{w t}{2}\right)\right) e^{i w z} d w \tag{2.5}
\end{equation*}
$$

where $g=(x, y, z) \in G$ and $\vec{x}=(x, y)$; see for example [53]. This formula for the heat kernel demonstrates the potential advantage of considering the infinite dimensional Wiener space representation of $p_{t}(g) d g$, since the Wiener space representation no longer involves the oscillatory integral that appears in (2.5).

We will need the following straightforward results.
Lemma 2.4. By the left invariance of $\nabla$ and $P_{t}$, the inequality $\left(I_{p}\right)$ holds for all $g \in G$, $f \in C_{p}^{\infty}(G)$, and $t>0$, if and only if,

$$
\begin{equation*}
\left|\nabla P_{t} f\right|^{p}(0) \leq K_{p}(t) P_{t}|\nabla f|^{p}(0), \tag{2.6}
\end{equation*}
$$

for all $f \in C_{p}^{\infty}(G)$ and $t>0$.

Proof. If (2.6) holds, then

$$
\begin{aligned}
\left|\nabla P_{t} f\right|^{p}(g) & =\left|\left(\nabla P_{t} f\right) \circ L_{g}\right|^{p}(0)=\left|\nabla\left(P_{t} f \circ L_{g}\right)\right|^{p}(0) \\
& =\left|\nabla\left(P_{t}\left(f \circ L_{g}\right)\right)\right|^{p}(0) \leq K_{p} P_{t}\left|\nabla\left(f \circ L_{g}\right)\right|^{p}(0) \\
& =K_{p} P_{t}\left|(\nabla f) \circ L_{g}\right|^{p}(0)=K_{p} P_{t}|\nabla f|^{p} \circ L_{g}(0) \\
& =K_{p} P_{t}|\nabla f|^{p}(g) .
\end{aligned}
$$

The converse is trivial.
Note that this result only depended on the left invariance of the $\nabla$ and $P_{t}$, and not the specific Heisenberg group structure. Thus, this result holds in general and will be used later in the general Lie group context.

Lemma 2.5. For $A \in \mathfrak{g}$,

$$
\tilde{A} P_{t} f(0)=P_{t} \hat{A} f(0)
$$

for all $f \in C_{p}^{\infty}(G)$ and $t>0$. More generally,

$$
\hat{A} P_{t} f=P_{t} \hat{A} f
$$

from which the previous equation follows, since $\hat{A}=\tilde{A}$ at 0 .
Proof. If we ignore domain issues, then $[\hat{A}, \tilde{B}]=0$ for all $B \in \mathfrak{g}$, so that $[\hat{A}, L]=0$, and thus we expect $\left[\hat{A}, e^{t L / 2}\right]=0$. Formally, we have

$$
\begin{aligned}
\tilde{A} P_{t} f(0)=\left.\frac{d}{d \epsilon}\right|_{0} P_{t} f\left(e^{\epsilon A}\right) & =\left.\frac{d}{d \epsilon}\right|_{0} \int_{G} f\left(e^{\epsilon A} g\right) p_{t}(g) d g \\
& =\left.\int_{G} \frac{d}{d \epsilon}\right|_{0} f\left(e^{\epsilon A} g\right) p_{t}(g) d g \\
& =\int_{G} \hat{A} f(g) p_{t}(g) d g=P_{t} \hat{A} f(0) .
\end{aligned}
$$

To differentiate under the integral, we have used the translation invariance of Haar measure (which is Lebesgue measure on $\mathbb{R}^{3}$ ) and the heat kernel bound

$$
p_{t}(g) \leq C t^{-2} e^{-\rho^{2}(g) / C t}
$$

where $\rho(g) \geq C^{\prime}\left(|x|+|y|+|z|^{1 / 2}\right)$ is the Carnot-Carathéodory distance on $G$, and $C$ and $C^{\prime}$ are some positive constants; see Theorem 5.4.3 in [51] and page 27 of 9 .

### 2.3 Scaling on $G$

So let $r>0$ and $g=(x, y, z)$, and define $\phi_{r}: G \rightarrow G$ by $\phi_{r}(x, y, z)=$ $\left(r x, r y, r^{2} z\right)$. Notice that

$$
\begin{aligned}
\phi_{r}((a, b, c) \cdot(x, y, z)) & =\phi_{r}\left(\left(a+x, b+y, c+z+\frac{1}{2}(a y-x b)\right)\right. \\
& =\phi_{r}\left(\left(r a+r x, r b+r y, r^{2} c+r^{2} z+\frac{r^{2}}{2}(a y-x b)\right)\right. \\
& =\phi_{r}(a, b, c) \phi_{r}(x, y, z) .
\end{aligned}
$$

Thus $\phi_{r}$ is a group isomorphism on $G$. (In fact, $\phi_{r}$ is a dilation on $G$ with generator $W$ given by,

$$
\begin{aligned}
W(x, y, z) & =\left.\frac{d}{d r}\right|_{r=1} \phi_{r}(x, y, z)=(x, y, 2 z)_{(x, y, z)} \\
& =x \partial_{x}+y \partial_{y}+2 z \partial_{z} \\
& =x\left(\tilde{X}+\frac{1}{2} y \partial_{z}\right)+y\left(\tilde{Y}-\frac{1}{2} x \partial_{z}\right)+2 z \partial_{z}=x \tilde{X}+y \tilde{Y}+2 z \tilde{Z}
\end{aligned}
$$

see Definition 4.8.) Using $e^{t \tilde{X}}(g)=g \cdot(t, 0,0)$ and

$$
\phi_{r *} \tilde{X} \circ \phi_{r}^{-1}(g)=\left.\frac{d}{d t}\right|_{0} \phi_{r}\left(e^{t \tilde{X}}\left(\phi_{r}^{-1}(g)\right)\right),
$$

along with similar formulas involving $\tilde{Y}$, one shows

$$
\begin{equation*}
\phi_{r *} \tilde{X} \circ \phi_{r}^{-1}=r \tilde{X} \text { and } \phi_{r *} \tilde{Y} \circ \phi_{r}^{-1}=r \tilde{Y} . \tag{2.7}
\end{equation*}
$$

The equations in (2.7) are equivalent to

$$
\tilde{X}\left(f \circ \phi_{r}\right)=r(\tilde{X} f) \circ \phi_{r} \text { and } \tilde{Y}\left(f \circ \phi_{r}\right)=r(\tilde{Y} f) \circ \phi_{r} .
$$

Therefore,

$$
\begin{align*}
& \nabla\left(f \circ \phi_{r}\right)=r(\nabla f) \circ \phi_{r} \text { and }  \tag{2.8}\\
& L\left(f \circ \phi_{r}\right)=r^{2}(L f) \circ \phi_{r} .
\end{align*}
$$

Also, from Equation (2.5), for $g=(x, y, z)$,

$$
\begin{align*}
p_{r^{2} t}(g) & =\frac{1}{8 \pi^{2}} \int_{\mathbb{R}} \frac{w}{\sinh \left(\frac{w r^{2} t}{2}\right)} \exp \left(-\frac{1}{4}|\vec{x}|^{2} w \operatorname{coth}\left(\frac{w r^{2} t}{2}\right)\right) e^{i w z} d w \\
& =\frac{1}{8 \pi^{2}} \int_{\mathbb{R}} \frac{w}{r^{2} \sinh \left(\frac{w t}{2}\right)} \exp \left(-\frac{1}{4 r^{2}}|\vec{x}|^{2} w \operatorname{coth}\left(\frac{w t}{2}\right)\right) e^{i w z / r^{2} r^{-2} d w} \\
& =r^{-4}\left(p_{t} \circ \phi_{r^{-1}}\right)(g) \tag{2.9}
\end{align*}
$$

through the change of variables $w \mapsto r^{-2} w$. Thus,

$$
\begin{aligned}
& P_{t}\left(f \circ \phi_{r}\right)(g)=\int_{G}\left(f \circ \phi_{r}\right)(g h) p_{t}(h) d h=\int_{G} f\left(\phi_{r}(g) \phi_{r}(h)\right) p_{t}(h) d h \\
& \quad=\int_{G} f\left(\phi_{r}(g) h\right)\left(p_{t} \circ \phi_{r^{-1}}\right)(h) r^{-4} d h=\int_{G} f\left(\phi_{r}(g) h\right) p_{r^{2} t}(h) d h=\left(P_{r^{2} t} f \circ \phi_{r}\right)(g) ;
\end{aligned}
$$

that is,

$$
\begin{equation*}
P_{t}\left(f \circ \phi_{r}\right)=e^{t L / 2}\left(f \circ \phi_{r}\right)=\left(e^{r^{2} t L / 2} f\right) \circ \phi_{r}=\left(P_{r^{2} t} f\right) \circ \phi_{r} . \tag{2.10}
\end{equation*}
$$

The above remarks lead to the following proposition.
Proposition 2.6. If $K_{p}$ is the best constant such that

$$
\begin{equation*}
\left|\nabla P_{1} f\right|^{p} \leq K_{p} P_{1}|\nabla f|^{p}, \tag{2.11}
\end{equation*}
$$

for all $f \in C_{p}^{\infty}(G)$, then $K_{p}(t)=K_{p}$ for all $t>0$, where $K_{p}(t)$ is the function introduced in Notation 1.9.

Proof. Fix $t>0$. Then by Equations (2.8) and 2.10,

$$
\begin{aligned}
\left|\nabla P_{t}\left(f \circ \phi_{t^{-1 / 2}}\right)\right|^{p} & =\left|\nabla\left[\left(P_{1} f\right) \circ \phi_{t^{-1 / 2}}\right]\right|^{p}=\left|t^{-1 / 2}\left(\nabla P_{1} f\right) \circ \phi_{t^{-1 / 2}}\right|^{p} \\
& \leq K_{p} t^{-p / 2}\left(P_{1}|\nabla f|^{p}\right) \circ \phi_{t^{-1 / 2}}=K_{p} t^{-p / 2} P_{t}\left(|\nabla f|^{p} \circ \phi_{t^{-1 / 2}}\right) \\
& =K_{p} P_{t}\left|\nabla\left(f \circ \phi_{t^{-1 / 2}}\right)\right|^{p},
\end{aligned}
$$

Replacing $f$ by $f \circ \phi_{t^{1 / 2}}$ in the above computations completes the proof. Moreover, reversing the above argument shows that $\left|\nabla P_{t} f\right|^{p} \leq K_{p} P_{t}|\nabla f|^{p}$ implies that $\left|\nabla P_{1} f\right|^{p} \leq$ $K_{p} P_{1}|\nabla f|^{p}$.

### 2.4 The constant $K_{p}>1$

Proposition 2.7. For $p \in[1, \infty)$, let $K_{p}$ be the best constant such that

$$
\begin{equation*}
\left|\nabla P_{t} f\right|^{p} \leq K_{p} P_{t}|\nabla f|^{p} \tag{2.12}
\end{equation*}
$$

for all $f \in C_{p}^{\infty}(G)$ and $t>0$. Then $K_{p}>1$. In particular, $K_{2} \geq 2$.
Proof. First consider the case $p=2 k$ for some positive integer $k$, and suppose that the constant $K_{2 k}=1$. Then

$$
\left|\nabla P_{t} f\right|^{2 k} \leq P_{t}|\nabla f|^{2 k}
$$

for all $t \geq 0$, and $|\nabla f|^{2 k}=\left|\nabla P_{0} f\right|^{2 k}=P_{0}|\nabla f|^{2 k}=|\nabla f|^{2 k}$, together would imply that

$$
\begin{equation*}
k|\nabla f|^{2(k-1)} \nabla f \cdot \nabla L f=\left.\frac{d}{d t}\right|_{0}\left|\nabla P_{t} f\right|^{2 k} \leq\left.\frac{d}{d t}\right|_{0} P_{t}|\nabla f|^{2 k}=\frac{1}{2} L|\nabla f|^{2 k} \tag{2.13}
\end{equation*}
$$

We now show that the function $f(x, y, z)=x+y z$ violates this inequality. Note that

$$
\begin{gathered}
L f=\nabla \cdot \nabla f=\binom{\tilde{X}}{\tilde{Y}} \cdot \nabla f=\binom{\tilde{X}}{\tilde{Y}} \cdot\binom{1-\frac{1}{2} y \cdot y}{z+\frac{1}{2} x \cdot y}=\frac{1}{2} x, \\
\nabla L f=\frac{1}{2}\binom{1}{0}, \nabla f \cdot \nabla L f=\frac{1}{2}\left(1-\frac{1}{2} y \cdot y\right), \text { and }|\nabla f|^{2}(0)=1 .
\end{gathered}
$$

Hence,

$$
\begin{equation*}
\left(k|\nabla f|^{2(k-1)} \nabla f \cdot \nabla L f\right)(0)=\frac{k}{2} . \tag{2.14}
\end{equation*}
$$

On the other hand,

$$
L \phi(g)=\phi^{\prime}(g) L g+\phi^{\prime \prime}(g)|\nabla g|^{2},
$$

and so setting $\phi(t)=t^{k}$ and $g=|\nabla f|^{2}$ gives

$$
L|\nabla f|^{2 k}=k|\nabla f|^{2(k-1)} L|\nabla f|^{2}+\left.\left.k(k-1)|\nabla f|^{2(k-2)}|\nabla| \nabla f\right|^{2}\right|^{2}
$$

From the above,

$$
\left.\left.|\nabla| \nabla f\right|^{2}\right|^{2}=\left|\binom{y z+\frac{1}{2} x y^{2}-\frac{1}{2} y(2 z+x y)}{-2 y+y^{3}+x z+\frac{1}{2} 2 x^{2} y+\frac{1}{2} x(2 z+x y)}\right|^{2},
$$

and hence

$$
\left.\left.|\nabla| \nabla f\right|^{2}\right|^{2}(0)=0,
$$

while $\left(L|\nabla f|^{2}\right)(0)=-2$. Therefore

$$
\begin{equation*}
\frac{1}{2}\left(L|\nabla f|^{2 k}\right)(0)=-k \tag{2.15}
\end{equation*}
$$

Inserting the results of Equations (2.14) and 2.15) into Equation (2.13) would imply that $\frac{k}{2} \leq-k$, which is absurd. Thus, $K_{2 k}>1$ for any positive integer $k$.

For any $p \in[1, \infty)$, there is some integer $k$ such that $p \leq 2 k$. Thus,

$$
\begin{align*}
\left|\nabla P_{t} f\right|^{2 k} & =\left(\left|\nabla P_{t} f\right|^{p}\right)^{2 k / p}  \tag{2.16}\\
& \leq K_{p}^{2 k / p}\left(P_{t}|\nabla f|^{p}\right)^{2 k / p} \leq K_{p}^{2 k / p} P_{t}|\nabla f|^{2 k}
\end{align*}
$$

Since $K_{2 k}$ is the optimal constant for which 2.16 holds and $K_{2 k}>1$,

$$
1<K_{2 k} \leq K_{p}^{2 k / p}
$$

implies that $K_{p}>1$.
We now quantify this estimate this estimate for $p=2$. Since

$$
K_{2}=\sup _{F \in C_{p}^{\infty}(G)} \frac{\left|\nabla P_{t} F\right|^{2}}{P_{t}|\nabla F|^{2}}(0) \geq \frac{\left|\nabla P_{t} f\right|^{2}}{P_{t}|\nabla f|^{2}}(0)=: C(t),
$$

where $f(x, y, z)=x+y z$, it follows that $K_{2} \geq \sup _{t>0} C(t)$. To finish the proof we compute $C(t)$ explicitly. Observe that $P_{t}$, when acting on polynomials, may be computed using

$$
P_{t}=e^{t L / 2}=\sum_{n=1}^{\infty} \frac{1}{n!}\left(\frac{t L}{2}\right)^{n}=I+\frac{t}{2} L+\frac{1}{2!} \cdot \frac{t^{2}}{4} L^{2}+\cdots .
$$

We then have

$$
P_{t} f=f+\frac{t}{2} L f=(x+y z)+\frac{t}{2} x, \quad \nabla P_{t} f=\binom{\left(1+\frac{t}{2}\right)-\frac{1}{2} y \cdot y}{z+\frac{1}{2} x \cdot y}, \text { and }
$$

$$
\begin{aligned}
\left|\nabla P_{t} f\right|^{2} & =\left(\left(1+\frac{t}{2}\right)-\frac{1}{2} y^{2}\right)^{2}+\left(z+\frac{1}{2} x y\right)^{2} \\
& =\left(1-y^{2}+\frac{1}{4} y^{4}+z^{2}+x y z+\frac{1}{4} x^{2} y^{2}\right)+\frac{t}{2}\left(2-y^{2}\right)+\frac{t^{2}}{8} \cdot 2 .
\end{aligned}
$$

Also, from before we have,

$$
\nabla f=\binom{1-\frac{1}{2} y \cdot y}{z+\frac{1}{2} x \cdot y}
$$

and so

$$
\begin{gathered}
|\nabla f|^{2}=\left(1-\frac{1}{2} y^{2}\right)^{2}+\left(z+\frac{1}{2} x y\right)^{2}=1-y^{2}+\frac{1}{4} y^{4}+z^{2}+x y z+\frac{1}{4} x^{2} y^{2}, \\
L|\nabla f|^{2}=\left(\frac{1}{2} y^{2}+\left(-2+3 y^{2}+\frac{1}{2} x^{2}\right)\right)+(x \cdot x-y \cdot y)+\frac{1}{4} r^{2} \cdot 2=-2+3 y^{2}+2 x^{2}, \\
L^{2}|\nabla f|^{2}=4+6=10 .
\end{gathered}
$$

Thus,

$$
P_{t}|\nabla f|^{2}=\left(1-y^{2}+\frac{1}{4} y^{4}+z^{2}+x y z+\frac{1}{4} x^{2} y^{2}\right)+\frac{t}{2}\left(-2+3 y^{2}+2 x^{2}\right)+\frac{t^{2}}{8} \cdot 10 .
$$

So evaluating at the identity, we have

$$
P_{t}|\nabla f|^{2}(0)=|\nabla f|^{2}(0)+\frac{t}{2} L|\nabla f|^{2}(0)+\frac{t^{2}}{8} L^{2}|\nabla f|^{2}(0)=1-t+\frac{5}{4} t^{2},
$$

and also

$$
\left|\nabla P_{t} f\right|^{2}(0)=1+t+\frac{1}{4} t^{2}
$$

We can find the maximum value of

$$
C(t)=\frac{1+t+\frac{1}{4} t^{2}}{1-t+\frac{5}{4} t^{2}}
$$

for $t \geq 0$ by taking derivatives in $t$ to show that $C(t)$ takes on maximum value 2 at $t=\frac{2}{3}$.

### 2.5 Review of calculus on Wiener space

This section contains a brief introduction to basic Wiener space definitions and notions of differentiability. For a more complete exposition, we refer the reader to $[28,29,30,32,33,38,39,40,52,53,54]$ and references contained therein. In particular, the first two chapters of Nualart [40] and Chapter V of Ikeda and Watanabe [30] are cited often here.

Let $\left(\mathscr{W}\left(\mathbb{R}^{k}\right), \mathcal{F}, \mu\right)$ denote classical k -dimensional Wiener space. That is, $\mathscr{W}=$ $\mathscr{W}\left(\mathbb{R}^{k}\right)$ is the Banach space of continuous paths $\omega:[0,1] \rightarrow \mathbb{R}^{k}$ such that $\omega_{0}=0$, equipped with the supremum norm

$$
\|\omega\|=\max _{t \in[0,1]}\left|\omega_{t}\right|,
$$

$\mu$ is standard Wiener measure, and $\mathcal{F}$ is the completion of the Borel $\sigma$-field on $\mathscr{W}$ with respect to $\mu$. By definition of $\mu$, the process

$$
b_{t}(\omega)=\left(b_{t}^{1}(\omega), \ldots, b_{t}^{k}(\omega)\right)=\omega_{t}
$$

is an $\mathbb{R}^{k}$ Brownian motion. For those $\omega \in \mathscr{W}$ which are absolutely continuous, let

$$
E(\omega):=\int_{0}^{1}\left|\dot{\omega}_{s}\right|^{2} d s
$$

denote the energy of $\omega$. The Cameron-Martin Hilbert space is the space of finite energy paths,

$$
\mathscr{H}=\mathscr{H}\left(\mathbb{R}^{k}\right):=\left\{\omega \in \mathscr{W}\left(\mathbb{R}^{k}\right): \omega \text { is absolutely continuous and } E(\omega)<\infty\right\},
$$

equipped with the inner product

$$
(h, k)_{\mathscr{H}}:=\int_{0}^{1} \dot{h}_{s} \cdot \dot{k}_{s} d s, \quad \text { for all } h, k \in \mathscr{H} .
$$

Definition 2.8. Denote by $\mathcal{S}$ the class of smooth cylinder functionals; that is, random variables $F: \mathscr{W} \rightarrow \mathbb{R}$ such that

$$
\begin{equation*}
F(\omega)=f\left(\omega_{t_{1}}, \ldots, \omega_{t_{n}}\right), \tag{2.17}
\end{equation*}
$$

for some $n \geq 1,0<t_{1}<\cdots<t_{n} \leq 1$, and function $f \in C_{p}^{\infty}\left(\left(\mathbb{R}^{k}\right)^{n}\right)$ (see Notation 1.1). For $E$ be a real separable Hilbert space, let $\mathcal{S}_{E}$ be the set of $E$-valued smooth cylinder functions $F: W \rightarrow E$ of the form

$$
\begin{equation*}
F=\sum_{j=1}^{m} F_{j} e_{j} \tag{2.18}
\end{equation*}
$$

for some $m \geq 1, e_{j} \in E$, and $F_{j} \in \mathcal{S}$.
Definition 2.9. Fix $h \in \mathscr{H}$. The directional derivative of a smooth cylinder functional $F \in \mathcal{S}$ of the form (2.17) along $h$ is given by

$$
\begin{equation*}
\partial_{h} F(\omega):=\left.\frac{d}{d \epsilon}\right|_{0} F(\omega+\epsilon h)=\sum_{i=1}^{n} \nabla^{i} f\left(\omega_{t_{1}}, \ldots, \omega_{t_{n}}\right) \cdot h_{t_{i}} \tag{2.19}
\end{equation*}
$$

where $\nabla^{i} f$ is the gradient of $f$ with respect to the $i^{\text {th }}$ variable.
The following integration by parts result is standard; see for example Theorem 8.2.2 of Hsu [27].

Proposition 2.10. Let $F, G \in \mathcal{S}$ and $h \in \mathscr{H}$. Then

$$
\left(\partial_{h} F, G\right)_{\mathscr{H}}=\left(F, \partial_{h}^{*} G\right)_{\mathscr{H}}
$$

where $\partial_{h}^{*}=-\partial_{h}+\int_{0}^{1} \dot{h}_{s} \cdot d b_{s}$.
Definition 2.11. The gradient of a smooth cylinder functional $F \in \mathcal{S}$ is the random process $D_{t} F$ taking values in $\mathscr{H}$ such that $(D F, h)_{\mathscr{H}}=\partial_{h} F$. It may be determined that, for $F$ of the form (2.17),

$$
D_{t} F=\sum_{i=1}^{n} \nabla^{i} f\left(\omega_{t_{1}}, \ldots, \omega_{t_{n}}\right)\left(t_{i} \wedge t\right)
$$

where $s \wedge t=\min \{s, t\}$. For $F \in \mathcal{S}_{E}$ of the form (2.18), we define the derivative $D_{t} F$ to be the random process taking values in $\mathscr{H} \otimes E$ given by

$$
D_{t} F:=\sum_{j=1}^{m} D_{t} F \otimes e_{j}
$$

Iterations of the derivative for smooth functionals $F \in \mathcal{S}$ are given by

$$
D_{t_{1}, \ldots, t_{k}}^{k} F=D_{t_{1}} \cdots D_{t_{k}} F \in \mathscr{H}^{\otimes k}
$$

for $k \in \mathbb{N}$. For $F \in \mathcal{S}_{E}$,

$$
D^{k} F=\sum_{j=1}^{m} D^{k} F_{j} \otimes e_{j}
$$

and these are measurable functions defined almost everywhere on $[0,1]^{k} \times \mathscr{W}$. The operator $D$ on $\mathcal{S}_{E}$ is closable, and there exist closed extensions $D^{k}$ to $L^{p}\left(\mathscr{W}, \mathscr{H}^{\otimes k} \otimes E\right)$; see, for example [40], Theorem 8.28 of [27], or Theorem 8.5 of [30]. We will denote the closure of the derivative operator also by $D$ and the domain of $D^{k}$ in $L^{p}\left([0,1]^{k} \times \mathscr{W}\right)$ by $\mathcal{D}^{k, p}$, which is the completion of the family of smooth Wiener functionals $\mathcal{S}$ with respect to the seminorm $\|\cdot\|_{k, p, E}$ on $\mathcal{S}_{E}$ given by

$$
\|F\|_{k, p, E}:=\left(\sum_{j=0}^{k} \mathbb{E}\left(\left\|D^{j} F\right\|_{\mathscr{H} \not{ }^{\otimes j \otimes E}}^{p}\right)\right)^{1 / p}
$$

for any $p \geq 1$. Let

$$
\mathcal{D}^{k, \infty}(E):=\bigcap_{p>1} \mathcal{D}^{k, p}(E) \text { and } \mathcal{D}^{\infty}(E):=\bigcap_{p>1} \bigcap_{k \geq 1} \mathcal{D}^{k, p}(E)
$$

When $E=\mathbb{R}$, we write $\mathcal{D}^{k, p}(\mathbb{R})=\mathcal{D}^{k, p}, \mathcal{D}^{k, \infty}(\mathbb{R})=\mathcal{D}^{k, \infty}$, and $\mathcal{D}^{\infty}(\mathbb{R})=\mathcal{D}^{\infty}$.
The operator $\partial_{h}$ on $\mathcal{S}$ is also closable, and there exists a closed extension of $\partial_{h^{1}} \cdots \partial_{h^{k}}$ to $L^{p}(\mu)$. We will denote the closure of $\partial_{h}$ also by $\partial_{h}$, with domain $\operatorname{Dom}\left(\partial_{h}\right)$. We will denote by $G^{k, p}$ the class of functions $F \in L^{p}(\mu)$ such that $\partial_{h^{1}} \cdots \partial_{h^{j}} F \in L^{p}(\mu)$, for all $h^{1}, \ldots, h^{j} \in \mathscr{H}, j=1, \ldots, k$. The norm on $G^{k, p}$ is given by

$$
\|F\|_{G^{k, p}}:=\sum_{j=0}^{k}\left(\mathbb{E}\left\|D^{j} F\right\|_{\left(\mathscr{C}^{\otimes j}\right)^{*}}^{p}\right)^{1 / p},
$$

where

$$
\left\|D^{j} F\right\|_{\left(\mathscr{H}^{\otimes j}\right)^{*}}=\sup \left\{\left|\partial_{h^{1}} \cdots \partial_{h^{j}} F\right|: h^{i} \in \mathscr{H},\left|h^{i}\right|_{\mathscr{H}} \leq 1\right\}
$$

is the operator norm on the space of continuous linear functionals $\left(\mathscr{H}^{\otimes j}\right)^{*}$.
The following result follows from Proposition 5.4.6 and Corollary 5.4.7 of Bogachev [10]. (Note well that here the space $\mathcal{D}^{k, p}$ corresponds to the space $W^{k, p}$ in that text.)

Theorem 2.12. For all $k \in \mathbb{N}$ and $p \in(1, \infty)$, $\mathcal{D}^{k, p} \subset G^{k, p}$. In particular, for $k=1$, $\mathcal{D}^{1, p}=G^{1, p}$, for all $p \in(1, \infty)$.

Definition 2.13. Let $D^{*}$ denote the $L^{2}(\mu)$-adjoint of the derivative operator $D$, which has domain in $L^{2}(\mathscr{W} \times[0,1], \mathscr{H})$ consisting of functions $G$ such that

$$
\left|\mathbb{E}\left[(D F, G)_{\mathscr{H}}\right]\right| \leq C\|F\|_{L^{2}(\mu)},
$$

for all $F \in \mathcal{D}^{1,2}$, where $C$ is a constant depending on $G$. For those functions $G$ in the domain of $D^{*}, D^{*} G$ is the element of $L^{2}(\mu)$ such that

$$
\mathbb{E}\left[F D^{*} G\right]=\mathbb{E}\left[(D F, G)_{\mathscr{H}}\right]
$$

It is known that $D$ is a continuous operator from $\mathcal{D}^{\infty}$ to $\mathcal{D}^{\infty}(\mathscr{H})$, and similarly, $D^{*}$ is continuous from $\mathcal{D}^{\infty}(\mathscr{H})$ to $\mathcal{D}^{\infty}$; see for example Theorem V-8.1 and its corollary in [30.

Malliavin [31, 36, 37] introduced the notion of derivatives of Wiener functionals and applied it to the regularity of probability laws induced by the solutions to stochastic differential equations at fixed times. The notion of Sobolev spaces of Wiener functionals was first introduced by Shigekawa 44 and Stroock [45, 46].

### 2.6 Heisenberg rolling map

For the rest of this chapter, let $\mathscr{W}=\mathscr{W}\left(\mathbb{R}^{2}\right)$ and $\mathscr{H}=\mathscr{H}\left(\mathbb{R}^{2}\right)$. Let $\xi:[0,1] \times$ $\mathscr{W}\left(\mathbb{R}^{2}\right) \rightarrow G$ denote the solution to the Stratonovich stochastic differential equation

$$
\begin{align*}
d \xi_{t} & =L_{\xi_{t}} X \circ d b_{t}^{1}+L_{\xi_{t}} Y \circ d b_{t}^{2} \\
& =\tilde{X}\left(\xi_{t}\right) \circ d b_{t}^{1}+\tilde{Y}\left(\xi_{t}\right) \circ d b_{t}^{2}  \tag{2.20}\\
\xi_{0} & =0 .
\end{align*}
$$

Remark 2.14. Since $\tilde{X}$ and $\tilde{Y}$ have smooth coefficients with bounded partial derivatives, Theorem 2.2.2 in Nualart [40] implies that $\xi_{t}^{i} \in \mathcal{D}^{\infty}$, for $i=1,2,3$ and all $t \in[0,1]$.

Because $G$ is a nilpotent Lie group, we may determine an "explicit" solution of (2.20),

$$
\begin{aligned}
d \xi_{t} & =\tilde{X}\left(\xi_{t}^{1}, \xi_{t}^{2}, \xi_{t}^{3}\right) \circ d b_{t}^{1}+\tilde{Y}\left(\xi_{t}^{1}, \xi_{t}^{2}, \xi_{t}^{3}\right) \circ d b_{t}^{2} \\
& =\left(\begin{array}{c}
1 \\
0 \\
-\frac{1}{2} \xi_{t}^{2}
\end{array}\right) \circ d b_{t}^{1}+\left(\begin{array}{c}
0 \\
1 \\
\frac{1}{2} \xi_{t}^{1}
\end{array}\right) \circ d b_{t}^{2} .
\end{aligned}
$$

Thus,

$$
d \xi_{t}^{1}=d b_{t}^{1}, d \xi_{t}^{2}=d b_{t}^{2}, \text { and } d \xi_{t}^{3}=-\frac{1}{2} \xi_{t}^{2} \circ d b_{t}^{1}+\frac{1}{2} \xi_{t}^{1} \circ d b_{t}^{2},
$$

and one may verify directly that

$$
\begin{equation*}
\xi_{t}=\left(b_{t}^{1}, b_{t}^{2}, \frac{1}{2} \int_{0}^{t}\left[b_{s}^{1} d b_{s}^{2}-b_{s}^{2} d b_{s}^{1}\right]\right) \tag{2.21}
\end{equation*}
$$

satisfies the required equation. Note that the third component of $\xi$ may be recognized as Lévy's stochastic area integral.

From Section 3.9 in Gīhman and Skorohod [21] and Theorem 1.22 in Bell [7, the solution $\xi$ is a time homogenous Markov process, and $P_{t}=e^{t L / 2}$ with $L=\tilde{X}^{2}+\tilde{Y}^{2}$ is the associated Markov diffusion semigroup to $\xi$; that is, $\nu_{t}:=\left(\xi_{t}\right)_{*} \mu=p_{t}(g) d g$ is the density of the transition probability of the diffusion process $\xi_{t}$, and

$$
\begin{equation*}
\left(P_{t} f\right)(0)=\mathbb{E}\left[f\left(\xi_{t}\right)\right], \tag{2.22}
\end{equation*}
$$

for any $f \in C^{\infty}(G)$, where the right hand side is expectation for $\xi_{t}$ started at 0 .
Theorem 2.15. The Malliavin covariance matrix of $\xi_{t}$,

$$
\sigma_{t}=\left(\left(D \xi_{t}^{i}, D \xi_{t}^{j}\right)_{\mathscr{H}}\right)_{1 \leq i, j \leq 3},
$$

is invertible a.s. for $t>0$, and

$$
(\operatorname{det} \sigma)^{-1} \in \cap_{p \geq 1} L^{p}(\mu)=: L^{\infty-}(\mu)
$$

This statement follows from the proof of Theorem 2.3.3 in Nualart [40] and is a special case of Theorem 3.5 in the Heisenberg setting, which will be proved in Section 3.1.2. however, we include the following self-contained proof.

Proof. From Equation (5.19) of the appendix, we have

$$
\operatorname{det} \sigma=t^{2} \int_{0}^{t}\left|b_{s}\right|^{2} d s-t\left(\int_{0}^{t} b_{s}^{1} d s\right)^{2}-t\left(\int_{0}^{t} b_{s}^{2} d s\right)^{2} .
$$

Thus $\operatorname{det} \sigma(\omega)=0$ if and only if $\omega=0$, since the Cauchy-Schwarz inequality implies that

$$
\operatorname{det} \sigma \geq t^{2} \int_{0}^{t}\left|b_{s}\right|^{2} d s-t^{2} \int_{0}^{t}\left(b_{s}^{1}\right)^{2} d s-t^{2} \int_{0}^{t}\left(b_{s}^{2}\right)^{2} d s=0
$$

with equality if and only if $b_{1}=c_{1}$ and $b_{2}=c_{2}$ for some constants $c_{1}$ and $c_{2}$. Since $b(0)=0$, then $b_{1}=0$ and $b_{2}=0$.

To simplify notation, we consider the quadratic form $Q$ on one-dimensional path space $\mathscr{W}([0, T], \mathbb{R})$ defined by

$$
\begin{equation*}
Q(\omega):=T^{2} \int_{0}^{T} \omega_{s}^{2} d s-T\left(\int_{0}^{T} \omega_{s} d s\right)^{2} \tag{2.23}
\end{equation*}
$$

The argument is similar on $\mathscr{W}\left([0, T], \mathbb{R}^{2}\right)$. We need to show that

$$
\int_{\mathscr{W}} \frac{1}{[Q(\omega)]^{p}} d \mu(\omega)<\infty
$$

for any $p \geq 1$. First note that

$$
\int_{0}^{\infty} e^{-\lambda t} t^{p} \frac{d t}{t}=\frac{1}{\lambda^{p}} \int_{0}^{\infty} e^{-t} t^{p} \frac{d t}{t}=: \frac{1}{\lambda^{p}} \Gamma(p),
$$

under the change of variables $t \mapsto \frac{t}{\lambda}$. Letting $\lambda=Q(\omega)$ then yields

$$
\frac{1}{[Q(\omega)]^{p}}=\frac{1}{\Gamma(p)} \int_{0}^{\infty} e^{-t Q(\omega)} t^{p-1} d t
$$

Thus,

$$
\int_{\mathscr{W}} \frac{1}{[Q(\omega)]^{p}} d \mu(\omega)=\frac{1}{\Gamma(p)} \int_{0}^{\infty} d t t^{p-1} \int_{W} d \mu(\omega) e^{-t Q(\omega)} .
$$

Define $F(t):=\int_{\mathscr{W}} e^{-t Q(\omega)} d \mu(\omega)$. To show that the above integral is finite for any $p \geq 1$, we must show that $F(t)$ decays faster than any polynomial as $t \rightarrow \infty$.

Consider first the finite dimensional analogue. That is, for any $a \in \mathbb{R}^{N}$, consider a quadratic form $Q$ such that

$$
Q(a):=\tilde{Q} a \cdot a
$$

for some $\tilde{Q}>0$ and $d \mu(a)=e^{-\frac{1}{2} a \cdot a} \frac{1}{(2 \pi)^{N / 2}} d a$. Since $\tilde{Q}$ is compact ( $\tilde{Q}$ is an operator on a finite dimensional vector space), there exists an orthonormal basis of eigenvectors $u_{1}, \ldots, u_{N}$ such that

$$
\tilde{Q} u_{i}=\lambda_{i} u_{i}, \quad \lambda_{i}>0 .
$$

In this basis, for any $a \in \mathbb{R}^{N}$, we may write $a=\sum_{i=1}^{N} \alpha_{i} u_{i}$, where $\alpha_{i}=a \cdot u_{i}$. Then

$$
\begin{aligned}
F(t) & =\int_{\mathbb{R}^{N}} e^{-t \tilde{Q} a \cdot a} e^{-\frac{1}{2} a \cdot a} \frac{1}{(2 \pi)^{N / 2}} d a \\
& =\int_{\mathbb{R}^{N}} e^{-t \sum_{i=1}^{N} \lambda \alpha_{i}^{2}} e^{-\frac{1}{2} \sum_{i=1}^{N} \alpha_{i}^{2}} \prod_{i=1}^{N} \frac{1}{\sqrt{2 \pi}} d \alpha_{i} \\
& =\prod_{i=1}^{N} \int_{\mathbb{R}} e^{-\frac{1}{2}\left(1+2 t \lambda_{i}\right) \alpha_{i}^{2}} \frac{1}{\sqrt{2 \pi}} d \alpha_{i} \\
& =\prod_{i=1}^{N} \sqrt{\frac{2 \pi}{1+2 t \lambda_{i}}} \cdot \frac{1}{\sqrt{2 \pi}}=\prod_{i=1}^{N} \frac{1}{\sqrt{1+2 t \lambda_{i}}} \leq C_{N} t^{-N / 2} .
\end{aligned}
$$

Returning then to the infinite dimensional case, Lemma 2.16 below implies that, for $Q: \mathscr{H} \rightarrow \mathbb{R}$ defined as in Equation (2.23),

$$
Q(h)=(\tilde{Q} h, h)_{\mathscr{H}(\mathbb{R})},
$$

where $\tilde{Q}: \mathscr{H}(\mathbb{R}) \rightarrow \mathscr{H}(\mathbb{R})$ is a positive compact operator. Thus, there exists an orthonormal basis of eigenvectors $\left\{h_{i}\right\}_{i=1}^{\infty}$ such that

$$
\tilde{Q} h_{i}=\lambda_{i} h_{i}, \quad \lambda_{i}>0 \text { and } \sum_{i=1}^{\infty} \lambda_{i}<\infty,
$$

and so, for any $h \in \mathscr{H}(\mathbb{R})$, we may write $h=\sum_{i=1}^{\infty}\left(h, h_{i}\right) \not{H} h_{i}$ and

$$
Q(h)=(\tilde{Q} h, h)_{\mathscr{H}}=\left(\sum_{i=1}^{\infty}\left(h, h_{i}\right)_{\mathscr{H}} \tilde{Q} h_{i}, \sum_{i=1}^{\infty}\left(h, h_{i}\right)_{\mathscr{H}} h_{i}\right)_{\mathscr{H}}=\sum_{i=1}^{\infty} \lambda_{i}\left(h, h_{i}\right)_{\mathscr{H}}^{2} .
$$

By Theorem 2.12 from Da Prato and Zabczyk [12], if $\left\{N_{i}\right\}_{i=1}^{\infty}$ is a sequence of independent, identically distributed normal random variables, and we define for $N \in \mathbb{N}$

$$
B_{N}:=\sum_{i=1}^{N} N_{i} h_{i}
$$

then

$$
\lim _{N \rightarrow \infty}\left\|B_{N}-\sum_{i=1}^{\infty} N_{i} h_{i}\right\|_{\infty}=0, \quad \text { a.s. }
$$

and $B:=\sum_{i=1}^{\infty} N_{i} h_{i}$ is Brownian motion. Then

$$
Q(B)=\lim _{N \rightarrow \infty} Q\left(B_{N}\right)=\lim _{N \rightarrow \infty} \sum_{i=1}^{N} \lambda_{i} N_{i}^{2}
$$

and

$$
\begin{aligned}
& \int_{W} e^{-t Q(\omega)} d \mu(\omega)=\mathbb{E}\left[e^{-t Q(B)}\right] \\
&= \mathbb{E}\left[\exp \left(-t \lim _{N \rightarrow \infty} \sum_{i=1}^{N} \lambda_{i} N_{i}^{2}\right)\right] \leq \mathbb{E}\left[\exp \left(-t \sum_{i=1}^{N} \lambda_{i} N_{i}^{2}\right)\right] \leq C_{N} t^{-N / 2}
\end{aligned}
$$

for all $N \in \mathbb{N}$.
Lemma 2.16. For $\tilde{Q}: \mathscr{H}(\mathbb{R}) \rightarrow \mathscr{H}(\mathbb{R})$ such that

$$
(\tilde{Q} h, h)_{\mathscr{H}(\mathbb{R})}=Q(h)=T^{2} \int_{0}^{T} \omega_{s}^{2} d s-T\left(\int_{0}^{T} \omega_{s} d s\right)^{2}
$$

$\tilde{Q}$ is a trace class operator, and so is compact.

Proof. Let $S$ be an orthonormal basis of $\mathscr{H}(\mathbb{R})$. Then

$$
\begin{aligned}
\operatorname{tr}(\tilde{Q})=\sum_{h \in S} Q(h) & =\sum_{h \in S}\left[T^{2} \int_{0}^{T} h^{2}(s) d s-T\left(\int_{0}^{T} h(s) d s\right)^{2}\right] \\
& =\sum_{h \in S}\left[T^{2} \int_{0}^{T} h^{2}(s) d s-T \int_{0}^{T} \int_{0}^{T} h(r) h(s) d r d s\right] \\
& =T^{2} \int_{0}^{T} s d s-T \int_{0}^{T} \int_{0}^{T} r \wedge s d r d s \\
& =T^{2} \cdot \frac{1}{2} T^{2}-\left.T \cdot\left(\frac{1}{2} T s^{2}-\frac{1}{6} s^{3}\right)\right|_{0} ^{T}=\frac{1}{6} T^{4}<\infty,
\end{aligned}
$$

where in the third equality, we have used that for any orthonormal basis $S$ of $\mathscr{H}(\mathbb{R})$,

$$
\sum_{h \in S} h(r) h(s)=r \wedge s
$$

Remark 2.17. By the general theory, Theorem 2.15 implies $\nu_{t}=\operatorname{Law}\left(\xi_{t}\right)$ is a smooth measure; see for example Theorem 2.12 and Remark 2.13 in Bell [7].

### 2.7 Lifted vector fields and their $L^{2}$-adjoints

Given $A \in \mathfrak{g}$, let $\tilde{A}=\left(\tilde{A}^{1}, \tilde{A}^{2}, \tilde{A}^{3}\right)$ be the associated left invariant vector field on G. In particular, we are interested in the vector fields $\tilde{X}=\left(1,0,-\frac{1}{2} y\right)$ and $\tilde{Y}=\left(0,1, \frac{1}{2} x\right)$. We define the "lifted vector field" $\mathbf{A}$ of $\tilde{A}$ as

$$
\begin{equation*}
\mathbf{A}=\mathbf{A}^{t}:=\sum_{i, j=1}^{3} \sigma_{i j}^{-1} \tilde{A}^{j}\left(\xi_{t}\right) D \xi_{t}^{i} \in \mathscr{H} \tag{2.24}
\end{equation*}
$$

acting on functions $F \in \mathcal{D}^{1,2}$ by

$$
\mathbf{A} F=(D F, \mathbf{A})_{\mathscr{H}} .
$$

Remark 2.18. Recall that $D$ is a continuous operator from $\mathcal{D}^{\infty}$ to $\mathcal{D}^{\infty}(\mathscr{H})$. Thus, Remark 2.14 implies that $\tilde{A}^{j}\left(\xi_{t}\right) \in \mathcal{D}^{\infty}$ and $D \xi_{t}^{i} \in \mathcal{D}^{\infty}(\mathscr{H})$, for all $t \in[0,1]$. So $\sigma_{i j} \in \mathcal{D}^{\infty}$ for $i, j=1,2,3$, and this along with Theorem 2.15 implies that $\sigma_{i j}^{-1} \in \mathcal{D}^{\infty}$. Hence, $\mathbf{A} \in \mathcal{D}^{\infty}(\mathscr{H})$.

Proposition 2.19. For all $f \in C_{p}^{\infty}(G)$,

$$
\mathbf{A}\left[f\left(\xi_{t}\right)\right]=(\tilde{A} f)\left(\xi_{t}\right)
$$

Proof. For any function $f \in C_{p}^{\infty}(G), f\left(\xi_{t}\right) \in \mathcal{D}^{\infty}$ and

$$
D\left[f\left(\xi_{t}\right)\right]=\sum_{k=1}^{3} \frac{\partial f}{\partial x_{k}}\left(\xi_{t}\right) D \xi_{t}^{k}
$$

see Proposition 1.2.3 from Nualart [40]. Then using Equation (2.24) and for the lifted vector field $\mathbf{A}$ and the Malliavin matrix $\sigma$, we have

$$
\begin{aligned}
\mathbf{A}\left[f\left(\xi_{t}\right)\right] & =\left(D f\left(\xi_{t}\right), \mathbf{A}\right)_{\mathscr{H}} \\
& =\sum_{i, j, k=1}^{3}\left(\frac{\partial f}{\partial x_{k}}\left(\xi_{t}\right) D \xi_{t}^{k}, \sigma_{i j}^{-1} \tilde{A}^{j}\left(\xi_{t}\right) D \xi_{t}^{i}\right)_{\mathscr{H}} \\
& =\sum_{i, j, k=1}^{3} \tilde{A}^{j}\left(\xi_{t}\right) \frac{\partial f}{\partial x_{k}}\left(\xi_{t}\right)\left(D \xi_{t}^{k}, D \xi_{t}^{i}\right)_{\mathscr{H}} \sigma_{i j}^{-1} \\
& =\sum_{j, k=1}^{3} \tilde{A}^{j}\left(\xi_{t}\right) \frac{\partial f}{\partial x_{k}}\left(\xi_{t}\right) \delta_{k j}=\sum_{j=1}^{3} \tilde{A}^{j}\left(\xi_{t}\right) \frac{\partial f}{\partial x_{j}}\left(\xi_{t}\right)=(\tilde{A} f)\left(\xi_{t}\right)
\end{aligned}
$$

as desired.
Definition 2.20. For a vector field $\mathbf{A}$ acting on functions of $\mathscr{W}$, we will denote the adjoint of $\mathbf{A}$ in the $L^{2}(\mu)$ inner product by $\mathbf{A}^{*}$, which has domain in $L^{2}(\mu)$ consisting of functions $G$ such that

$$
|\mathbb{E}[(\mathbf{A} F) G]| \leq C\|F\|_{L^{2}(\mu)}
$$

for all $F \in \mathcal{D}^{1,2}$, for some constant $C$. For functions $G$ in the domain of $\mathbf{A}^{*}$,

$$
\mathbb{E}\left[F\left(\mathbf{A}^{*} G\right)\right]=\mathbb{E}[(\mathbf{A} F) G]
$$

for all $F \in \mathcal{D}^{1,2}$.
Note that for any $F \in \mathcal{D}^{1,2}$,

$$
\mathbb{E}[\mathbf{A} F]=\mathbb{E}\left[(D F, \mathbf{A})_{\mathscr{H}}\right]=\mathbb{E}\left[F D^{*} \mathbf{A}\right] .
$$

Thus, we must have that $\mathbf{A}^{*} 1=D^{*} \mathbf{A}$ a.s. Recall that $D^{*}$ is a continuous operator from $\mathcal{D}^{\infty}(\mathscr{H})$ into $\mathcal{D}^{\infty}$. Thus, for $\mathbf{A}$ a vector field on $\mathscr{W}$ as defined in Equation (2.24), Remark 2.18 implies that

$$
D^{*} \mathbf{A}=\sum_{i, j=1}^{3} D^{*}\left(\sigma_{i j}^{-1} \tilde{A}^{j}\left(\xi_{t}\right) D \xi_{t}^{i}\right) \in \mathcal{D}^{\infty}
$$

Thus we have the following proposition.
Proposition 2.21. Let $\tilde{A}$ be a left invariant vector field on $G$ with lifted vector field $\mathbf{A}$ on $\mathscr{W}$ as defined by Equation (2.24). Then $\mathbf{A}^{*} 1$, where $\mathbf{A}^{*}$ is the $L^{2}(\mu)$-adjoint of $\mathbf{A}$, is an element of $\mathcal{D}^{\infty}$.

### 2.8 An $L^{p}$-type gradient estimate $(p>1)$ and a Poincaré inequality

Theorem 2.22. For all $p>1$,

$$
\begin{equation*}
\left|\nabla P_{t} f\right|^{p} \leq K_{p} P_{t}|\nabla f|^{p}, \tag{2.25}
\end{equation*}
$$

for all $f \in C_{p}^{\infty}(G)$ and $t>0$, where

$$
\begin{equation*}
K_{p}:=2^{p / q}+2^{p\left(\frac{1}{q}+\frac{1}{2}\right)}\left[\left\|\mathbf{X}^{*} \xi_{1}^{1}\right\|_{L^{q}(\mu)}^{2}+\left\|\mathbf{X}^{*} \xi_{1}^{2}\right\|_{L^{q}(\mu)}^{2}\right]^{p / 2}<\infty \tag{2.26}
\end{equation*}
$$

with $\mathbf{X}^{*}$ the adjoint of the lifted vector field $\mathbf{X}$ as in Equation (2.24) with $t=1$, and $q=\frac{p}{p-1}$.

Proof. By Proposition 2.6, we know the constants $K_{p}$ are independent of $t$. Also, Lemma 2.4 states that the inequality is translation invariant. Thus, the proof is reduced to verifying the inequality at the identity for $t=1$; that is, we must find finite constants $K_{p}$ such that

$$
\begin{equation*}
\left|\nabla P_{1} f\right|^{p}(0) \leq K_{p} P_{1}|\nabla f|^{p}(0), \tag{2.27}
\end{equation*}
$$

for all $f \in C_{p}^{\infty}(G)$. So applying Remark 2.1 and Lemma 4.2, consider

$$
\begin{aligned}
\tilde{X} P_{1} f(0) & =P_{1} \hat{X} f(0) \\
& =P_{1}(\tilde{X}+y \tilde{Z}) f(0)=P_{1}(\tilde{X} f)(0)+P_{1}(y \tilde{Z} f)(0) .
\end{aligned}
$$

Similarly,

$$
\tilde{Y} P_{1} f(0)=P_{1}(\tilde{Y} f)(0)-P_{1}(x \tilde{Z} f)(0) .
$$

Thus,

$$
\begin{align*}
\left|\nabla P_{1} f\right|^{p}(0) & =\left|P_{1} \nabla f+P_{1}\left(\binom{y}{-x} \tilde{Z} f\right)\right|^{p}(0) \\
& \leq\left(\left|P_{1} \nabla f\right|+\left|P_{1}\left(\binom{y}{-x} \tilde{Z} f\right)\right|\right)^{p}(0) \\
& \leq 2^{p / q}\left(\left|P_{1} \nabla f\right|^{p}(0)+\left|P_{1}\left(\binom{y}{-x} \tilde{Z} f\right)\right|^{p}(0)\right) \tag{2.28}
\end{align*}
$$

where

$$
\left|P_{1}\left(\binom{y}{-x} \tilde{Z} f\right)\right|^{p}(0)=\left[\left|P_{1}(y \tilde{Z} f)\right|^{2}(0)+\left|P_{1}(x \tilde{Z} f)\right|^{2}(0)\right]^{p / 2}
$$

and $q=\frac{p}{p-1}$ is the conjugate exponent to $p$. Expanding $\tilde{Z}$ in terms of the Lie bracket gives $\tilde{Z}=\tilde{X} \tilde{Y}-\tilde{Y} \tilde{X}$. So let $F=\left(F^{1}, F^{2}, F^{3}\right):=\xi_{1}$, and using Equation 2.22, consider

$$
\begin{align*}
P_{1}(y \tilde{Z} f)(0) & =P_{1}(y \tilde{X} \tilde{Y} f)(0)-P_{1}(y \tilde{Y} \tilde{X} f)(0) \\
& =\mathbb{E}\left[F^{2}(\tilde{X} \tilde{Y} f)(F)\right]-\mathbb{E}\left[F^{2}(\tilde{Y} \tilde{X} f)(F)\right] \\
& =\mathbb{E}\left[F^{2} \mathbf{X}((\tilde{Y} f)(F))\right]-\mathbb{E}\left[F^{2} \mathbf{Y}((\tilde{X} f)(F))\right] \\
& =\mathbb{E}\left[\mathbf{X}^{*} F^{2} \cdot(\tilde{Y} f)(F)\right]-\mathbb{E}\left[\mathbf{Y}^{*} F^{2} \cdot(\tilde{X} f)(F)\right], \tag{2.29}
\end{align*}
$$

where $\mathbf{X}$ and $\mathbf{Y}$ are the lifted vector fields of $\tilde{X}$ and $\tilde{Y}$, as in Equation 2.24 , with $t=1$. Hence,

$$
\begin{aligned}
\left|P_{1}(y \tilde{Z} f)\right|^{2}(0) \leq & \left(\left|\mathbb{E}\left[\mathbf{X}^{*} F^{2}(\tilde{Y} f)(F)\right]\right|+\left|\mathbb{E}\left[\mathbf{Y}^{*} F^{2}(\tilde{X} f)(F)\right]\right|\right)^{2} \\
\leq & 2\left(\left|\mathbb{E}\left[\mathbf{X}^{*} F^{2}(\tilde{Y} f)(F)\right]\right|^{2}+\left|\mathbb{E}\left[\mathbf{Y}^{*} F^{2}(\tilde{X} f)(F)\right]\right|^{2}\right) \\
\leq & 2\left[\left(\mathbb{E}\left|\mathbf{X}^{*} F^{2}\right|^{q}\right)^{2 / q}\left(P_{1}|\tilde{Y} f|^{p}\right)^{2 / p}(0)\right. \\
& \left.+\left(\mathbb{E}\left|\mathbf{Y}^{*} F^{2}\right|^{q}\right)^{2 / q}\left(P_{1}|\tilde{X} f|^{p}\right)^{2 / p}(0)\right]
\end{aligned}
$$

by Hölder's inequality. Similarly,

$$
\left|P_{1}(x \tilde{Z} f)\right|^{2}(0) \leq 2\left[\left(\mathbb{E}\left|\mathbf{X}^{*} F^{1}\right|^{q}\right)^{2 / q}\left(P_{1}|\tilde{Y} f|^{p}\right)^{2 / p}(0)+\left(\mathbb{E}\left|\mathbf{Y}^{*} F^{1}\right|^{q}\right)^{2 / q}\left(P_{1}|\tilde{X} f|^{p}\right)^{2 / p}(0)\right] .
$$

Combining this with Equation (2.28), we have

$$
\begin{aligned}
\left|\nabla P_{1} f\right|^{p}(0) \leq & 2^{p / q}\left(\left|P_{1} \nabla f\right|^{p}(0)+\left[2\left(\mathbb{E}\left|\mathbf{X}^{*} F^{2}\right|^{q}\right)^{2 / q}\left(P_{1}|\tilde{Y} f|^{p}\right)^{2 / p}(0)\right.\right. \\
& +2\left(\mathbb{E}\left|\mathbf{Y}^{*} F^{2}\right|^{q}\right)^{2 / q}\left(P_{1}|\tilde{X} f|^{p}\right)^{2 / p}(0) \\
& +2\left(\mathbb{E}\left|\mathbf{X}^{*} F^{1}\right|^{q}\right)^{2 / q}\left(P_{1}|\tilde{Y} f|^{p}\right)^{2 / p}(0) \\
& \left.\left.+2\left(\mathbb{E}\left|\mathbf{Y}^{*} F^{1}\right|^{q}\right)^{2 / q}\left(P_{1}|\tilde{X} f|^{p}\right)^{2 / p}(0)\right]^{p / 2}\right) \\
\leq & 2^{p / q}\left(P_{1}|\nabla f|^{p}(0)\right. \\
& +2^{p / 2}\left[\left(P_{1}|\tilde{X} f|^{p}\right)^{2 / p}(0)\left[\left(\mathbb{E}\left|\mathbf{Y}^{*} F^{1}\right|^{q}\right)^{2 / q}+\left(\mathbb{E}\left|\mathbf{Y}^{*} F^{2}\right|^{q}\right)^{2 / q}\right]\right. \\
& \left.+\left(P_{1}|\tilde{Y} f|^{p}\right)^{2 / p}(0)\left[\left(\mathbb{E}\left|\mathbf{X}^{*} F^{1}\right|^{q}\right)^{2 / q}+\left(\mathbb{E}\left|\mathbf{X}^{*} F^{2}\right|^{q}\right)^{2 / q}\right]^{p / 2}\right)
\end{aligned}
$$

where we use Hölder's inequality and that $p_{1}(g) d g$ is a probability measure to get

$$
\left|P_{1} \nabla f\right|^{p}(0) \leq P_{1}|\nabla f|^{p}(0)
$$

So let

$$
C_{p}:=\left(\mathbb{E}\left|\mathbf{X}^{*} F^{1}\right|^{q}\right)^{2 / q}+\left(\mathbb{E}\left|\mathbf{X}^{*} F^{2}\right|^{q}\right)^{2 / q} .
$$

Note that $C_{p}$ is a finite constant for all $p>1$ by Hölder's inequality, Remark 2.14, and Proposition 2.21, since

$$
\mathbf{A}^{*} F=D^{*}(F \mathbf{A})
$$

for any vector field $\mathbf{A}$ on $\mathscr{W}$ and $F \in \mathcal{D}^{\infty}$. By symmetry,

$$
C_{p}=\left(\mathbb{E}\left|\mathbf{Y}^{*} F^{1}\right|^{q}\right)^{2 / q}+\left(\mathbb{E}\left|\mathbf{Y}^{*} F^{2}\right|^{q}\right)^{2 / q} .
$$

Thus,

$$
\begin{aligned}
\left|\nabla P_{1} f\right|^{p}(0) & \leq 2^{p / q} P_{1}|\nabla f|^{p}(0)+\left(2 C_{p}\right)^{p / 2}\left[\left(P_{1}|\tilde{X} f|^{p}\right)^{2 / p}(0)+\left(P_{1}|\tilde{Y} f|^{p}\right)^{2 / p}(0)\right]^{p / 2} \\
& \leq\left(2^{p / q}+2^{p\left(\frac{1}{q}+\frac{1}{2}\right)} C_{p}^{p / 2}\right) P_{1}|\nabla f|^{p}(0),
\end{aligned}
$$

which proves Equation (2.27), and hence, the theorem.
Theorem 2.23 (Poincaré Inequality). Let $K_{2}$ be given by Equation (2.26) with $p=2$;
that is

$$
K_{2}=2+4\left(\left\|\mathbf{X}^{*} \xi_{1}^{1}\right\|_{L^{2}(\mu)}^{2}+\left\|\mathbf{X}^{*} \xi_{1}^{2}\right\|_{L^{2}(\mu)}^{2}\right)<\infty
$$

Let $p_{t}(g) d g$ be the hypoelliptic heat kernel. Then

$$
\begin{equation*}
\int_{\mathbb{R}^{3}} f^{2}(g) p_{t}(g) d g-\left(\int_{\mathbb{R}^{3}} f(g) p_{t}(g) d g\right)^{2} \leq K_{2} t \int_{\mathbb{R}^{3}}|\nabla f|^{2}(g) p_{t}(g) d g, \tag{2.30}
\end{equation*}
$$

for all $f \in C_{p}^{\infty}(G)$ and $t>0$.
Proof. Let $F_{t}(g)=\left(P_{t} f\right)(g)$. Then

$$
\frac{d}{d t} P_{t-s} F_{s}^{2}=P_{t-s}\left(-\frac{1}{2} L F_{s}^{2}+F_{s} L F_{s}\right)=-P_{t-s}\left|\nabla F_{s}\right|^{2}
$$

Integrating this equation on $t$ implies that

$$
\begin{aligned}
P_{t} f^{2}-\left(P_{t} f\right)^{2} & =\int_{0}^{t} P_{t-s}\left|\nabla F_{s}\right|^{2} d s \\
& =\int_{0}^{t} P_{t-s}\left|\nabla P_{s} f\right|^{2} d s \\
& \leq \int_{0}^{t} K_{2} P_{t-s} P_{s}|\nabla f|^{2} d s=K_{2} t P_{t}|\nabla f|^{2},
\end{aligned}
$$

wherein we have made use of Theorem 2.22. Evaluating the above at 0 gives the desired result.

### 2.9 The $L^{1}$-type gradient estimate

### 2.9.1 Method fails for the $p=1$ case

In this section, we show that the argument in the proof of Theorem 2.22 can not be used to prove the inequality $(2.25)$ for $p=1$.

Proposition 2.24. Let $F=\left(F^{1}, F^{2}, F^{3}\right):=\xi_{1}$. Then

$$
\left\|\mathbf{X}^{*} F^{1}\right\|_{L^{\infty}(\mu)}+\left\|\mathbf{X}^{*} F^{2}\right\|_{L^{\infty}(\mu)}=\infty
$$

Proof. Let $\sigma(F)$ denote the $\sigma$ - algebra generated by $F: \mathscr{W} \rightarrow G$ and $p_{t}(g) d g$ denote the Heisenberg group heat kernel. Then for $f \in C_{c}^{1}\left(\mathbb{R}^{3}\right)$

$$
\begin{aligned}
\mathbb{E}\left[\mathbf{X}^{*} F^{1} f(F)\right] & =\mathbb{E}\left[F^{1}(\tilde{X} f)(F)\right]=P_{1}(x \tilde{X} f)(0) \\
& =\int_{G} x \tilde{X} f(g) p_{1}(g) d g \\
& =-\int_{G} f(g) \tilde{X}\left(x p_{1}(g)\right) d g \\
& =-\int_{G} f(g)\left(1+x \tilde{X} \ln p_{1}(g)\right) p_{1}(g) d g=-\mathbb{E}\left[f(F)\left(1+x \tilde{X} \ln p_{1}\right)(F)\right],
\end{aligned}
$$

where in the third line we have applied standard integration by parts. Consequently, we have shown

$$
\mathbb{E}\left[\mathbf{X}^{*} F^{1} \mid \sigma(F)\right]=-\left(1+x \tilde{X} \ln p_{1}\right)(F) .
$$

By a similar computation one also shows

$$
\mathbb{E}\left[\mathbf{X}^{*} F^{2} \mid \sigma(F)\right]=-\left(y \tilde{X} \ln p_{1}\right)(F) .
$$

Since conditional expectation is $L^{p}$-contractive and the law of $F$ is absolutely continuous relative to Lebesgue measure, it now follows that

$$
\begin{aligned}
\left\|\mathbf{X}^{*} F^{1}\right\|_{L^{\infty}(\mu)}+\left\|\mathbf{X}^{*} F^{2}\right\|_{L^{\infty}(\mu)} & \geq\left\|\mathbb{E}\left[\mathbf{X}^{*} F^{1} \mid \sigma(F)\right]\right\|_{L^{\infty}(\mu)}+\left\|\mathbb{E}\left[\mathbf{X}^{*} F^{2} \mid \sigma(F)\right]\right\|_{L^{\infty}(\mu)} \\
& =\left\|1+x \tilde{X} \ln p_{1}\right\|_{L^{\infty}\left(\mathbb{R}^{3}, m\right)}+\left\|y \tilde{X} \ln p_{1}\right\|_{L^{\infty}\left(\mathbb{R}^{3}, m\right)},
\end{aligned}
$$

where $m$ is Lebesgue measure. Hence, it suffices to show that either $x \tilde{X} \ln p_{1}$ or $y \tilde{X} \ln p_{1}$ is unbounded. We will show $x \tilde{X} \ln p_{1}$ is unbounded by making use the formula for $p_{t}(g)$ in Equation (2.5). Letting $t=1$ in Equation (2.5) and making the change of variables $w \mapsto 2 w$, we have

$$
p_{1}(g)=\frac{1}{2 \pi^{2}} \int_{\mathbb{R}} \frac{w}{\sinh w} \exp \left(-\frac{1}{2}|\vec{x}|^{2} w \operatorname{coth} w\right) e^{2 i w z} d w .
$$

Then applying $\tilde{X}=\partial_{x}-\frac{1}{2} y \partial_{z}$ yields

$$
\tilde{X} p_{1}(g)=-\frac{1}{2 \pi^{2}} \int_{\mathbb{R}}(x w \operatorname{coth} w+i y w) \frac{w}{\sinh w} \exp \left(-\frac{1}{2}|\vec{x}|^{2} w \operatorname{coth} w\right) e^{2 i w z} d w .
$$

Setting $y=z=0$, we then have

$$
\tilde{X} \ln p_{1}(x, 0,0)=-x \int_{\mathbb{R}} w \operatorname{coth} w d \nu_{x}(w),
$$

where

$$
\begin{equation*}
d \nu_{x}(w):=\frac{1}{z_{x}} \frac{w}{\sinh w} \exp \left(-\frac{1}{2} x^{2} w \operatorname{coth} w\right) d w \tag{2.31}
\end{equation*}
$$

and $z_{x}$ is the normalizing constant

$$
z_{x}:=\int_{\mathbb{R}} \frac{w}{\sinh w} \exp \left(-\frac{1}{2} x^{2} w \operatorname{coth} w\right) d w .
$$

By Lemma 2.25 below,

$$
\lim _{x \rightarrow \infty} \int_{\mathbb{R}} w \operatorname{coth} w d \nu_{x}(w)=1
$$

and so

$$
\lim _{x \rightarrow \infty} \tilde{X} \ln p_{1}(x, 0,0)=\lim _{x \rightarrow \infty}\left(-x \int_{\mathbb{R}} w \operatorname{coth} w d \nu_{x}(w)\right)=-\infty .
$$

Lemma 2.25. Let $\psi(w)=w \operatorname{coth} w-1$ and $\nu_{x}$ be as in Equation 2.31. Then

$$
\begin{equation*}
\lim _{x \rightarrow \infty} \int \psi d \nu_{x}=\psi(0)=0 \tag{2.32}
\end{equation*}
$$

Proof. Since $\psi(0)=0$ and $\psi$ is continuous, to prove Equation (2.32) it suffices to show by the usual approximation of $\delta$ - function arguments that

$$
\lim _{x \rightarrow \infty} \int_{|w| \geq \epsilon} \psi(w) d \nu_{x}(w)=0
$$

holds for every $\epsilon>0$. We begin by rewriting Equation 2.31) as

$$
d \nu_{x}(w)=\frac{1}{Z_{x}} \frac{w}{\sinh w} \exp \left(-\frac{1}{2} x^{2} \psi(w)\right) d w
$$

where

$$
Z_{x}:=\int_{\mathbb{R}} \frac{w}{\sinh w} \exp \left(-\frac{1}{2} x^{2} \psi(w)\right) d w .
$$

A glance at the graph of $\psi$ will convince the reader that there are constants $\alpha, \beta>0$ (depending on $\epsilon>0$ ) such that $\alpha|w| \leq \psi(w) \leq \beta|w|$ for all $|w| \geq \epsilon$. (In fact, one could take $\beta=1$ independent of $\epsilon$.) Thus

$$
\begin{aligned}
\int_{|w| \geq \epsilon} \psi(w) \frac{w}{\sinh w} \exp \left(-\frac{1}{2} x^{2} \psi(w)\right) d w & \leq 2 \int_{w \geq \epsilon} \beta w e^{-\alpha x^{2} w / 2} d w \\
& =\frac{4 \beta}{x^{2} \alpha}\left(\epsilon+\frac{2}{x^{2} \alpha}\right) e^{-\alpha x^{2} \epsilon / 2}
\end{aligned}
$$

where in the inequality we have also used that $\frac{w}{\sinh w} \leq 1$.
Now consider the constant $Z_{x}$. We know that for $w$ small, there exists a constant $\gamma>0$ such that $\psi(w) \leq \gamma w^{2}$. So letting $\varphi(w)=\frac{w}{\sinh w}$, we have

$$
\begin{aligned}
Z_{x} & \geq \int_{|w| \leq \epsilon} \varphi(w) \exp \left(-\frac{1}{2} x^{2} \psi(w)\right) d w \\
& \geq \int_{-\epsilon}^{\epsilon} \varphi(w) e^{-\gamma x^{2} w^{2} / 2} d w=\frac{1}{x} \int_{-\epsilon x}^{\epsilon x} \varphi\left(\frac{w}{x}\right) e^{-\gamma w^{2} / 2} d w
\end{aligned}
$$

where we have made the change of variables $w \mapsto \frac{w}{x}$. So, by the dominated convergence theorem,

$$
\liminf _{x \rightarrow \infty}\left(x Z_{x}\right) \geq \liminf _{x \rightarrow \infty} \int_{-\epsilon x}^{\epsilon x} \varphi\left(\frac{w}{x}\right) e^{-\gamma w^{2} / 2} d w=\varphi(0) \int_{-\infty}^{\infty} e^{-\gamma w^{2} / 2} d w=\sqrt{\frac{2 \pi}{\gamma}}
$$

Thus, $Z_{x} \geq \frac{1}{2} \sqrt{\frac{2 \pi}{\gamma}} \frac{1}{x}$ for $x$ sufficiently large, and so

$$
\begin{aligned}
\lim _{x \rightarrow \infty} \int_{|w| \geq \epsilon} \psi(w) d \nu_{x}(w) & =\lim _{x \rightarrow \infty} \frac{1}{Z_{x}} \int_{|w| \geq \epsilon} \psi(w) \frac{w}{\sinh w} \exp \left(-\frac{1}{2} x^{2} \psi(w)\right) d w \\
& \leq 2 \lim _{x \rightarrow \infty} \frac{\frac{4 \beta}{x^{2} \alpha}\left(\epsilon+\frac{2}{x^{2} \alpha}\right) e^{-\alpha x^{2} \epsilon / 2}}{\sqrt{\frac{2 \pi}{\gamma}} \frac{1}{x}}=0
\end{aligned}
$$

as desired.

### 2.9.2 Function classes for which the $p=1$ inequality holds

Although we have not been able to adapt this method of proof to show the $p=1$ case, we have shown using ad hoc methods that several large classes of functions on $G$ do satisfy the $p=1$ inequality.

## Even and odd functions in $z$.

Proposition 2.26. Suppose $f \in C_{p}^{\infty}(G)$ is an even or odd function in $z$ on the Heisenberg group $G$. Then

$$
\left|\nabla P_{t} f\right| \leq \sqrt{2} P_{t}|\nabla f|,
$$

for all $t>0$.
Proof. Suppose $f \in C_{p}^{\infty}(G)$ is an even function on $G$. Let $g=(x, y, z)$. Then

$$
\begin{aligned}
\tilde{X} P_{t} f(0)=P_{t}(\hat{X} f)(0) & =\int_{G}\left(f_{x}+\frac{1}{2} y f_{z}\right)(g) p_{t}(g) d g \\
& =\int_{G}\left(f_{x}-\frac{1}{2} y f_{z}\right)(g) p_{t}(g) d g=P_{t}(\tilde{X} f)(0)
\end{aligned}
$$

where the third equality follows from

$$
\int_{\mathbb{R}} f_{z}(g) p_{t}(g) d z=0
$$

since $f_{z} p_{t}$ is an odd function in $z$. Thus,

$$
\left|\tilde{X} P_{t} f\right|(0)=\left|P_{t} \tilde{X} f\right|(0) \leq P_{t}|\tilde{X} f|(0) .
$$

Similarly, it may be shown that $\tilde{Y} P_{t} f(0)=P_{t} \tilde{Y} f(0)$, which implies that

$$
\left|\tilde{Y} P_{t} f\right|(0) \leq P_{t}|\tilde{Y} f|(0)
$$

Thus, for all $f \in C_{c}^{\infty}(G)$ which are even in $z$,

$$
\begin{aligned}
\left|\nabla P_{t} f\right|(0) & =\left(\left|\tilde{X} P_{t} f\right|^{2}+\left|\tilde{Y} P_{t} f\right|^{2}\right)^{1 / 2} \\
& \leq\left(\left(P_{t}|\tilde{X} f|\right)^{2}+\left(P_{t}|\tilde{Y} f|\right)^{2}\right)^{1 / 2}(0) \leq P_{t}|\tilde{X} f|(0)+P_{t}|\tilde{Y} f|(0) \leq \sqrt{2} P_{t}|\nabla f|(0)
\end{aligned}
$$

since $(a+b)^{1 / 2} \leq a^{1 / 2}+b^{1 / 2}$.
Now suppose $f \in C_{p}^{\infty}(G)$ is odd in $z$. As in the argument above, write

$$
\begin{aligned}
\tilde{X} P_{t} f(0)=P_{t}(\hat{X} f)(0) & =\int_{G}\left(f_{x}+\frac{1}{2} y f_{z}\right)(g) p_{t}(g) d g \\
& =\int_{G}\left(-f_{x}+\frac{1}{2} y f_{z}\right)(g) p_{t}(g) d g=-P_{t}(\tilde{X} f)(0)
\end{aligned}
$$

where in the third equality we have used that

$$
\int_{\mathbb{R}} f_{x}(g) p_{t}(g) d z=0
$$

as $f_{x}$ is odd in $z$ and $p_{t}$ is even in $z$. Thus,

$$
\left|\tilde{X} P_{t} f\right|(0)=\left|P_{t} \tilde{X} f\right|(0) \leq P_{t}|\tilde{X} f|(0)
$$

A parallel argument shows that $\left|\tilde{Y} P_{t} f\right|(0) \leq P_{t}|\tilde{Y} f|(0)$, and thus

$$
\left|\nabla P_{t} f\right|(0) \leq \sqrt{2} P_{t}|\nabla f|(0)
$$

for all $f \in C_{c}^{\infty}(G)$ which are odd functions in $z$.

Functions $f(x, y, z)=e^{i \lambda z} h(x, y)$.
Proposition 2.27. Suppose $f(x, y, z)=e^{i \lambda z} h(x, y)$ for some $h \in C_{p}^{\infty}\left(\mathbb{R}^{2}, \mathbb{C}\right)$. Then

$$
\left|\nabla P_{t} f\right| \leq(\sqrt{2}+4) P_{t}|\nabla f|
$$

for all $t>0$.

Proof. In the case $\lambda=0, f=f(x, y)$ depends only on $\vec{x}=(x, y)$ and

$$
P_{t} f=e^{t L / 2} f=e^{t \Delta / 2} f,
$$

where $\Delta$ is the standard two-dimensional Laplacian, and this reduces to the Gaussian measure case, for which the inequality holds.

So suppose $\lambda \neq 0$, and let $g=(x, y, z)$. For any function $f \in C_{p}^{\infty}(G)$, the computations in 2.28) imply that

$$
\begin{align*}
\left|\nabla P_{t} f\right|(0) \leq\left|P_{t} \nabla f\right|(0)+\left|\int_{G}\binom{y}{-x}(\tilde{Z} f)(g) p_{t}(g) d g\right| \\
\leq\left|P_{t} \nabla f\right|(0)+\left|\int_{G} y(\tilde{Z} f)(g) p_{t}(g) d g\right|+\left|\int_{G} x(\tilde{Z} f)(g) p_{t}(g) d g\right| \tag{2.33}
\end{align*}
$$

For $f(x, y, z)=e^{i \lambda z} h(x, y)$, define the ratios

$$
\Lambda_{1}(f):=\frac{\left|\int_{G} y \tilde{Z} f(g) p_{2}(g) d g\right|}{P_{2}|\tilde{Y} f|(0)} \quad \text { and } \quad \Lambda_{2}(f):=\frac{\left|\int_{G} x \tilde{Z} f(g) p_{2}(g) d g\right|}{P_{2}|\tilde{X} f|(0)}
$$

For $\lambda>0$ and $h$ a real-valued function,

$$
\Lambda_{1}(f)=\frac{\left|\int_{G} y \lambda e^{i \lambda z} h(x, y) p_{2}(g) d g\right|}{\int_{G}\left|g_{y}(x, y)+\frac{i}{2} \lambda x h(x, y)\right| p_{2}(g) d g} .
$$

The numerator of this expression may be written as

$$
\begin{aligned}
& \int_{G} y \lambda e^{i \lambda z} h(x, y) p_{2}(g) d g \\
&=\frac{1}{8 \pi^{2}} \int_{\mathbb{R}} d x \int_{\mathbb{R}} d y \lambda y h(x, y) \int_{\mathbb{R}} d z e^{i \lambda z} \int_{\mathbb{R}} d w \frac{w}{\sinh w} \exp \left(-|\vec{x}|^{2} \frac{w \operatorname{coth} w}{4}\right) e^{i w z} \\
&=\frac{1}{8 \pi^{2}} \int_{\mathbb{R}^{2}} \lambda y h(x, y) \frac{\lambda}{\sinh \lambda} \exp \left(-|\vec{x}|^{2} \frac{\lambda \operatorname{coth} \lambda}{4}\right) d x d y \\
&=\frac{1}{8 \pi^{2}} \frac{\lambda^{2}}{\sinh \lambda} \int_{\mathbb{R}^{2}} h(x, y) \frac{-2}{\lambda \operatorname{coth} \lambda} \partial_{y}\left[\exp \left(-|\vec{x}|^{2} \frac{\lambda \operatorname{coth} \lambda}{4}\right)\right] d x d y \\
&=\frac{1}{8 \pi^{2}} \frac{-2 \lambda}{\cosh \lambda} \int_{\mathbb{R}^{2}} h_{y}(x, y) \exp \left(-|\vec{x}|^{2} \frac{\lambda \operatorname{coth} \lambda}{4}\right) d x d y,
\end{aligned}
$$

where in the second equality we have evaluated the Fourier transform in $w$ and the inverse Fourier transform in $-z$. Thus,

$$
\begin{aligned}
\Lambda_{1}(f) & =\frac{\frac{1}{8 \pi^{2}}\left|\frac{2 \lambda}{\cosh \lambda} \int_{\mathbb{R}^{2}} h_{y}(x, y) \exp \left(-|\vec{x}|^{2} \frac{\lambda \operatorname{coth} \lambda}{4}\right) d x d y\right|}{\int_{G}\left|h_{y}(x, y)+\frac{i \lambda}{2} x h(x, y)\right| p_{2}(g) d g} \\
& \leq\left|\frac{2 \lambda}{\cosh \lambda}\right| \frac{\frac{1}{8 \pi^{2}}\left|\int_{\mathbb{R}^{2}} h_{y}(x, y) \exp \left(-|\vec{x}|^{2} \frac{\lambda \operatorname{coth} \lambda}{4}\right) d x d y\right|}{\int_{G}\left|h_{y}(x, y)\right| p_{2}(g) d g} \\
& =\left|\frac{2 \lambda}{\cosh \lambda}\right| \frac{\left|\int_{\mathbb{R}^{2}} h_{y}(x, y) \exp \left(-|\vec{x}|^{2} \frac{\lambda \operatorname{coth} \lambda}{4}\right) d x d y\right|}{\int_{\mathbb{R}^{2}}\left|h_{y}(x, y)\right| e^{-\frac{|\overrightarrow{\mid}|^{2}}{4}} d x d y} \\
& \leq\left|\frac{2 \lambda}{\cosh \lambda}\right| \frac{\int_{\mathbb{R}^{2}}\left|h_{y}(x, y)\right| \exp \left(-|\vec{x}|^{2} \frac{\lambda \operatorname{coth} \lambda}{4}\right) d x d y}{\int_{\mathbb{R}^{2}}\left|h_{y}(x, y)\right| e^{-\frac{|\vec{x}|^{2}}{4}} d x d y} \leq\left|\frac{2 \lambda}{\cosh \lambda}\right| \leq 2,
\end{aligned}
$$

where in the second equality we have evaluated

$$
\int_{\mathbb{R}} p_{2}(g) d z=\frac{1}{8 \pi^{2}} e^{-\frac{|\vec{x}|^{2}}{4}}
$$

and in the penultimate inequality we have used that $\lambda \operatorname{coth} \lambda \geq 1$ and

$$
\left|\frac{\int_{\mathbb{R}^{2}} h(x, y) e^{k|\vec{x}|^{2}} d x d y}{\int_{\mathbb{R}^{2}} h(x, y) e^{|\vec{x}|^{2}} d x d y}\right| \leq 1
$$

for any function $h \in C_{p}^{\infty}\left(\mathbb{R}^{2}\right)$ and $k \geq 1$. Thus it has been shown that

$$
\left|\int_{G} y \tilde{Z} f(g) p_{2}(g) d g\right| \leq 2 P_{2}|\tilde{Y} f|(0)
$$

Similarly, one may show that $\Lambda_{2}(f) \leq 2$ by performing integration by parts in the $x$ variable and comparing with the appropriate denominator term. Thus,

$$
\left|\int_{G} x \tilde{Z} f(g) p_{2}(g) d g\right| \leq 2 P_{2}|\tilde{X} f|(0)
$$

Combining these inequalities with 2.33 ,

$$
\begin{aligned}
& \left|\nabla P_{2} f\right|(0) \leq\left|P_{2} \nabla f\right|(0)+2\left(P_{2}|\tilde{Y} f|(0)+P_{2}|\tilde{X} f|(0)\right) \\
& \quad \leq P_{2}|\nabla f|(0)+2 \sqrt{2} P_{2}|\nabla f|(0)=(1+2 \sqrt{2}) P_{2}|\nabla f|(0)
\end{aligned}
$$

for all functions $f(x, y, z)=e^{i \lambda z} g(x, y)$ where $g$ is real valued. With obvious modifications to the statement and proof of Proposition 2.6, this is then sufficient proof that

$$
\left|\nabla P_{t} f\right|(0) \leq K P_{t}|\nabla f|(0)
$$

with $K:=1+2 \sqrt{2}$, for all $t>0$.
Clearly the previous argument holds when $g$ is strictly imaginary. So consider $g=u+i v$ for $u, v$ real functions.

$$
\begin{aligned}
\left|\nabla P_{t} f\right|(x, y, z) & =\left|\nabla P_{t} e^{i \lambda z}(u(x, y)+i v(x, y))\right| \\
& \leq K\left(P_{t}\left|\nabla e^{i \lambda z} u(x, y)\right|+P_{t}\left|\nabla e^{i \lambda z} v(x, y)\right|\right) \\
& \leq \sqrt{2} K P_{t}\left|\nabla e^{i \lambda z}(u(x, y)+i v(x, y))\right|=\sqrt{2} K P_{t}|\nabla f|(x, y, z)
\end{aligned}
$$

as desired.

Radial functions $f=f(r, z)$.
Proposition 2.28. Suppose $f(x, y, z)=g(r, z)$ for some function $g \in C_{c}^{\infty}\left(\mathbb{R}^{2}\right)$, where $r^{2}=x^{2}+y^{2}$. Then

$$
|\hat{\nabla} f|=|\nabla f|
$$

where $\hat{\nabla}=(\hat{X}, \hat{Y})$ is the right invariant gradient on $G$. In particular, this implies that

$$
\left|\nabla P_{t} f\right| \leq P_{t}|\nabla f|
$$

for all $t>0$.
Proof. For $f(x, y, z)=g(r, z)$,

$$
\tilde{X} f=\frac{x}{r} g_{r}-\frac{1}{2} y g_{z} \text { and } \tilde{Y} f=\frac{y}{r} g_{r}+\frac{1}{2} x g_{z} .
$$

Thus,

$$
|\nabla f|^{2}=\left(\frac{x^{2}}{r^{2}} g_{r}^{2}-\frac{x y}{r} g_{r} g_{z}+\frac{1}{4} y^{2} g_{z}^{2}\right)+\left(\frac{y^{2}}{r^{2}} g_{r}^{2}+\frac{x y}{r} g_{r} g_{z}+\frac{1}{4} x^{2} g_{z}^{2}\right)=g_{r}^{2}+\frac{1}{4} r^{2} g_{z}^{2}
$$

Recall from Equation 2.4 that

$$
\hat{X}=\partial_{x}+\frac{1}{2} y \partial_{z} \text { and } \hat{Y}=\partial_{y}-\frac{1}{2} x \partial_{z} .
$$

Therefore,

$$
\hat{X} f=\frac{x}{r} g_{r}+\frac{1}{2} y g_{z} \text { and } \hat{Y} f=\frac{y}{r} g_{r}-\frac{1}{2} x g_{z},
$$

implies that

$$
|\hat{\nabla} f|^{2}=\left(\frac{x^{2}}{r^{2}} g_{r}^{2}+\frac{x y}{r} g_{r} g_{z}+\frac{1}{4} y^{2} g_{z}^{2}\right)+\left(\frac{y^{2}}{r^{2}} g_{r}^{2}-\frac{x y}{r} g_{r} g_{z}+\frac{1}{4} x^{2} g_{z}^{2}\right)=g_{r}^{2}+\frac{1}{4} r^{2} g_{z}^{2} .
$$

Hence, $|\hat{\nabla} f|=|\nabla f|$, and it follows that

$$
\left|\nabla P_{t} f\right|(0)=\left|P_{t} \hat{\nabla} f\right|(0) \leq P_{t}|\hat{\nabla} f|(0)=P_{t}|\nabla f|(0)
$$

### 2.9.3 A Green's function inequality

Recall that for a negative linear operator $L$ on a Banach space, the set $\rho(L)$ of all $\lambda \in \mathbb{C}$ such that $\lambda I+L$ is invertible is called the resolvent set of $L$, and the family of bounded linear operators

$$
\left\{R(\lambda):=(\lambda I+L)^{-1}: \lambda \in \rho(L)\right\}
$$

is called the resolvent of $L$; see for example [25, 41].
Proposition 2.29. If

$$
\begin{equation*}
\left|\nabla P_{t} f\right| \leq K P_{t}|\nabla f|, \tag{2.34}
\end{equation*}
$$

for all $f \in C_{c}^{\infty}(G)$ and $t>0$, then

$$
\left|\nabla R(\lambda)^{m} f\right| \leq K R(\operatorname{Re} \lambda)^{m}|\nabla f|,
$$

for all $f \in C_{c}^{\infty}(G), \operatorname{Re} \lambda>0$, and $m \in \mathbb{N}$.

Proof. For all $m \in \mathbb{N}$ and $\lambda \in \rho(L)$,

$$
R(\lambda)^{m}=\frac{1}{(m-1)!} \int_{0}^{\infty} t^{m-1} e^{-\lambda t} P_{t} f d t
$$

see [41]. Suppose (2.34) holds. Then,

$$
\begin{aligned}
\left|\nabla R(\lambda)^{m} f\right| & =\left|\nabla \frac{1}{(m-1)!} \int_{0}^{\infty} t^{m-1} e^{-\operatorname{Re} \lambda t} e^{-i \operatorname{Im} \lambda t} P_{t} f d t\right| \\
& \leq \frac{1}{(m-1)!} \int_{0}^{\infty} t^{m-1} e^{-\operatorname{Re} \lambda t}\left|\nabla P_{t} f\right| d t \\
& \leq K \int_{0}^{\infty} t^{m-1} e^{-\operatorname{Re} \lambda t} P_{t}|\nabla f| d t=R(\operatorname{Re} \lambda)^{m}|\nabla f|,
\end{aligned}
$$

where we have also used the inequality (2.34) to justify differentiating through the integral.

It is shown in Folland [18] that, on the Heisenberg group, for $\lambda=0$ and $n=1$, $R(0) f=L^{-1} f=f * h$, where

$$
h(x, y, z):=\frac{C}{\sqrt{\left(x^{2}+y^{2}\right)^{2}+16 z^{2}}},
$$

for $C$ is a positive constant; that is, $u=f * h$ is a solution to the equation $L u=f$. Then the following result then is a mild indication that the Heisenberg $L^{1}$ gradient inequality might hold.

Proposition 2.30. For any $f \in C_{c}^{\infty}(G)$,

$$
|\nabla(f * h)| \leq 9|\nabla f| * h .
$$

Proof. By the translation invariance of Haar measure, it is again sufficient to verify the inequality at the identity. Considering $\nabla(f * g)(0)$ componentwise,

$$
\begin{align*}
\tilde{X}(f * h)(0)=\hat{X} & (f * h)(0)=((\hat{X} f) * h)(0) \\
& =\left(\left(\left(\tilde{X}+\frac{1}{2} y \tilde{Z}\right) f\right) * h\right)(0)=((\tilde{X} f) * h)(0)+\frac{1}{2}((y \tilde{Z} f) * h)(0) . \tag{2.35}
\end{align*}
$$

Writing $\tilde{Z}$ as the commutator of $\tilde{X}$ and $\tilde{Y}$ and letting $g=(x, y, z)$, the second term in Equation (2.35) may be expanded as

$$
\begin{aligned}
\frac{1}{2}((y \tilde{Z} f) * h)(0)= & \frac{1}{2} \int_{G} y \tilde{Z} f(g) h(g) d g=\frac{1}{2} \int_{G} y[\tilde{X}, \tilde{Y}] f(g) h(g) d g \\
= & \frac{1}{2} \int_{G} \tilde{X} \tilde{Y} f(g) y h(g) d g-\frac{1}{2} \int_{G} \tilde{Y} \tilde{X} f(g) y h(g) d g \\
= & -\frac{1}{2} \int_{G} \tilde{Y} f(g) \tilde{X}(y h(g)) d g+\frac{1}{2} \int_{G} \tilde{X} f(g) \tilde{Y}(y h(g)) d g \\
= & -\frac{1}{2} \int_{G} \tilde{Y} f(g)\left[\frac{-2 x y\left(x^{2}+y^{2}\right)+8 x y z}{\left(x^{2}+y^{2}\right)^{2}+16 z^{2}}\right] h(g) d g \\
& +\frac{1}{2} \int_{G} \tilde{X} f(g)\left[1-\frac{2 y^{2}\left(x^{2}+y^{2}\right)+8 x y z}{\left(x^{2}+y^{2}\right)^{2}+16 z^{2}}\right] h(g) d g .
\end{aligned}
$$

The following inequalities are trivially true.

$$
\begin{aligned}
&\left|\frac{2 x y\left(x^{2}+y^{2}\right)}{\left(x^{2}+y^{2}\right)^{2}+16 z^{2}}\right| \leq 1,\left|\frac{8 y^{2} z}{\left(x^{2}+y^{2}\right)^{2}+16 z^{2}}\right| \leq 1,\left|\frac{2 y^{2}\left(x^{2}+y^{2}\right)}{\left(x^{2}+y^{2}\right)^{2}+16 z^{2}}\right| \leq 2, \text { and } \\
&\left|\frac{8 x y z}{\left(x^{2}+y^{2}\right)^{2}+16 z^{2}}\right| \leq \frac{1}{2} .
\end{aligned}
$$

Thus, $|(y \tilde{Z} f) * h|(0) \leq 2(|\tilde{Y} f| * h)(0)+\frac{7}{2}(|\tilde{X} f| * h)(0)$. Combining this with Equation (2.35), we have

$$
|\tilde{X}(f * h)|(0) \leq 2(|\tilde{Y} f| * h)(0)+\frac{9}{2}(|\tilde{X} f| * h)(0) .
$$

In a similar manner, one may determine that

$$
|\tilde{Y}(f * h)|(0) \leq 2(|\tilde{X} f| * h)(0)+\frac{9}{2}(|\tilde{Y} f| * h)(0) .
$$

Hence,

$$
\begin{align*}
|\nabla(f * h)|^{2}(0) & =|\tilde{X}(f * h)|^{2}(0)+|\tilde{Y}(f * h)|^{2}(0) \\
& \leq\left[\left(2|\tilde{Y} f|+\frac{9}{2}|\tilde{X} f|\right) * h\right]^{2}(0)+\left[\left(2|\tilde{X} f|+\frac{9}{2}|\tilde{Y} f|\right) * h\right]^{2}(  \tag{0}\\
& \leq 2\left[\frac{9}{2}(|\tilde{X} f|+|\tilde{Y} f|) * h\right]^{2}(0) \\
& \leq 2\left[\frac{9}{2}\left(\sqrt{2} \sqrt{|\tilde{X} f|^{2}+|\tilde{Y} f|^{2}}\right) * h\right]^{2}(0)=(9|\nabla f| * h)^{2}(0),
\end{align*}
$$

as desired.

## Chapter 3

## Wiener calculus over $G$

We now return to the case of a general Lie group $G$ with identity $e$ and Lie algebra $\operatorname{Lie}(G)=\mathfrak{g}$, and suppose $\left\{X_{i}\right\}_{i=1}^{k} \subset \mathfrak{g}$ is a Lie generating set, in the sense of Equation (1.7). Recall that we have defined the inner product on $\mathfrak{g}$ such that $\left\{X_{i}\right\}_{i=1}^{k}$ is an orthonormal basis of the hypoelliptic subspace $\mathfrak{g}_{0}=\operatorname{span}\left(\left\{X_{i}\right\}_{i=1}^{k}\right)$.

Notation 3.1. Let $\operatorname{Ad}: G \rightarrow \operatorname{End}(\mathfrak{g})$ denote the adjoint representation of $G$ with differential ad $:=d(\operatorname{Ad}): \mathfrak{g} \rightarrow \operatorname{End}(\mathfrak{g})$. That is, $\operatorname{Ad}(g)=\operatorname{Ad}_{g}=L_{g *} R_{g^{-1} *}$, for all $g \in G$, and $\operatorname{ad}(X)=\operatorname{ad}_{X}=[X, \cdot]$, for all $X \in \mathfrak{g}$. For any function $\varphi \in C^{1}(G)$, define $\hat{\nabla} \varphi, \tilde{\nabla} \varphi: G \rightarrow \mathfrak{g}$ such that, for any $g \in G$ and $X \in \mathfrak{g}$,

$$
\langle\hat{\nabla} \varphi(g), X\rangle:=\left\langle d \varphi(g), R_{g *} X\right\rangle=(\hat{X} \varphi)(g)
$$

and

$$
\langle\tilde{\nabla} \varphi(g), X\rangle:=\left\langle d \varphi(g), L_{g *} X\right\rangle=(\tilde{X} \varphi)(g)
$$

Then,

$$
\langle\tilde{\nabla} \hat{\nabla} \varphi(g), X \otimes Y\rangle=\left.\left.\frac{d}{d s}\right|_{0} \frac{d}{d t}\right|_{0} \varphi\left(e^{s X} g e^{t Y}\right)
$$

for all $X, Y \in \mathfrak{g}$, and similarly for $\hat{\nabla} \tilde{\nabla} \varphi$.

We will use the following facts in the sequel:

$$
\begin{align*}
\langle\hat{\nabla} \varphi(g), X\rangle & =\left\langle d \varphi(g), L_{g *} L_{g^{-1} *} R_{g *} X\right\rangle  \tag{3.1}\\
& =\left\langle d \varphi(g), L_{g *} \operatorname{Ad}_{g^{-1}} X\right\rangle=\left\langle\tilde{\nabla} \varphi(g), \operatorname{Ad}_{g^{-1}} X\right\rangle
\end{align*}
$$

and similarly

$$
\begin{equation*}
\langle\tilde{\nabla} \varphi(g), X\rangle=\left\langle\hat{\nabla} \varphi(g), \operatorname{Ad}_{g} X\right\rangle . \tag{3.2}
\end{equation*}
$$

We will also use the notation developed in Section 2.5 for calculus on the path space $\mathscr{W}\left(\mathbb{R}^{k}\right)$.

### 3.1 The main theorems

### 3.1.1 The rolling map on $G$

Now suppose $\left\{b_{t}^{1}, \ldots, b_{t}^{k}\right\}$ are $k$ independent real-valued Brownian motions. Then

$$
\vec{b}_{t}:=X_{i} b_{t}^{i}:=\sum_{i=1}^{k} X_{i} b_{t}^{i}
$$

is a $\left(\mathfrak{g}_{0},\langle\cdot, \cdot\rangle\right)$ Brownian motion. In the sequel, we will always observe the convention of summing over repeated upper and lower indices. Let $\xi:[0,1] \times \mathscr{W}\left(\mathbb{R}^{k}\right) \rightarrow G$ denote the solution to the Stratonovich stochastic differential equation

$$
\begin{equation*}
d \xi_{t}=\xi_{t} \circ d \vec{b}_{t}:=L_{\xi_{t^{*}}} \circ d \vec{b}_{t}=L_{\xi_{t^{*}}} X_{i} \circ d b_{t}^{i}=\tilde{X}_{i}\left(\xi_{t}\right) \circ d b_{t}^{i}, \text { with } \xi_{0}=e . \tag{3.3}
\end{equation*}
$$

The solution $\xi$ exists by the standard theory; see, for example, Theorem V-1.1 of 30]. Additionally, Remark V-10.3 of [30] implies that $P_{t}=e^{t L / 2}$, with $L=\sum_{i=1}^{k} \tilde{X}_{i}^{2}$, is the associated Markov diffusion semigroup to $\xi$, where $P_{t}$ is as defined in Definition 1.8; that is, $\nu_{t}:=\left(\xi_{t}\right)_{*} \mu=p_{t}(g) d g$ is the density of the transition probability of the diffusion process $\xi_{t}$, where $d g$ here denotes right Haar measure, and

$$
\begin{equation*}
\left(P_{t} f\right)(0)=\mathbb{E}\left[f\left(\xi_{t}\right)\right], \tag{3.4}
\end{equation*}
$$

for any $f \in C_{c}^{\infty}(G)$, where the right hand side is expectation conditioned on $\xi_{0}=e$.

Theorem 3.2. For any $f \in C_{c}^{\infty}(G), f\left(\xi_{t}\right) \in \mathcal{D}^{\infty}$ for all $t \in[0,1]$. In particular, $D\left[f\left(\xi_{t}\right)\right] \in \mathscr{H}\left(\mathbb{R}^{k}\right) \otimes \mathbb{R}^{k}$ and

$$
\begin{equation*}
\left(D\left[f\left(\xi_{t}\right)\right]\right)^{i}=\left\langle\hat{\nabla} f\left(\xi_{t}\right), \int_{0}^{\cdot \wedge t} \operatorname{Ad}_{\xi_{\tau}} X_{i} d \tau\right\rangle, \tag{3.5}
\end{equation*}
$$

for $i=1, \ldots, k$, componentwise in $\mathscr{H}\left(\mathbb{R}^{k}\right)$.
Notation 3.3. In the sequel, $\mathscr{W}=\mathscr{W}\left(\mathbb{R}^{k}\right)$ and $\mathscr{H}=\mathscr{H}\left(\mathbb{R}^{k}\right)$. For $h=\left(h^{1}, \ldots, h^{k}\right) \in \mathscr{H}$, let $\vec{h}:=X_{i} h^{i} \in \mathfrak{g}_{0}$. We also let

$$
\bigcap_{p>1} L^{p}(\mu)=: L^{\infty-}(\mu)
$$

Proof. This proof goes through a series of convergence arguments for solutions to cutoff versions of Equation (3.3). Let $|g|$ denote the distance from a point $g \in G$ to $e$ with respect to the right invariant metric. Let $\left\{\varphi_{m}\right\}_{m=1}^{\infty} \subset C_{c}^{\infty}(G,[0,1])$ be a sequence of functions, $\varphi_{m} \uparrow 1$, such that $\varphi_{m}(g)=1$ when $|g| \leq m$ and $\sup _{m} \sup _{g \in G}\left|\tilde{\nabla}^{k} \varphi_{m}(g)\right|<\infty$ for $k=0,1,2, \ldots$; see Lemma 3.6 of Driver and Gross [16]. Let $\psi \in C_{c}^{\infty}(\operatorname{End}(\mathfrak{g}),[0,1])$ such that $\psi=1$ in a neighborhood of $I$, and $\psi(x)=0$ if $|x| \geq 2$, where here $|\cdot|$ is the distance from $I$ with respect to any metric on $\operatorname{End}(\mathfrak{g})$. Define $v(g):=\psi\left(\operatorname{Ad}_{g}\right)$ and $u_{m}(g):=\varphi_{m}(g) v(g)$. Let $\eta^{m}:[0,1] \times \mathscr{W} \rightarrow G$ denote the solution to the equation

$$
d \eta^{m}=u_{m}\left(\eta^{m}\right) \eta^{m} \circ d \vec{b}=u_{m}\left(\eta^{m}\right) \tilde{X}_{i}\left(\eta^{m}\right) \circ d b^{i}
$$

with $\eta_{0}^{m}=e$, and let $\eta:[0,1] \times \mathscr{W} \rightarrow G$ denote the solution to

$$
d \eta=v(\eta) \eta \circ d \vec{b}=v(\eta) \tilde{X}_{i}(\eta) \circ d b^{i}
$$

with $\eta_{0}=e$. By Lemma 3.19, these solutions exist for all time $t$, and by Proposition 3.23. for any $f \in C_{c}^{\infty}(G)$,

$$
\begin{equation*}
\lim _{m \rightarrow \infty} \mathbb{E} \sup _{\tau \leq 1}\left|f\left(\eta_{\tau}^{m}\right)-f\left(\eta_{\tau}\right)\right|^{p}=0 \tag{3.6}
\end{equation*}
$$

for all $p \in(1, \infty)$.

Since the vector fields $u_{m} \tilde{X}_{i}$ have compact support, we may embed $G$ as a Euclidean submanifold in such a way that the embedded vector fields will be bounded with bounded derivatives. Thus, by Theorem 2.1 of Taniguchi [47, for $f \in C_{c}^{\infty}(G)$ and $t \in[0,1], f\left(\eta_{t}^{m}\right) \in \mathcal{D}^{\infty}$. Now let $\theta^{m}:[0,1] \times \mathscr{W} \rightarrow \mathfrak{g}$ be the solution to the equation

$$
d \theta^{m}=\left\langle\hat{\nabla} u_{m}\left(\eta^{m}\right), \theta^{m}\right\rangle \operatorname{Ad}_{\eta^{m}} \circ d \vec{b}+u_{m}\left(\eta^{m}\right) \operatorname{Ad}_{\eta^{m}} d \vec{h}
$$

with $\theta_{0}^{m}=0$, and $\theta_{t}: \mathscr{W} \rightarrow \mathfrak{g}$ be the solution to

$$
d \theta=\langle\hat{\nabla} v(\eta), \theta\rangle \operatorname{Ad}_{\eta} \circ d \vec{b}+v(\eta) \operatorname{Ad}_{\eta} d \vec{h},
$$

with $\theta_{0}=0$. In Proposition 3.25, we show that

$$
\lim _{m \rightarrow \infty} \mathbb{E} \sup _{\tau \leq 1}\left|\theta_{\tau}^{m}-\theta_{\tau}\right|^{p}=0
$$

for all $p \in(1, \infty)$, and Proposition 3.20 implies that, for any $h \in \mathscr{H}$,

$$
\partial_{h}\left[f\left(\eta_{t}^{m}\right)\right]=\left\langle\hat{\nabla} f\left(\eta_{t}^{m}\right), \theta_{t}^{m}\right\rangle .
$$

It is then easily shown in Proposition 3.26 that these two facts give

$$
\begin{equation*}
\lim _{m \rightarrow \infty} \mathbb{E} \sup _{\tau \leq t}\left|\partial_{h}\left[f\left(\eta_{\tau}^{m}\right)\right]-\left\langle\hat{\nabla} f\left(\eta_{\tau}\right), \theta_{\tau}\right\rangle\right|^{p}=0 \tag{3.7}
\end{equation*}
$$

for all $p \in(1, \infty)$. Since $\partial_{h}$ is a closable operator, Equations (3.6) and (3.7) imply that $f\left(\eta_{t}\right) \in \operatorname{Dom}\left(\partial_{h}\right)$, and $\partial_{h}\left[f\left(\eta_{t}\right)\right]=\left\langle\hat{\nabla} f\left(\eta_{t}\right), \theta_{t}\right\rangle \in L^{\infty-}(\mu)$. It then follows from Theorem 2.12 that $f\left(\eta_{t}\right) \in \mathcal{D}^{1, \infty}$.

Now set $\psi_{n}(x):=\psi\left(n^{-1} x\right)$ and $v_{n}(g):=\psi_{n}\left(\operatorname{Ad}_{g}\right)=\psi\left(n^{-1} \operatorname{Ad}_{g}\right)$. Let $\xi^{n}:$ $[0,1] \times \mathscr{W} \rightarrow G$ denote the solution to

$$
d \xi^{n}=v_{n}\left(\xi^{n}\right) \xi^{n} \circ d \vec{b}=v_{n}\left(\xi^{n}\right) \tilde{X}_{i}\left(\xi^{n}\right) \circ d b^{i},
$$

with $\xi_{0}^{n}=e$, and let $\xi$ denote the solution to Equation (3.3). By the previous argument, for any $f \in C_{c}^{\infty}(G)$ and $t \in[0,1], f\left(\xi_{t}^{n}\right) \in \mathcal{D}^{1, \infty}$, and

$$
\partial_{h}\left[f\left(\xi_{t}^{n}\right)\right]=\left\langle\hat{\nabla} f\left(\xi_{t}^{n}\right), \Theta_{t}^{n}\right\rangle,
$$

where $\Theta_{t}^{n}: \mathscr{W} \rightarrow \mathfrak{g}$ satisfies the equation

$$
d \Theta^{n}=\left\langle\hat{\nabla} v_{n}\left(\xi^{n}\right), \Theta^{n}\right\rangle \operatorname{Ad}_{\xi^{n}} \circ d \vec{b}+v_{n}\left(\xi^{n}\right) \operatorname{Ad}_{\xi^{n}} d \vec{h},
$$

with $\Theta_{0}^{n}=0$. In Theorem 3.28, we show that

$$
\lim _{n \rightarrow \infty} \mathbb{E} \sup _{\tau \leq 1}\left|f\left(\xi_{\tau}^{n}\right)-f\left(\xi_{\tau}\right)\right|^{p}=0
$$

and

$$
\lim _{n \rightarrow \infty} \mathbb{E} \sup _{\tau \leq 1}\left|\Theta_{\tau}^{n}-R_{\xi_{\tau^{*}}} \int_{0}^{\tau} \operatorname{Ad}_{\xi} X_{i} \dot{h}^{i} d s\right|^{p}=0
$$

for all $p \in(1, \infty)$. Thus, for any $h \in \mathscr{H}$,

$$
\lim _{n \rightarrow \infty} \mathbb{E} \sup _{\tau \leq 1}\left|\partial_{h} f\left(\xi_{\tau}^{n}\right)-\left\langle\hat{\nabla} f\left(\xi_{\tau}\right), \int_{0}^{\tau} \operatorname{Ad}_{\xi} X_{i} \dot{h}^{i} d s\right\rangle\right|^{p}=0
$$

So $f\left(\xi_{t}\right) \in \operatorname{Dom}\left(\partial_{h}\right)$, and

$$
\begin{equation*}
\partial_{h} f\left(\xi_{t}\right)=\left\langle\hat{\nabla} f\left(\xi_{t}\right), \int_{0}^{t} \operatorname{Ad}_{\xi} X_{i} \dot{h}^{i} d s\right\rangle \in L^{\infty-}(\mu) \tag{3.8}
\end{equation*}
$$

From this equality, we may then show that $\left\|f\left(\xi_{t}\right)\right\|_{G^{1, p}}<\infty$, for all $p \in(1, \infty)$. (The norm $\|\cdot\|_{G^{k, p}}$ was defined in Section 2.5.) It then follows from Theorem 2.12 that $f\left(\xi_{t}\right) \in \mathcal{D}^{1, \infty}$, and

$$
D_{s}\left[f\left(\xi_{t}\right)\right]^{i}=\left\langle\hat{\nabla} f\left(\xi_{t}\right), \int_{0}^{s \wedge t} \operatorname{Ad}_{\xi_{\tau}} X_{i} d \tau\right\rangle
$$

componentwise in $\mathscr{H}$.
Finally, by Proposition 3.15, $\bar{W}=\int_{0} \operatorname{Ad}_{\xi_{\tau}} d \tau \in \mathcal{D}^{\infty}(\mathscr{H}(\operatorname{End}(\mathfrak{g})))$. Since $\mathcal{D}^{\infty}$ is an algebra, this and Equation (3.5) are sufficient to show that $f\left(\xi_{t}\right) \in \mathcal{D}^{\infty}$, for all $f \in C_{c}^{\infty}(G)$ and $t \in[0,1]$.

The supporting propositions and theorems cited in the above proof may be found in Section 3.2.

### 3.1.2 Covariance matrix of $\xi_{t}$

The Malliavin covariance matrix of $\xi$ is the matrix $\sigma_{t}(\omega):=\xi_{t}^{\prime}(\omega) \xi_{t}^{\prime}(\omega)^{*}$ : $T_{\xi_{t}(\omega)} G \rightarrow T_{\xi_{t}(\omega)} G$, where $\xi_{t}^{\prime}(\omega): \mathscr{H} \rightarrow T_{\xi_{t}(\omega)} G$ is the Frechet derivative given by

$$
\xi_{t}^{\prime}(\omega) h:=\left.\frac{d}{d \epsilon}\right|_{0} \xi_{t}(\omega+\epsilon h),
$$

for all $h \in \mathscr{H}$, and its adjoint $\xi_{t}^{\prime}(\omega)^{*}: T_{\xi_{t}(\omega)} \rightarrow \mathscr{H}$ is computed relative to the CameronMartin inner product on $\mathscr{H}$ and the chosen metric on $G$.

Notation 3.4. In the following, let $\operatorname{Ad}_{\xi_{t}}^{\dagger}$ denote the adjoint of $\operatorname{Ad}_{\xi_{t}}$ as an operator on $\mathfrak{g}$, and let $P: \mathfrak{g} \rightarrow \mathfrak{g}_{0}$ be orthogonal projection onto the hypoelliptic subspace $\mathfrak{g}_{0}$.

Theorem 3.5. The Malliavin covariance matrix of $\xi$ is

$$
\begin{equation*}
\sigma_{t}:=\xi_{t}^{\prime}(\omega) \xi_{t}^{\prime}(\omega)^{*}=R_{\xi_{t} *}\left(\int_{0}^{t} \operatorname{Ad}_{\xi_{s}} P \operatorname{Ad}_{\xi_{s}}^{\dagger} d s\right) R_{\xi_{t} *}^{\mathrm{tr}}, \tag{3.9}
\end{equation*}
$$

Let $\bar{\sigma}_{t}=\int_{0}^{t} \operatorname{Ad}_{\xi_{s}} P \operatorname{Ad}_{\xi_{s}}^{\dagger} d s$, and $\Delta_{t}:=\operatorname{det} \bar{\sigma}_{t}$. Then $\Delta_{t}>0$ a.e., and so $\bar{\sigma}_{t}$ is invertible a.e. for $t>0$. Moreover,

$$
\Delta_{t}^{-1} \in L^{\infty-}(\mu)
$$

Proof. We now compute $\xi_{t}^{\prime}(\omega)^{*}: T_{\xi_{t}(\omega)} G \rightarrow \mathscr{H}$, the adjoint in $\xi_{t}^{\prime}(\omega)$ with respect to the Cameron-Martin inner product and the right invariant metric on $T G$. By Equation (3.8), for any $X \in \mathfrak{g}$,

$$
\begin{aligned}
\left(\xi_{t}^{\prime}(\omega)^{*}\left(R_{\xi_{t} *} X\right), h\right)_{\mathscr{H}} & =\left\langle R_{\xi_{t^{*}}} X, \xi_{t}^{\prime}(\omega) h\right\rangle \\
& =\left\langle R_{\xi_{t} *} X, R_{\xi_{t^{*}}} \int_{0}^{t} \operatorname{Ad}_{\xi_{s}} \dot{h}_{s} d s\right\rangle \\
& =\left\langle X, \int_{0}^{t} \operatorname{Ad}_{\xi_{s}} \dot{h}_{s} d s\right\rangle=\int_{0}^{t}\left\langle\operatorname{Ad}_{\xi_{s}}^{\dagger} X, X_{i}\right\rangle \dot{h}_{s}^{i} d s,
\end{aligned}
$$

where the penultimate equality follows from the right invariance of the metric on $G$. It then follows that

$$
\begin{equation*}
\dot{k}_{s}^{i}=\frac{d}{d s}\left[\xi_{t}^{\prime}(\omega)^{*}\left(R_{\xi_{t} *} X\right)\right]_{s}^{i}=1_{s \leq t}\left\langle\operatorname{Ad}_{\xi_{s}}^{\dagger} X, X_{i}\right\rangle, \tag{3.10}
\end{equation*}
$$

as an element $\left(k^{1}, \ldots, k^{k}\right)$ of the Cameron-Martin space $\mathscr{H}$. Combining Equations (3.8) and (3.10) for $k:=X_{i} k^{i}$, we have

$$
\begin{aligned}
\xi_{t}^{\prime}(\omega) \xi_{t}^{\prime}(\omega)^{*}\left(R_{\xi_{t} *} X\right) & =\sum_{i=1}^{k} R_{\xi_{t} *} \int_{0}^{t} \operatorname{Ad}_{\xi_{s}} X\left\langle\operatorname{Ad}_{\xi_{s}}^{\dagger} X, X_{i}\right\rangle d s \\
& =R_{\xi_{t^{*}}} \int_{0}^{t} \operatorname{Ad}_{\xi_{s}} P \operatorname{Ad}_{\xi_{s}}^{\dagger} X d s,
\end{aligned}
$$

and Equation (3.9) follows from the above.
The proof that $\Delta_{t}>0$ and $\Delta_{t}^{-1} \in L^{\infty-}(\mu)$ is by now standard and relies on satisfaction of the Hörmander bracket condition, $\operatorname{Lie}\left(\left\{\tilde{X}_{i}\right\}_{i=1}^{k}\right)=\mathfrak{g}$; see, for example, the proof of Theorem 8.6 in Driver [15]. For completeness, we have included the adapted proof in the appendix.

Remark 3.6. By the general theory, Theorem 3.5 implies $\nu_{t}=\operatorname{Law}\left(\xi_{t}\right)$ is a smooth measure; see for example Remark V-10.3 of [30].

### 3.1.3 Lifted vector fields and their $L^{2}$-adjoints

Throughout this section, $t \in[0,1]$ will be fixed.
Definition 3.7. Given $X \in \mathfrak{g}$, let $\tilde{X}$ be the associated left invariant vector field on $G$. Define the "lifted vector field" $\mathbf{X}$ of $\tilde{X}$ as

$$
\begin{equation*}
\mathbf{X}=\mathbf{X}^{t}:=\xi_{t}^{\prime}(\omega)^{*}\left[\xi_{t}^{\prime}(\omega) \xi_{t}^{\prime}(\omega)^{*}\right]^{-1} \tilde{X}\left(\xi_{t}\right)=\xi_{t}^{\prime}(\omega)^{*} \sigma_{t}^{-1} \tilde{X}\left(\xi_{t}\right) \in \mathscr{H}, \tag{3.11}
\end{equation*}
$$

acting on functions $F \in \mathcal{D}^{1,2}$ by

$$
\mathbf{X} F=(D F, \mathbf{X})_{\mathscr{H}}
$$

Proposition 3.8. For any $X \in \mathfrak{g}, \mathbf{X} \in \mathcal{D}^{\infty}(\mathscr{H})$, and

$$
\mathbf{X}\left[f\left(\xi_{t}\right)\right]=(\tilde{X} f)\left(\xi_{t}\right),
$$

for any $f \in C^{\infty}(G)$,

Proof. Combining Equations (3.9) and (3.10), we have that

$$
\frac{d}{d s} \mathbf{X}_{s}^{i}=1_{s \leq t}\left\langle\operatorname{Ad}_{\xi_{s}}^{\dagger}\left(\int_{0}^{t} \operatorname{Ad}_{\xi_{r}} P \operatorname{Ad}_{\xi_{r}}^{\dagger} d r\right)^{-1} \operatorname{Ad}_{\xi_{t}} X, X_{i}\right\rangle
$$

Thus, we may rewrite Equation (3.11) explicitly as

$$
\begin{align*}
\mathbf{X}^{i} & =\int_{0}^{\cdot \wedge t}\left\langle\operatorname{Ad}_{\xi_{s}}^{\dagger}\left(\int_{0}^{t} \operatorname{Ad}_{\xi_{r}} P \operatorname{Ad}_{\xi_{r}}^{\dagger} d r\right)^{-1} \operatorname{Ad}_{\xi_{t}} X, X_{i}\right\rangle d s \\
& =\left\langle\left(\int_{0}^{\wedge} \operatorname{Ad}_{\xi_{s}}^{\dagger} d s\right)\left(\int_{0}^{t} \operatorname{Ad}_{\xi_{s}} P \operatorname{Ad}_{\xi_{s}}^{\dagger} d s\right)^{-1} \operatorname{Ad}_{\xi_{t}} X, X_{i}\right\rangle \tag{3.12}
\end{align*}
$$

By Proposition 3.15. $\bar{W}=\int_{0}^{\sim} \operatorname{Ad}_{\xi_{\tau}} d \tau \in \mathcal{D}^{\infty}(\mathscr{H}(\operatorname{End}(\mathfrak{g})))$. Note that $W_{t}^{\dagger}=\operatorname{Ad}_{\xi_{t}}^{\dagger}: \mathscr{W} \rightarrow$ $\operatorname{End}(\mathfrak{g})$ satisfies the differential equation

$$
d W_{t}^{\dagger}=\operatorname{ad}_{X_{i}}^{\dagger} W_{t}^{\dagger} \circ d b_{t}^{i}, \text { with } W_{0}^{\dagger}=I
$$

which is linear with smooth coefficients. Then by a similar argument to that in Proposition 3.15, we may show that

$$
\bar{W}^{\dagger}:=\int_{0} \operatorname{Ad}_{\xi_{\tau}}^{\dagger} d \tau \in \mathcal{D}^{\infty}(\mathscr{H}(\operatorname{End}(\mathfrak{g}))),
$$

for all $t \in[0,1]$. Also, Theorem 3.5 implies that

$$
\bar{\sigma}_{t}^{-1}=\left(\int_{0}^{t} \operatorname{Ad}_{\xi_{r}} P \operatorname{Ad}_{\xi_{r}}^{\dagger} d r\right)^{-1}
$$

exists and is in $L^{\infty-}(\mu)$ componentwise. Thus the Equation 3.12 implies that $\mathbf{X} \in$ $\mathcal{D}^{\infty}(\mathscr{H})$, since $\mathcal{D}^{\infty}(\mathscr{H})$ is an algebra.

For $f \in C^{\infty}(G)$ and $\left(h^{1}, \ldots, h^{k}\right) \in \mathscr{H}$, by Equation (3.53),

$$
\partial_{h}\left[f\left(\xi_{t}\right)\right]=\left(D\left[f\left(\xi_{t}\right)\right], h\right)_{\mathscr{H}}=\left\langle\hat{\nabla} f\left(\xi_{t}\right), \int_{0}^{t} \operatorname{Ad}_{\xi_{s}} X_{i} \dot{h}_{s}^{i} d s\right\rangle,
$$

and so

$$
\begin{aligned}
\mathbf{X}\left[f\left(\xi_{t}\right)\right] & =\left(D\left[f\left(\xi_{t}\right)\right], \mathbf{X}\right)_{\mathscr{H}} \\
& =\left\langle\hat{\nabla} f\left(\xi_{t}\right), \int_{0}^{t} \operatorname{Ad}_{\xi_{s}} X_{i}\left\langle\operatorname{Ad}_{\xi_{s}}^{\dagger}\left(\int_{0}^{t} \operatorname{Ad}_{\xi_{r}} P \operatorname{Ad}_{\xi_{r}}^{\dagger} d r\right)^{-1} \operatorname{Ad}_{\xi_{t}} X, X_{i}\right\rangle d s\right\rangle \\
& =\left\langle\hat{\nabla} f\left(\xi_{t}\right), \int_{0}^{t} \operatorname{Ad}_{\xi_{s}} P \operatorname{Ad}_{\xi_{s}}^{\dagger}\left(\int_{0}^{t} \operatorname{Ad}_{\xi_{r}} P \operatorname{Ad}_{\xi_{r}}^{\dagger} d r\right)^{-1} \operatorname{Ad}_{\xi_{t}} X d s\right\rangle \\
& =\left\langle\hat{\nabla} f\left(\xi_{t}\right), \operatorname{Ad}_{\xi_{t}} X\right\rangle=\left\langle\tilde{\nabla} f\left(\xi_{t}\right), \operatorname{Ad}_{\xi_{t}^{-1}} \operatorname{Ad}_{\xi_{t}} X\right\rangle \\
& =(\tilde{X} f)\left(\xi_{t}\right),
\end{aligned}
$$

where we have used Equation (3.1) in the penultimate equality.
Definition 3.9. For a vector field $\mathbf{X}$ acting on functions of $\mathscr{W}$, we will denote the adjoint of $\mathbf{X}$ in the $L^{2}(\mu)$ inner product by $\mathbf{X}^{*}$, which has domain in $L^{2}(\mu)$ consisting of functions $G$ such that for all $F \in \mathcal{D}^{1,2}$,

$$
\mathbb{E}[(\mathbf{X} F) G] \leq c\|F\|_{L^{2}(\mu)}
$$

for some constant $c$. For functions $G$ in the domain of $\mathbf{X}^{*}$,

$$
\begin{equation*}
\mathbb{E}\left[F\left(\mathbf{X}^{*} G\right)\right]=\mathbb{E}[(\mathbf{X} F) G] \tag{3.13}
\end{equation*}
$$

for all $F \in \mathcal{D}^{1,2}$.
Note that for any lifted vector field $\mathbf{X}$ acting on function $F \in \mathcal{D}^{1,2}$ as defined in Definition 3.7.

$$
\mathbb{E}[\mathbf{X} F]=\mathbb{E}\left[(D F, \mathbf{X})_{\mathscr{H}}\right]=\mathbb{E}\left[F D^{*} \mathbf{X}\right] .
$$

Thus, we must have that $\mathbf{X}^{*}=\mathbf{X}^{*} 1=D^{*} \mathbf{X}$ a.s. Recall that $D^{*}$ is a continuous operator from $\mathcal{D}^{\infty}(\mathscr{H})$ into $\mathcal{D}^{\infty}$; see for example Theorem V-8.1 and its corollary in 30]. Thus, for $\mathbf{X}$ a vector field on $W$ as defined in Equation (3.11), Proposition 3.8 implies that $D^{*} \mathbf{X} \in \mathcal{D}^{\infty}$. We have then proven the following proposition.

Proposition 3.10. Let $\tilde{X}$ be a left invariant vector field on $G$. Then for the vector field on $\mathscr{W}$ defined by

$$
\mathbf{X}=\xi_{t}^{\prime}(\omega)^{*}\left[\xi_{t}^{\prime}(\omega) \xi_{t}^{\prime}(\omega)^{*}\right]^{-1} \tilde{X}\left(\xi_{t}(\omega)\right),
$$

$\mathbf{X}^{*} \in \mathcal{D}^{\infty}$, where $\mathbf{X}^{*}$ is the $L^{2}(\mu)$-adjoint of $\mathbf{X}$.
This completes the main results of this chapter, and if desired the reader may now continue to Chapter 4 and skip the next technical section which gives details of the proof of Theorem 3.2.

### 3.2 Results needed in the proof of Theorem 3.2

To prove Theorem 3.2, we will repeatedly use the following standard proposition; see for example Driver [15].

Proposition 3.11. Suppose $p \in[2, \infty)$, $\alpha_{t}$ is a predictable $\mathbb{R}^{d}$-valued process, $A_{t}$ is a predictable $\operatorname{Hom}\left(\mathfrak{g}_{0}, \mathbb{R}^{d}\right)$-valued process, and

$$
\begin{equation*}
Y_{t}:=\int_{0}^{t} A_{\tau} d \vec{b}_{\tau}+\int_{0}^{t} \alpha_{\tau} d \tau=\int_{0}^{t} A_{\tau} X_{i} d b_{\tau}^{i}+\int_{0}^{t} \alpha_{\tau} d \tau \tag{3.14}
\end{equation*}
$$

where $\left\{b^{1}, \ldots, b^{k}\right\}$ are $k$ independent real Brownian motions. Then

$$
\begin{equation*}
\mathbb{E} \sup _{\tau \leq t}\left|Y_{\tau}\right|^{p} \leq C_{p}\left\{\mathbb{E}\left(\int_{0}^{t}\left|A_{\tau}\right|^{2} d \tau\right)^{p / 2}+\mathbb{E}\left(\int_{0}^{t}\left|\alpha_{\tau}\right| d \tau\right)^{p}\right\}, \tag{3.15}
\end{equation*}
$$

where

$$
|A|^{2}=\operatorname{tr}\left(A A^{*}\right)=\sum_{i=1}^{n}\left(A A^{*}\right)_{i i}=\sum_{i, j} A_{i j} A_{i j}=\operatorname{tr}\left(A^{*} A\right) .
$$

Notation 3.12. Here and in the sequel, we will let $\delta_{n}$ denote constants such that $\lim _{n \rightarrow \infty} \delta_{n}=0$. Also, we will write $f \lesssim g$, if there is a positive constant $C$ so that $f \leq C g$.

This section is divided into two parts. Section 3.2 .1 gives convergence results in a matrix group setting which are necessary to resolve certain convergence issues on the Lie group. These Lie group issues are addressed in Section 3.2.2.

### 3.2.1 Matrix group computations

In this section, let $M=\operatorname{End}(\mathfrak{g})$ and $A_{i}=\operatorname{ad}_{X_{i}} \in M$, for $i=1, \ldots, k$. Further, let $B:=A_{i} b^{i}=\operatorname{ad}_{X_{i}} b^{i}=\operatorname{ad}_{\vec{b}}$, where $\left\{b^{i}\right\}_{i=1}^{k}$ is a set of independent real-valued Brownian motions, and let $H:=A_{i} h^{i}=\operatorname{ad}_{X_{i}} h^{i}$, where $h=\left(h^{1}, \ldots, h^{k}\right)$ is a fixed element of $\mathscr{H}\left(\mathbb{R}^{k}\right)$.

Proposition 3.13. Let $W:[0,1] \times \mathscr{W} \rightarrow M$ denote the solution to the Stratonovich differential equations

$$
\begin{equation*}
d W=W \circ d B=W d B+\frac{1}{2} \sum_{i=1}^{k} W A_{i}^{2} d t, \text { with } W_{0}=I, \tag{3.16}
\end{equation*}
$$

and let $W^{s}:[0,1]^{2} \times \mathscr{W} \rightarrow M$ denote the solution to the equation

$$
\begin{equation*}
d W^{s}=W^{s}(\circ d B+s d H), \text { with } W_{0}^{s}=I \tag{3.17}
\end{equation*}
$$

Then

$$
\lim _{s \downarrow 0} \mathbb{E} \sup _{\tau \leq 1}\left|W_{\tau}^{s}-W_{\tau}\right|^{p}=0,
$$

for all $p \in(1, \infty)$.
Proof. Writing Equation (3.17) in Itô form gives

$$
d W^{s}=W^{s} d B+\frac{1}{2} \sum_{i=1}^{k} W^{s} A_{i}^{2} d t+s W^{s} d H
$$

Then by Proposition 3.11, for any $s \in[0,1]$,

$$
\begin{aligned}
\mathbb{E} \sup _{\tau \leq 1}\left|W_{\tau}^{s}\right|^{p} & \lesssim 1+\mathbb{E}\left(\int_{0}^{t}\left|W_{\tau}^{s}\right|^{2} d \tau\right)^{p / 2}+\mathbb{E}\left(\int_{0}^{t}\left|W_{\tau}^{s}\right| d \tau\right)^{p} \\
& \lesssim 1+\mathbb{E} \int_{0}^{t}\left|W_{\tau}^{s}\right|^{p} d \tau
\end{aligned}
$$

for all $t \in[0,1]$. An application of Gronwall's inequality then shows that

$$
\mathbb{E} \sup _{\tau \leq 1}\left|W_{\tau}^{s}\right|^{p} \leq C e^{C},
$$

where these constants are independent of $s$; that is, there exists some finite constant $C_{p}$ such that

$$
\sup _{s \in[0,1]}\left(\mathbb{E} \sup _{\tau \leq 1}\left|W_{\tau}^{s}\right|^{p}\right) \leq C_{p}
$$

for all $p \in(1, \infty)$.
Now let $\varepsilon^{s}:=W^{s}-W$. Then from Equation (3.16), $\varepsilon^{s}$ satisfies

$$
\begin{aligned}
d \varepsilon^{s} & =\left(W^{s}-W\right) d B+\frac{1}{2} \sum_{i=1}^{k}\left(W^{s}-W\right) A_{i}^{2} d t+s W^{s} d H \\
& =\varepsilon^{s} d B+\frac{1}{2} \sum_{i=1}^{k} \varepsilon^{s} A_{i}^{2} d t+s W^{s} d H .
\end{aligned}
$$

Then applying Proposition 3.11, we have

$$
\mathbb{E} \sup _{\tau \leq 1}\left|\varepsilon_{\tau}^{s}\right|^{p} \lesssim \int_{0}^{t}\left|\varepsilon^{s}\right|^{p} d \tau+\delta_{s}
$$

for all $t \in[0,1]$, where

$$
\delta_{s}=\sum_{i=1}^{k} \int_{0}^{t} s^{p}\left|W^{s} \dot{h}^{i}\right|^{p} d \tau \rightarrow 0,
$$

as $s \downarrow 0$, by the dominated convergence theorem. Thus, by Gronwall's inequality,

$$
\mathbb{E} \sup _{\tau \leq 1}\left|W_{\tau}^{s}-W_{\tau}\right|^{p}=\mathbb{E} \sup _{\tau \leq 1}\left|\varepsilon_{\tau}^{s}\right|^{p} \leq \delta_{s} e^{C} \rightarrow 0,
$$

as $s \downarrow 0$, for all $p \in(1, \infty)$.
Proposition 3.14. Let $W$ be the solution to Equation (3.16) and $W^{s}$ be the solution to Equation 3.17). Then $W_{t} \in \mathcal{D}^{\infty}(\operatorname{End}(\mathfrak{g}))$, for all $t \in[0,1]$, and $\partial_{h} W:[0,1] \times \mathscr{W} \rightarrow M$ solves the equation

$$
\begin{equation*}
\partial_{h} W_{t}=\left(\int_{0}^{t} W_{\tau} \dot{H} W_{\tau}^{-1} d \tau\right) W_{t} \tag{3.18}
\end{equation*}
$$

Furthermore,

$$
\lim _{s \downarrow 0} \mathbb{E} \sup _{\tau \leq 1}\left|\frac{W_{\tau}^{s}-W_{\tau}}{s}-\partial_{h} W_{\tau}\right|^{p}=0
$$

for all $p \in(1, \infty)$.

Proof. Note that $W=\operatorname{Ad}_{\xi}$ satisfies the Stratonovich stochastic differential equation

$$
\begin{equation*}
d \operatorname{Ad}_{\xi}=\operatorname{Ad}_{\xi} \circ \operatorname{ad}_{d b}=\operatorname{Ad}_{\xi} \operatorname{ad}_{X_{i}} \circ d b^{i} \tag{3.19}
\end{equation*}
$$

a linear differential equation with smooth coefficients. Then by Theorem V-10.1 of Ikeda and Watanabe [30, $W_{t}=\operatorname{Ad}_{\xi_{t}} \in \mathcal{D}^{\infty}(\operatorname{End}(\mathfrak{g}))$ componentwise with respect to some basis.

Now let $\Psi:[0,1] \times \mathscr{W} \rightarrow M$ denote the solution to the equation

$$
d \Psi=\Psi W \circ d B+W d H=\Psi W d B+\frac{1}{2} \sum_{i=1}^{k} \Psi A_{i}^{2} d t+W d H
$$

with $\Psi_{0}=0$. Let $\varepsilon^{s}:=\left(\frac{W^{s}-W}{s}-\Psi\right)$. Then we have that

$$
\begin{aligned}
d \varepsilon^{s} & =\left(\frac{W^{s}-W}{s}-\Psi\right) d B+\frac{1}{2} \sum_{i=1}^{k}\left(\frac{W^{s}-W}{s}-\Psi\right) A_{i}^{2} d t+\left(W^{s}-W\right) d H \\
& =\varepsilon^{s} d B+\frac{1}{2} \sum_{i=1}^{k} \varepsilon^{s} A_{i}^{2} d t+\left(W^{s}-W\right) d H
\end{aligned}
$$

By Proposition 3.11, this implies that

$$
\mathbb{E} \sup _{\tau \leq t}\left|\varepsilon_{\tau}^{s}\right|^{p} \lesssim \int_{0}^{t}\left|\varepsilon^{s}\right|^{p} d \tau+\delta_{s}
$$

for all $t \in[0,1]$, where

$$
\delta_{s}=\sum_{i=1}^{k} \int_{0}^{t}\left|\left(W^{s}-W\right) \dot{h}^{i}\right|^{p} d \tau \rightarrow 0
$$

as $s \downarrow 0$, by Proposition 3.13 and the dominated convergence theorem. An application of Gronwall's inequality then gives

$$
\mathbb{E} \sup _{\tau \leq 1}\left|\frac{W_{\tau}^{s}-W_{\tau}}{s}-\Psi_{\tau}\right|^{p}=\mathbb{E} \sup _{\tau \leq 1}\left|\varepsilon_{\tau}^{s}\right|^{p} \leq \delta_{s} e^{C} \rightarrow 0
$$

as $s \downarrow 0$, for all $p \in(1, \infty)$.

By Theorem VIII.2B of Elworthy [17], there exists a modification of $W_{t}^{s}$ so that the mapping $s \mapsto W_{t}^{s}$ is smooth. Let $F \in \mathcal{S}$ be a smooth cylinder function on $\mathscr{W}$. By the above convergence, we have that

$$
\left.\frac{d}{d s}\right|_{0} \mathbb{E}\left[W_{t}^{s} F\right]=\mathbb{E}\left[\Psi_{t} F\right]
$$

Consider also

$$
\begin{aligned}
\left.\frac{d}{d s}\right|_{0} \mathbb{E}\left[W_{t}^{s} F\right] & =\left.\int \frac{d}{d s}\right|_{0} W_{t}(b+s h) F(b) d \mu(b) \\
& =\left.\int W_{t}(b) \frac{d}{d s}\right|_{0} F(b-s h) d \mu(b-s h) \\
& =\int W_{t}(b)\left[-\partial_{h} F(b)+\left(\int_{0}^{1} \dot{h}_{s} \cdot d b_{s}\right) F(b)\right] d \mu(b)=\mathbb{E}\left[W_{t} \partial_{h}^{*} F\right] ;
\end{aligned}
$$

where the third equality follows from differentiating the shifted measure (see for example Theorem 8.1.1 of Hsu [27]), and the final equality follows from Proposition 2.10. This then implies that $\mathbb{E}\left[\Psi_{t} F\right]=\mathbb{E}\left[W_{t} \partial_{h}^{*} F\right]$, and so $\partial_{h} W_{t}=\Psi_{t}$. Thus, $\partial_{h} W$ satisfies the differential equation

$$
d\left(\partial_{h} W\right)=\partial_{h} W \circ d B+W d H,
$$

with $\partial_{h} W_{0}=0$. Equation (3.18) then follows from an application of Duhamel's principle.

Proposition 3.15. Let $W$ be the solution to Equation (3.16). For $\bar{W}:[0,1] \times \mathscr{W} \rightarrow$ $\operatorname{End}(\mathfrak{g})$ given by

$$
\bar{W}_{t}:=\int_{0}^{t} W_{\tau} d \tau,
$$

we have $\bar{W} \in \mathcal{D}^{\infty}(\mathscr{H}(\operatorname{End}(\mathfrak{g})))$.
Proof. Let $V=W^{-1}:[0,1] \times \mathscr{W} \rightarrow \operatorname{End}(\mathfrak{g})$. By differentiating the identity $W_{t} W_{t}^{-1}=I$, one may verify that $V$ satisfies the differential equation

$$
d V=-\circ d B \quad V=-A_{i} V \circ d b^{i}, \text { with } V_{0}=I
$$

Let $V_{t}^{s}:[0,1]^{2} \times \mathscr{W} \rightarrow \operatorname{End}(\mathfrak{g})$ denote the solution to the equation

$$
\begin{equation*}
d V_{t}^{s}=-(\circ d B+s d H) V_{t}^{s}=-A_{i} V_{t}^{s} \circ d b_{t}^{i}-s A_{i} V_{t}^{s} \dot{h}_{t}^{i} d t, \text { with } V_{0}^{s}=I \tag{3.20}
\end{equation*}
$$

By the same arguments as in Propositions 3.13 and 3.14, we may determine that

$$
\begin{gather*}
\lim _{s \downarrow 0} \mathbb{E} \sup _{\tau \leq 1}\left|V_{\tau}^{s}-V_{\tau}\right|^{p}=0  \tag{3.21}\\
\partial_{h} V_{t}=-V_{t} \int_{0}^{t} V_{\tau}^{-1} \dot{H} V_{\tau} d \tau=-V_{t} \int_{0}^{t} W_{\tau} \dot{H} V_{\tau} d \tau \tag{3.22}
\end{gather*}
$$

and

$$
\begin{equation*}
\lim _{s \downarrow 0} \mathbb{E} \sup _{\tau \leq 1}\left|\frac{V_{\tau}^{s}-V_{\tau}}{s}-\partial_{h} V_{\tau}\right|^{p}=0 \tag{3.23}
\end{equation*}
$$

for all $p \in(1, \infty)$. From the proof of Proposition 3.13. we know that there exists a finite constant $C_{p}$ such that

$$
\sup _{s \in[0,1]}\left(\sup _{\tau \leq 1} \mathbb{E}\left|W_{\tau}^{s}\right|^{p}\right) \leq C_{p}
$$

for all $p \in(1, \infty)$. We may similarly show that

$$
\sup _{s \in[0,1]}\left(\sup _{\tau \leq 1} \mathbb{E}\left|V_{\tau}^{s}\right|^{p}\right) \leq C_{p}
$$

for all $p \in(1, \infty)$.
Now let $W_{s, t}:=V_{s} W_{t}$. Note that $W_{s, t}^{-1}=W_{t, s}$, and $W_{s, t} W_{t, u}=W_{s, u}$. By the above bounds on $W$ and $V$, there exist finite constants $C_{p}$ such that

$$
\begin{equation*}
\sup _{\tau_{1}, \tau_{2} \leq t} \mathbb{E}\left|W_{\tau_{1}, \tau_{2}}\right|^{p} \leq C_{p} \tag{3.24}
\end{equation*}
$$

for all $p \in(1, \infty)$. Using this notation and Equation (3.18) we may write

$$
\begin{align*}
\partial_{h} W_{t} & =\left(\int_{0}^{t} W_{\tau} A_{i} W_{\tau}^{-1} \dot{h}_{\tau}^{i} d \tau\right) W_{t}  \tag{3.25}\\
& =\int_{\{0 \leq \tau \leq t\}} W_{\tau} A_{i} W_{\tau, t} \dot{h}_{\tau}^{i} d \tau,
\end{align*}
$$

and so, for $t_{1} \leq t$ and $\alpha_{1}=1 \ldots, k$,

$$
\begin{equation*}
\left[D_{t_{1}} W_{t}\right]^{\alpha_{1}}=\int_{\left\{0 \leq \tau \leq t_{1}\right\}} W_{\tau} A_{\alpha_{1}} W_{\tau, t} d \tau \tag{3.26}
\end{equation*}
$$

where, for $F \in \operatorname{End}(\mathfrak{g}), F=\sum_{i=1}^{k} F^{i} \otimes A_{i}$. For $t_{1}>t, D_{t_{1}} W_{t}=0$.
Let $r \in \mathbb{N}$. For $\left\{t_{1}, \ldots, t_{r}\right\} \subset[0,1]$ and multi-index $\alpha=\left(\alpha_{1}, \ldots, \alpha_{r}\right) \subset$ $\{1, \ldots, k\}^{r}$, let $D_{t_{r}, \ldots, t_{1}}^{r, \alpha} W_{t}$ denote the $\alpha^{t h}$ component of $D_{t_{r}, \ldots, t_{1}}^{r} W_{t}$; that is,

$$
D_{t_{r}, \ldots, t_{1}}^{r, \alpha} W_{t}:=\left[D_{t_{r}}\left[\cdots\left[D_{t_{2}}\left[D_{t_{1}} W_{t}\right]^{\alpha_{1}}\right]^{\alpha_{2}} \cdots\right]^{\alpha_{r}} .\right.
$$

We will now show that the following properties hold for any integer $r \geq 1$ :
(P1) For any $p \in(1, \infty)$ and multi-index $\alpha, W_{t} \in \mathcal{D}^{r, p}(\operatorname{End}(\mathfrak{g}))$ and

$$
\sup _{\left\{t_{1}, \ldots, t_{r}\right\} \in[0,1]} \mathbb{E}\left|D_{t_{r}, \ldots, t_{1}}^{r, \alpha} W_{t}\right|^{p}<\infty .
$$

(P2) Let $T:=\min \left\{t_{1}, \ldots, t_{r}\right\}$. If $T \leq t$, then the $r^{t h}$ derivative of $W_{t}$ satisfies the linear differential equation

$$
\begin{equation*}
D_{t_{r}, \ldots, t_{1}}^{r, \alpha} W_{t}=\int_{[0, T]^{r}} W_{T_{1}} A_{J_{1}} W_{T_{1}, T_{2}} A_{J_{2}} \cdots W_{T_{r-1}, T_{r}} A_{J_{r}} W_{T_{r}, t} d \tau_{r} \cdots d \tau_{1} \tag{3.27}
\end{equation*}
$$

where $T_{i}$ denotes the $i^{\text {th }}$ smallest element of the set $\left\{\tau_{1}, \ldots, \tau_{r}\right\}$, and $J_{i}$ denotes the index corresponding to $T_{i}$ (that is, $J_{i}:=\sum_{l=1}^{r} 1_{\left\{T_{i}=\tau_{l}\right\}} \alpha_{l}$ ). If $T>t$, then $D_{t_{1}, \ldots, t_{r}}^{r} W_{t}=0$.

By Equation (3.26), the above holds for $r=1$. Now assume that these properaties hold up to and including order $r$. Note that Equations (3.22) and (3.25) imply that

$$
\begin{aligned}
\partial_{h} W_{s, t} & =\left(\partial_{h} V_{s}\right) W_{t}+V_{s}\left(\partial_{h} W_{t}\right) \\
& =-V_{s}\left(\int_{0}^{s} W_{\tau} A_{i} V_{\tau} \dot{h}_{\tau}^{i} d \tau\right) W_{t}+V_{s}\left(\int_{0}^{t} W_{\tau} A_{i} V_{\tau} \dot{h}_{\tau}^{i} d \tau\right) W_{t} \\
& =\int_{\{s \leq \tau \leq t\}} W_{s, \tau} A_{i} W_{\tau, t} \dot{h}_{\tau}^{i} d \tau .
\end{aligned}
$$

Let $W^{s}$ denote the solution to Equation (3.17), and $V^{s}$ denote the solution to (3.20). Then, for $W_{\tau_{1}, \tau_{2}}^{s}=V_{\tau_{1}}^{s} W_{\tau_{2}}^{s}$,

$$
\begin{aligned}
\left\lvert\, \frac{W_{\tau_{1}, \tau_{2}}^{s}-W_{\tau_{1}, \tau_{2}}}{s}-\right. & \partial_{h} W_{\tau_{1}, \tau_{2}} \mid \\
= & \left|V_{\tau_{1}}^{s} \frac{W_{\tau_{2}}^{s}-W_{\tau_{2}}}{s}-V_{\tau_{1}}\left(\partial_{h} W_{\tau_{2}}\right)+\frac{V_{\tau_{1}}^{s}-V_{\tau_{1}}}{s} W_{\tau_{2}}-\left(\partial_{h} V_{\tau_{1}}\right) W_{\tau_{2}}\right| \\
\leq & \left|V_{\tau_{1}}^{s}-V_{\tau_{1}}\right|\left|\frac{W_{\tau_{2}}^{s}-W_{\tau_{2}}}{s}\right|+\left|V_{\tau_{1}}\right|\left|\frac{W_{\tau_{2}}^{s}-W_{\tau_{2}}}{s}-\left(\partial_{h} W_{\tau_{2}}\right)\right| \\
& \quad+\left|\frac{V_{\tau_{1}}^{s}-V_{\tau_{1}}}{s}-\left(\partial_{h} V_{\tau_{1}}\right)\right|\left|W_{\tau_{2}}\right|
\end{aligned}
$$

and thus Proposition 3.14 and Equations (3.21) and (3.23) imply that

$$
\lim _{s \downarrow 0} \sup _{\tau_{1}, \tau_{2} \leq t} \mathbb{E}\left|\frac{W_{\tau_{1}, \tau_{2}}^{s}-W_{\tau_{1}, \tau_{2}}}{s}-\partial_{h} W_{\tau_{1}, \tau_{2}}\right|^{p}=0 .
$$

So for $a_{1}, a_{2}, b_{1}, b_{2} \in[0,1]$, we have

$$
\begin{aligned}
\lim _{s \downarrow 0} \mathbb{E} \mid \int_{a_{2}}^{b_{2}} \int_{a_{1}}^{b_{1}} & \frac{W_{\tau_{1}, \tau_{2}}^{s}-W_{\tau_{1}, \tau_{2}}}{s} d \tau_{1} d \tau_{2}-\left.\int_{a_{2}}^{b_{2}} \int_{a_{1}}^{b_{1}} \partial_{h} W_{\tau_{1}, \tau_{2}} d \tau_{1} d \tau_{2}\right|^{p} \\
& \leq \lim _{s \downarrow 0} \mathbb{E} \int_{a_{2}}^{b_{2}} \int_{a_{1}}^{b_{1}}\left|\frac{W_{\tau_{1}, \tau_{2}}^{s}-W_{\tau_{1}, \tau_{2}}}{s}-\partial_{h} W_{\tau_{1}, \tau_{2}}\right|^{p} d \tau_{1} d \tau_{2} \\
& =\mathbb{E} \int_{a_{2}}^{b_{2}} \int_{a_{1}}^{b_{1}} \lim _{s \downarrow 0}\left|\frac{W_{\tau_{1}, \tau_{2}}^{s}-W_{\tau_{1}, \tau_{2}}}{s}-\partial_{h} W_{\tau_{1}, \tau_{2}}\right|^{p} d \tau_{1} d \tau_{2}=0
\end{aligned}
$$

by dominated convergence. Thus,

$$
\partial_{h}\left(\int_{a_{2}}^{b_{2}} \int_{a_{1}}^{b_{1}} W_{\tau_{1}, \tau_{2}} d \tau_{1} d \tau_{2}\right)=\int_{a_{2}}^{b_{2}} \int_{a_{1}}^{b_{1}}\left(\partial_{h} W_{\tau_{1}, \tau_{2}}\right) d \tau_{1} d \tau_{2} .
$$

Performing similar estimates, we may then show that

$$
\begin{aligned}
& \partial_{h} D_{t_{r}, \ldots, t_{1}}^{r, \alpha} \\
& =\int_{[0, T]^{r}} \\
& \quad+W_{T_{1}} A_{J_{1}}\left(\partial_{h} W_{T_{1}} W_{T_{1}, T_{2}}\right) A_{J_{1}} W_{T_{1}, T_{2}} A_{J_{2}} \cdots W_{T_{r-1}, T_{r}} A_{J_{r}} W_{T_{r}, t} \\
& \quad+\cdots+W_{T_{r-1}, T_{r}} A_{J_{r}} W_{T_{r}, t} \\
& \quad+\cdots A_{J_{1}} W_{T_{1}, T_{2}} A_{J_{2}} \cdots\left(\partial_{h} W_{T_{r-1}, T_{r}}\right) A_{J_{r}} W_{T_{r}, t} \\
& \left.\quad+W_{T_{1}} A_{J_{1}} W_{T_{1}, T_{2}} A_{J_{2}} \cdots W_{T_{r-1}, T_{r}} A_{J_{r}}\left(\partial_{h} W_{T_{r}, t}\right)\right] d \tau_{r} \cdots d \tau_{1} \\
& \quad+\int_{\left\{0 \leq \tau_{r+1} \leq T_{1}\right\}} W_{\tau_{r+1}} A_{\alpha_{r+1}} W_{\tau_{r+1}, T_{1}} A_{J_{1}} W_{T_{1}, T_{2}} A_{J_{2}} \cdots W_{T_{r-1}, T_{r}} A_{J_{r}} W_{T_{r}, t} \\
& \quad+\int_{\left\{T_{1} \leq \tau_{r+1} \leq T_{2}\right\}} W_{T_{1}} A_{J_{1}} W_{T_{1}, \tau_{r+1}} A_{\alpha_{r+1}} W_{\tau_{r+1}, T_{2}} A_{J_{2}} \cdots W_{T_{r-1}, T_{r}} A_{J_{r}} W_{T_{r}, t} \\
& \quad+\cdots+\int_{\left\{T_{r-1} \leq \tau_{r+1} \leq T_{r}\right\}} W_{T_{1}} A_{J_{1}} W_{T_{1}, T_{2}} A_{J_{2}} \cdots W_{T_{r-1}, \tau_{r+1}} A_{\alpha_{r+1}} W_{\tau_{r+1}, T_{r}} A_{J_{r}} W_{T_{r}, t} \\
& \left.\quad+\int_{\left\{T_{r} \leq \tau_{r+1} \leq t\right\}} W_{T_{1}} A_{J_{1}} W_{T_{1}, T_{2}} A_{J_{2}} \cdots W_{T_{r-1}, T_{r}} A_{J_{r}} W_{T_{r}, \tau_{r+1}} A_{\alpha_{r+1}} W_{\tau_{r+1}, t}\right] \\
& \\
& \quad \times \dot{h}_{\tau_{r+1}}^{\alpha_{r+1}} d \tau_{r+1} d \tau_{r} \cdots d \tau_{1},
\end{aligned}
$$

which implies exactly that (P2) holds for $r+1$, since

$$
\left\{0 \leq \tau_{r+1} \leq T_{1}\right\},\left\{T_{1} \leq \tau_{r+1} \leq T_{2}\right\}, \cdots,\left\{T_{r-1} \leq \tau_{r+1} \leq T_{r}\right\},\left\{T_{r} \leq \tau_{r+1} \leq t\right\}
$$

partitions the set $[0, t]$. Clearly, this also implies that, for all $t_{r+1} \in[0,1]$ and $\alpha_{r+1}=$ $1, \ldots, k,\left[D_{t_{r+1}} D_{t_{r}, \ldots, t_{1}}^{r, \alpha}\right]^{\alpha_{r+1}} \in L^{p}(\mu)$ for all $p \in(1, \infty)$, by Equation (3.24), and so $W_{t} \in \mathcal{D}^{r+1, p}(\operatorname{End}(\mathfrak{g}))$.

Now, for $\bar{W}_{t}=\int_{0}^{t} W_{\tau} d \tau$, the above arguments imply that

$$
\partial_{h} \bar{W}_{t}=\int_{0}^{t} \partial_{h} W_{\tau} d \tau=\int_{0}^{t} \int_{\left\{0 \leq \tau_{1} \leq \tau\right\}} W_{\tau_{1}} A_{\alpha_{1}} W_{\tau_{1}, \tau} \dot{h}_{\tau_{1}}^{\alpha_{1}} d \tau_{1} d \tau
$$

and, for $t_{1} \leq t$ and $\alpha_{1}=1, \ldots, k$,

$$
\left[D_{t_{1}} \bar{W}_{t}\right]^{\alpha_{1}}=\int_{0}^{t} \int_{\left\{0 \leq \tau_{1} \leq \tau \wedge t_{1}\right\}} W_{\tau_{1}} A_{\alpha_{1}} W_{\tau_{1}, \tau} d \tau_{1} d \tau .
$$

Let $r \in \mathbb{N}$. For any $\left\{t_{1}, \ldots, t_{r}\right\} \subset[0,1]$ and multindex $\alpha=\left(\alpha_{1}, \ldots, \alpha_{r}\right)$, let $D_{t_{1}, \ldots, t_{r}}^{r, \alpha} \bar{W}_{t}$ be the $\alpha^{t h}$ component of $D_{t_{1}, \ldots, t_{r}}^{r} \bar{W}_{t}$; that is,

$$
D_{t_{r}, \ldots, t_{1}}^{r, \bar{W}_{t}}:=\left[D_{t_{r}}\left[\cdots\left[D_{t_{2}}\left[D_{t_{1}} \bar{W}_{t}\right]^{\alpha_{1}}\right]^{\alpha_{2}}\right] \cdots\right]^{\alpha_{r}} .
$$

Then working as above, we may show that

$$
\begin{aligned}
D_{t_{r}, \ldots, t_{1}}^{r, \alpha} \bar{W}_{t} & =\int_{0}^{t} \int_{[0, T \wedge \tau]^{r}} W_{T_{1}} A_{J_{1}} W_{T_{1}, T_{2}} A_{J_{2}} \cdots W_{T_{r-1}, T_{r}} A_{J_{r}} W_{T_{r}, \tau} d \tau_{r} \cdots d \tau_{1} d \tau \\
& =\int_{0}^{t} D_{t_{r}, \ldots, t_{1}}^{r, \alpha} W_{\tau} d \tau
\end{aligned}
$$

and since $W_{\tau}$ satisfies (P1) for all $\tau \in[0, t]$, we have $\bar{W} \in \mathcal{D}^{r, p}(\mathscr{H}(\operatorname{End}(\mathfrak{g})))$, for all $r \in \mathbb{N}$ and $p \in(1, \infty)$.

Proposition 3.16. Let $\psi \in C_{c}^{\infty}(M)$ such that $\psi=1$ near $I$ and $\psi(x)=0$ if $|x| \geq 2$, where $|\cdot|$ is the distance from I with respect to any metric on $M$. For any $n \in \mathbb{N}$, define $\psi_{n}(x):=\psi(x / n)$, and, for any $A \in M$, define $\left\langle\psi^{\prime}(x), A\right\rangle:=\left.\frac{d}{d t}\right|_{0} \psi(x+t A)$. Let $W$ denote the solution to Equation (3.16), and let $W^{n}:[0,1] \times \mathscr{W} \rightarrow M$ denote the solution to the Stratonovich differential equation

$$
\begin{equation*}
d W^{n}=\psi_{n}\left(W^{n}\right) W^{n} \circ d B, \text { with } W_{0}=I \tag{3.28}
\end{equation*}
$$

Then

$$
\lim _{n \rightarrow \infty} \mathbb{E} \sup _{\tau \leq 1}\left|W_{\tau}^{n}-W_{\tau}\right|^{p}=0
$$

for all $p \in(1, \infty)$.
Remark 3.17. Notice that $\psi_{n}^{\prime}(x)=n^{-1} \psi^{\prime}(x / n)$, and therefore

$$
\left|\psi_{n}^{\prime}(x)\right||x| \leq n^{-1} C 2 n=2 C
$$

where $C$ is a bound on $\psi^{\prime}$. Similarly, we may show

$$
\left|\psi_{n}^{\prime \prime}(x)\right||x|^{2} \leq C,
$$

where $C$ is determined by a bound on $\psi^{\prime \prime}$. These bounds will be used repeatedly in the sequel without further mention.

Proof. Equation (3.28) in Itô form is

$$
\begin{align*}
& d W^{n}= \psi_{n}\left(W^{n}\right) W^{n} d B  \tag{3.29}\\
&+\frac{1}{2}\left[\left\langle\psi_{n}^{\prime}\left(W^{n}\right), \psi_{n}\left(W^{n}\right) W^{n} d B\right\rangle W^{n}+\psi_{n}\left(W^{n}\right) \psi_{n}\left(W^{n}\right) W^{n} d B\right] d B \\
&=\psi_{n}\left(W^{n}\right) W^{n} d B  \tag{3.30}\\
&+\frac{1}{2} \sum_{i=1}^{k}\left[\psi_{n}\left(W^{n}\right)\left\langle\psi_{n}^{\prime}\left(W^{n}\right), W^{n} A_{i}\right\rangle W^{n} A_{i}+\psi_{n}^{2}\left(W^{n}\right) W^{n} A_{i}^{2}\right] d t . \tag{3.31}
\end{align*}
$$

By Proposition 3.11, Equation (3.31) implies that

$$
\begin{aligned}
\mathbb{E} \sup _{\tau \leq t}\left|W_{\tau}^{n}\right|^{p} & \lesssim 1+\mathbb{E}\left(\int_{0}^{t}\left|W^{n}\right|^{2} d \tau\right)^{p / 2}+\mathbb{E}\left(\int_{0}^{t}\left|W^{n}\right| d \tau\right)^{p} \\
& \lesssim 1+\mathbb{E} \int_{0}^{t}\left|W^{n}\right|^{p} d \tau
\end{aligned}
$$

for all $t \in[0,1]$. An application of Gronwall's inequality then shows that

$$
\mathbb{E} \sup _{\tau \leq 1}\left|W_{\tau}^{n}\right|^{p} \leq C e^{C},
$$

where these constants are independent of $n$. Thus, there exists some finite constants $C_{p}$ so that

$$
\begin{equation*}
\sup _{n \in \mathbb{N}}\left(\mathbb{E} \sup _{\tau \leq 1}\left|W_{\tau}^{n}\right|^{p}\right) \leq C_{p}, \tag{3.32}
\end{equation*}
$$

for all $p \in(1, \infty)$.
Now let $\varepsilon^{n}:=W^{n}-W$ (so $W^{n}=W+\varepsilon^{n}$ ). Then we have

$$
\begin{align*}
d \varepsilon^{n}= & {\left[\psi_{n}\left(W^{n}\right) W^{n}-W\right] d B } \\
& +\frac{1}{2} \sum_{i=1}^{k}\left[\psi_{n}\left(W^{n}\right)\left\langle\psi_{n}^{\prime}\left(W^{n}\right), W^{n} A_{i}\right\rangle W^{n} A_{i}+\left(\psi_{n}^{2}\left(W^{n}\right) W^{n}-W\right) A_{i}^{2}\right] d t . \tag{3.33}
\end{align*}
$$

Applying Proposition 3.11 to each term in Equation (3.33), we may bound $\mathbb{E}\left|\varepsilon^{n}\right|^{p}$. For
the first term,

$$
\begin{aligned}
& \mathbb{E}\left|\int_{0}^{t}\left[\psi_{n}\left(W^{n}\right) W^{n}-W\right] d B\right|^{p}=\mathbb{E}\left|\int_{0}^{t}\left[\psi_{n}\left(W^{n}\right)\left(\varepsilon^{n}+W\right)-W\right] d B\right|^{p} \\
& \lesssim \mathbb{E} \int_{0}^{t}\left|\psi_{n}\left(W^{n}\right) \varepsilon^{n}+\left(\psi_{n}\left(W^{n}\right)-1\right) W\right|^{p} d \tau \\
& \lesssim \mathbb{E} \int_{0}^{t}\left|\psi_{n}\left(W^{n}\right)\right|^{p}\left|\varepsilon^{n}\right|^{p} d \tau+\mathbb{E} \int_{0}^{t}\left|\psi_{n}\left(W^{n}\right)-1\right|^{p}|W|^{p} d \tau
\end{aligned}
$$

Similarly, for the second term,

$$
\begin{aligned}
\mathbb{E} \mid \int_{0}^{t} \psi_{n} & \left.\left(W^{n}\right)\left\langle\psi_{n}^{\prime}\left(W^{n}\right), W^{n} A_{i}\right\rangle W^{n} A_{i} d \tau\right|^{p} \\
& =\mathbb{E}\left|\int_{0}^{t} \psi_{n}\left(W^{n}\right)\left\langle\psi_{n}^{\prime}\left(W^{n}\right), W^{n} A_{i}\right\rangle\left(\varepsilon^{n}+W\right) A_{i} d \tau\right|^{p} \\
& \lesssim \mathbb{E} \int_{0}^{t}\left[\left|\psi_{n}\left(W^{n}\right)\right|\left|\psi_{n}^{\prime}\left(W^{n}\right)\right|\left|W^{n} \| \varepsilon^{n}+W\right|\right]^{p} d \tau \\
& \lesssim \mathbb{E} \int_{0}^{t}\left[\left|\psi_{n}^{\prime}\left(W^{n}\right)\right|\left|W^{n}\right|\right]^{p}\left|\varepsilon^{n}\right|^{p} d \tau+\mathbb{E} \int_{0}^{t}\left[\left|\psi_{n}^{\prime}\left(W^{n}\right)\right|\left|W^{n}\right||W|\right]^{p} d \tau
\end{aligned}
$$

And finally,

$$
\begin{aligned}
\mathbb{E} \mid \int_{0}^{t}\left[\psi_{n}^{2}\left(W^{n}\right)\right. & \left.W^{n}-W\right]\left.A_{i}^{2} d \tau\right|^{p}=\mathbb{E}\left|\int_{0}^{t}\left[\psi_{n}^{2}\left(W^{n}\right)\left(\varepsilon^{n}+W\right)-W\right] A_{i}^{2} d \tau\right|^{p} \\
& \lesssim \mathbb{E} \int_{0}^{t}\left|\psi_{n}^{2}\left(W^{n}\right)\left(\varepsilon^{n}+W\right)-W\right|^{p} d \tau \\
& \lesssim \mathbb{E} \int_{0}^{t}\left|\psi_{n}^{2}\left(W^{n}\right)\right|^{p}\left|\varepsilon^{n}\right|^{p} d \tau+\mathbb{E} \int_{0}^{t}\left|\psi_{n}^{2}\left(W^{n}\right)-1\right|^{p}|W|^{p} d \tau .
\end{aligned}
$$

Bringing all of this together together, we have

$$
\mathbb{E} \sup _{\tau \leq t}\left|\varepsilon_{\tau}^{n}\right|^{p} \leq C \mathbb{E} \int_{0}^{t}\left|\varepsilon^{n}\right|^{p} d \tau+\delta_{n}
$$

for all $t \in[0,1]$, where (up to constant multiple)

$$
\begin{aligned}
& \delta_{n}=\mathbb{E} \int_{0}^{t}\left|\psi_{n}\left(W^{n}\right)-1\right|^{p}|W|^{p} d \tau \\
&+\mathbb{E} \int_{0}^{t}\left[\left|\psi_{n}^{\prime}\left(W^{n}\right)\right|\left|W^{n}\right||W|\right]^{p} d \tau+\mathbb{E} \int_{0}^{t}\left|\psi_{n}^{2}\left(W^{n}\right)-1\right|^{p}|W|^{p} d \tau
\end{aligned}
$$

Since $\left|\psi_{n}^{\prime}\left(W^{n}\right)\right|\left|W^{n}\right|$ and $|W|$ remain bounded, the dominated convergence theorem implies that $\lim _{n \rightarrow \infty} \delta_{n}=0$. Thus, by Gronwall again, we have

$$
\mathbb{E} \sup _{\tau \leq 1}\left|W_{\tau}^{n}-W_{\tau}\right|^{p}=\mathbb{E} \sup _{\tau \leq 1}\left|\varepsilon_{\tau}^{n}\right|^{p} \leq \delta_{n} e^{C} \rightarrow 0
$$

as $n \rightarrow \infty$.
Proposition 3.18. Let $W$ be the solution to Equation (3.16), and let $W^{n}$ be the solution to Equation (3.28). Then $W^{n} \in \operatorname{Dom}\left(\partial_{h}\right)$ and

$$
\lim _{n \rightarrow \infty} \mathbb{E} \sup _{\tau \leq 1}\left|\partial_{h} W_{\tau}^{n}-\partial_{h} W_{\tau}\right|^{p}=0
$$

for all $p \in(1, \infty)$.
Proof. As in Proposition 3.14, we may show that $\partial_{h} W^{n}$ satisfies the Itô equation

$$
\begin{aligned}
d\left(\partial_{h} W^{n}\right)= & {\left[\left\langle\psi_{n}^{\prime}\left(W^{n}\right), \partial_{h} W^{n}\right\rangle W^{n}+\psi_{n}\left(W^{n}\right)\left(\partial_{h} W^{n}\right)\right] d B+\psi_{n}\left(W^{n}\right) W^{n} d H } \\
& +\frac{1}{2} \sum_{i=1}^{k}\left[\begin{array}{c}
\left\langle\psi_{n}^{\prime}\left(W^{n}\right), \partial_{h} W^{n}\right\rangle\left\langle\psi_{n}^{\prime}\left(W^{n}\right), W^{n} A_{i}\right\rangle W^{n} A_{i} \\
+\psi_{n}\left(W^{n}\right)\left\langle\psi_{n}^{\prime \prime}\left(W^{n}\right), \partial_{h} W^{n} \otimes W^{n} A_{i}\right\rangle W^{n} A_{i} \\
+\psi_{n}\left(W^{n}\right)\left\langle\psi_{n}^{\prime}\left(W^{n}\right),\left(\partial_{h} W^{n}\right) A_{i}\right\rangle W^{n} A_{i} \\
+\psi_{n}\left(W^{n}\right)\left\langle\psi_{n}^{\prime}\left(W^{n}\right), W^{n} A_{i}\right\rangle\left(\partial_{h} W^{n}\right) A_{i} \\
+2 \psi_{n}\left(W^{n}\right)\left\langle\psi_{n}^{\prime}\left(W^{n}\right), \partial_{h} W^{n}\right\rangle W^{n} A_{i}^{2}+\psi_{n}^{2}\left(W^{n}\right)\left(\partial_{h} W^{n}\right) A_{i}^{2}
\end{array}\right] d t .
\end{aligned}
$$

Recall also that

$$
d\left(\partial_{h} W\right)=\left(\partial_{h} W\right) d B+W d H+\frac{1}{2}\left(\partial_{h} W\right) A_{i}^{2} d t .
$$

Let $\varepsilon^{n}:=\partial_{h} W^{n}-\partial_{h} W$. Then

$$
\begin{gathered}
d \varepsilon^{n}=\left[\left\langle\psi_{n}^{\prime}\left(W^{n}\right), \partial_{h} W^{n}\right\rangle W^{n}+\left(\psi_{n}\left(W^{n}\right)\left(\partial_{h} W^{n}\right)-\left(\partial_{h} W\right)\right)\right] d B \\
+\left[\psi_{n}\left(W^{n}\right) W^{n}-W\right] d H \\
+\frac{1}{2} \sum_{i=1}^{k}\left[\begin{array}{c}
\left\langle\psi_{n}^{\prime}\left(W^{n}\right), \partial_{h} W^{n}\right\rangle\left\langle\psi_{n}^{\prime}\left(W^{n}\right), W^{n} A_{i}\right\rangle W^{n} A_{i} \\
+\psi_{n}\left(W^{n}\right)\left\langle\psi_{n}^{\prime \prime}\left(W^{n}\right), \partial_{h} W^{n} \otimes W^{n} A_{i}\right\rangle W^{n} A_{i} \\
+\psi_{n}\left(W^{n}\right)\left\langle\psi_{n}^{\prime}\left(W^{n}\right),\left(\partial_{h} W^{n}\right) A_{i}\right\rangle W^{n} A_{i} \\
+\psi_{n}\left(W^{n}\right)\left\langle\psi_{n}^{\prime}\left(W^{n}\right), W^{n} A_{i}\right\rangle\left(\partial_{h} W^{n}\right) A_{i} \\
+2 \psi_{n}\left(W^{n}\right)\left\langle\psi_{n}^{\prime}\left(W^{n}\right), \partial_{h} W^{n}\right\rangle W^{n} A_{i}^{2}+\left[\psi_{n}^{2}\left(W^{n}\right) \partial_{h} W^{n}-\partial_{h} W\right] A_{i}^{2}
\end{array}\right] d t .
\end{gathered}
$$

That is,

$$
\begin{align*}
& d \varepsilon^{n}=\left\langle\psi_{n}^{\prime}\left(W^{n}\right), \partial_{h} W^{n}\right\rangle W^{n} d B+\left[\psi_{n}\left(W^{n}\right)\left(\partial_{h} W^{n}\right)-\left(\partial_{h} W\right)\right] d B+\left[\psi_{n}\left(W^{n}\right) W^{n}-W\right] d H \\
&+\frac{1}{2} \sum_{i=1}^{k}\left[\Gamma_{n}\left(\partial_{h} W^{n}\right)+\left[\psi_{n}^{2}\left(W^{n}\right) \partial_{h} W^{n}-\partial_{h} W\right] A_{i}^{2}\right] d t, \tag{3.34}
\end{align*}
$$

where

$$
\begin{aligned}
\Gamma_{n}(A)= & \sum_{i=1}^{k}\left[\left\langle\psi_{n}^{\prime}\left(W^{n}\right), A\right\rangle\left\langle\psi_{n}^{\prime}\left(W^{n}\right), W^{n} A_{i}\right\rangle W^{n} A_{i}\right. \\
+ & \psi_{n}\left(W^{n}\right)\left\langle\psi_{n}^{\prime \prime}\left(W^{n}\right), A \otimes W^{n} A_{i}\right\rangle W^{n} A_{i}+\psi_{n}\left(W^{n}\right)\left\langle\psi_{n}^{\prime}\left(W^{n}\right), A A_{i}\right\rangle W^{n} A_{i} \\
& \left.\quad+\psi_{n}\left(W^{n}\right)\left\langle\psi_{n}^{\prime}\left(W^{n}\right), W^{n} A_{i}\right\rangle A A_{i}+2 \psi_{n}\left(W^{n}\right)\left\langle\psi_{n}^{\prime}\left(W^{n}\right), A\right\rangle W^{n} A_{i}^{2}\right]
\end{aligned}
$$

satisfies the bound

$$
\begin{equation*}
\left|\Gamma_{n}\right| \lesssim\left|\psi_{n}^{\prime \prime}\left(W^{n}\right)\right|\left|W^{n}\right|^{2}+\left|\psi_{n}^{\prime}\left(W^{n}\right)\right|\left|W^{n}\right|+\left|\psi_{n}^{\prime}\left(W^{n}\right)\right|^{2}\left|W^{n}\right|^{2} . \tag{3.35}
\end{equation*}
$$

Again using Proposition (3.11), we may work through (3.34) term by term to bound $\mathbb{E}\left|\varepsilon^{n}\right|^{p}$ :
For the first term in the sum,

$$
\begin{aligned}
\mathbb{E}\left|\int_{0}^{t}\left\langle\psi_{n}^{\prime}\left(W^{n}\right), \partial_{h} W^{n}\right\rangle W^{n} d B\right|^{p} & \lesssim \mathbb{E} \int_{0}^{t}\left|\left\langle\psi_{n}^{\prime}\left(W^{n}\right), \varepsilon^{n}+\partial_{h} W\right\rangle W^{n}\right|^{p} d \tau \\
& \lesssim \mathbb{E} \int_{0}^{t}\left|\varepsilon^{n}\right|^{p} d \tau+\mathbb{E} \int_{0}^{t}\left[\left|\psi_{n}^{\prime}\left(W^{n}\right)\right|\left|\partial_{h} W\right|\left|W^{n}\right|\right]^{p} d \tau
\end{aligned}
$$

Considering the second term, we have

$$
\begin{aligned}
\mathbb{E}\left|\int_{0}^{t}\left[\psi_{n}\left(W^{n}\right) \partial_{h} W^{n}-\partial_{h} W\right] d B\right|^{p} & \lesssim \mathbb{E} \int_{0}^{t}\left|\psi_{n}\left(W^{n}\right)\left(\varepsilon^{n}+\partial_{h} W\right)-\partial_{h} W\right|^{p} d \tau \\
& \lesssim \mathbb{E} \int_{0}^{t}\left|\varepsilon^{n}\right|^{p} d \tau+\mathbb{E} \int_{0}^{t}\left|\psi_{n}\left(W^{n}\right)-1\right|^{p}\left|\partial_{h} W\right|^{p} d \tau .
\end{aligned}
$$

For the third term, note that

$$
\mathbb{E}\left|\int_{0}^{t}\left[\psi_{n}\left(W^{n}\right) W^{n}-W\right] d H\right|^{p}=C\|H\|_{\mathscr{H}}^{p} \mathbb{E} \sup _{\tau \leq 1}\left|\psi_{n}\left(W_{\tau}^{n}\right) W_{\tau}^{n}-W_{\tau}\right|^{p}
$$

and

$$
\begin{aligned}
\mathbb{E} \sup _{\tau \leq 1}\left|\psi_{n}\left(W_{\tau}^{n}\right) W_{\tau}^{n}-W_{\tau}\right|^{p} & =\mathbb{E} \sup _{\tau \leq 1}\left|\psi_{n}\left(W_{t}^{n}\right)\left(W_{t}+\left(W_{t}^{n}-W_{t}\right)\right)-W_{t}\right|^{p} \\
& \leq \mathbb{E} \sup _{\tau \leq 1}\left(\left|\psi_{n}\left(W_{\tau}^{n}\right)-1\right|^{p}\left|W_{\tau}\right|^{p}+\left|\psi_{n}\left(W_{\tau}^{n}\right)\right|^{n}\left|W_{\tau}^{n}-W_{\tau}\right|^{p}\right) \rightarrow 0,
\end{aligned}
$$

as $n \rightarrow \infty$, by Proposition 3.16 and dominated convergence. Using the bound in 3.35) on the fourth term, we have

$$
\begin{aligned}
& \mathbb{E}\left|\int_{0}^{t} \Gamma_{n}\left(\partial_{h} W^{n}\right) d \tau\right|^{p}=\mathbb{E}\left|\int_{0}^{t} \Gamma_{n}\left(\varepsilon^{n}+\partial_{h} W\right) d \tau\right|^{p} \\
& \quad \lesssim \mathbb{E} \int_{0}^{t}\left|\varepsilon^{n}\right|^{p} d \tau \\
& \quad+\mathbb{E} \int_{0}^{t}\left(\left|\psi_{n}^{\prime \prime}\left(W^{n}\right)\right|\left|W^{n}\right|^{2}+\left|\psi_{n}^{\prime}\left(W^{n}\right)\right|\left|W^{n}\right|+\left|\psi_{n}^{\prime}\left(W^{n}\right)\right|^{2}\left|W^{n}\right|^{2}\right)^{p}\left|\partial_{h} W\right|^{p} d \tau .
\end{aligned}
$$

Finally, for the last term,

$$
\begin{aligned}
\mathbb{E}\left|\int_{0}^{t}\left[\psi_{n}^{2}\left(W^{n}\right) \partial_{h} W^{n}-\partial_{h} W\right] A_{i}^{2} d \tau\right|^{p} & \lesssim \mathbb{E} \int_{0}^{t}\left|\psi_{n}^{2}\left(W^{n}\right)\left(\varepsilon^{n}+\partial_{h} W\right)-\partial_{h} W\right|^{p} d \tau \\
& \lesssim \mathbb{E} \int_{0}^{t}\left|\varepsilon^{n}\right|^{p} d \tau+\mathbb{E} \int_{0}^{t}\left|\psi_{n}^{2}\left(W^{n}\right)-1\right|^{p}\left|\partial_{h} W\right|^{p} d \tau
\end{aligned}
$$

Putting this all together shows

$$
\mathbb{E} \sup _{\tau \leq t}\left|\varepsilon_{t}^{n}\right|^{p} \leq C \mathbb{E} \int_{0}^{t}\left|\varepsilon^{n}\right|^{p} d \tau+\delta_{n},
$$

for all $t \in[0,1]$, where

$$
\begin{aligned}
\delta_{n}= & \mathbb{E} \int_{0}^{t}\left|\psi_{n}^{\prime}\left(W^{n}\right)\right|^{p}\left|\partial_{h} W\right|^{p}\left|W^{n}\right|^{p} d \tau+\mathbb{E} \int_{0}^{t}\left|\psi_{n}\left(W^{n}\right)-1\right|^{p}\left|\partial_{h} W\right|^{p} d \tau \\
& +\mathbb{E}\left|\int_{0}^{t}\left[\psi_{n}\left(W^{n}\right) W^{n}-W\right] d H\right|^{p} \\
& +\mathbb{E} \int_{0}^{t}\left(\left|\psi_{n}^{\prime \prime}\left(W^{n}\right)\right|\left|W^{n}\right|^{2}+\left|\psi_{n}^{\prime}\left(W^{n}\right)\right|\left|W^{n}\right|+\left|\psi_{n}^{\prime}\left(W^{n}\right)\right|^{2}\left|W^{n}\right|^{2}\right)^{p}\left|\partial_{h} W\right|^{p} d \tau \\
& +\mathbb{E} \int_{0}^{t}\left|\psi_{n}^{2}\left(W^{n}\right)-1\right|^{p}\left|\partial_{h} W\right|^{p} d \tau .
\end{aligned}
$$

Again using Remark 3.17, $\lim _{n \rightarrow \infty} \delta_{n}=0$, by the dominated convergence theorem. Thus, another application of Gronwall's inequality shows that

$$
\mathbb{E} \sup _{\tau \leq 1}\left|\partial_{h} W_{\tau}^{n}-\partial_{h} W_{\tau}\right|^{p}=\mathbb{E} \sup _{\tau \leq 1}\left|\varepsilon_{\tau}^{n}\right|^{p} \leq \delta_{n} e^{C} \rightarrow 0
$$

as $n \rightarrow \infty$.

### 3.2.2 Lie group computations

Lemma 3.19. Let $u \in C^{1}(G)$ such that $u$ and $\tilde{X} u$ are bounded for all $X \in \mathfrak{g}_{0}$. Then the solution $\eta:[0,1] \times \mathscr{W}\left(\mathbb{R}^{k}\right) \rightarrow G$ to the stochastic differential equation

$$
d \eta=u(\eta) \eta \circ d \vec{b}:=u(\eta) L_{\eta *} \circ d \vec{b}=u(\eta) \tilde{X}_{i}(\eta) \circ d b^{i}, \text { with } \eta_{0}=e
$$

exists for all time; that is, $\eta$ has no explosion.
Proof. Let $\zeta$ be the life-time of $\eta$ and $\varphi \in C_{c}^{\infty}(G)$. Then on $\{t<\zeta\}$,

$$
\begin{align*}
d \varphi(\eta)= & u(\eta) \tilde{X}_{i} \varphi(\eta) \circ d b^{i}=u(\eta)\left\langle\tilde{\nabla} \varphi(\eta), X_{i}\right\rangle \circ d b^{i}=u(\eta)\langle\tilde{\nabla} \varphi(\eta), \circ d \vec{b}\rangle \\
= & u(\eta)\langle\tilde{\nabla} \varphi(\eta), d \vec{b}\rangle+\frac{1}{2} d[u(\eta)\langle\tilde{\nabla} \varphi(\eta), \cdot\rangle] d \vec{b} \\
= & u(\eta)\langle\tilde{\nabla} \varphi(\eta), d \vec{b}\rangle+\frac{1}{2}\langle\tilde{\nabla}[u(\eta)\langle\tilde{\nabla} \varphi(\eta), \cdot\rangle], d \vec{b}\rangle d \vec{b} \\
= & u(\eta)\langle\tilde{\nabla} \varphi(\eta), d \vec{b}\rangle+\frac{1}{2}\left[\langle\tilde{\nabla} u(\eta), d \vec{b}\rangle\langle\tilde{\nabla} \varphi(\eta), d \vec{b}\rangle+u(\eta)\left\langle\tilde{\nabla}^{2} \varphi(\eta), d \vec{b} \otimes d \vec{b}\right\rangle\right] \\
= & u(\eta)\langle\tilde{\nabla} \varphi(\eta), d \vec{b}\rangle \\
& \quad+\frac{1}{2} \sum_{i=1}^{k}\left[\left\langle\tilde{\nabla} u(\eta), X_{i}\right\rangle\left\langle\tilde{\nabla} \varphi(\eta), X_{i}\right\rangle+u(\eta)\left\langle\tilde{\nabla}^{2} \varphi(\eta), X_{i} \otimes X_{i}\right\rangle\right] d t . \tag{3.36}
\end{align*}
$$

Let $|g|$ denote the distance from a point $g \in G$ to $e$ with respect to the right invariant metric. Let $\left\{\varphi_{m}\right\}_{m=1}^{\infty} \subset C_{c}^{\infty}(G,[0,1])$ be a sequence of functions, $\varphi_{m} \uparrow 1$, such that $\varphi_{m}(g)=1$ when $|g| \leq m$ and $\sup _{m} \sup _{g \in G}\left|\tilde{\nabla}^{k} \varphi_{m}(g)\right|<\infty$ for $k=0,1,2, \ldots$; see Lemma 3.6 of Driver and Gross [16]. Also, let $\left\{K_{n}\right\}_{n=1}^{\infty}$ be a sequence of nested compact sets in $G$, such that $K_{n} \uparrow G$, and take $\tau_{n}$ to be the exit time from $K_{n}$,

$$
\tau_{n}:=\inf \left\{t>0: \eta_{t} \notin K_{n}\right\} .
$$

Let $\eta_{t}^{n}:=\eta_{t \wedge \tau_{n}}$. Then from Equation (3.36) (using the convention that $\varphi_{m}\left(\eta_{t}\right)=0$, $\tilde{\nabla} \varphi_{m}\left(\eta_{t}\right)=0$, and $\tilde{\nabla}^{2} \varphi_{m}\left(\eta_{t}\right)=0$ on $\left.\{t>\zeta\}\right)$ we have

$$
\begin{aligned}
\varphi_{m}\left(\eta_{t}^{n}\right) & =1+\int_{0}^{t \wedge \tau_{n}} u(\eta)\left\langle\tilde{\nabla} \varphi_{m}(\eta), d \vec{b}\right\rangle \\
& +\frac{1}{2} \sum_{i=1}^{k} \int_{0}^{t \wedge \tau_{n}}\left[\left\langle\tilde{\nabla} u(\eta), X_{i}\right\rangle\left\langle\tilde{\nabla} \varphi_{m}(\eta), X_{i}\right\rangle+u(\eta)\left\langle\tilde{\nabla}^{2} \varphi_{m}(\eta), X_{i} \otimes X_{i}\right\rangle\right] d \tau
\end{aligned}
$$

Taking the expectation of this equation then gives

$$
\mathbb{E}\left[\varphi_{m}\left(\eta_{t}^{n}\right)\right]=1+\frac{1}{2} \delta_{m},
$$

where

$$
\delta_{m}:=\sum_{i=1}^{k} \mathbb{E} \int_{0}^{t \wedge \tau_{n}}\left[\left\langle\tilde{\nabla} u(\eta), X_{i}\right\rangle\left\langle\tilde{\nabla} \varphi_{m}(\eta), X_{i}\right\rangle+u(\eta)\left\langle\tilde{\nabla}^{2} \varphi_{m}(\eta), X_{i} \otimes X_{i}\right\rangle\right] d \tau
$$

Now by construction of the $\varphi_{m}$ and $\tau_{m}$ and the assumptions on $u$, there is a constant $M<\infty$ such that

$$
1_{\left\{\tau \leq t \wedge \tau_{m}\right\}}\left|\left\langle\tilde{\nabla} u\left(\eta_{\tau}\right), X_{i}\right\rangle\left\langle\tilde{\nabla} \varphi_{m}\left(\eta_{\tau}\right), X_{i}\right\rangle+u\left(\eta_{\tau}\right)\left\langle\tilde{\nabla}^{2} \varphi_{m}\left(\eta_{\tau}\right), X_{i} \otimes X_{i}\right\rangle\right| \leq M
$$

Moreover, $\lim _{m \rightarrow \infty}\left|\tilde{\nabla} \varphi_{m}\right|=0=\lim _{m \rightarrow \infty}\left|\tilde{\nabla}^{2} \varphi_{m}\right|$. So it follows by the dominated convergence theorem that $\lim _{m \rightarrow \infty} \delta_{m}=0$, and we have proved

$$
1=\lim _{m \rightarrow \infty} \mathbb{E}\left[\varphi_{m}\left(\eta_{t}^{n}\right)\right]=\mathbb{E}\left[\lim _{m \rightarrow \infty} \varphi_{m}\left(\eta_{t}^{n}\right)\right]=\mathbb{E} 1_{\left\{t \wedge \tau_{n}<\zeta\right\}}=P\left(t \wedge \tau_{n}<\zeta\right)
$$

for all $n=1,2, \ldots$ Thus, $1=\lim _{n \rightarrow \infty} P\left(t \wedge \tau_{n}<\zeta\right)=P(t<\zeta)$, and so $\zeta=\infty \mu$-a.s.
Now let $u \in C_{c}^{\infty}(G)$, and suppose $\eta:[0,1] \times \mathscr{W}\left(\mathbb{R}^{k}\right) \rightarrow G$ is a solution to the stochastic differential equation

$$
\begin{equation*}
d \eta=u(\eta) \eta \circ d \vec{b}:=u(\eta) L_{\eta *} \circ d \vec{b}=u(\eta) \tilde{X}_{i}(\eta) \circ d b^{i} \tag{3.37}
\end{equation*}
$$

Since the vector fields $u \tilde{X}_{i}$ have compact support, we may embed $G$ as a Euclidean submanifold in a "nice" way so that the embedded vector fields are bounded with bounded derivatives. Then Theorem 2.1 of Taniguchi 47] implies that, for any $f \in C_{c}^{\infty}(G)$, we have $f\left(\eta_{t}\right) \in \mathcal{D}^{\infty}$.

Proposition 3.20. Fix $h \in \mathscr{H}$, and let $\eta$ be the solution to Equation (3.37). Then, for $f \in C_{c}^{\infty}(G)$,

$$
\partial_{h} f\left(\eta_{t}\right)=\left\langle\hat{\nabla} f\left(\eta_{t}\right), \theta_{t}\right\rangle,
$$

where $\theta_{t}: \mathscr{W} \rightarrow \mathfrak{g}$ solves

$$
\begin{equation*}
\theta_{t}=\int_{0}^{t}\left(\langle\hat{\nabla} u(\eta), \theta\rangle \operatorname{Ad}_{\eta} \circ d \vec{b}+u(\eta) \operatorname{Ad}_{\eta} d \vec{h}\right), \tag{3.38}
\end{equation*}
$$

where $d \vec{h}=X_{i} \dot{h}^{i} d t$; writing the above in Itô form gives

$$
\begin{equation*}
\theta_{t}=\int_{0}^{t}\left(\langle\hat{\nabla} u(\eta), \theta\rangle \operatorname{Ad}_{\eta} d \vec{b}+u(\eta) \operatorname{Ad}_{\eta} d \vec{h}+\frac{1}{2} \sum_{i=1}^{k}\left\langle\tilde{\nabla} \hat{\nabla} u(\eta), X_{i} \otimes \theta\right\rangle \operatorname{Ad}_{\eta} X_{i} d \tau\right) . \tag{3.39}
\end{equation*}
$$

Proof. Let $\eta_{t}^{s}:=\eta_{t}(\cdot+s h):[0,1]^{2} \times \mathscr{W} \rightarrow G$. Then $\eta_{t}^{s}$ satisfies the Stratonovich equation in $t$,

$$
\begin{equation*}
d \eta_{t}^{s}=u\left(\eta_{t}^{s}\right) L_{\eta_{t}^{s} *}\left(o d \vec{b}_{t}+s d \vec{h}_{t}\right) . \tag{3.40}
\end{equation*}
$$

By Corollary 4.3 of Driver [14], there exists a modification of $\eta_{t}^{s}$ so that the mapping $s \mapsto \eta_{t}^{s}$ is smooth in the sense that, for any function $f \in C_{c}^{\infty}(G), s \mapsto f\left(\eta_{t}^{s}\right)$ is smooth, and, furthermore, for any one-form $\vartheta$ acting on $T_{\eta_{t}^{s}} G$,

$$
\begin{equation*}
\left.\frac{\partial}{\partial s}\right|_{0} \int_{0}^{t} \vartheta\left(d \eta_{\tau}^{s}\right)=\int_{0}^{t} d \vartheta\left(\left.\frac{\partial}{\partial s}\right|_{0} \eta_{\tau}^{s}, d \eta_{\tau}\right)+\left.\vartheta\left(\left.\frac{\partial}{\partial s}\right|_{0} \eta_{\tau}^{s}\right)\right|_{\tau=0} ^{t} \tag{3.41}
\end{equation*}
$$

Let $\vartheta$ be the $\mathfrak{g}$-valued one-form such that $\vartheta(\hat{X})=X$, for all $X \in \mathfrak{g}$. Since $\vartheta$ is right invariant, the two-form $d \vartheta$ satisfies the identity

$$
\begin{equation*}
d \vartheta=\vartheta \wedge \vartheta \tag{3.42}
\end{equation*}
$$

where $\vartheta \wedge \vartheta(X, Y):=[\vartheta(X), \vartheta(Y)]$ for any $X, Y \in \mathfrak{g}$; see for example [1]. Let $\theta_{t}:=$ $\vartheta\left(\left.\frac{\partial}{\partial s}\right|_{0} \eta_{t}^{s}\right)$, so that $\left.\frac{\partial}{\partial s}\right|_{0} \eta_{t}^{s}=R_{\eta_{t} *} \theta_{t}$. Thus,

$$
\left.\frac{\partial}{\partial s}\right|_{0} \operatorname{Ad}_{\eta_{t}^{s}}=d \operatorname{Ad}\left(\theta_{t}\right) \operatorname{Ad}_{\eta_{t}}=\operatorname{ad}_{\theta_{t}} \operatorname{Ad}_{\eta_{t}}
$$

By Equation (3.41), we then have

$$
\begin{aligned}
\theta_{t}= & \vartheta\left(\left.\frac{\partial}{\partial s}\right|_{0} ^{\eta_{t}^{s}}\right)=\left.\frac{\partial}{\partial s}\right|_{0} \int_{0}^{t} \vartheta\left(d \eta_{\tau}^{s}\right)-\int_{0}^{t}\left[\vartheta\left(\left.\frac{\partial}{\partial s}\right|_{0} \eta_{t}^{s}\right), \vartheta\left(d \eta_{\tau}\right)\right] \\
= & \left.\frac{\partial}{\partial s}\right|_{0} \int_{0}^{t} u\left(\eta_{t}^{s}\right) \operatorname{Ad}_{\eta_{t}^{s}}\left(\circ d \vec{b}_{\tau}+s d \vec{h}_{\tau}\right)-\int_{0}^{t}\left[\theta_{\tau}, u\left(\eta_{t}\right) \operatorname{Ad}_{\eta_{t}} \circ d \vec{b}_{\tau}\right] \\
= & \int_{0}^{t}\left(\left\langle\hat{\nabla} u\left(\eta_{\tau}\right), \theta_{\tau}\right\rangle \operatorname{Ad}_{\eta_{\tau}} \circ d \vec{b}_{\tau}+u\left(\eta_{\tau}\right) \operatorname{ad}_{\theta_{\tau}} \operatorname{Ad}_{\eta_{\tau}} \circ d \vec{b}_{\tau}+u\left(\eta_{t}\right) \operatorname{Ad}_{\eta_{t}} d \vec{h}_{\tau}\right) \\
& -\int_{0}^{t}\left[\theta_{\tau}, u\left(\eta_{t}\right) \operatorname{Ad}_{\eta_{t}} \circ d \vec{b}_{\tau}\right] \\
= & \int_{0}^{t}\left(\left\langle\hat{\nabla} u\left(\eta_{\tau}\right), \theta_{\tau}\right\rangle \operatorname{Ad}_{\eta_{\tau}} \circ d \vec{b}_{\tau}+u\left(\eta_{\tau}\right) \operatorname{Ad}_{\eta_{\tau}} d \vec{h}_{\tau}\right),
\end{aligned}
$$

and for any $f \in C_{c}^{\infty}(G)$,

$$
\partial_{h} f\left(\eta_{t}\right)=\left.\frac{\partial}{\partial s}\right|_{0} f\left(\eta_{t}^{s}\right)=\left\langle\hat{\nabla} f\left(\eta_{t}\right), \vartheta\left(\left.\frac{\partial}{\partial s}\right|_{0} \eta_{t}^{s}\right)\right\rangle=\left\langle\hat{\nabla} f\left(\eta_{t}\right), \theta_{t}\right\rangle .
$$

Now, to write this equation in Itô form, first note that

$$
\begin{align*}
d \operatorname{Ad}_{\eta} & =u(\eta) \operatorname{Ad}_{\eta} \circ \operatorname{ad}_{d \vec{b}} \\
& =u(\eta) \operatorname{Ad}_{\eta} \operatorname{ad}_{d \vec{b}}+\frac{1}{2}\left[\langle\tilde{\nabla} u(\eta), d \vec{b}\rangle \operatorname{Ad}_{\eta}+u^{2}(\eta) \operatorname{Ad}_{\eta} \operatorname{ad}_{d \vec{b}}\right] \cdot \operatorname{ad}_{d \vec{b}} \\
& =u(\eta) \operatorname{Ad}_{\eta} \operatorname{ad}_{d \vec{b}}+\frac{1}{2} \sum_{i=1}^{k}\left[\left\langle\tilde{\nabla} u(\eta), X_{i}\right\rangle \operatorname{Ad}_{\eta} \operatorname{ad}_{X_{i}}+u^{2}(\eta) \operatorname{Ad}_{\eta} \operatorname{ad}_{X_{i}}^{2}\right] d t \tag{3.43}
\end{align*}
$$

where $^{\operatorname{ad}_{d \vec{b}}}=\operatorname{ad}_{X_{i}} d b^{i}$. This then implies

$$
\begin{aligned}
d\left[\langle\hat{\nabla} u(\eta), \theta\rangle \operatorname{Ad}_{\eta}\right] & \cdot d \vec{b}=\left[\langle\tilde{\nabla} \hat{\nabla} u(\eta), d \vec{b} \otimes \theta\rangle \operatorname{Ad}_{\eta}+\langle\hat{\nabla} u(\eta), \theta\rangle u(\eta) \operatorname{Ad}_{\eta} \operatorname{ad}_{d \vec{b}}\right] d \vec{b} \\
= & \sum_{i=1}^{k}\left[\left\langle\tilde{\nabla} \hat{\nabla} u(\eta), X_{i} \otimes \theta\right\rangle \operatorname{Ad}_{\eta}+\langle\hat{\nabla} u(\eta), \theta\rangle u(\eta) \operatorname{Ad}_{\eta} \operatorname{ad}_{X_{i}}\right] X_{i} d t \\
& =\sum_{i=1}^{k}\left\langle\tilde{\nabla} \hat{\nabla} u(\eta), X_{i} \otimes \theta\right\rangle \operatorname{Ad}_{\eta} X_{i} d t
\end{aligned}
$$

wherein we have used $\operatorname{ad}_{X} X=[X, X]=0$. Therefore,

$$
\begin{aligned}
\langle\hat{\nabla} u(\eta), \theta\rangle \operatorname{Ad}_{\eta} \circ d \vec{b} & =\langle\hat{\nabla} u(\eta), \theta\rangle \operatorname{Ad}_{\eta} d \vec{b}+\frac{1}{2} d\left[\langle\hat{\nabla} u(\eta), \theta\rangle \operatorname{Ad}_{\eta}\right] \cdot d \vec{b} \\
& =\langle\hat{\nabla} u(\eta), \theta\rangle \operatorname{Ad}_{\eta} d \vec{b}+\frac{1}{2} \sum_{i=1}^{k}\left\langle\tilde{\nabla} \hat{\nabla} u(\eta), X_{i} \otimes \theta\right\rangle \operatorname{Ad}_{\eta} X_{i} d t,
\end{aligned}
$$

and $\theta$ satisfies the Itô differential equation

$$
\begin{equation*}
d \theta=\langle\hat{\nabla} u(\eta), \theta\rangle \operatorname{Ad}_{\eta} d \vec{b}+u(\eta) \operatorname{Ad}_{\eta} d \vec{h}+\frac{1}{2} \sum_{i=1}^{k}\left\langle\tilde{\nabla} \hat{\nabla} v(\eta), X_{i} \otimes \theta\right\rangle \operatorname{Ad}_{\eta} X_{i} d t . \tag{3.44}
\end{equation*}
$$

Notation 3.21. Let us now fix $\psi \in C_{c}^{\infty}(\operatorname{End}(\mathfrak{g}))$ with

$$
\left\langle\psi^{\prime}(x), A\right\rangle:=\left.\frac{d}{d t}\right|_{0} \psi(x+t A),
$$

for any $A \in \operatorname{End}(\mathfrak{g})$. Let $|g|$ denote the distance from a point $g \in G$ to $e$ with respect to the right invariant metric, and let $\left\{\varphi_{m}\right\}_{m=1}^{\infty} \in C_{c}^{\infty}(G)$ be chosen so that $\varphi_{m} \uparrow 1$ as $m \rightarrow \infty$ such that $\varphi_{m}(g)=1$ when $|g| \leq m$ and $\sup _{m} \sup _{g \in G}\left|\hat{\nabla}^{k} \varphi_{m}(g)\right|<\infty$ for $k=0,1,2, \ldots$ Take $v(g):=\psi\left(\operatorname{Ad}_{g}\right)$ and $u_{m}(g):=\varphi_{m}(g) v(g)$.

Remark 3.22. For $u_{m}=\varphi_{m} v$, we have the derivative formulae

$$
\hat{\nabla} u_{m}=\hat{\nabla} v \cdot \varphi_{m}+v \cdot \hat{\nabla} \varphi_{m}
$$

and

$$
\begin{aligned}
\tilde{\nabla} \hat{\nabla} u_{m} & =\tilde{\nabla}\left[\hat{\nabla} v \cdot \varphi_{m}+v \cdot \hat{\nabla} \varphi_{m}\right] \\
& =\tilde{\nabla} \hat{\nabla} v \cdot \varphi_{m}+\tilde{\nabla} v \otimes \hat{\nabla} \varphi_{m}+\tilde{\nabla} \varphi_{m} \otimes \hat{\nabla} v+v \cdot \tilde{\nabla} \hat{\nabla} \varphi_{m} .
\end{aligned}
$$

Proposition 3.23. Let $\eta^{m}:[0,1] \times \mathscr{W} \rightarrow G$ denote the solution to the equation

$$
\begin{align*}
d \eta^{m}=u_{m}\left(\eta^{m}\right) \eta^{m} \circ d \vec{b} & =\varphi_{m}\left(\eta^{m}\right) v\left(\eta^{m}\right) \eta^{m} \circ d \vec{b} \\
& =\varphi_{m}\left(\eta^{m}\right) \psi\left(\operatorname{Ad}_{\eta^{m}}\right) \eta^{m} \circ d \vec{b}, \text { with } \eta_{0}^{m}=e, \tag{3.45}
\end{align*}
$$

and $\eta:[0,1] \times \mathscr{W} \rightarrow G$ denote the solution to

$$
\begin{equation*}
d \eta=v(\eta) \eta \circ d \vec{b}=\psi\left(\operatorname{Ad}_{\eta}\right) \eta \circ d \vec{b}, \text { with } \eta_{0}=e . \tag{3.46}
\end{equation*}
$$

Then for all $f \in C_{c}^{\infty}(G)$,

$$
\lim _{m \rightarrow \infty} \mathbb{E} \sup _{\tau \leq 1}\left|f\left(\eta_{\tau}^{m}\right)-f\left(\eta_{\tau}\right)\right|^{p}=0
$$

for all $p \in(1, \infty)$.
Proof. By Lemma 3.19, Equation (3.45) has a global solution. Notice also that $\tilde{X} v(\eta)=$ $\left\langle\psi^{\prime}\left(\operatorname{Ad}_{\eta}\right), \operatorname{ad}_{X}\right\rangle$ is bounded, so that Equation(3.46) also has a global solution. Then by Equation (3.36),

$$
\begin{aligned}
d f\left(\eta^{m}\right) & =u_{m}\left(\eta^{m}\right)\left\langle\tilde{\nabla} f\left(\eta^{m}\right), d \vec{b}\right\rangle \\
& +\frac{1}{2} \sum_{i=1}^{k}\left[\left\langle\tilde{\nabla} u_{m}\left(\eta^{m}\right), X_{i}\right\rangle\left\langle\tilde{\nabla} f\left(\eta^{m}\right), X_{i}\right\rangle+u_{m}\left(\eta^{m}\right)\left\langle\tilde{\nabla}^{2} f\left(\eta^{m}\right), X_{i} \otimes X_{i}\right\rangle\right] d t,
\end{aligned}
$$

and similarly,

$$
\begin{aligned}
& d f(\eta)=v(\eta)\langle\tilde{\nabla} f(\eta), d \vec{b}\rangle \\
& \quad+\frac{1}{2} \sum_{i=1}^{k}\left[\left\langle\tilde{\nabla} v(\eta), X_{i}\right\rangle\left\langle\tilde{\nabla} f(\eta), X_{i}\right\rangle+v(\eta)\left\langle\tilde{\nabla}^{2} f(\eta), X_{i} \otimes X_{i}\right\rangle\right] d t .
\end{aligned}
$$

Thus,

$$
\begin{aligned}
d\left[f\left(\eta^{m}\right)-f(\eta)\right] & =u_{m}\left(\eta^{m}\right)\left\langle\tilde{\nabla} f\left(\eta^{m}\right), d \vec{b}\right\rangle-v(\eta)\langle\tilde{\nabla} f(\eta), d \vec{b}\rangle \\
& +\frac{1}{2} \sum_{i=1}^{k}\left[\left\langle\tilde{\nabla} u_{m}\left(\eta^{m}\right), X_{i}\right\rangle\left\langle\tilde{\nabla} f\left(\eta^{m}\right), X_{i}\right\rangle-\left\langle\tilde{\nabla} v(\eta), X_{i}\right\rangle\left\langle\tilde{\nabla} f(\eta), X_{i}\right\rangle\right. \\
& \left.+u_{m}\left(\eta^{m}\right)\left\langle\tilde{\nabla}^{2} f\left(\eta^{m}\right), X_{i} \otimes X_{i}\right\rangle-v(\eta)\left\langle\tilde{\nabla}^{2} f(\eta), X_{i} \otimes X_{i}\right\rangle\right] d t .
\end{aligned}
$$

We bound this expression by applying Proposition 3.11 to each term. For the first term, note that $u_{m} \rightarrow v$ boundedly, as $m \rightarrow \infty$, and this implies that

$$
\begin{aligned}
\mathbb{E} \mid \int_{0}^{t}\left[u _ { m } ( \eta ^ { m } ) \left\langle\tilde{\nabla} f\left(\eta^{m}\right)\right.\right. & , d \vec{b}\rangle-v(\eta)\langle\tilde{\nabla} f(\eta), d \vec{b}\rangle]\left.\right|^{p} \\
& \lesssim \sum_{i=1}^{k} \mathbb{E} \int_{0}^{t}\left|u_{m}\left(\eta^{m}\right)\left\langle\tilde{\nabla} f\left(\eta^{m}\right), X_{i}\right\rangle-v(\eta)\left\langle\tilde{\nabla} f(\eta), X_{i}\right\rangle\right|^{p} d \tau \\
& \lesssim\|\tilde{\nabla} f\|_{\infty} \mathbb{E} \int_{0}^{t}\left|u_{m}\left(\eta^{m}\right)-v(\eta)\right|^{p} d \tau \rightarrow 0,
\end{aligned}
$$

as $m \rightarrow \infty$, by the dominated convergence theorem. Similarly, for the second term, $\tilde{\nabla} u_{m} \rightarrow \tilde{\nabla} v$ boundedly, as $m \rightarrow \infty$, implies that

$$
\begin{aligned}
\mathbb{E} \mid \int_{0}^{t}[ & \left.\left\langle\tilde{\nabla} u_{m}\left(\eta^{m}\right), X_{i}\right\rangle\left\langle\tilde{\nabla} f\left(\eta^{m}\right), X_{i}\right\rangle-\left\langle\tilde{\nabla} v(\eta), X_{i}\right\rangle\left\langle\tilde{\nabla} f(\eta), X_{i}\right\rangle\right]\left.d \tau\right|^{p} \\
& \lesssim \mathbb{E} \int_{0}^{t}\left|\left\langle\tilde{\nabla} u_{m}\left(\eta^{m}\right), X_{i}\right\rangle\left\langle\tilde{\nabla} f\left(\eta^{m}\right), X_{i}\right\rangle-\left\langle\tilde{\nabla} v(\eta), X_{i}\right\rangle\left\langle\tilde{\nabla} f(\eta), X_{i}\right\rangle\right|^{p} d \tau \\
& \lesssim\|\tilde{\nabla} f\|_{\infty} \mathbb{E} \int_{0}^{t}\left|\tilde{\nabla} u_{m}\left(\eta^{m}\right)-\tilde{\nabla} v(\eta)\right|^{p} d \tau \rightarrow 0,
\end{aligned}
$$

as $m \rightarrow \infty$, by the dominated convergence theorem. Finally,

$$
\begin{aligned}
\mathbb{E} \mid \int_{0}^{t}\left[u_{m}\left(\eta^{m}\right)\right. & \left.\left\langle\tilde{\nabla}^{2} f\left(\eta^{m}\right), X_{i} \otimes X_{i}\right\rangle-v(\eta)\left\langle\tilde{\nabla}^{2} f(\eta), X_{i} \otimes X_{i}\right\rangle\right]\left.d \tau\right|^{p} \\
& =\mathbb{E} \int_{0}^{t}\left|u_{m}\left(\eta^{m}\right)\left\langle\tilde{\nabla}^{2} f\left(\eta^{m}\right), X_{i} \otimes X_{i}\right\rangle-v(\eta)\left\langle\tilde{\nabla}^{2} f(\eta), X_{i} \otimes X_{i}\right\rangle\right|^{p} d \tau \\
& \lesssim\left\|\tilde{\nabla}^{2} f\right\|_{\infty} \mathbb{E} \int_{0}^{t}\left|u_{m}\left(\eta^{m}\right)-v(\eta)\right|^{p} d \tau \rightarrow 0,
\end{aligned}
$$

as $m \rightarrow \infty$, again by dominated convergence. Thus,

$$
\lim _{m \rightarrow \infty} \mathbb{E} \sup _{\tau \leq 1}\left|f\left(\eta_{\tau}^{m}\right)-f\left(\eta_{\tau}\right)\right|^{p}=0
$$

as desired.
Proposition 3.24. Let $U_{t}^{m}=\operatorname{Ad}_{\eta_{t}^{m}}: \mathscr{W} \rightarrow \operatorname{End}(\mathfrak{g})$ and $U_{t}=\operatorname{Ad}_{\eta_{t}}: \mathscr{W} \rightarrow \operatorname{End}(\mathfrak{g})$, which satisfy the stochastic differential equations

$$
\begin{equation*}
d U^{m}=u_{m}\left(\eta^{m}\right) U^{m} \circ \operatorname{ad}_{d \vec{b}}=\varphi_{m}\left(\eta^{m}\right) \psi\left(U^{m}\right) U^{m} \circ \operatorname{ad}_{d \vec{b}} \text {, with } U_{0}^{m}=I \text {, } \tag{3.47}
\end{equation*}
$$

and

$$
\begin{equation*}
d U=v(\eta) U \circ \operatorname{ad}_{d \vec{b}}=\psi(U) U \circ \operatorname{ad}_{d \vec{b}}, \text { with } U_{0}=I, \tag{3.48}
\end{equation*}
$$

where $\operatorname{ad}_{d \vec{b}}=\operatorname{ad}_{X_{i}} d b^{i}$. Then

$$
\lim _{m \rightarrow \infty} \mathbb{E} \sup _{\tau \leq 1}\left|U_{\tau}^{m}-U_{\tau}\right|^{p}=0
$$

for all $p \in(1, \infty)$.
Proof. From Equation (3.43), we may rewrite (3.47) and (3.48) in Itô form as

$$
\begin{aligned}
d U^{m}= & u_{m}\left(\eta^{m}\right) U^{m} \operatorname{ad}_{d \vec{b}}+\frac{1}{2} \sum_{i=1}^{k}\left[\left\langle\tilde{\nabla} u_{m}\left(\eta^{m}\right), X_{i}\right\rangle U^{m} \operatorname{ad}_{X_{i}}+u_{m}^{2}\left(\eta^{m}\right) U^{m} \operatorname{ad}_{X_{i}}^{2}\right] d t \\
= & \varphi_{m}\left(\eta^{m}\right) \psi\left(U^{m}\right) U^{m} \operatorname{ad}_{d \vec{b}}+\frac{1}{2} \sum_{i=1}^{k}\left[\left\langle\tilde{\nabla} \varphi_{m}\left(\eta^{m}\right), X_{i}\right\rangle \psi\left(U^{m}\right) U^{m} \operatorname{ad}_{X_{i}}\right. \\
& \left.+\varphi_{m}\left(\eta^{m}\right) \psi\left(U^{m}\right)\left\langle\psi^{\prime}\left(U^{m}\right), U^{m} \operatorname{ad}_{X_{i}}\right\rangle U^{m} \operatorname{ad}_{X_{i}}+\varphi_{m}^{2}\left(\eta^{m}\right) \psi^{2}\left(U^{m}\right) U^{m} \operatorname{ad}_{X_{i}}^{2}\right] d t
\end{aligned}
$$

and

$$
\begin{aligned}
d U & =v(\eta) U \operatorname{ad}_{d \vec{b}}+\frac{1}{2} \sum_{i=1}^{k}\left[\left\langle\tilde{\nabla} v(\eta), X_{i}\right\rangle U \operatorname{ad}_{X_{i}}+v^{2}(\eta) U \operatorname{ad}_{X_{i}}^{2}\right] d t \\
& =\psi(U) U \operatorname{ad}_{d \vec{b}}+\frac{1}{2} \sum_{i=1}^{k}\left[\psi(U)\left\langle\psi^{\prime}(U), U \operatorname{ad}_{X_{i}}\right\rangle U \operatorname{ad}_{X_{i}}+\psi^{2}(U) U \operatorname{ad}_{X_{i}}^{2}\right] d t .
\end{aligned}
$$

Let $\varepsilon^{m}:=U^{m}-U$. Then by the above,

$$
\begin{array}{r}
d \varepsilon^{m}=\left(\varphi_{m}\left(\eta^{m}\right) \psi\left(U^{m}\right) U^{m}-\psi(U) U\right) \operatorname{ad}_{d \vec{b}}+\frac{1}{2} \sum_{i=1}^{k}\left[\left\langle\tilde{\nabla} \varphi_{m}\left(\eta^{m}\right), X_{i}\right\rangle \psi\left(U^{m}\right) U^{m} \operatorname{ad}_{X_{i}}\right. \\
+\left(\varphi_{m}\left(\eta^{m}\right) \psi\left(U^{m}\right)\left\langle\psi^{\prime}\left(U^{m}\right), U^{m} \operatorname{ad}_{X_{i}}\right\rangle U^{m}-\psi(U)\left\langle\psi^{\prime}(U), U \operatorname{ad}_{X_{i}}\right\rangle U\right) \operatorname{ad}_{X_{i}} \\
 \tag{3.49}\\
\left.+\left(\varphi_{m}^{2}\left(\eta^{m}\right) \psi^{2}\left(U^{m}\right) U^{m}-\psi^{2}(U) U\right) \operatorname{ad}_{X_{i}}^{2}\right] d t .
\end{array}
$$

Again applying Proposition 3.11, we work term by term to bound the above expression. Note first that, since $\psi$ has compact support, $U$ and $U^{m}$ always remain in a fixed
compact subset of $\operatorname{End}(\mathfrak{g})$. Thus,

$$
\begin{aligned}
\mathbb{E} \mid \int_{0}^{t}\left[\varphi_{m}\left(\eta^{m}\right)\right. & \left.\psi\left(U^{m}\right) U^{m}-\psi(U) U\right]\left.\mathrm{ad}_{d \vec{b}}\right|^{p} \\
& \lesssim \mathbb{E} \int_{0}^{t}\left|\varphi_{m}\left(\eta^{m}\right) \psi\left(U^{m}\right) U^{m}-\psi(U) U\right|^{p} d \tau \\
& \lesssim \mathbb{E} \int_{0}^{t}\left|\psi\left(U^{m}\right) U^{m}-\psi(U) U\right|^{p} d \tau+\mathbb{E} \int_{0}^{t}\left|\varphi_{m}\left(\eta^{m}\right)-1\right|^{p}\left|\psi\left(U^{m}\right)\right|^{p} d \tau \\
& \lesssim \mathbb{E} \int_{0}^{t}\left|\varepsilon^{m}\right|^{p} d \tau+\delta_{m}
\end{aligned}
$$

wherein we have applied the mean value inequality to $x \mapsto \psi(x) x$ to learn

$$
\left|\psi\left(U^{m}\right) U^{m}-\psi(U) U\right| \leq C(\psi)\left|U^{m}-U\right|=C\left|\varepsilon^{m}\right|,
$$

and

$$
\delta_{m}=\mathbb{E} \int_{0}^{t}\left|\varphi_{m}\left(\eta^{m}\right)-1\right|^{p}\left|\psi\left(U^{m}\right)\right|^{p} d \tau \rightarrow 0
$$

as $m \rightarrow \infty$, by the dominated convergence theorem. Similarly, for the last term of the sum in (3.49),

$$
\begin{aligned}
\mathbb{E} \mid \int_{0}^{t} & \left.\left(\varphi_{m}^{2}\left(g^{n}\right) \psi^{2}\left(U^{m}\right) U^{m}-\psi^{2}(U) U\right) \operatorname{ad}_{X_{i}}^{2} d \tau\right|^{p} \\
& \lesssim \mathbb{E} \int_{0}^{t}\left|\psi^{2}\left(U^{m}\right) U^{m}-\psi^{2}(U) U\right|^{p} d \tau+\mathbb{E} \int_{0}^{t}\left|\varphi_{m}^{2}\left(\eta^{m}\right)-1\right|^{p}\left|\psi\left(U^{m}\right)\right|^{2 p} d \tau \\
& \lesssim \mathbb{E} \int_{0}^{t}\left|\varepsilon^{m}\right|^{p} d \tau+\delta_{m},
\end{aligned}
$$

where the mean value inequality has now been applied to the function $x \mapsto \psi^{2}(x) x$, and

$$
\delta_{m}=\mathbb{E} \int_{0}^{t}\left|\varphi_{m}^{2}\left(\eta^{m}\right)-1\right|^{p}\left|\psi\left(U^{m}\right)\right|^{2 p} d \tau \rightarrow 0,
$$

as $m \rightarrow \infty$. For the second term,

$$
\begin{aligned}
& \mathbb{E}\left|\int_{0}^{t}\left\langle\tilde{\nabla} \varphi_{m}\left(\eta^{m}\right), X_{i}\right\rangle \psi\left(U^{m}\right) U^{m} \operatorname{ad}_{X_{i}} d \tau\right|^{p} \lesssim \mathbb{E} \int_{0}^{t}\left|\left\langle\tilde{\nabla} \varphi_{m}\left(\eta^{m}\right), X_{i}\right\rangle \psi\left(U^{m}\right) U^{m}\right|^{p} d \tau \\
&=\mathbb{E} \int_{0}^{t}\left|\left\langle\tilde{\nabla} \varphi_{m}\left(\eta^{m}\right), X_{i}\right\rangle \psi\left(U^{m}\right)\left(\varepsilon^{m}+U\right)\right|^{p} d \tau \\
& \lesssim \mathbb{E} \int_{0}^{t}\left|\varepsilon^{m}\right|^{p} d \tau+\delta_{m},
\end{aligned}
$$

where

$$
\delta_{m}=\mathbb{E} \int_{0}^{t}\left|\left\langle\tilde{\nabla} \varphi_{m}\left(\eta^{m}\right), X_{i}\right\rangle \psi\left(U^{m}\right) U\right|^{p} d \tau \rightarrow 0
$$

as $m \rightarrow \infty$, since $\lim _{m \rightarrow \infty}\left|\tilde{\nabla} \varphi_{m}\right|=0$. Finally, for the the third term, note first that

$$
\begin{aligned}
\varphi_{m}\left(g^{n}\right)\left\langle\psi^{\prime}\right. & \left.\left(U^{m}\right), U^{m} \operatorname{ad}_{X_{i}}\right\rangle U^{m} \\
= & \varphi_{m}\left(\eta^{m}\right)\left\langle\psi^{\prime}\left(U^{m}\right), U^{m} \operatorname{ad}_{X_{i}}\right\rangle\left(\varepsilon^{m}+U\right) \\
= & \varphi_{m}\left(\eta^{m}\right)\left\langle\psi^{\prime}\left(U^{m}\right), U^{m} \operatorname{ad}_{X_{i}}\right\rangle \varepsilon^{m}+\varphi_{m}\left(\eta^{m}\right)\left\langle\psi^{\prime}\left(U^{m}\right),\left(\varepsilon^{m}+U\right) \operatorname{ad}_{X_{i}}\right\rangle U \\
= & \varphi_{m}\left(\eta^{m}\right)\left\langle\psi^{\prime}\left(U^{m}\right), U^{m} \operatorname{ad}_{X_{i}}\right\rangle \varepsilon^{m}+\varphi_{m}\left(\eta^{m}\right)\left\langle\psi^{\prime}\left(U^{m}\right), \varepsilon^{m} \operatorname{ad}_{X_{i}}\right\rangle U \\
& \quad+\varphi_{m}\left(\eta^{m}\right)\left\langle\psi^{\prime}\left(U^{m}\right), U \operatorname{ad}_{X_{i}}\right\rangle U
\end{aligned}
$$

Thus,

$$
\begin{aligned}
& \mathbb{E}\left|\int_{0}^{t}\left(\varphi_{m}\left(\eta^{m}\right) \psi\left(U^{m}\right)\left\langle\psi^{\prime}\left(U^{m}\right), U^{m} \operatorname{ad}_{X_{i}}\right\rangle U^{m}-\psi(U)\left\langle\psi^{\prime}(U), U \operatorname{ad}_{X_{i}}\right\rangle U\right) \operatorname{ad}_{X_{i}} d \tau\right|^{p} \\
& \lesssim \mathbb{E} \int_{0}^{t} \mid \varphi_{m}\left(\eta^{m}\right)\left\langle\psi^{\prime}\left(U^{m}\right), U^{m} \operatorname{ad}_{X_{i}}\right\rangle \varepsilon^{m}+\varphi_{m}\left(\eta^{m}\right)\left\langle\psi^{\prime}\left(U^{m}\right), \varepsilon^{m} \operatorname{ad}_{X_{i}}\right\rangle U \\
& \quad+\varphi_{m}\left(\eta^{m}\right)\left\langle\psi^{\prime}\left(U^{m}\right), U \operatorname{ad}_{X_{i}}\right\rangle U-\left.\left\langle\psi^{\prime}(U), U \operatorname{ad}_{X_{i}}\right\rangle U\right|^{p} d \tau \\
& \lesssim \mathbb{E} \int_{0}^{t}\left|\varepsilon^{m}\right|^{p} d \tau+\delta_{m}
\end{aligned}
$$

where

$$
\delta_{m}=\mathbb{E} \int_{0}^{t}\left|\varphi_{m}\left(\eta^{m}\right)\left\langle\psi^{\prime}\left(U^{m}\right), U \operatorname{ad}_{X_{i}}\right\rangle-\left\langle\psi^{\prime}(U), U \operatorname{ad}_{X_{i}}\right\rangle\right|^{p}|U|^{p} d \tau \rightarrow 0
$$

as $m \rightarrow \infty$, since $\varphi_{m}\left(\eta^{m}\right) \psi^{\prime}\left(U^{m}\right) \rightarrow \psi^{\prime}(U)$ boundedly. These bounds then imply that

$$
\mathbb{E} \sup _{\tau \leq t}\left|\varepsilon_{\tau}^{m}\right|^{p} \leq C \int_{0}^{t}\left|\varepsilon^{m}\right|^{p} d \tau+\delta_{m}
$$

for all $t \in[0,1]$. Thus, by Gronwall's inequality, we have

$$
\mathbb{E} \sup _{\tau \leq 1}\left|U_{\tau}^{m}-U_{\tau}\right|^{p}=\mathbb{E} \sup _{\tau \leq 1}\left|\varepsilon_{\tau}^{m}\right|^{p} \leq \delta_{m} e^{C} \rightarrow 0
$$

as $m \rightarrow \infty$.

Proposition 3.25. Let $\theta_{t}^{m}: \mathscr{W} \rightarrow \mathfrak{g}$ be as in Equation (3.38) with $u$ replaced by $u_{m}$; that is,

$$
\begin{equation*}
\theta_{t}^{m}=\int_{0}^{t}\left(\left\langle\hat{\nabla} u_{m}\left(\eta^{m}\right), \theta\right\rangle \operatorname{Ad}_{\eta^{m}} \circ d \vec{b}+u_{m}\left(\eta^{m}\right) \operatorname{Ad}_{\eta^{m}} d \vec{h}\right) \tag{3.50}
\end{equation*}
$$

Then

$$
\lim _{m \rightarrow \infty} \mathbb{E} \sup _{\tau \leq 1}\left|\theta_{\tau}^{m}-\theta_{\tau}\right|^{p}=0
$$

for all $p \in(1, \infty)$, where $\theta_{t}: \mathscr{W} \rightarrow \mathfrak{g}$ is the solution to

$$
\begin{equation*}
\theta_{t}=\int_{0}^{t}\left(\langle\hat{\nabla} v(\eta), \theta\rangle \operatorname{Ad}_{\eta} \circ d \vec{b}+v(\eta) \operatorname{Ad}_{\eta} d \vec{h}\right) \tag{3.51}
\end{equation*}
$$

Proof. Let $U_{t}^{m}=\operatorname{Ad}_{\eta_{t}^{m}}$ and $U_{t}=\operatorname{Ad}_{\eta_{t}}$. We rewrite Equation (3.50) in Itô form as $d \theta^{m}=\left\langle\hat{\nabla} u_{m}\left(\eta^{m}\right), \theta^{m}\right\rangle U^{m} d \vec{b}+u_{m}\left(\eta^{m}\right) U^{m} d \vec{h}+\frac{1}{2} \sum_{i=1}^{k}\left\langle\tilde{\nabla} \hat{\nabla} u_{m}\left(\eta^{m}\right), X_{i} \otimes \theta^{m}\right\rangle U^{m} X_{i} d t$.
Note that, formally, $\theta$ is the solution to Equation (3.38) with $u$ replaced by $v$ (although $v$ is not a function with compact support), and we may rewrite Equation (3.51) in Itô form as

$$
\begin{aligned}
d \theta & =\langle\hat{\nabla} v(\eta), \theta\rangle \operatorname{Ad}_{\eta} d \vec{b}+v(\eta) \operatorname{Ad}_{\eta} d \vec{h}+\frac{1}{2} \sum_{i=1}^{k}\left\langle\tilde{\nabla} \hat{\nabla} v(\eta), X_{i} \otimes \theta\right\rangle \operatorname{Ad}_{\eta} X_{i} d t \\
& =\langle\hat{\nabla} v(\eta), \theta\rangle U d \vec{b}+v(\eta) U d \vec{h}+\frac{1}{2} \sum_{i=1}^{k}\left\langle\tilde{\nabla} \hat{\nabla} v(\eta), X_{i} \otimes \theta\right\rangle U X_{i} d t
\end{aligned}
$$

Let $\varepsilon^{m}:=\theta^{m}-\theta$ (so that $\theta^{m}=\varepsilon^{m}+\theta$ ). Then

$$
\begin{aligned}
& d \varepsilon^{m}=\left[\left\langle\hat{\nabla} u_{m}\left(\eta^{m}\right), \theta^{m}\right\rangle U^{m}-\langle\hat{\nabla} v(\eta), \theta\rangle U\right] d \vec{b}+\left[u_{m}\left(\eta^{m}\right) U^{m}-v(\eta) U\right] d \vec{h} \\
& \quad+\frac{1}{2} \sum_{i=1}^{k}\left[\left\langle\tilde{\nabla} \hat{\nabla} u_{m}\left(\eta^{m}\right), X_{i} \otimes \theta^{m}\right\rangle U^{m}-\left\langle\tilde{\nabla} \hat{\nabla} v(\eta), X_{i} \otimes \theta\right\rangle U\right] X_{i} d t .
\end{aligned}
$$

Considering the first term of this expression, we have

$$
\begin{aligned}
\left\langle\hat{\nabla} u_{m}\left(\eta^{m}\right), \theta^{m}\right\rangle U^{m}- & \langle\hat{\nabla} v(\eta), \theta\rangle U \\
& =\left\langle\hat{\nabla} u_{m}\left(\eta^{m}\right), \theta+\varepsilon^{m}\right\rangle U^{m}-\langle\hat{\nabla} v(\eta), \theta\rangle U \\
& =\left\langle\hat{\nabla} u_{m}\left(\eta^{m}\right), \varepsilon^{m}\right\rangle U^{m}+\left\langle\hat{\nabla} u_{m}\left(\eta^{m}\right), \theta\right\rangle U^{m}-\langle\hat{\nabla} v(\eta), \theta\rangle U
\end{aligned}
$$

Using again the $U$ and $U^{m}$ remain in a fixed compact subset of $\operatorname{End}(\mathfrak{g})$ and the fact that $\hat{\nabla} u_{\left(\eta^{m}\right)} \rightarrow \hat{\nabla} v(\eta)$ boundedly, we have

$$
\mathbb{E}\left|\int_{0}^{t}\left[\left\langle\hat{\nabla} u_{m}\left(\eta^{m}\right), \theta^{m}\right\rangle U^{m}-\langle\hat{\nabla} v(\eta), \theta\rangle U\right] d \vec{b}\right|^{p} \lesssim \mathbb{E} \int_{0}^{t}\left|\varepsilon^{m}\right|^{p} d \tau+\delta_{m},
$$

where

$$
\delta_{m}=\mathbb{E} \int_{0}^{t}\left|\left\langle\hat{\nabla} u_{m}\left(\eta^{m}\right), \theta\right\rangle U^{m}-\langle\hat{\nabla} v(\eta), \theta\rangle U\right| d \tau \rightarrow 0
$$

as $m \rightarrow \infty$. The second term converges to 0 since

$$
\mathbb{E}\left|\int_{0}^{t}\left[u_{m}\left(\eta^{m}\right) U^{m}-v(\eta) U\right] d \vec{h}\right|^{p}=\mathbb{E}\left|\int_{0}^{t}\left[\varphi_{m}\left(\eta^{m}\right) v\left(\eta^{m}\right) U^{m}-v(\eta) U\right] d \vec{h}\right|^{p} \rightarrow 0
$$

as $m \rightarrow \infty$, by the dominated convergence theorem. For the third term, note that

$$
\begin{aligned}
\mid\left[\left\langle\tilde{\nabla} \hat{\nabla} u_{m}\left(\eta^{m}\right),\right.\right. & \left.\left.X_{i} \otimes\left(\theta+\varepsilon^{m}\right)\right\rangle U^{m}-\left\langle\tilde{\nabla} \hat{\nabla} v(\eta), X_{i} \otimes \theta\right\rangle U\right] X_{i} \mid \\
& \lesssim\left|\varepsilon^{m}\right|+\left|\left[\left\langle\tilde{\nabla} \hat{\nabla} u_{m}\left(\eta^{m}\right), X_{i} \otimes \theta\right\rangle U^{m}-\left\langle\tilde{\nabla} \hat{\nabla} v(\eta), X_{i} \otimes \theta\right\rangle U\right] X_{i}\right|,
\end{aligned}
$$

and so

$$
\begin{aligned}
& \mathbb{E}\left|\int_{0}^{t}\left[\left\langle\tilde{\nabla} \hat{\nabla} u_{m}\left(\eta^{m}\right), X_{i} \otimes \theta^{m}\right\rangle U^{m}-\left\langle\tilde{\nabla} \hat{\nabla} v(\eta), X_{i} \otimes \theta\right\rangle U\right] X_{i} d \tau\right|^{p} \\
& \lesssim \mathbb{E} \int_{0}^{t}\left|\varepsilon^{m}\right|^{p} d \tau+\delta_{m}
\end{aligned}
$$

where

$$
\delta_{m}=\mathbb{E}\left|\int_{0}^{t}\left[\left\langle\tilde{\nabla} \hat{\nabla} u_{m}\left(\eta^{m}\right), X_{i} \otimes \theta\right\rangle U^{m}-\left\langle\tilde{\nabla} \hat{\nabla} v(\eta), X_{i} \otimes \theta\right\rangle U\right] X_{i} d \tau\right|^{p} \rightarrow 0
$$

as $m \rightarrow \infty$, where we have used that $\tilde{\nabla} \hat{\nabla} u_{m}\left(\eta^{m}\right) \rightarrow \tilde{\nabla} \hat{\nabla} v(\eta)$ boundedly to apply the dominated convergence theorem.

Putting these bounds together then shows

$$
\mathbb{E} \sup _{\tau \leq t}\left|\varepsilon_{\tau}^{m}\right|^{p} \leq C \mathbb{E} \int_{0}^{t}\left|\varepsilon^{m}\right|^{p} d \tau+\delta_{m}
$$

for all $t \in[0,1]$, and again applying Gronwall's inequality gives

$$
\mathbb{E} \sup _{\tau \leq 1}\left|\theta_{\tau}^{m}-\theta_{\tau}\right|^{p}=\mathbb{E} \sup _{\tau \leq 1}\left|\varepsilon_{\tau}^{m}\right|^{p} \leq \delta_{m} e^{C} \rightarrow 0
$$

as $m \rightarrow \infty$, finishes the proof.

Proposition 3.26. Let $\eta^{m}$ be the solution to Equation (3.45) and $\theta^{m}$ be the solution to Equation (3.50). Then, for any $h \in \mathscr{H}, f \in C_{c}^{\infty}(G)$, and $t \in[0,1]$,

$$
\partial_{h} f\left(\eta_{t}^{m}\right)=\left\langle\hat{\nabla} f\left(\eta_{t}^{m}\right), \theta_{t}^{m}\right\rangle .
$$

Furthermore, for $\eta$ the solution to (3.46) and $\theta$ the solution to (3.51),

$$
\lim _{m \rightarrow \infty} \mathbb{E}\left|\partial_{h} f\left(\eta_{t}^{m}\right)-\left\langle\hat{\nabla} f\left(\eta_{t}\right), \theta_{t}\right\rangle\right|^{p}=0
$$

for all $p \in(1, \infty)$.
Proof. The first claim follows immediately from Proposition 3.20. Now note that

$$
\begin{aligned}
\mid \partial_{h} f\left(\eta^{m}\right)- & \langle\hat{\nabla} f(\eta), \theta\rangle\left|=\left|\left\langle\hat{\nabla} f\left(\eta^{m}\right), \theta^{m}\right\rangle-\langle\hat{\nabla} f(\eta), \theta\rangle\right|\right. \\
& \leq\left|\left\langle\hat{\nabla} f\left(\eta^{m}\right), \theta^{m}\right\rangle-\left\langle\hat{\nabla} f\left(\eta^{m}\right), \theta\right\rangle\right|+\left|\left\langle\hat{\nabla} f\left(\eta^{m}\right), \theta\right\rangle-\langle\hat{\nabla} f(\eta), \theta\rangle\right| \\
\leq & |\hat{\nabla} f|\left|\theta^{m}-\theta\right|+\left|\hat{\nabla} f\left(\eta^{m}\right)-\hat{\nabla} f(\eta)\right||\theta|
\end{aligned}
$$

Thus,

$$
\lim _{m \rightarrow \infty} \mathbb{E}\left|\partial_{h} f\left(\eta^{m}\right)-\langle\hat{\nabla} f(\eta), \theta\rangle\right|^{p} \leq \lim _{m \rightarrow \infty} \mathbb{E}\left[|\hat{\nabla} f|\left|\theta^{m}-\theta\right|+\left|\hat{\nabla} f\left(\eta^{m}\right)-\hat{\nabla} f(\eta)\right||\theta|\right]^{p}=0
$$

by Propositions 3.23 and 3.25 and the dominated convergence theorem.
Corollary 3.27. For any $h \in \mathscr{H}, f \in C_{c}^{\infty}(G)$, and $t \in[0,1], f\left(\eta_{t}\right) \in \operatorname{Dom}\left(\partial_{h}\right)$ and

$$
\partial_{h} f\left(\eta_{t}\right)=\left\langle\hat{\nabla} f\left(\eta_{t}\right), \theta_{t}\right\rangle \in L^{\infty-}(\mu)
$$

This corollary follows from $\partial_{h}$ being a closed operator (we are thinking here that $\partial_{h}=\bar{\partial}_{h}$ ). Now we are able to remove the cutoff functions and prove the primary result of this section.

Theorem 3.28. Let $\xi:[0,1] \times \mathscr{W} \rightarrow G$ denote the solution to Equation (3.3), and let $\Theta: \mathscr{W} \rightarrow \mathfrak{g}$ be the solution to

$$
\begin{equation*}
\Theta_{t}:=\int_{0}^{t} \operatorname{Ad}_{\xi} d \vec{h}=\int_{0}^{t} \operatorname{Ad}_{\xi} X_{i} \dot{h}^{i} d \tau \tag{3.52}
\end{equation*}
$$

Then, for any $f \in C_{c}^{\infty}(G)$ and $t \in[0,1], f\left(\xi_{t}\right) \in \mathcal{D}^{1, \infty}$, and for any $h \in \mathscr{H}$,

$$
\begin{equation*}
\partial_{h}\left[f\left(\xi_{t}\right)\right]=\int_{0}^{t}\left\langle\hat{\nabla} f\left(\xi_{t}\right), \operatorname{Ad}_{\xi_{\tau}} X_{i}\right\rangle \dot{h}_{\tau}^{i} d \tau . \tag{3.53}
\end{equation*}
$$

Thus,

$$
\left(D\left[f\left(\xi_{t}\right)\right]\right)^{i}=\int_{0}^{\cdot \wedge t}\left\langle\hat{\nabla} f\left(\xi_{t}\right), \operatorname{Ad}_{\xi_{\tau}} X_{i}\right\rangle d \tau
$$

componentwise in $\mathscr{H}$.
Proof. Let $\psi \in C_{c}^{\infty}(\operatorname{End}(\mathfrak{g}),[0,1])$ be chosen so that $\psi=1$ near $I$ and $\psi(x)=0$ if $|x| \geq 2$, where $|\cdot|$ is the distance from $I$ with respect to any metric on $\operatorname{End}(\mathfrak{g})$. Define $v_{n}(x):=\psi\left(n^{-1} \operatorname{Ad}_{x}\right)$, and let $\xi^{n}:[0,1] \times \mathscr{W} \rightarrow G$ be the solution to the Stratonovich equation

$$
d \xi^{n}=v_{n}\left(\xi^{n}\right) \xi^{n} \circ d \vec{b}, \text { with } \xi_{0}^{n}=e .
$$

By Lemma 3.19, $\xi^{n}$ exists for all time $t \in[0,1]$. Noting that $v_{n} \rightarrow 1$ and $\tilde{\nabla} v_{n} \rightarrow 0$ boundedly as $n \rightarrow \infty$, by an argument identical to that in Proposition 3.23, we may show that, for all $f \in C_{c}^{\infty}(G)$,

$$
\lim _{n \rightarrow \infty} \mathbb{E} \sup _{\tau \leq 1}\left|f\left(\xi_{\tau}^{n}\right)-f\left(\xi_{\tau}\right)\right|^{p}=0
$$

for all $p \in(1, \infty)$. Note that this convergence implies that $f\left(\xi_{t}\right) \in L^{\infty-}(\mu)$.
Now let $\Theta_{t}^{n}: \mathscr{W} \rightarrow \mathfrak{g}$ denote the solution to the Itô equation

$$
\begin{align*}
& d \Theta^{n}=\left\langle\hat{\nabla} v_{n}\left(\xi^{n}\right), \Theta^{n}\right\rangle \operatorname{Ad}_{\xi^{n}} d \vec{b}+v_{n}\left(\xi^{n}\right) \operatorname{Ad}_{\xi^{n}} d \vec{h} \\
&+\frac{1}{2} \sum_{i=1}^{k}\left\langle\tilde{\nabla} \hat{\nabla} v_{n}\left(\xi^{n}\right), X_{i} \otimes \Theta^{n}\right\rangle \operatorname{Ad}_{\xi^{n}} X_{i} d t, \tag{3.54}
\end{align*}
$$

with $\Theta_{0}^{n}=0$. We will now show that

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \mathbb{E} \sup _{\tau \leq 1}\left|\Theta_{\tau}^{n}-\Theta_{\tau}\right|^{p}=0 \tag{3.55}
\end{equation*}
$$

for all $p \in(1, \infty)$. So let $W^{n}=\operatorname{Ad}_{\xi^{n}}$ and $W=\operatorname{Ad}_{\xi}$. Then $W^{n}, W: \mathscr{W} \rightarrow \operatorname{End}(\mathfrak{g})$ satisfy the equations

$$
d W^{n}=v_{n}\left(\xi^{n}\right) W^{n} \circ \operatorname{ad}_{d \vec{b}}, \text { with } W_{0}^{n}=I,
$$

and

$$
d W=W \circ \operatorname{ad}_{d \vec{b}}, \text { with } W_{0}=I
$$

Rewriting Equation (3.54), $\Theta^{n}$ solves

$$
d \Theta^{n}=\left\langle\hat{\nabla} v_{n}\left(\xi^{n}\right), \Theta^{n}\right\rangle W^{n} d \vec{b}+v_{n}\left(\xi^{n}\right) W^{n} d \vec{h}+\frac{1}{2} \sum_{i=1}^{k}\left\langle\tilde{\nabla} \hat{\nabla} v_{n}\left(\xi^{n}\right), X_{i} \otimes \Theta^{n}\right\rangle W^{n} X_{i} d t
$$

and Equation (3.52) implies that

$$
d \Theta=W d \vec{h}
$$

Thus, for $\varepsilon^{n}:=\Theta^{n}-\Theta$,

$$
\begin{align*}
d \varepsilon^{n}=\left\langle\hat{\nabla} v_{n}\left(\xi^{n}\right), \Theta^{n}\right\rangle W^{n} d \vec{b}+\left(v_{n}\left(\xi^{n}\right)\right. & \left.W^{n}-W\right) d \vec{h} \\
& +\frac{1}{2} \sum_{i=1}^{k}\left\langle\tilde{\nabla} \hat{\nabla} v_{n}\left(\xi^{n}\right), X_{i} \otimes \Theta^{n}\right\rangle W^{n} X_{i} d t \tag{3.56}
\end{align*}
$$

Since $W$ and $W^{n}$ are the solutions to Equations (3.16) and (3.28), which have smooth, bounded coefficients, Theorem V-10.1 of Ikeda and Watanabe [30] implies that $W$ and $W^{n}$ are (componentwise) in the domain $\mathcal{D}^{\infty}$. By Propositions 3.16 and 3.18 ,

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \mathbb{E} \sup _{\tau \leq 1}\left|W_{\tau}^{n}-W_{\tau}\right|^{p} \tag{3.57}
\end{equation*}
$$

and

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \mathbb{E} \sup _{\tau \leq 1}\left|\partial_{h} W_{\tau}^{n}-\partial_{h} W_{\tau}\right|^{p} \tag{3.58}
\end{equation*}
$$

for all $p \in(1, \infty)$. Furthermore,

$$
\partial_{h} W^{n}=\partial_{h} \operatorname{Ad}_{\xi^{n}}=\operatorname{ad}_{\Theta^{n}} \operatorname{Ad}_{\xi^{n}}=\operatorname{ad}_{\Theta^{n}} W^{n} .
$$

This then implies that

$$
\begin{aligned}
\left\langle\hat{\nabla} v_{n}\left(\xi^{n}\right), \Theta^{n}\right\rangle & =\left.\frac{d}{d t}\right|_{0} \psi\left(n^{-1} \operatorname{Ad}_{e^{t \Theta^{n}} \xi^{n}}\right) \\
& =\frac{1}{n}\left\langle\psi^{\prime}\left(n^{-1} \operatorname{Ad}_{\xi^{n}}\right), \operatorname{ad}_{\Theta^{n}} \operatorname{Ad}_{\xi^{n}}\right\rangle=\frac{1}{n}\left\langle\psi^{\prime}\left(n^{-1} W^{n}\right), \partial_{h} W^{n}\right\rangle .
\end{aligned}
$$

Thus, for the first term of Equation (3.56),

$$
\begin{aligned}
& \mathbb{E}\left|\int_{0}^{t}\left\langle\hat{\nabla} v_{n}\left(\xi^{n}\right), \Theta^{n}\right\rangle W^{n} d \vec{b}\right|^{p} \lesssim \mathbb{E} \int_{0}^{t} \frac{1}{n}\left|\left\langle\psi^{\prime}\left(n^{-1} W^{n}\right), \partial_{h} W^{n}\right\rangle W^{n}\right|^{p} d \tau \\
& \lesssim \mathbb{E} \int_{0}^{t}\left(\frac{1}{n}\left|\psi^{\prime}\left(n^{-1} W^{n}\right)\right|\left|W^{n}\right|\right)^{p}\left|\partial_{h} W^{n}\right|^{p} d \tau \\
& \lesssim \mathbb{E} \int_{0}^{t}\left(\frac{1}{n}\left|\psi^{\prime}\left(n^{-1} W^{n}\right)\right|\left|W^{n}\right|\right)^{p}\left|\partial_{h} W\right|^{p} d \tau \\
& \quad+\mathbb{E} \int_{0}^{t}\left(\frac{1}{n}\left|\psi^{\prime}\left(n^{-1} W^{n}\right)\right|\left|W^{n}\right|\right)^{p}\left|\partial_{h} W^{n}-\partial_{h} W\right|^{p} d \tau \rightarrow 0,
\end{aligned}
$$

as $n \rightarrow \infty$, by the dominated convergence theorem. Similarly,

$$
\begin{aligned}
\left\langle\tilde{\nabla} \hat{\nabla} v_{n}\left(\xi^{n}\right), X_{i} \otimes \Theta^{n}\right\rangle & =\frac{1}{n^{2}}\left\langle\psi^{\prime \prime}\left(n^{-1} W^{n}\right), W^{n} \operatorname{ad}_{X_{i}} \otimes \operatorname{ad}_{\Theta^{n}} \operatorname{Ad}_{\xi^{n}}\right\rangle \\
& =\frac{1}{n^{2}}\left\langle\psi^{\prime \prime}\left(n^{-1} W^{n}\right), W^{n} \operatorname{ad}_{X_{i}} \otimes \partial_{h} W^{n}\right\rangle,
\end{aligned}
$$

so that

$$
\begin{aligned}
& \mathbb{E} \mid \int_{0}^{t}\left\langle\tilde{\nabla} \hat{\nabla} v_{n}\left(\xi^{n}\right),\right.\left.X_{i} \otimes \Theta^{n}\right\rangle\left. W^{n} X_{i} d \tau\right|^{p} \\
&=\mathbb{E}\left|\int_{0}^{t} \frac{1}{n^{2}}\left\langle\psi^{\prime \prime}\left(n^{-1} W^{n}\right), W^{n} \operatorname{ad}_{X_{i}} \otimes \partial_{h} W^{n}\right\rangle W^{n} X_{i}\right|^{p} d \tau \\
& \lesssim \mathbb{E} \int_{0}^{t}\left(\frac{1}{n^{2}}\left|\psi^{\prime \prime}\left(n^{-1} W^{n}\right)\right|\left|W^{n}\right|^{2}\right)^{p}\left|\partial_{h} W^{n}\right|^{p} d \tau \\
& \lesssim \mathbb{E} \int_{0}^{t}\left(\frac{1}{n^{2}}\left|\psi^{\prime \prime}\left(n^{-1} W^{n}\right)\right|\left|W^{n}\right|^{2}\right)^{p}\left|\partial_{h} W\right|^{p} d \tau \\
& \quad+\mathbb{E} \int_{0}^{t}\left(\frac{1}{n^{2}}\left|\psi^{\prime \prime}\left(n^{-1} W^{n}\right)\right|\left|W^{n}\right|^{2}\right)^{p}\left|\partial_{h} W^{n}-\partial_{h} W\right|^{p} d \tau \rightarrow 0
\end{aligned}
$$

as $n \rightarrow \infty$, again by dominated convergence. Finally,

$$
\begin{aligned}
& \mathbb{E}\left|\int_{0}^{t}\left(v_{n}\left(\xi^{n}\right) W^{n}-W\right) d \vec{h}\right|^{p} \\
& \lesssim \mathbb{E} \int_{0}^{t}\left|v_{n}\left(\xi^{n}\right)-1\right|^{p}|W|^{p} d \vec{h}+\mathbb{E} \int_{0}^{t}\left|v_{n}\left(\xi^{n}\right)\right|^{p}\left|W^{n}-W\right|^{p} d \vec{h} \rightarrow 0
\end{aligned}
$$

as $n \rightarrow \infty$, by the dominated convergence theorem. Thus, Proposition 3.11 completes the proof of Equation (3.55).

Now, by Corollary 3.27, for any $f \in C_{c}^{\infty}(G)$ and $t \in[0,1], \partial_{h}\left[f\left(\xi_{t}^{n}\right)\right] \in L^{\infty-}(\mu)$, and

$$
\partial_{h}\left[f\left(\xi_{t}^{n}\right)\right]=\left\langle\hat{\nabla} f\left(\xi_{t}^{n}\right), \Theta_{t}^{n}\right\rangle .
$$

Thus, by the same argument as in Proposition 3.26. Equation (3.55) implies that

$$
\lim _{n \rightarrow \infty} \mathbb{E}\left|\partial_{h} f\left(\xi_{t}^{n}\right)-\left\langle\hat{\nabla} f\left(\xi_{t}\right), \Theta_{t}\right\rangle\right|^{p}=0,
$$

for all $p \in(1, \infty)$. Since $\partial_{h}$ is a closed operator, this and $f\left(\xi_{t}\right) \in L^{\infty-}(\mu)$ imply that $f\left(\xi_{t}\right) \in \operatorname{Dom}\left(\partial_{h}\right)$ and

$$
\partial_{h}\left[f\left(\xi_{t}\right)\right]=\left\langle\hat{\nabla} f\left(\xi_{t}\right), \Theta_{t}\right\rangle=\left\langle\hat{\nabla} f\left(\xi_{t}\right), \int_{0}^{t} \operatorname{Ad}_{\xi} X_{i} \dot{h}^{i} d \tau\right\rangle \in L^{\infty-}(\mu) .
$$

In particular, for any $h \in \mathscr{H}$ such that $\|h\|_{\mathscr{H}}=1$,

$$
\left.\mathbb{E}\left|\partial_{h}\left[f\left(\xi_{t}\right)\right]^{p} \leq\|\hat{\nabla} f\|_{\infty} \mathbb{E}\right| \int_{0}^{t} \operatorname{Ad}_{\xi} X_{i} \dot{h}^{i} d \tau\right|^{p} \leq\|\hat{\nabla} f\|_{\infty} \sum_{i=1}^{k} \mathbb{E} \int_{0}^{t}\left|\operatorname{Ad}_{\xi} X_{i}\right|^{p} d \tau .
$$

By the proof of Proposition 3.13,

$$
\mathbb{E} \sup _{\tau \leq 1}\left|\operatorname{Ad}_{\xi_{\tau}}\right|^{p} \leq C,
$$

for some finite constant $C$, and thus $\left\|f\left(\xi_{t}\right)\right\|_{G^{1, p}}<\infty$ for all $p \in(1, \infty)$, where

$$
\|F\|_{G^{1, p}}=\|F\|_{L^{p}(\mu)}+\left(\mathbb{E} \sup _{\|h\|_{\mathscr{C}}=1}\left|\partial_{h} F\right|^{p}\right)^{1 / p} ;
$$

see Section 2.5. Then by Theorem 2.12, $f\left(\xi_{t}\right) \in \mathcal{D}^{1, \infty}$, and

$$
\mathbb{E}\left[f\left(\xi_{t}\right) D^{*} h\right]=\mathbb{E}\left[\partial_{h}\left[f\left(\xi_{t}\right)\right]\right]=\mathbb{E}\left[\int_{0}^{t}\left\langle\hat{\nabla} f\left(\xi_{t}\right), \operatorname{Ad}_{\xi} X_{i}\right\rangle \dot{h}^{i} d \tau\right],
$$

implies that

$$
\left(D_{s}\left[f\left(\xi_{t}\right)\right]\right)^{i}=\left\langle\hat{\nabla} f\left(\xi_{t}\right), \int_{0}^{s \wedge t} \operatorname{Ad}_{\xi} X_{i} d \tau\right\rangle
$$

componentwise in $\mathscr{H}$.

## Chapter 4

## Lie group inequalities

Again, let $G$ be a Lie group with identity $e$ and Lie algebra $\operatorname{Lie}(G)=\mathfrak{g}$, and suppose $\left\{X_{i}\right\}_{i=1}^{k} \subset \mathfrak{g}$ is a Lie generating set, in the sense of Equation 1.7). We have the gradient $\nabla=\left(\tilde{X}_{1}, \ldots, \tilde{X}_{k}\right)$ and the subLaplacian $L=\sum_{i=1}^{k} \tilde{X}_{i}^{2}$ as operators on smooth functions of $G$ with compact support. Let $L$ also denote the self-adjoint extension of the subLaplacian and $P_{t}=e^{t L / 2}$ be the heat semigroup as in Definition 1.8.

We recall the following lemmas, proved in Chapter 2 in the Heisenberg context (Lemmas 2.4 and 2.5). The proofs are identical in the general Lie group case.

Lemma 4.1. By the left invariance of $\nabla$ and $P_{t}$, the inequality $I_{p}$ holds for all $g \in G$, $f \in C_{c}^{\infty}(G)$, and $t>0$, if and only if,

$$
\begin{equation*}
\left|\nabla P_{t} f\right|^{p}(e) \leq K_{p}(t) P_{t}|\nabla f|^{p}(e), \tag{4.1}
\end{equation*}
$$

for all $f \in C_{p}^{\infty}(G)$ and $t>0$, where $e \in G$ is the identity element.
Lemma 4.2. For $X \in \mathfrak{g}$,

$$
\tilde{X} P_{t} f(e)=P_{t} \hat{X} f(e) .
$$

for all $f \in C_{c}^{\infty}(G)$. More generally,

$$
\hat{X} P_{t} f=P_{t} \hat{X} f
$$

from which the previous equation follows, since $\hat{X}=\tilde{X}$ at e.
(The proof of this statement is slightly easier than its analogue Lemma 2.5 , since working with functions with compact support - versus functions with polynomial growth - requires only the invariance of Haar measure to justify passing the derivative through the integral.)

### 4.1 An $L^{p}$-type gradient estimate $(p>1)$

Notation 4.3. For each $r \in\{1, \ldots, m\}$, let $\Lambda^{r}=\Lambda^{k, r}$ be the set of multi-indices $\alpha=\left(\alpha_{0}, \alpha_{1}, \ldots, \alpha_{r}\right)$, with $\alpha_{l} \in\{1, \ldots, k\}$ for $l \in\{0, \ldots, r\}$. For any $\alpha \in \Lambda^{r}$, define

$$
\begin{gathered}
\alpha^{\prime}:=\left(\alpha_{1}, \ldots, \alpha_{r}\right) \text { and } \\
\bar{\alpha}:=\left(\alpha_{r}, \ldots, \alpha_{0}\right)=\alpha \text { reversed } .
\end{gathered}
$$

We define the order of $\alpha$ by $|\alpha|:=r+1$. Let

$$
\begin{gathered}
X_{\alpha}=\left[X_{\alpha_{r}},\left[\cdots,\left[X_{\alpha_{1}}, X_{\alpha_{0}}\right] \cdots\right]\right]=\operatorname{ad}_{X_{\alpha_{r}}} \cdots \operatorname{ad}_{X_{\alpha_{1}}} X_{\alpha_{0}} \text { and } \\
X^{\alpha}=X_{\alpha_{r}} \cdots X_{\alpha_{0}} .
\end{gathered}
$$

When $|\alpha|=1$, that is, $\alpha=\left(\alpha_{0}\right)$, then $X^{\alpha}=X_{\alpha_{0}}=X_{\alpha}$. For each $\alpha \in \Lambda^{r}$, there exist $\epsilon_{\beta, \alpha} \in\{-1,0,1\}$ such that

$$
X_{\alpha}=\sum_{\beta \in \Lambda^{r}} \epsilon_{\beta, \alpha} X^{\beta}
$$

For example, returning the Heisenberg group case and letting $X_{1}=X, X_{2}=Y$, and $X_{3}=Z$, with commutation relations given as in Section 2.1, we have

$$
X_{(1,2)}=\left[X_{1}, X_{2}\right]=X_{3}, X_{(\overline{1,2})}=X_{(2,1)}=\left[X_{2}, X_{1}\right]=-X_{3}, \quad \text { and } \quad X^{(1,2)^{\prime}}=X^{(2)}=X_{2} .
$$

Proposition 4.4. For each $i=1, \ldots, k, \hat{X}_{i}$ may be written as a linear combination,

$$
\begin{equation*}
\hat{X}_{i}=\sum_{r=0}^{m} \sum_{\alpha \in \Lambda^{r}} c_{i, \alpha} \tilde{X}^{\alpha} \tag{4.2}
\end{equation*}
$$

with coefficients $c_{i, \alpha}: G \rightarrow \mathbb{R}$ (some of these are 0) such that $c_{i, \alpha}\left(\xi_{t}\right) \in \mathcal{D}^{\infty}$, for all $t \in[0,1]$.

Proof. Recall from Notation 1.4 that we defined

$$
\Sigma_{r}=\left\{\left[X_{i_{1}},\left[\cdots,\left[X_{i_{r-1}}, X_{i_{r}}\right] \cdots\right]: i_{l} \in\{1, \ldots, k\}, l \in\{1, \ldots, r\}\right\}=\left\{X_{\alpha}: \alpha \in \Lambda^{r}\right\},\right.
$$

for $r=0, \ldots, m$. Recall also from Notation 1.4 that we have an orthonormal basis $\left\{X_{i}, Y_{j}: i \in\{1, \ldots, k\}, j \in\{1, \ldots, d-k\}\right\}$ of $\mathfrak{g}$, where $d=\operatorname{dim}(G)$ and, for each $j \in\{1, \ldots, d-k\}, Y_{j}=X_{\alpha(j)} \in \Sigma_{r(j)}$ for some $\alpha(j) \in \Lambda^{r(j)}, r(j) \in\{1, \ldots, m\}$. Thus, for any $g \in G$ and $X \in \mathfrak{g}$,

$$
\begin{aligned}
\hat{X}(g) & =R_{g *} X=L_{g *} L_{g^{-1}} R_{g *} X=L_{g *} \operatorname{Ad}_{g^{-1}} X \\
& =L_{g *}\left(\sum_{i=1}^{k}\left\langle\operatorname{Ad}_{g^{-1}} X, X_{i}\right\rangle X_{i}+\sum_{j=1}^{d-k}\left\langle\operatorname{Ad}_{g^{-1}} X, Y_{j}\right\rangle Y_{j}\right) \\
& =L_{g *}\left(\sum_{i=1}^{k}\left\langle\operatorname{Ad}_{g^{-1}} X, X_{i}\right\rangle X_{i}+\sum_{j=1}^{d-k} \sum_{\alpha \in \Lambda^{r(j)}} \epsilon_{\alpha, \alpha(j)}\left\langle\operatorname{Ad}_{g^{-1}} X, Y_{j}\right\rangle X^{\alpha}\right) \\
& =\sum_{i=1}^{k}\left\langle\operatorname{Ad}_{g^{-1}} X, X_{i}\right\rangle \tilde{X}_{i}+\sum_{j=1}^{d-k} \sum_{\alpha \in \Lambda^{r(j)}} \epsilon_{\alpha, \alpha(j)}\left\langle\operatorname{Ad}_{g^{-1}} X, Y_{j}\right\rangle \tilde{X}^{\alpha}
\end{aligned}
$$

where $\epsilon_{\alpha, \alpha(j)} \in\{-1,0,1\}$. So, for each $l=1, \ldots, k$,

$$
\begin{aligned}
\hat{X}_{l}(g) & =\sum_{i=1}^{k}\left\langle\operatorname{Ad}_{g^{-1}} X_{l}, X_{i}\right\rangle \tilde{X}_{i}(g)+\sum_{j=1}^{d-k} \sum_{\alpha \in \Lambda^{r(j)}} \epsilon_{\alpha, \alpha(j)}\left\langle\operatorname{Ad}_{g^{-1}} X_{l}, Y_{j}\right\rangle \tilde{X}^{\alpha}(g) \\
& =\sum_{r=0}^{m} \sum_{\alpha \in \Lambda^{r}} c_{l, \alpha} \tilde{X}^{\alpha}(g)
\end{aligned}
$$

where

$$
c_{l, \alpha}(g)=\left\langle\operatorname{Ad}_{g^{-1}} X_{l}, X_{i}\right\rangle, \quad \text { when } r=0 \text { and } \alpha=(i)
$$

and

$$
c_{l, \alpha}(g)=\epsilon\left\langle\operatorname{Ad}_{g^{-1}} X_{l}, Y_{j}\right\rangle, \epsilon \in\{-1,0,1\}, \text { when } r \in\{1, \ldots, m\}
$$

Now recall that $\operatorname{Ad}_{\xi_{t}} \in \mathcal{D}^{\infty}(\operatorname{End}(\mathfrak{g}))$ by Proposition 3.14, and $u: \operatorname{End}(\mathfrak{g}) \rightarrow \mathbb{R}$ given by $u(W)=\langle W X, Y\rangle$ is a smooth function for any fixed $X, Y \in \mathfrak{g}$. Thus, $u\left(\operatorname{Ad}_{\xi_{t}}\right) \in$ $\mathcal{D}^{\infty}$ for all $t \in[0,1]$. Since $c_{l, \alpha}\left(\xi_{t}\right)=\epsilon u\left(\operatorname{Ad}_{\xi_{t}}\right)$, with $X=X_{l}$ and $Y=X_{i}$ or $Y_{j}$, this implies that $c_{l, \alpha}\left(\xi_{t}\right) \in \mathcal{D}^{\infty}$, for all $\alpha \in \Lambda^{r}$ and $l \in\{1, \ldots, k\}$.

Theorem 4.5. Let $G$ be a Lie group with Lie algebra $\mathfrak{g}$, and suppose $\left\{X_{1}, \ldots, X_{k}\right\} \subset \mathfrak{g}$ is a Lie generating set. Then, for all $p \in(1, \infty), K_{p}(t)<\infty$, where $K_{p}(t)$ are the functions defined in Notation 1.9; that is, $K_{p}(t)$ is the best function so that

$$
\left|\nabla P_{t} f\right|^{p} \leq K_{p}(t) P_{t}|\nabla f|^{p}
$$

for all $f \in C_{c}^{\infty}(G)$ and $t>0$.
Proof. Lemma 4.1 implies that the inequality $\left(\overline{I_{p}}\right)$ is translation invariant on groups. Thus, we need only determine finite $K_{p}(t)$ such that the inequality holds at the identity. Using Lemma 4.2 and Equation (4.2), we have

$$
\left|\nabla P_{t} f\right|^{2}(e)=\sum_{i=1}^{k}\left|\tilde{X}_{i} P_{t} f\right|^{2}(e)=\sum_{i=1}^{k}\left|P_{t} \hat{X}_{i} f\right|^{2}(e) \leq C \sum_{i=1}^{k} \sum_{r=0}^{m} \sum_{\alpha \in \Lambda^{r}}\left|P_{t} c_{i, \alpha} \tilde{X}^{\alpha} f\right|^{2}(e),
$$

for a constant $C=C(k, m)$. Recall that Equation (3.4) implies that, for any $f \in C_{c}^{\infty}(G)$, $P_{t} f(e)=\mathbb{E}\left[f\left(\xi_{t}\right)\right]$, where $\xi$ is the solution to the Stratonovich equation (3.3). Thus, for any $\alpha \in \Lambda^{r}$,

$$
\begin{aligned}
\left|P_{t} c \tilde{X}^{\alpha} f\right|(e) & \leq \mathbb{E}\left|c\left(\xi_{t}\right)\left(\tilde{X}^{\alpha} f\right)\left(\xi_{t}\right)\right|=\mathbb{E}\left|c\left(\xi_{t}\right) \mathbf{X}^{\alpha^{\prime}}\left[\left(\tilde{X}_{\alpha_{0}} f\right)\left(\xi_{t}\right)\right]\right| \\
& =\mathbb{E}\left|\left(\mathbf{X}^{\overline{\alpha^{\prime}}}\right)^{*}\left[c\left(\xi_{t}\right)\right]\left(\tilde{X}_{\alpha_{0}} f\right)\left(\xi_{t}\right)\right| \\
& \leq\left(\mathbb{E}\left|\left(\mathbf{X}^{\overline{\alpha^{\prime}}}\right)^{*}\left[c\left(\xi_{t}\right)\right]\right|^{q}\right)^{1 / q}\left(\mathbb{E}\left|\left(\tilde{X}_{\alpha_{0}} f\right)\left(\xi_{t}\right)\right|^{p}\right)^{1 / p} \\
& =\left(\mathbb{E}\left|\left(\mathbf{X}^{\overline{\alpha^{\prime}}}\right)^{*}\left[c\left(\xi_{t}\right)\right]\right|^{q}\right)^{1 / q}\left(P_{t}\left|\tilde{X}_{\alpha_{0}} f\right|^{p}(e)\right)^{1 / p}
\end{aligned}
$$

by Hölder's inequality, where $q$ is the conjugate exponent to $p, \mathbf{X}^{\alpha}$ is the lifted vector field on $W$ of the vector field $\tilde{X}^{\alpha}$, as defined in Equation 3.11), and $\left(\mathbf{X}^{\alpha}\right)^{*}=\mathbf{X}_{\alpha_{r}}^{*} \cdots \mathbf{X}_{\alpha_{0}}^{*}$
(so $\left.\left(\mathbf{X}^{\overline{\alpha^{\prime}}}\right)^{*}=\mathbf{X}_{\alpha_{1}}^{*} \cdots \mathbf{X}_{\alpha_{r}}^{*}\right)$. Thus,

$$
\begin{aligned}
\left|\nabla P_{t} f\right|^{p}(e) & \leq C\left[\sum_{i=1}^{k} \sum_{r=0}^{m} \sum_{\alpha \in \Lambda^{r}}\left(\mathbb{E}\left|\left(\mathbf{X}^{\overline{\alpha^{\prime}}}\right)^{*}\left[c_{i, \alpha}\left(\xi_{t}\right)\right]\right|^{q}\right)^{2 / q}\left(P_{t}\left|\tilde{X}_{\alpha_{0}} f\right|^{p}(e)\right)^{2 / p}\right]^{p / 2} \\
& \leq C\left[\sum_{i=1}^{k} \sum_{r=0}^{m} \sum_{\alpha \in \Lambda^{r}}\left(\mathbb{E}\left|\left(\mathbf{X}^{\overline{\alpha^{\prime}}}\right)^{*}\left[c_{i, \alpha}\left(\xi_{t}\right)\right]\right|^{q}\right)^{p / q}\right] P_{t}|\nabla f|^{p}(e),
\end{aligned}
$$

where $C=C(k, m, p)$ and $q=\frac{p-1}{p}$. Thus, the inequality $I_{p}$ holds with

$$
\begin{equation*}
C_{p}(t)=C(k, m, p) \sum_{i=1}^{k} \sum_{r=0}^{m} \sum_{\alpha \in \Lambda^{r}}\left(\mathbb{E}\left|\left(\mathbf{X}^{\overline{\alpha^{\prime}}}\right)^{*}\left[c_{i, \alpha}\left(\xi_{t}\right)\right]\right|^{q}\right)^{p / q} \tag{4.3}
\end{equation*}
$$

Propositions 3.10 and 4.4 imply that $\left(\mathbf{X}^{\overline{\alpha^{\prime}}}\right)^{*}\left[c_{i, \alpha}\left(\xi_{t}\right)\right] \in L^{\infty-}(\mu)$, for all $i=1, \ldots, k$ and $\alpha \in \Lambda^{r}$. Therefore, $K_{p}(t) \leq C_{p}(t)<\infty$ for all $t>0$ and $p \in(1, \infty)$.

It is important to note that there is currently no good control over the behavior of the functions $C_{p}$ in Equation 4.3 with respect to $t$. In fact, from certain scaling arguments, it is expected that $C_{p}(t) \rightarrow \infty$ and $t \rightarrow 0$; see for example [8, 33]. However, we do not claim these are the best constants for which the inequality holds.

### 4.2 A Poincaré inequality

The following result is a direct corollary to Theorem 4.5. The proof is completely analogous to the proof of Theorem 2.23 .

Theorem 4.6 (Poincaré Inequality). Let $K_{2}(t)$ be the best function for which Equation $I_{p}$ holds for $p=2$, and let $p_{t}(g) d g$ be the hypoelliptic heat kernel. Then

$$
\begin{equation*}
\int_{G} f^{2}(g) p_{t}(g) d g-\left(\int_{G} f(g) p_{t}(g) d g\right)^{2} \leq \Lambda(t) \int_{G}|\nabla f|^{2}(g) p_{t}(g) d g \tag{4.4}
\end{equation*}
$$

for all $f \in C_{c}^{\infty}(G)$ and $t>0$, where

$$
\Lambda(t)=\int_{0}^{t} K_{2}(s) d s
$$

Proof. Let $F_{t}(g)=\left(P_{t} f\right)(g)$. Then

$$
\frac{d}{d t} P_{t-s} F_{s}^{2}=P_{t-s}\left(-\frac{1}{2} L F_{s}^{2}+F_{s} L F_{s}\right)=-P_{t-s}\left|\nabla F_{s}\right|^{2} .
$$

Integrating this equation on $t$ implies that

$$
\begin{aligned}
P_{t} f^{2}-\left(P_{t} f\right)^{2} & =\int_{0}^{t} P_{t-s}\left|\nabla F_{s}\right|^{2} d s \\
& =\int_{0}^{t} P_{t-s}\left|\nabla P_{s} f\right|^{2} d s \\
& \leq \int_{0}^{t} K_{2}(s) P_{t-s} P_{s}|\nabla f|^{2} d s=\left(\int_{0}^{t} K_{2}(s) d s\right) \cdot P_{t}|\nabla f|^{2}
\end{aligned}
$$

wherein we have made use of Theorem 4.5. Evaluating the above at $e \in G$ gives the desired result.

This Poincaré inequality is less useful in the general Lie group case because nothing is known about the integrability of $K_{p}(t)$. However, in the next two sections, we will carry out a generalization of the argument that was made in Section 2.3. That is, we will show that when $G$ is a nilpotent Lie group, $K_{p}(t)$ is a bounded function for all $p \in(1, \infty)$. In particular, when $p=2$, this implies the Poincaré inequality holds with $\Lambda(t)<\infty$, for all $t>0$.

### 4.2.1 Stratified nilpotent Lie groups

Definition 4.7. A Lie algebra $\mathfrak{g}$ is said to be nilpotent if $\mathrm{ad}_{X}$ is a nilpotent endomorphism of $\mathfrak{g}$ for all $X \in \mathfrak{g}$, that is, if there exists $m \in \mathbb{N}$ such that

$$
\operatorname{ad}_{Y_{1}} \cdots \operatorname{ad}_{Y_{m-1}} Y_{m}=\left[Y_{1},\left[\cdots,\left[Y_{m-1}, Y_{m}\right] \cdots\right]=0\right.
$$

for any $Y_{1}, \ldots, Y_{m} \in \mathfrak{g}$. We say that $\mathfrak{g}$ is nilpotent of step $m$. A Lie group $G$ is nilpotent if $\mathfrak{g}=\operatorname{Lie}(G)$ is a nilpotent Lie algebra.

Definition 4.8. A family of dilations on a Lie algebra $\mathfrak{g}$ is a family of algebra automorphisms $\left\{\Phi_{r}\right\}_{r>0}$ on $\mathfrak{g}$ of the form $\Phi_{r}=\exp (W \log r)$, where $W$ is a diagonalizable linear operator on $\mathfrak{g}$ with positive eigenvalues.

Definition 4.9. A stratified group $G$ is a simply connected nilpotent group for which there exists a subset of the Lie algebra $V_{1} \subset \mathfrak{g}$, such that $\mathfrak{g}=\oplus_{j=1}^{m} V_{j}$ with $V_{j+1}=\left[V_{1}, V_{j}\right]$, for $j=1, \ldots, m-1$, and $V_{m+1}=\left[V_{1}, V_{m}\right]=\{0\}$.

For a general exposition on nilpotent Lie groups and dilations, see [19, 22] and references contained therein. If $G$ is a stratified Lie group, a natural family of dilations may be defined on $\mathfrak{g}$ by setting $\Phi_{r}(X)=r^{j} X$, for all $X \in V_{j}$. The generator $W$ of this dilation acts on parts of the vector space decomposition by $W V_{j}=j V_{j}$, for each $j=$ $1, \ldots, m$. The automorphism $\Phi_{r}$ induces a group dilation $\phi_{r}$ via the exponential maps, $\phi_{r}=\exp \circ \Phi_{r} \circ \exp ^{-1}$. Since $G$ is a simply connected nilpotent group, the exponential map is in fact a global diffeomorphism on $\mathfrak{g}$, and $\exp ^{-1}$ exists everywhere on $G$; see for example Theorem 3.6.2 of Varadarajan [49]. Then for each $X \in V_{1}$,

$$
\begin{align*}
\tilde{X}\left(f \circ \phi_{r}\right)(g)=\left.\frac{d}{d \epsilon}\right|_{0} & \left(f \circ \phi_{r}\right)\left(g e^{\epsilon X}\right)=\left.\frac{d}{d \epsilon}\right|_{0} f\left(\phi_{r}(g) \phi_{r}\left(e^{\epsilon X}\right)\right) \\
& =\left.\frac{d}{d \epsilon}\right|_{0} f\left(\phi_{r}(g) e^{r \epsilon X}\right)=\left.\frac{d}{d \epsilon}\right|_{0} r f\left(\phi_{r}(g) e^{\epsilon X}\right)=r\left(\tilde{X} f \circ \phi_{r}\right)(g), \tag{4.5}
\end{align*}
$$

for all $f \in C^{1}(G)$, where in the second equality we have used that $\phi_{r}$ is a homomorphism. Let $\left\{X_{i}\right\}_{i=1}^{k} \subset V_{1}$ be a basis of $V_{1}$, and consider the operators $\nabla=\left(\tilde{X}_{1}, \ldots, \tilde{X}_{k}\right)$ and $L=\sum_{i=1}^{k}$. Equation (4.5) implies that

$$
\nabla\left(f \circ \phi_{r}\right)=r(\nabla f) \circ \phi_{r},
$$

and we have the following proposition.
Proposition 4.10. Let $L$ denote the self-adjoint extension of $\sum_{i=1}^{k} \tilde{X}_{i}^{2}$, and $P_{t}=e^{t L / 2}$ be as in Definition 1.8. Then

$$
L\left(f \circ \phi_{r}\right)=r^{2}(L f) \circ \phi_{r}
$$

and

$$
P_{t}\left(f \circ \phi_{r}\right)=e^{t L / 2}\left(f \circ \phi_{r}\right)=\left(e^{r^{2} t L / 2} f\right) \circ \phi_{r}=P_{r^{2} t} \circ \phi_{r},
$$

for any $f \in C_{c}^{\infty}(G)$.

Proof. Let $\mathcal{E}^{0}(f, h):=\sum_{i=1}^{k}\left(\tilde{X}_{i} f, \tilde{X}_{i} h\right)_{L^{2}(G)}$ be a Dirichlet form associated to $L$. Recall from Chapter 1 that $\mathcal{E}^{0}$ has a closed extension $\mathcal{E}$. By definition,

$$
f_{1} \in C_{c}^{\infty}(G) \text { and } L f_{1}=h \Longleftrightarrow \mathcal{E}\left(f_{1}, f_{2}\right)=\left(h, f_{2}\right), \forall f_{2} \in \operatorname{Dom}(\mathcal{E}) .
$$

Now note that

$$
\begin{aligned}
\mathcal{E}^{0}\left(f \circ \phi_{r}, f \circ \phi_{r}\right) & =\sum_{i=1}^{k} \int_{G}\left|\tilde{X}_{i}\left(f \circ \phi_{r}\right)\right|^{2}(g) d g \\
& =\sum_{i=1}^{k} r^{2} \int\left|\left(\tilde{X}_{i} f\right) \circ \phi_{r}\right|^{2}(g) d g \\
& =\sum_{i=1}^{k} r^{2} \int\left|\tilde{X}_{i} f\right|^{2}(g) J\left(r^{-1}\right) d g=r^{2} J\left(r^{-1}\right) \mathcal{E}^{0}(f, f),
\end{aligned}
$$

where $J(r)$ is the Jacobian of the transformation $\phi_{r}$,

$$
J(r)=\prod_{j=1}^{m}\left(r^{j}\right)^{d_{j}}
$$

with $d_{j}=\operatorname{dim}\left(V_{j}\right)$. Thus, $J\left(r^{-1}\right)=J(r)^{-1}$. So $f \in \operatorname{Dom}(\mathcal{E})$ implies that $f \circ \phi_{r} \in$ $\operatorname{Dom}(\mathcal{E})$, and, in general, $\mathcal{E}\left(f \circ \phi_{r}, h \circ \phi_{r}\right)=r^{2} J\left(r^{-1}\right) \mathcal{E}(f, h)$, for $f, h \in \operatorname{Dom}(\mathcal{E})$. Replacing $h$ here by $h \circ \phi_{r^{-1}}$ gives

$$
\begin{aligned}
\mathcal{E}\left(f \circ \phi_{r}, h\right) & =r^{2} J\left(r^{-1}\right) \mathcal{E}\left(f, h \circ \phi_{r^{-1}}\right) \\
& =r^{2} J\left(r^{-1}\right)\left(L f, h \circ \phi_{r^{-1}}\right)_{L^{2}(G)} \\
& =r^{2} J\left(r^{-1}\right) J(r)\left(L f \circ \phi_{r}, h\right)_{L^{2}(G)}=r^{2}\left(L f \circ \phi_{r}, h\right)_{L^{2}(G)},
\end{aligned}
$$

implies that if $f \in \operatorname{Dom}(L)$, then $f \circ \phi_{r} \in \operatorname{Dom}(L)$ and $L\left(f \circ \phi_{r}\right)=r^{2} L f \circ \phi_{r}$.
Now, for $r>0$, let $U_{r}: L^{2}(G) \rightarrow L^{2}(G)$ be the unitary operator given by $U_{r} f=\frac{1}{\sqrt{J\left(r^{-1}\right)}} f \circ \phi_{r}$. Then

$$
L U_{r}=r^{2} U_{r} L=U_{r}\left(r^{2} L\right)
$$

as operators, and thus $U_{r}^{-1} L U_{r}=r^{2} L$. Then we have that

$$
U_{r}^{-1} e^{t L / 2} U_{r}=e^{t U_{r}^{-1} L U_{r} / 2}=e^{r^{2} t L / 2}
$$

from which it follows that

$$
r^{2} e^{t L / 2}\left(f \circ \phi_{r}\right)=e^{t L / 2} U_{r} f=U_{r} e^{r^{2} t L / 2} f=r^{2}\left(e^{r^{2} t L / 2} f\right) \circ \phi_{r} .
$$

Given these equations, an argument identical to the proof of Proposition 2.6 proves the following proposition.

Proposition 4.11. Suppose $G$ is a stratified Lie group with vector space decomposition $\oplus_{j=1}^{m} V_{j}$. Let $\left\{X_{i}\right\}_{i=1}^{k} \subset V_{1}, \nabla$, and $L$ be as above, and let $p \in(1, \infty)$. If $K_{p}$ is the best constant such that

$$
\left|\nabla P_{1} f\right|^{p} \leq K_{p} P_{1}|\nabla f|^{p},
$$

for all $f \in C_{c}^{\infty}(G)$, then $K_{p}(t)=K_{p}$ for all $t>0$, where $K_{p}(t)$ is the function defined in Notation 1.9

### 4.2.2 Nilpotent Lie groups

Now let $G$ be a general nilpotent Lie group. Because not all nilpotent Lie groups admit dilations, we can not show that the functions $K_{p}(t)$ should be scale invariant in this context. However, by covering $G$ with a group which has a family of dilations adapted to its structure, we may show that there exists some constant $K_{p}<\infty$ for which $K_{p}(t)<K_{p}$ for all $t>0$.

Definition 4.12. Let $\mathcal{L}=\mathcal{L}(k, m)$ be the free nilpotent Lie algebra of step $m$ with $k$ generators $\left\{e_{i}\right\}_{i=1}^{k}$. Then $\mathcal{L}$ is the unique (up to isomorphism) nilpotent Lie algebra of rank $m$ such that, for every nilpotent Lie algebra $\mathfrak{g}$ of rank $m$ and map $\tilde{\Pi}:\left\{e_{1}, \ldots, e_{k}\right\} \rightarrow$ $\mathfrak{g}$, there exists a unique homomorphism $\Pi: \mathcal{L} \rightarrow \mathfrak{g}$ which extends $\tilde{\Pi}$. Let $\mathcal{N}=\mathcal{N}(k, m)$ be the free nilpotent Lie group of rank $m$ with $k$ generators, which is the simply connected group of $\mathcal{L}(k, m)$.

The Lie algebra $\mathcal{L}(k, m)$ admits a vector space decomposition by setting $V_{1}=$ $\operatorname{span}\left\{e_{1}, \ldots, e_{k}\right\}$. Thus, $\mathcal{N}$ is a stratified Lie group with Hörmander set $\left\{e_{i}\right\}_{i=1}^{k} \subset \mathcal{L}$;
for definitions and further details, see [51]. Let $\nabla_{\mathcal{L}}=\left(\tilde{e}_{1}, \ldots, \tilde{e}_{k}\right), \mathscr{L}=\sum_{i=1}^{k} \tilde{e}_{i}^{2}$, and $\mathscr{P}_{t}=e^{t \mathscr{L} / 2}$. Theorem 4.5 and Proposition 4.11 imply that, for all $p \in(1, \infty)$, there exist constants $K_{p}^{\mathcal{L}}<\infty$ such that

$$
\begin{equation*}
\left|\nabla_{\mathcal{L}} \mathscr{P}_{t} f\right|^{p} \leq K_{p}^{\mathcal{L}} \mathscr{P}_{t}\left|\nabla_{\mathcal{L}} f\right|^{p} \tag{4.6}
\end{equation*}
$$

for all $f \in C_{c}^{\infty}(\mathcal{N})$ and $t>0$.
Proposition 4.13. Let $G$ be a nilpotent group of step $m$ with Hörmander set $\left\{X_{i}\right\}_{i=1}^{k}$. Then $K_{p}(t) \leq K_{p}^{\mathcal{L}}$ for all $t>0$, where $K_{p}(t)$ is the function defined in Notation 1.9.

Proof. By definition of $\mathcal{L}=\mathcal{L}(k, m)$, there exists a unique Lie algebra homomorphism $\Pi: \mathcal{L} \rightarrow \mathfrak{g}$ such that $\Pi\left(e_{i}\right)=X_{i}$. Then $\Pi$ induces a group homomorphism $\pi: \mathcal{N} \rightarrow G$ via the exponential maps,

$$
\pi=\exp _{G} \circ \Pi \circ \exp _{\mathcal{N}}^{-1}
$$

Again, because $\mathcal{N}$ is a simply connected nilpotent Lie group, the exponential map on $\mathcal{L}$ is a global diffeomorphism. Note that $d \pi=\Pi$,

and the vector fields $\tilde{X}_{i}$ and $\tilde{e}_{i}$ are $\pi$-related; that is,

$$
\tilde{e}_{\alpha}(f \circ \pi)=\left(\tilde{X}_{\alpha} f\right) \circ \pi
$$

for any multi-index $\alpha \in \Lambda^{r}$ and $f \in C_{c}^{\infty}(G)$. Since $\pi$ is a Lie group homomorphism, $f \circ \pi \in C_{c}^{\infty}(\mathcal{N})$, and thus,

$$
\left|\nabla P_{t} f\right|^{p}(e)=\left|\nabla_{\mathcal{L}} \mathscr{P}_{t}(f \circ \pi)\right|^{p}\left(e_{\mathcal{N}}\right) \leq K_{p}^{\mathcal{L}} \mathscr{P}_{t}\left|\nabla_{\mathcal{L}}(f \circ \pi)\right|^{p}\left(e_{\mathcal{N}}\right)=K_{p}^{\mathcal{L}} P_{t}|\nabla f|^{p}(e)
$$

where $e_{\mathcal{N}}$ is the identity element of $\mathcal{N}$. Since $K_{p}(t)$ is the best constant for which

$$
\left|\nabla P_{t} f\right|^{p}(e) \leq K_{p}(t) P_{t}|\nabla f|^{p}(e)
$$

holds, the above implies that $K_{p}(t) \leq K_{p}^{\mathcal{L}}$ for all $t>0$.

This method of lifting the vector fields to a free nilpotent Lie algebra was learned from [50, 51]. A generalization of this procedure may be found in [43].

Remark 4.14. Note that the above argument is independent of the minimality of the Hörmander set $\left\{X_{i}\right\}_{i=1}^{k}$. So suppose that the collection $\left\{X_{i}\right\}_{i=1}^{k}$ spans the Lie algebra $\mathfrak{g}$. Since $G$ is a nilpotent Lie group (and thus unimodular) it is then well known that the operator $L=\sum_{i=1}^{k} \tilde{X}_{i}^{2}$ is in fact the Laplace-Beltrami operator on the Riemannian manifold $(G,\langle\cdot, \cdot\rangle)$. Then by Theorem 1.2 ,

$$
\left|\nabla P_{t} f\right|^{p} \leq e^{p k t} P_{t}|\nabla f|^{p},
$$

where $-k$ is a lower bound on the Ricci curvature. Proposition 4.13 improves this result for large $t$ by implying that there exists a $K_{p}<\infty$ independent of $t$ such that

$$
\left|\nabla P_{t} f\right|^{p} \leq K_{p} P_{t}|\nabla f|^{p},
$$

for all $f \in C_{p}^{\infty}(G)$ and $t>0$. Thus, we have the following corollary.
Corollary 4.15. Let $G$ be a nilpotent Lie group of step $m$ and $\left\{X_{i}\right\}_{i=1}^{k} \subset \mathfrak{g}$ such that $\left\{X_{i}\right\}_{i=1}^{k}$ spans the Lie algebra $\mathfrak{g}$. Then, for $K_{p}(t)$ as in Notation 1.9.

$$
K_{p}(t) \leq \min \left\{K_{p}^{\mathcal{L}}, e^{p k t}\right\},
$$

where $K_{p}^{\mathcal{L}}$ is the best constant so that $I_{p}$ holds on $\mathcal{L}(k, m)$ and $-k$ is a lower bound on the Ricci curvature associated to the Riemannian metric determined by $L=\sum_{i=1}^{k} \tilde{X}_{i}^{2}$.

Corollary 4.16 (Poincaré Inequality for nilpotent Lie groups). Suppose $G$ is a nilpotent Lie group, and let $K_{2}$ be a finite constant for which the inequality (In) holds for $p=2$. Then the inequality (4.4) holds with $\Lambda(t)=K_{2} t$, for all $t>0$.

## Chapter 5

## Appendix

### 5.1 Proof of Theorem 1.2

Proof. Comparing Equations (1.1) and (1.2), it is clear that a lower bound on the Ricci curvature immediately gives

$$
\Gamma_{2}(f)=\left|\nabla^{2} f\right|^{2}+(\operatorname{Ric} \nabla f, \nabla f) \geq(\operatorname{Ric} \nabla f, \nabla f) \geq-2 k|\nabla f|^{2}=\Gamma_{1}(f),
$$

and so (1) implies (2). For the converse, let $m \in M, v \in T_{m} M$, and let $f \in C^{\infty}(M)$ such that $f\left(\exp _{m}(w)\right)=(v, w)$, for all $w \in T_{m} M$ in a neighborhood of $0 \in T_{m} M$. Since $f\left(\exp _{m}(t w)\right)=t v \cdot w$, it follows that

$$
d f\left(w_{m}\right)=\left.\frac{d}{d t}\right|_{0} f\left(\exp _{m}(t w)\right)=(v, w)
$$

and

$$
\nabla^{2} f\left(w_{m}, w_{m}\right)=\left.\frac{d^{2}}{d t^{2}}\right|_{0} f\left(\exp _{m}(t w)\right)=0
$$

Thus, $\left(\nabla^{2} f\right)_{m}=0, \nabla f(m)=v$, and so

$$
-2 k|v|^{2}=-2 k \Gamma_{1}(f) \leq \Gamma_{2}(f)=(\operatorname{Ric} v, v) ;
$$

that is, (2) implies (1).

An application of the Cauchy-Schwarz inequality shows that (3) implies (4), and (5) is a trivial consequence of (4). Assuming (5) holds, differentiating the inequality

$$
\left|\nabla P_{t} f\right|^{2} \leq K(t) P_{t}|\nabla f|^{2}
$$

at $t=0$, implies that

$$
\begin{aligned}
2\left(\nabla f, \frac{1}{2} \nabla L f\right) & =\left.\frac{d}{d t}\right|_{0}\left|\nabla P_{t} f\right|^{2} \\
& \leq\left.\frac{d}{d t}\right|_{0} K(t) P_{t}|\nabla f|^{2}=\dot{K}(0)|\nabla f|^{2}+\frac{1}{2} L|\nabla f|^{2} .
\end{aligned}
$$

Hence,

$$
\Gamma_{2}(f)=\frac{1}{2} L|\nabla f|^{2}-(\nabla f, \nabla L f) \geq \dot{K}(0)|\nabla f|^{2}=\dot{K}(0) \Gamma_{1}(f)
$$

and thus $\sqrt{22}$ follows with $k=\frac{1}{2} \dot{K}(0)$.
Finally, to verify that (1) implies (3), we may use the Feynman-Kac formula,

$$
e^{t \Delta / 2} f(x)=\mathbb{E}\left[Q_{t} / /_{t}^{-1} f\left(B_{t}^{x}\right)\right],
$$

where $B_{t}^{x}$ is Brownian motion on $M$ such that $B_{0}^{x}=x, / / t$ is parallel translation, and $Q_{t}$ satisfies

$$
\frac{d}{d t} Q_{t}=-\operatorname{Ric}_{/ / t} Q_{t}=0 \text { with } Q_{0}=I d
$$

see for example Theorem 7.2.1 of [27]. Thus,

$$
\nabla e^{t \Delta / 2} f(x)=\mathbb{E}\left[Q_{t} / /_{t}^{-1} \nabla f\left(B_{t}^{x}\right)\right] .
$$

So if Ric $\geq-k g$, then $\left|Q_{t}\right| \leq e^{k t}$, and hence

$$
\begin{aligned}
&\left|\nabla e^{t \Delta / 2} f\right|(x) \leq \mathbb{E} \mid Q_{t} / / /_{t}^{-1} \nabla f\left(B_{t}^{x}\right) \mid \\
& \leq e^{k t} \mathbb{E} \mid / / t \\
&-1 \\
&
\end{aligned}
$$

and the proof is complete.

### 5.2 Proof of Theorem 3.5

Notation 5.1. Let $S:=\{v \in \mathfrak{g}:\langle v, v\rangle=1\}$, that is, $S$ is the unit sphere in $\mathfrak{g}$.
Proof. To show $\bar{\sigma}_{t}^{-1} \in L^{\infty-}(\mu)$, it suffices to show that

$$
\mu\left(\inf _{v \in S}\left\langle\bar{\sigma}_{t} v, v\right\rangle<\varepsilon\right)=O\left(\varepsilon^{\infty-}\right)
$$

To verify this claim, notice that $\lambda_{0}:=\inf _{v \in S}\left\langle\bar{\sigma}_{t} v, v\right\rangle$ is the smallest eigenvalue of $\bar{\sigma}_{t}$. Since det $\bar{\sigma}_{t}$ is the product of the eigenvalues of $\bar{\sigma}_{t}$ it follows that $\Delta_{t}:=\operatorname{det} \bar{\sigma}_{t} \geq \lambda_{0}^{n}$ and so $\left\{\operatorname{det} \bar{\sigma}_{t}<\varepsilon^{n}\right\} \subset\left\{\lambda_{0}<\varepsilon\right\}$ and hence

$$
\mu\left(\operatorname{det} \bar{\sigma}_{t}<\varepsilon^{n}\right) \leq \mu\left(\lambda_{0}<\varepsilon\right)=O\left(\varepsilon^{\infty-}\right) .
$$

By replacing $\varepsilon$ by $\varepsilon^{1 / n}$ above this implies $\mu\left(\Delta_{t}<\varepsilon\right)=O\left(\varepsilon^{\infty-}\right)$. From this estimate it then follows that

$$
\begin{aligned}
& \mathbb{E}\left[\Delta_{t}^{-q}\right]=\mathbb{E} \int_{\Delta_{t}}^{\infty} q \tau^{-q-1} d \tau=q \mathbb{E} \int_{0}^{\infty} 1_{\Delta_{t} \leq \tau} \tau^{-q-1} d \tau \\
&=q \int_{0}^{\infty} \mu\left(\Delta_{t} \leq \tau\right) \tau^{-q-1} d \tau=q \int_{0}^{\infty} O\left(\tau^{p}\right) \tau^{-q-1} d \tau
\end{aligned}
$$

which is seen to be finite by taking $p \geq q+1$.
More generally, if $T$ is any stopping time with $T \leq t$, since $\left\langle\bar{\sigma}_{T} v, v\right\rangle \leq\left\langle\bar{\sigma}_{t} v, v\right\rangle$ for all $v \in S$, it suffices to prove

$$
\begin{equation*}
\mu\left(\inf _{v \in S}\left\langle\bar{\sigma}_{T} v, v\right\rangle<\varepsilon\right)=O\left(\varepsilon^{\infty-}\right) \tag{5.1}
\end{equation*}
$$

According to Lemma 5.2 and Proposition 5.3 below, Equation (5.1) holds with

$$
\begin{equation*}
T=T_{\delta}:=\inf \left\{t>0: \max \left\{\left|\operatorname{Ad}_{\xi_{t}}-I_{\mathfrak{g}}\right|,\left|\xi_{t}\right|\right\}>\delta\right\} \tag{5.2}
\end{equation*}
$$

provided $\delta>0$ is chosen sufficiently small.
The rest of this section is now devoted to the proof of Lemma 5.2 and Proposition 5.3 below. In what follows we will make repeated use of the identity,

$$
\begin{equation*}
\left\langle\bar{\sigma}_{T} v, v\right\rangle=\sum_{i=1}^{k} \int_{0}^{T}\left\langle\operatorname{Ad}_{\xi_{\tau}} X_{i}, v\right\rangle^{2} d \tau \tag{5.3}
\end{equation*}
$$

To prove this, recall that $\left\{X_{i}\right\}_{i=1}^{k}$ is an orthonormal basis of $\mathfrak{g}_{0}$. Thus,

$$
P \operatorname{Ad}_{\xi_{\tau}}^{\dagger} v=\sum_{i=1}^{k}\left\langle\operatorname{Ad}_{\xi_{\tau}}^{\dagger} v, X_{i}\right\rangle X_{i}=\sum_{i=1}^{k}\left\langle\operatorname{Ad}_{\xi_{\tau}} X_{i}, v\right\rangle X_{i}
$$

so that

$$
\left\langle\operatorname{Ad}_{\xi_{\tau}} P \operatorname{Ad}_{\xi_{\tau}}^{\dagger} v, v\right\rangle=\sum_{i=1}^{k}\left\langle\operatorname{Ad}_{\xi_{\tau}}\left\langle\operatorname{Ad}_{\xi_{\tau}} X_{i}, v\right\rangle X_{i}, v\right\rangle=\sum_{i=1}^{k}\left\langle\operatorname{Ad}_{\xi_{\tau}} X_{i}, v\right\rangle^{2}
$$

which upon integrating in $\tau$ gives Equation (5.3).
Lemma 5.2 (Compactness Argument). Let $T_{\delta}$ be as in Equation (5.2), and suppose for all $v \in S$ there exists $i \in\{1, \ldots, k\}$ and an open neighborhood $N \subset_{o} S$ of $v$ such that

$$
\begin{equation*}
\sup _{u \in N} \mu\left(\int_{0}^{T_{\delta}}\left\langle\operatorname{Ad}_{\xi_{\tau}} X_{i}, u\right\rangle^{2} d \tau<\varepsilon\right)=O\left(\varepsilon^{\infty-}\right), \tag{5.4}
\end{equation*}
$$

then Equation (5.1) holds provided $\delta>0$ is sufficiently small.
Proof. By compactness of $S$, it follows from Equation (5.4) that

$$
\begin{equation*}
\sup _{u \in S} \mu\left(\int_{0}^{T_{\delta}}\left\langle\operatorname{Ad}_{\xi_{\tau}} X_{i}, u\right\rangle^{2} d \tau<\varepsilon\right)=O\left(\varepsilon^{\infty-}\right) . \tag{5.5}
\end{equation*}
$$

For $w \in \mathfrak{g}$, let $\partial_{w}$ denote the directional derivative acting on functions $f(v)$ with $v \in \mathfrak{g}$. For all $v, w \in \mathbb{R}^{n}$ with $|v| \leq 1$ and $|w| \leq 1$ (using Equation (5.3)),

$$
\begin{aligned}
\left|\partial_{w}\left\langle\bar{\sigma}_{T_{\delta}} v, v\right\rangle\right| & \leq 2 \sum_{i=1}^{k} \int_{0}^{T_{\delta}}\left|\left\langle\operatorname{Ad}_{\xi_{\tau}} X_{i}, v\right\rangle\left\langle\operatorname{Ad}_{\xi_{\tau}} X_{i}, w\right\rangle\right| d \tau \\
& \leq 2 \sum_{i=1}^{k} \int_{0}^{T_{\delta}}\left|\operatorname{Ad}_{\xi_{\tau}} X_{i}\right|_{\operatorname{Hom}\left(\mathbb{R}^{n}, \mathfrak{g}\right)}^{2} d \tau .
\end{aligned}
$$

Thus, by choosing $\delta>0$ in Equation (5.2) sufficiently small, we may assume there is a non-random constant $\theta<\infty$ such that

$$
\sup _{|v|,|w| \leq 1}\left|\partial_{w}\left\langle\bar{\sigma}_{T_{\delta}} v, v\right\rangle\right| \leq \theta<\infty .
$$

With this choice of $\delta$, if $v, w \in S$ satisfy $|v-w|<\theta / \varepsilon$, then

$$
\begin{equation*}
\left|\left\langle\bar{\sigma}_{T_{\delta}} v, v\right\rangle-\left\langle\bar{\sigma}_{T_{\delta}} w, w\right\rangle\right|<\varepsilon . \tag{5.6}
\end{equation*}
$$

There exists $D<\infty$ such that ,for any $\varepsilon>0$, there is an open cover of $S$ with at most $D \cdot(\theta / \varepsilon)^{n}$ balls of the form $B\left(v_{j}, \varepsilon / \theta\right)$. From Equation (5.6), for any $v \in S$ there exists $j$ such that $v \in B\left(v_{j}, \varepsilon / \theta\right) \cap S$ and

$$
\left|\left\langle\bar{\sigma}_{T_{\delta}} v, v\right\rangle-\left\langle\bar{\sigma}_{T_{\delta}} v_{j}, v_{j}\right\rangle\right|<\varepsilon .
$$

So if $\inf _{v \in S}\left\langle\bar{\sigma}_{T_{\delta}} v, v\right\rangle<\varepsilon$, then $\min _{j}\left\langle\bar{\sigma}_{T_{\delta}} v_{j}, v_{j}\right\rangle<2 \varepsilon$; that is,

$$
\left\{\inf _{v \in S}\left\langle\bar{\sigma}_{T_{\delta}} v, v\right\rangle<\varepsilon\right\} \subset\left\{\min _{j}\left\langle\bar{\sigma}_{T_{\delta}} v_{j}, v_{j}\right\rangle<2 \varepsilon\right\} \subset \bigcup_{j}\left\{\left\langle\bar{\sigma}_{T_{\delta}} v_{j}, v_{j}\right\rangle<2 \varepsilon\right\}
$$

Therefore,

$$
\begin{aligned}
\mu\left(\inf _{v \in S}\left\langle\bar{\sigma}_{T_{\delta}} v, v\right\rangle\right. & <\varepsilon) \leq \sum_{j} \mu\left(\left\langle\bar{\sigma}_{T_{\delta}} v_{j}, v_{j}\right\rangle<2 \varepsilon\right) \\
& \leq D \cdot(\theta / \varepsilon)^{n} \cdot \sup _{v \in S} \mu\left(\left\langle\bar{\sigma}_{T_{\delta}} v, v\right\rangle<2 \varepsilon\right) \leq D \cdot(\theta / \varepsilon)^{n} O\left(\varepsilon^{\infty-}\right)=O\left(\varepsilon^{\infty-}\right)
\end{aligned}
$$

Recall from Equation (3.19), that for any $X \in \mathfrak{g}$,

$$
d \operatorname{Ad}_{\xi_{s}} X=\operatorname{Ad}_{\xi_{s}} \circ \operatorname{ad}_{d b_{s}} X=\operatorname{Ad}_{\xi_{s}} \operatorname{ad}_{X_{i}} X \circ d b_{s}^{i}=\operatorname{Ad}_{\xi_{s}}\left[X_{i}, X\right] \circ d b_{s}^{i}
$$

which may be rewritten in Itô form as

$$
\begin{equation*}
d\left[\operatorname{Ad}_{\xi_{s}} X\right]=\operatorname{Ad}_{\xi_{s}}\left[X_{i}, X\right] d b_{s}^{i}+\frac{1}{2} \operatorname{Ad}_{\xi_{s}} \operatorname{ad}_{X_{i}}^{2} X d s \tag{5.7}
\end{equation*}
$$

This gives the first hint that Hörmander's condition is relevant to showing $\Delta_{t}^{-1} \in$ $L^{\infty-}(\mu)$, or equivalently that $\bar{\sigma}_{t}^{-1} \in L^{\infty-}(\mu)$.

Proposition 5.3. Let $T_{\delta}$ be as in Equation (5.2). If Hörmander's restricted bracket condition holds and $v \in S$ is given, there exists $i \in\{1, \ldots, k\}$ and an open neighborhood $U \subset_{o} S$ of $v$ such that

$$
\sup _{u \in U} \mu\left(\int_{0}^{T_{\delta}}\left\langle\operatorname{Ad}_{\xi_{\tau}} X_{i}, u\right\rangle^{2} d \tau \leq \varepsilon\right)=O\left(\varepsilon^{\infty-}\right)
$$

Proof. Hörmander's condition implies that there exist $m \in \mathbb{N}$ and $\beta>0$ such that

$$
\sum_{r=0}^{m} \frac{1}{\left|\Sigma_{r}\right|} \sum_{V \in \Sigma_{r}} V V^{\operatorname{tr}} \geq 3 \beta I
$$

Equivalently, for all $v \in S$,

$$
3 \beta \leq \sum_{r=0}^{m} \frac{1}{\left|\Sigma_{r}\right|} \sum_{V \in \Sigma_{r}}\langle V, v\rangle^{2} \leq \max _{\substack{V \in \Sigma_{r} \\ r \in\{0, \ldots, m\}}}\langle V, v\rangle^{2} .
$$

By choosing $\delta>0$ in Equation (5.2) sufficiently small we may assume that

$$
\max _{\substack{V \in \Sigma_{r} \\ r \in\{0, \ldots, m\}}} \inf _{\tau \leq T_{\delta}}\left\langle\operatorname{Ad}_{\xi_{\tau}} V, v\right\rangle^{2} \geq 2 \beta \text {, for all } v \in S \text {. }
$$

Fix a $v \in S$ and $V \in \cup_{r=0}^{m} \Sigma_{r}$ such that

$$
\inf _{\tau \leq T_{\delta}}\left\langle\operatorname{Ad}_{\xi_{\tau}} V, v\right\rangle^{2} \geq 2 \beta,
$$

and choose an open neighborhood $U \subset S$ of $v$ such that

$$
\inf _{\tau \leq T_{\delta}}\left\langle\operatorname{Ad}_{\xi_{\tau}} V, u\right\rangle^{2} \geq \beta \text {, for all } u \in U .
$$

Then, by Equation (8.19) of Driver [15], $\mu\left(T_{\delta} \leq \tau\right)=O\left(\tau^{\infty-}\right)$, and so

$$
\begin{equation*}
\sup _{u \in U} \mu\left(\int_{0}^{T_{\delta}}\left\langle\operatorname{Ad}_{\xi_{\tau}} V, u\right\rangle^{2} d \tau \leq \varepsilon\right) \leq \mu\left(\int_{0}^{T_{\delta}} \beta d t \leq \varepsilon\right)=\mu\left(T_{\delta} \leq \varepsilon / \beta\right)=O\left(\varepsilon^{\infty-}\right) . \tag{5.8}
\end{equation*}
$$

Write $V=\operatorname{ad}_{X_{i_{r}}} \cdots \operatorname{ad}_{X_{i_{2}}} X_{i_{1}}$ with $r \leq m$. When $r=1$, Equation (5.8) becomes

$$
\sup _{u \in U} \mu\left(\left\langle\bar{\sigma}_{T_{\delta}} u, u\right\rangle \leq \varepsilon\right) \leq \sup _{u \in U} \mu\left(\int_{0}^{T_{\delta}}\left\langle\operatorname{Ad}_{\xi_{\tau}} X_{i_{1}}, u\right\rangle^{2} d t \leq \varepsilon\right)=O\left(\varepsilon^{\infty-}\right),
$$

and we are done. So now suppose $r>1$. Set

$$
V_{j}=\operatorname{ad}_{X_{i_{j}}} \ldots \operatorname{ad}_{X_{i_{2}}} X_{i_{1}} \text { for } j=1,2, \ldots, r,
$$

so that $V_{r}=V$. We will now show by (decreasing) induction on $j$ that

$$
\begin{equation*}
\sup _{u \in U} \mu\left(\int_{0}^{T_{\delta}}\left\langle\operatorname{Ad}_{\xi_{\tau}} V_{j}, u\right\rangle^{2} d t \leq \varepsilon\right)=O\left(\varepsilon^{\infty-}\right) . \tag{5.9}
\end{equation*}
$$

From Equation (5.7), we have

$$
d\left[\operatorname{Ad}_{\xi_{\tau}} V_{j-1}\right]=\operatorname{Ad}_{\xi_{\tau}} \sum_{i=1}^{k}\left[X_{i}, V_{j-1}\right] d b_{\tau}^{i}+\frac{1}{2} \operatorname{Ad}_{\xi_{\tau}} \sum_{i=1}^{k} \operatorname{ad}_{X_{i}}^{2} V_{j-1} d \tau
$$

which upon integrating in $\tau$ gives

$$
\left\langle\operatorname{Ad}_{\xi_{t}} V_{j-1}, u\right\rangle=\left\langle V_{j-1}, u\right\rangle+\int_{0}^{t}\left\langle\operatorname{Ad}_{\xi_{\tau}}\left[X_{i}, V_{j-1}\right], u\right\rangle d B_{\tau}^{i}+\frac{1}{2} \int_{0}^{t}\left\langle\operatorname{Ad}_{\xi_{\tau}} \operatorname{ad}_{X_{i}}^{2} V_{j-1}, u\right\rangle d \tau .
$$

Applying Proposition 5.4 of the appendix with $T=T_{\delta}$,

$$
\begin{aligned}
Y_{t} & :=\left\langle\operatorname{Ad}_{\xi_{t}} V_{j-1}, u\right\rangle, y=\left\langle V_{j-1}, u\right\rangle, \\
M_{t} & :=\int_{0}^{t}\left\langle\operatorname{Ad}_{\xi_{\tau}}\left[X_{i}, V_{j-1}\right], u\right\rangle d b_{\tau}^{i}, \text { and } \\
A_{t} & :=\frac{1}{2} \int_{0}^{t}\left\langle\operatorname{Ad}_{\xi_{\tau}} \operatorname{ad}_{X_{i}}^{2} V_{j-1}, u\right\rangle d \tau,
\end{aligned}
$$

implies that

$$
\begin{equation*}
\sup _{u \in U} \mu\left(\Omega_{1}(u) \cap \Omega_{2}(u)\right)=O\left(\varepsilon^{\infty-}\right), \tag{5.10}
\end{equation*}
$$

where

$$
\begin{aligned}
& \Omega_{1}(u):=\left\{\int_{0}^{T_{\delta}}\left\langle\operatorname{Ad}_{\xi_{\tau}} V_{j-1}, u\right\rangle^{2} d t<\varepsilon^{q}\right\}, \\
& \Omega_{2}(u):=\left\{\int_{0}^{T_{\delta}} \sum_{i=1}^{k}\left\langle\operatorname{Ad}_{\xi_{\tau}}\left[X_{i}, V_{j-1}\right], u\right\rangle^{2} d \tau \geq \varepsilon\right\},
\end{aligned}
$$

and $q>4$. Since

$$
\begin{aligned}
\sup _{u \in U} \mu\left(\left[\Omega_{2}(u)\right]^{c}\right) & =\sup _{u \in U} \mu\left(\int_{0}^{T_{\delta}} \sum_{i=1}^{k}\left\langle\operatorname{Ad}_{\xi_{\tau}}\left[X_{i}, V_{j-1}\right], u\right\rangle^{2} d \tau<\varepsilon\right) \\
& \leq \sup _{u \in U} \mu\left(\int_{0}^{T_{\delta}}\left\langle\operatorname{Ad}_{\xi_{\tau}} V_{j}, u\right\rangle^{2} d \tau<\varepsilon\right),
\end{aligned}
$$

we may apply the induction hypothesis to learn

$$
\begin{equation*}
\sup _{u \in U} \mu\left(\left[\Omega_{2}(u)\right]^{c}\right)=O\left(\varepsilon^{\infty-}\right) . \tag{5.11}
\end{equation*}
$$

It now follows from Equations (5.10) and (5.11) that

$$
\begin{aligned}
\sup _{u \in U} \mu\left(\Omega_{1}(u)\right) & \leq \sup _{u \in U} \mu\left(\Omega_{1}(u) \cap \Omega_{2}(u)\right)+\sup _{u \in U} \mu\left(\Omega_{1}(u) \cap\left[\Omega_{2}(u)\right]^{c}\right) \\
& \leq \sup _{u \in U} \mu\left(\Omega_{1}(u) \cap \Omega_{2}(u)\right)+\sup _{u \in U} \mu\left(\left[\Omega_{2}(u)\right]^{c}\right) \\
& =O\left(\varepsilon^{\infty-}\right)+O\left(\varepsilon^{\infty-}\right)=O\left(\varepsilon^{\infty-}\right),
\end{aligned}
$$

that is,

$$
\sup _{u \in U} \mu\left(\int_{0}^{T_{\delta}}\left\langle\operatorname{Ad}_{\xi_{\tau}} V_{j-1}, u\right\rangle^{2} d t<\varepsilon^{q}\right)=O\left(\varepsilon^{\infty-}\right) .
$$

Replacing $\varepsilon$ by $\varepsilon^{1 / q}$ in the previous equation and using $O\left(\left(\varepsilon^{1 / q}\right)^{\infty-}\right)=O\left(\varepsilon^{\infty-}\right)$, completes the induction argument and hence the proof.

The following proposition is contained in the appendix to [15] and is included here for completeness.

Proposition 5.4 (A martingale inequality). Let $T$ be a stopping time bounded by $t_{0}<\infty$, and let $Y=y+M+A$, where $M$ is a continuous martingale and $A$ is a process of bounded variation such that $M_{0}=A_{0}=0$. Further assume that, on the set $\{t \leq T\}$, $\langle M\rangle_{t}$ and $|A|_{t}$ are absolutely continuous functions and there exist finite positive constants $c_{1}$ and $c_{2}$ such that

$$
\frac{d\langle M\rangle_{t}}{d t} \leq c_{1} \text { and } \frac{d|A|_{t}}{d t} \leq c_{2} .
$$

Then, for all $\nu>0$ and $q>\nu+4$, there exist constants $c=c\left(t_{0}, q, \nu, c_{1}, c_{2}\right)>0$ and $\varepsilon_{0}=\varepsilon_{0}\left(t_{0}, q, \nu, c_{1}, c_{2}\right)>0$ such that

$$
\begin{equation*}
P\left(\int_{0}^{T} Y_{t}^{2} d t<\varepsilon^{q},\langle Y\rangle_{T}=\langle M\rangle_{T} \geq \varepsilon\right) \leq 2 \exp \left(-\frac{1}{2 c \varepsilon^{\nu}}\right)=O\left(\varepsilon^{-\infty}\right) \tag{5.12}
\end{equation*}
$$

for all $\varepsilon \in\left(0, \varepsilon_{0}\right]$.

### 5.3 Spectral decomposition of the Heisenberg heat kernel

In this section, we choose our notation to correspond to that in Thangavelu [48]; any uncited assertions in this section may be found in this text. Let $G$ denote the

3 -dimensional Heisenberg Lie group, and let $g=(z, t) \in G$. Then for any $k=0,1,2, \ldots$, let

$$
L_{k}(t):=\left.\frac{1}{k!}\left(\frac{\partial}{\partial x}\right)^{k}\right|_{x=0} \frac{e^{-t x /(1-x)}}{1-x}
$$

and

$$
e_{k}^{\lambda}(z, t):=e^{i \lambda t} L_{k}\left(\frac{1}{2}|\lambda \| z|^{2}\right) e^{-\left.\frac{1}{4}|\lambda| z\right|^{2}} .
$$

The functions $L_{k}$ are Laguerre polynomials of type 0 , and the following identity is easy to verify,

$$
\begin{equation*}
\sum_{k=0}^{\infty} L_{k}(t) x^{k}=\frac{e^{-t x /(1-x)}}{1-x} . \tag{5.13}
\end{equation*}
$$

The functions $e_{k}^{\lambda}$ are eigenfunctions of $L$ with corresponding eigenvalues $(2 k+1)|\lambda|$; since $L$ is left invariant, this implies that for any $f \in C_{c}^{\infty}(G)$

$$
\begin{equation*}
L\left(f * e_{k}^{\lambda}\right)=f * L e_{k}^{\lambda}=(2 k+1)|\lambda|\left(f * e_{k}^{\lambda}\right) . \tag{5.14}
\end{equation*}
$$

Using these definitions, we have the following proposition.
Proposition 5.5. The Heisenberg heat kernel may be written as the eigenfunction expansion,

$$
p_{\tau}(z, t)=\sum_{k=0}^{\infty} \int_{\mathbb{R}} e^{-\tau(2 k+1)|\lambda|} e_{k}^{\lambda}(z, t) d \mu(\lambda),
$$

where $d \mu(\lambda)=(2 \pi)^{-2}|\lambda| d \lambda$ is Plancherel measure for $G$.
Proof. Theorem 2.1.1 of Thangavelu [48] states that, for any $f \in L^{2}(G)$, we have the expansion

$$
\begin{equation*}
f(z, t)=\sum_{k=0}^{\infty} \int_{\mathbb{R}} f * e_{k}^{\lambda}(z, t) d \mu(\lambda) \tag{5.15}
\end{equation*}
$$

For each $k=0,1,2, \ldots$,

$$
\int_{\mathbb{R}} f * e_{k}^{\lambda}(z, t) d \mu(\lambda)=f * \int_{\mathbb{R}} e_{k}^{\lambda}(z, t) d \mu(\lambda)
$$

where the right hand side should be interpreted in the principal value sense; see Section 2.2 of 48. Additionally, since $f, e_{k}^{\lambda} \in L^{2}(G)$ and $e_{k}^{\lambda}$ is summable in $k$, then

$$
\sum_{k=0}^{\infty} f * \int_{\mathbb{R}} e_{k}^{\lambda} d \mu(\lambda)=f * \sum_{k=0}^{\infty} \int_{\mathbb{R}} e_{k}^{\lambda} d \mu(\lambda)
$$

Then, by Equation 5.14 and the fact that $\left\{e^{-\tau(2 k+1)|\lambda|}\right\}_{k=0}^{\infty} \subset L^{2}(d \mu)$ is summable in $k$,

$$
e^{-\tau L / 2} f=e^{-\tau L / 2}\left(f * \sum_{k=0}^{\infty} \int_{\mathbb{R}} e_{k}^{\lambda} d \mu(\lambda)\right)=f * \sum_{k=0}^{\infty} \int_{\mathbb{R}} e^{-\tau(2 k+1)|\lambda|} e_{k}^{\lambda} d \mu(\lambda),
$$

and thus we may write the heat kernel of $P_{\tau}=e^{\tau L / 2}$ as

$$
p_{\tau}(z, t)=\sum_{k=0}^{\infty} \int_{\mathbb{R}} e^{-\tau(2 k+1)|\lambda|} e_{k}^{\lambda}(z, t) d \mu(\lambda),
$$

as desired.
We would like to reassure ourselves that this decomposition is equivalent to the expression we already have for the Heisenberg heat kernel in Equation 2.5). Thus, using the definitions of $L_{k}$ and $e_{k}^{\lambda}$ given at the beginning of this section, we have

$$
\begin{aligned}
p_{\tau}(z, t)= & \sum_{k=0}^{\infty} \int_{\mathbb{R}} e^{-\tau(2 k+1)|\lambda|} e_{k}^{\lambda}(z, t) d \mu(\lambda) \\
= & \frac{1}{4 \pi^{2}} \sum_{k=0}^{\infty} \int_{\mathbb{R}} e^{i \lambda t} e^{-\tau(2 k+1)|\lambda|} e^{-\frac{1}{4}|\lambda||z|^{2}} L_{k}\left(\frac{1}{2}|\lambda||z|^{2}\right)|\lambda| d \lambda \\
= & \frac{1}{4 \pi^{2}} \int_{\mathbb{R}} e^{i \lambda t} e^{-\tau|\lambda|} e^{-\frac{1}{4}|\lambda||z|^{2}}|\lambda| \\
& \times\left[\left.\sum_{k=0}^{\infty} \frac{\left(e^{-2 \tau|\lambda|}\right)^{k}}{k!}\left(\frac{\partial}{\partial x}\right)^{k}\right|_{x=0} \frac{e^{-\frac{1}{2}|\lambda||z|^{2} x /(1-x)}}{1-x}\right] d \lambda \\
= & \frac{1}{4 \pi^{2}} \int_{\mathbb{R}} e^{i \lambda t} e^{-\tau|\lambda|} e^{-\frac{1}{4}|\lambda||z|^{2}}|\lambda| \frac{\exp \left(-\frac{1}{2}|\lambda||z|^{2} \frac{e^{-2 \tau|\lambda|}}{1-e^{-2 \tau|\lambda|}}\right)}{1-e^{-2 \tau|\lambda|}} d \lambda \\
= & \frac{1}{4 \pi^{2}} \int_{\mathbb{R}} e^{i \lambda t} \frac{|\lambda|}{e^{\tau|\lambda|}-e^{-\tau|\lambda|}} \exp \left(-\frac{1}{2}|\lambda||z|^{2}\left(\frac{1}{2}+\frac{e^{-\tau|\lambda|}}{e^{\tau|\lambda|}-e^{-\tau|\lambda|}}\right)\right) d \lambda \\
= & \frac{1}{8 \pi^{2}} \int_{\mathbb{R}} \frac{\lambda}{\sinh \tau \lambda} \exp \left(-\frac{1}{4}|\lambda||z|^{2} \operatorname{coth} \tau \lambda\right) e^{i \lambda t} d \lambda,
\end{aligned}
$$

where we have applied the identity in Equation (5.13). Thus the expressions are equivalent, and we are satisfied.

### 5.4 Covariance matrix of $\xi_{t}$ on the Heisenberg group

In this section, we explicitly compute the Malliavin covariance matrix of $\xi$ in the Heisenberg case,

$$
\sigma_{t}:=\xi_{t}^{\prime}(b) \xi_{t}^{\prime}(b)^{*}: \mathbb{R}^{3} \rightarrow \mathbb{R}^{3},
$$

where the adjoint here is taken relative to the Cameron-Martin and Euclidean metrics. Using Equation 2.21, we calculate $\xi_{t}^{\prime}(b): T_{b} W \cong H \rightarrow T_{\xi_{t}(b)} G \cong \mathbb{R}^{3}$ as follows.

$$
\left.\left.\begin{array}{rl}
\xi_{t}^{\prime}(b) h & =\partial_{h} \xi_{t}(b)=\left.\frac{d}{d \epsilon}\right|_{0} \xi_{t}(b+\epsilon h) \\
& =\left.\frac{d}{d \epsilon}\right|_{0}\left(\begin{array}{c}
b_{t}^{1}+\epsilon h_{t}^{1} \\
b_{t}^{2}+\epsilon h_{t}^{2} \\
\frac{1}{2} \int_{0}^{t}\left[\left(b_{s}^{1}+\epsilon h_{s}^{1}\right)\left(d b_{s}^{2}+\epsilon d h_{s}^{2}\right)-\left(b_{s}^{2}+\epsilon h_{s}^{2}\right)\left(d b_{s}^{1}+\epsilon d h_{s}^{1}\right)\right]
\end{array}\right) \\
& h_{t}^{1}  \tag{5.16}\\
h_{t}^{2} \\
\frac{1}{2} \int_{0}^{t}\left[h_{s}^{1} d b_{s}^{2}+b_{s}^{1} d h_{s}^{2}-h_{s}^{2} d b_{s}^{1}-b_{s}^{2} d h_{s}^{1}\right]
\end{array}\right), \begin{array}{c}
\int_{0}^{t} \dot{h}_{s}^{1} d s \\
\int_{0}^{t} \dot{h}_{s}^{2} d s \\
\int_{0}^{t}\left[\left(\frac{1}{2} b_{t}^{2}-b_{s}^{2}\right) \dot{h}_{s}^{1}-\left(\frac{1}{2} b_{t}^{1}-b_{s}^{1}\right) \dot{h}_{s}^{2}\right] d s
\end{array}\right), ~\left(\begin{array}{c} 
\\
\end{array}\right),
$$

by integration by parts. To determine the action of the adjoint $\xi_{t}^{\prime}(b)^{*}: \mathbb{R}^{3} \rightarrow H$ on an element $a \in \mathbb{R}^{3}$, fix $h \in H$, and consider the following computations.

$$
\begin{align*}
\left(\xi_{t}^{\prime}(b)^{*}\right. & a, h)_{H}=a \cdot \xi_{t}^{\prime}(b) h \\
& =a_{1} \int_{0}^{t} \dot{h}_{s}^{1} d s+a_{2} \int_{0}^{t} \dot{h}_{s}^{2} d s+a_{3} \int_{0}^{t}\left[\left(\frac{1}{2} b_{t}^{2}-b_{s}^{2}\right) \dot{h}_{s}^{1}-\left(\frac{1}{2} b_{t}^{1}-b_{s}^{1}\right) \dot{h}_{s}^{2}\right] d s  \tag{5.17}\\
& =\int_{0}^{t}\left[\left(a_{1}+\frac{1}{2} a_{3} b_{t}^{2}-a_{3} b_{s}^{2}\right) \dot{h}_{s}^{1}+\left(a_{2}-\frac{1}{2} a_{3} b_{t}^{1}+a_{3} b_{s}^{1}\right) \dot{h}_{s}^{2}\right] d s .
\end{align*}
$$

By definition of the Cameron-Martin inner product,

$$
\left(\xi_{t}^{\prime}(b)^{*} a, h\right)_{H}=\int_{0}^{1} \frac{d}{d s}\left[\xi_{t}^{\prime}(w)^{*} a\right]_{s} \cdot \dot{h}_{s} d s
$$

and so Equation (5.17) then implies

$$
\begin{equation*}
\frac{d}{d s}\left[\xi_{t}^{\prime}(b)^{*} a\right]_{s}=\binom{a_{1}+\frac{1}{2} a_{3} b_{t}^{2}-a_{3} b_{s}^{2}}{a_{2}-\frac{1}{2} a_{3} b_{t}^{1}+a_{3} b_{s}^{1}} \tag{5.18}
\end{equation*}
$$

Use (5.16) to write the vector $\xi_{t}^{\prime}(b) \xi_{t}^{\prime}(b)^{*} a \in \mathbb{R}^{3}$ componentwise:

$$
\begin{gathered}
{\left[\xi_{t}^{\prime}(b) \xi_{t}^{\prime}(b)^{*} a\right]_{1}=\int_{0}^{t}\left(a_{1}+\frac{1}{2} a_{3} b_{t}^{2}-a_{3} b_{s}^{2}\right) d s=a_{1} t+\frac{1}{2} a_{3} t b_{t}^{2}-a_{3} \overline{b_{t}^{2}},} \\
{\left[\xi_{t}^{\prime}(b) \xi_{t}^{\prime}(b)^{*} a\right]_{2}=\int_{0}^{t}\left(a_{2}-\frac{1}{2} a_{3} b_{t}^{1}+a_{3} b_{s}^{1}\right) d s=a_{2} t-\frac{1}{2} a_{3} t b_{t}^{1}+a_{3} \overline{b_{t}^{1}}, \text { and }} \\
{\left[\xi_{t}^{\prime}(b) \xi_{t}^{\prime}(b)^{*} a\right]_{3}} \\
=\int_{0}^{t}\left[\left(\frac{1}{2} b_{t}^{2}-b_{s}^{2}\right)\left(a_{1}+\frac{1}{2} a_{3} b_{t}^{2}-a_{3} b_{s}^{2}\right)-\left(\frac{1}{2} b_{t}^{1}-b_{s}^{1}\right)\left(a_{2}-\frac{1}{2} a_{3} b_{t}^{1}+a_{3} b_{s}^{1}\right)\right] d s \\
=a_{1}\left(\frac{1}{2} t b_{t}^{2}-\overline{b_{t}^{2}}\right)-a_{2}\left(\frac{1}{2} t b_{t}^{1}-\overline{b_{t}^{1}}\right) \\
\quad+a_{3}\left(\frac{1}{4} t\left(\left(b_{t}^{1}\right)^{2}+\left(b_{t}^{2}\right)^{2}\right)-b_{t}^{1} \overline{b_{t}^{1}}-b_{t}^{2} \overline{b_{t}^{2}}+\int_{0}^{t}\left(\left(b_{s}^{1}\right)^{2}+\left(b_{s}^{2}\right)^{2}\right) d s\right),
\end{gathered}
$$

where $\overline{b^{i}}=\int_{0}^{t} b_{s}^{i} d s$ for $i=1,2$. So the Malliavin covariance matrix may be written as

$$
\xi_{t}^{\prime}(b) \xi_{t}^{\prime}(b)^{*}=\left(\begin{array}{ccc}
t & 0 & \frac{1}{2} t b_{t}^{2}-\overline{b_{t}^{2}}  \tag{5.19}\\
0 & t & -\frac{1}{2} t b_{t}^{1}+\overline{b_{t}^{1}} \\
\frac{1}{2} t b_{t}^{2}-\overline{b_{t}^{2}} & -\frac{1}{2} t b_{t}^{1}+\overline{b_{t}^{1}} & \Xi(t, b)
\end{array}\right),
$$

where

$$
\Xi(t, b)=\frac{1}{4} t|b|^{2}-b_{t}^{1} \overline{b_{t}^{1}}-b_{t}^{2} \overline{b_{t}^{2}}+\|b\|_{2}^{2},
$$

with $\left|b_{t}\right|^{2}=\left(b_{t}^{1}\right)^{2}+\left(b_{t}^{2}\right)^{2}$ and $\|b\|_{2}=\left(\int_{0}^{t}\left|b_{s}\right|^{2} d s\right)^{1 / 2}$ the norm in $L^{2}([0, t])$.

Alternately, we may use the more general formula (3.9) for the Malliavin covariance matrix to determine $\sigma_{t}$ in the Heisenberg case. For $g=(x, y, z) \in G$ and $(a, b, c) \in \mathbb{R}^{3}$,

$$
\begin{aligned}
\operatorname{Ad}_{g}(a, b, c) & =(x, y, z)(a, b, c)(x, y, z)^{-1} \\
& =\left(x+a, y+b, z+c+\frac{1}{2}(x b-a y)\right)(-x,-y,-z)=(a, b, c+x b-a y) .
\end{aligned}
$$

Thus, for $X, Y, Z \in \mathfrak{g} \cong \mathbb{R}^{3}$,

$$
\operatorname{Ad}_{\xi_{t}} X=\left(1,0,-b_{t}^{2}\right), \operatorname{Ad}_{\xi_{t}} Y=\left(0,1, b_{t}^{1}\right), \text { and } \operatorname{Ad}_{\xi_{t}} Z=(0,0,1)
$$

implies that

$$
\operatorname{Ad}_{\xi_{t}}=\left(\begin{array}{ccc}
1 & 0 & 0 \\
0 & 1 & 0 \\
-b_{t}^{2} & b_{t}^{1} & 1
\end{array}\right) .
$$

(Equivalently, $\operatorname{Ad}_{g}=R_{g^{-1} *} L_{g *}$, and so, from Equations (2.3) and (2.4),

$$
\left.\operatorname{Ad}_{\xi_{t}}=R_{\xi_{t}^{-1} *} L_{\xi_{t *}}=\left(\begin{array}{ccc}
1 & 0 & 0 \\
0 & 1 & 0 \\
-\frac{1}{2} b_{t}^{2} & \frac{1}{2} b_{t}^{1} & 1
\end{array}\right)\left(\begin{array}{ccc}
1 & 0 & 0 \\
0 & 1 & 0 \\
-\frac{1}{2} b_{t}^{2} & \frac{1}{2} b_{t}^{1} & 1
\end{array}\right)=\left(\begin{array}{ccc}
1 & 0 & 0 \\
0 & 1 & 0 \\
-b_{t}^{2} & b_{t}^{1} & 1
\end{array}\right) .\right)
$$

Thus,

$$
\operatorname{Ad}_{\xi_{s}} P \operatorname{Ad}_{\xi_{s}}^{\dagger}=\left(\begin{array}{ccc}
1 & 0 & 0 \\
0 & 1 & 0 \\
-b_{s}^{2} & b_{s}^{1} & 1
\end{array}\right)\left(\begin{array}{ccc}
1 & 0 & 0 \\
0 & 1 & 0 \\
0 & 0 & 0
\end{array}\right)\left(\begin{array}{ccc}
1 & 0 & -b_{s}^{2} \\
0 & 1 & b_{s}^{1} \\
0 & 0 & 1
\end{array}\right)=\left(\begin{array}{ccc}
1 & 0 & -b_{s}^{2} \\
0 & 1 & b_{s}^{1} \\
-b_{s}^{2} & b_{s}^{1} & \left|b_{s}\right|^{2}
\end{array}\right),
$$

and combining this with Equation (3.9) gives

$$
\begin{aligned}
\xi_{t}^{\prime}(b) \xi_{t}^{\prime}(b)^{*} & =R_{\xi_{t} *}\left(\int_{0}^{t} \operatorname{Ad}_{\xi_{s}} P \operatorname{Ad}_{\xi_{s}}^{\dagger} d s\right) R_{\xi_{t} *}^{\operatorname{tr}} \\
& =\left(\begin{array}{ccc}
1 & 0 & 0 \\
0 & 1 & 0 \\
\frac{1}{2} b_{t}^{2} & -\frac{1}{2} b_{t}^{1} & 1
\end{array}\right)\left(\begin{array}{ccc}
t & 0 & -\overline{b_{t}^{2}} \\
0 & t & \overline{b_{t}^{1}} \\
-\overline{b_{t}^{2}} & \overline{b_{t}^{1}} & \|b\|_{2}^{2}
\end{array}\right)\left(\begin{array}{ccc}
1 & 0 & \frac{1}{2} b_{t}^{2} \\
0 & 1 & -\frac{1}{2} b_{t}^{1} \\
0 & 0 & 1
\end{array}\right) \\
& =\left(\begin{array}{ccc}
t & 0 & \frac{1}{2} t b_{t}^{2}-\overline{b_{t}^{2}} \\
0 & t & -\frac{1}{2} t b_{t}^{1}+\overline{b_{t}^{1}} \\
\frac{1}{2} t b_{t}^{2}-\overline{b_{t}^{2}} & -\frac{1}{2} t b_{t}^{1}+\overline{b_{t}^{1}} & \Xi(t, b)
\end{array}\right),
\end{aligned}
$$

as in Equation 5.19).
The determinant of the covariance matrix may be determined from these computations,

$$
\begin{aligned}
\operatorname{det} \sigma_{t} & =t^{2}\left[\frac{1}{4} t\left|b_{t}\right|^{2}-\frac{1}{2} b_{t}^{1} \overline{b_{t}^{1}}-\frac{1}{2} b_{t}^{1} \overline{b_{t}^{1}}+\|b\|_{2}^{2}\right]-t\left(-\frac{1}{2} t b_{t}^{1}+\overline{b_{t}^{1}}\right)^{2}-t\left(-\frac{1}{2} t b_{t}^{2}+\overline{b_{t}^{2}}\right)^{2} \\
& =t^{2}\|b\|_{2}^{2}-t\left(\overline{b_{t}^{1}}\right)^{2}-t\left(\overline{b_{t}^{2}}\right)^{2} .
\end{aligned}
$$

Writing $\left|\overline{b_{t}}\right|^{2}=\left(\overline{b_{t}^{1}}\right)^{2}+\left(\overline{b_{t}^{2}}\right)^{2}$ and $\Delta=\operatorname{det} \sigma$,

$$
\Delta_{t}=t^{2}\|b\|_{2}^{2}-t\left|\overline{b_{t}}\right|^{2}
$$

Proposition 5.6. Let $B_{s}=\frac{1}{\sqrt{t}} b_{t s}$. Then

$$
\Delta_{t}=\operatorname{det} \sigma_{t}=t^{4}\left(\int_{0}^{1}\left|B_{s}\right|^{2} d s-\left|\int_{0}^{1} B_{s} d s\right|^{2}\right) .
$$

Proof. Note that

$$
\|b\|_{2}^{2}=\int_{0}^{t}\left|b_{r}\right|^{2} d r=t \int_{0}^{1}\left|b_{t s}\right|^{2} d s=t^{2} \int_{0}^{1}\left|\frac{1}{\sqrt{t}} b_{t s}\right|^{2} d s=t^{2} \int_{0}^{1}\left|B_{s}\right|^{2} d s
$$

and for $i=1,2$

$$
\overline{b_{t}^{i}}=\int_{0}^{t} b_{r}^{i} d r=t \int_{0}^{1} b_{t s}^{i} d s=t^{3 / 2} \int_{0}^{1} \frac{1}{\sqrt{t}} b_{t s}^{i} d s=t^{3 / 2} \int_{0}^{1} B_{s}^{i} d s
$$

Thus,

$$
t^{2}\|b\|_{2}^{2}-t\left|\overline{b_{t}}\right|^{2}=t^{4}\left(\int_{0}^{1}\left|B_{s}\right|^{2} d s-\left|\int_{0}^{1} B_{s} d s\right|^{2}\right)
$$

### 5.5 Heisenberg lifted vector fields

Here we compute the lift of the Heisenberg vector fields,

$$
\tilde{X}=\partial_{x}-\frac{1}{2} y \partial_{x}, \quad \tilde{Y}=\partial_{y}+\frac{1}{2} x \partial_{z}, \quad \text { and } \tilde{Z}=\partial_{z}
$$

For now, let $A=\left(a_{1}, a_{2}, a_{3}\right)$ be any element of $\mathfrak{g}$. Then

$$
\tilde{A}\left(\xi_{t}\right)=L_{\xi_{t^{*}}} A=\left(a_{1}, a_{2}, a_{3}+\frac{1}{2}\left(a_{2} b_{t}^{1}-a_{1} b_{t}^{2}\right)\right)
$$

is an element of $T_{\xi_{t}} G \cong \mathbb{R}^{3}$. The inverse of the matrix $\sigma$ may be computed directly from Equation (5.19), so that for $t>0$

$$
\begin{aligned}
& {\left[\left(\xi_{t}^{\prime}(b) \xi_{t}^{\prime}(b)^{*}\right)^{-1}\right]_{11}=\Delta_{t}^{-1}\left(t\|b\|_{2}^{2}-\left(\overline{b_{t}^{1}}\right)^{2}-t b_{t}^{2} \overline{\bar{b}_{t}^{2}}+\frac{1}{4}\left(b_{t}^{2}\right)^{2}\right)} \\
& {\left[\left(\xi_{t}^{\prime}(b) \xi_{t}^{\prime}(b)^{*}\right)^{-1}\right]_{12}=\Delta_{t}^{-1}\left(-\overline{b_{t}^{1} b_{t}^{2}}+\frac{1}{2} t b_{t}^{2} \overline{b_{t}^{1}}+\frac{1}{2} t b_{t}^{1} \overline{b_{t}^{2}}-\frac{1}{4} t^{2} b_{t}^{1} b_{t}^{2}\right)} \\
& {\left[\left(\xi_{t}^{\prime}(b) \xi_{t}^{\prime}(b)^{*}\right)^{-1}\right]_{13}=\Delta_{t}^{-1}\left(t \overline{b_{t}^{2}}-\frac{1}{2} t^{2} b_{t}^{2}\right)} \\
& {\left[\left(\xi_{t}^{\prime}(b) \xi_{t}^{\prime}(b)^{*}\right)^{-1}\right]_{21}=\Delta_{t}^{-1}\left(-\overline{b_{t}^{1} b_{t}^{2}}+\frac{1}{2} t b_{t}^{2} \overline{b_{t}^{1}}+\frac{1}{2} t b_{t}^{1} \overline{b_{t}^{2}}-\frac{1}{4} t^{2} b_{t}^{1} b_{t}^{2}\right)} \\
& {\left[\left(\xi_{t}^{\prime}(b) \xi_{t}^{\prime}(b)^{*}\right)^{-1}\right]_{22}=\Delta_{t}^{-1}\left(t\|b\|_{2}^{2}-\left(\overline{b_{t}^{2}}\right)^{2}-t b_{t}^{1} \overline{b_{t}^{1}}+\frac{1}{4}\left(b_{t}^{1}\right)^{2}\right)} \\
& {\left[\left(\xi_{t}^{\prime}(b) \xi_{t}^{\prime}(b)^{*}\right)^{-1}\right]_{23}=\Delta_{t}^{-1}\left(-t \overline{b_{t}^{1}}+\frac{1}{2} t^{2} b_{t}^{1}\right)} \\
& {\left[\left(\xi_{t}^{\prime}(b) \xi_{t}^{\prime}(b)^{*}\right)^{-1}\right]_{31}=\Delta_{t}^{-1}\left(t \overline{b_{t}^{2}}-\frac{1}{2} t^{2} b_{t}^{2}\right)} \\
& {\left[\left(\xi_{t}^{\prime}(b) \xi_{t}^{\prime}(b)^{*}\right)^{-1}\right]_{32}=\Delta_{t}^{-1}\left(-t \overline{b_{t}^{1}}+\frac{1}{2} t^{2} b_{t}^{1}\right)} \\
& {\left[\left(\xi_{t}^{\prime}(b) \xi_{t}^{\prime}(b)^{*}\right)^{-1}\right]_{33}=\Delta_{t}^{-1} t^{2} .}
\end{aligned}
$$

Thus, writing $\xi_{t}^{\prime}(b) \xi_{t}^{\prime}(b)^{*} \tilde{A}\left(\xi_{t}\right)$ componentwise,

$$
\begin{aligned}
& {\left[\left(\xi_{t}^{\prime}(b) \xi_{t}^{\prime}(b)^{*}\right)^{-1} \tilde{A}\left(\xi_{t}\right)\right]_{1}} \\
& \begin{array}{l}
=a_{1}\left(t\|b\|_{2}^{2}-\left(\overline{b_{t}^{1}}\right)^{2}-t b_{t}^{2} \overline{b^{2}}+\frac{1}{4}\left(b_{t}^{2}\right)^{2}\right)+a_{2}\left(-\overline{b_{t}^{1} b_{t}^{2}}+\frac{1}{2} t b_{t}^{2} \overline{b_{t}^{1}}+\frac{1}{2} t b_{t}^{1} \overline{b_{t}^{2}}-\frac{1}{4} t^{2} b_{t}^{1} b_{t}^{2}\right) \\
\quad+\left(a_{3}+\frac{1}{2}\left(a_{2} b_{t}^{1}-a_{1} b_{t}^{2}\right)\right)\left(t \overline{b_{t}^{2}}-\frac{1}{2} t^{2} b_{t}^{2}\right) \\
=a_{1}\left(t\|b\|_{2}^{2}-\left(\overline{b_{t}^{1}}\right)^{2}-\frac{3}{2} t b_{t}^{2} \overline{b_{t}^{2}}+\frac{1}{4}\left(1+t^{2}\right)\left(b_{t}^{2}\right)^{2}\right) \\
\quad+a_{2}\left(-\overline{b_{t}^{1} b_{t}^{2}}+\frac{1}{2} t b_{t}^{2} \overline{b_{t}^{1}}+t b_{t}^{1} \overline{b^{2}}-\frac{1}{2} t^{2} b_{t}^{1} b_{t}^{2}\right)+a_{3}\left(t \overline{b_{t}^{2}}-\frac{1}{2} t^{2} b_{t}^{2}\right) \\
\left.\quad \begin{array}{l}
{\left[\left(\xi_{t}^{\prime}(b) \xi_{t}^{\prime}(b)^{*}\right)^{-1} \tilde{A}\left(\xi_{t}\right)\right]_{2}} \\
=a_{1}\left(-\overline{b_{t}^{1} b_{t}^{2}}+\frac{1}{2} t b_{t}^{2} \overline{b_{t}^{1}}+\frac{1}{2} t b_{t}^{1} \overline{b_{t}^{2}}-\frac{1}{4} t^{2} b_{t}^{1} b_{t}^{2}\right)+a_{2}\left(t\|b\|_{2}^{2}-\bar{b}_{2}^{2}-t b_{t}^{1} \overline{b_{t}^{1}}+\frac{1}{4}\left(b_{t}^{1}\right)^{2}\right) \\
\quad+\left(a_{3}+\frac{1}{2}\left(a_{2} b_{t}^{1}-a_{1} b_{t}^{2}\right)\right)\left(-t \overline{b_{t}^{1}}+\frac{1}{2} t^{2} b_{t}^{1}\right) \\
=a_{1}\left(-\overline{b_{t}^{1} b_{t}^{2}}+t b_{t}^{2} \overline{b_{t}^{1}}+\frac{1}{2} t b_{t}^{1} \overline{b_{t}^{2}}-\frac{1}{2} t^{2} b_{t}^{1} b_{t}^{2}\right) \\
\quad+a_{2}\left(t\|b\|_{2}^{2}-\left(\overline{b_{t}^{2}}\right)^{2}-\frac{3}{2} t b_{t}^{1} \bar{b}_{t}^{1}\right.
\end{array} \frac{1}{4}\left(1+t^{2}\right)\left(b_{t}^{1}\right)^{2}\right)+a_{3}\left(-t \overline{b_{t}^{1}}+\frac{1}{2} t^{2} b_{t}^{1}\right)
\end{array} \\
& {\left[\left(\xi_{t}^{\prime}(b) \xi_{t}^{\prime}(b)^{*}\right)^{-1} \tilde{A}\left(\xi_{t}\right)\right]_{3}} \\
& \quad=a_{1}\left(t \overline{t b_{t}^{2}}-\frac{1}{2} t^{2} b_{t}^{2}\right)+a_{2}\left(-t \overline{b_{t}^{1}}+\frac{1}{2} t^{2} b_{t}^{1}\right)+\left(a_{3}+\frac{1}{2}\left(a_{2} b_{t}^{1}-a_{1} b_{t}^{2}\right)\right) t^{2} \\
& \quad=a_{1}\left(t \overline{b_{t}^{2}}-t^{2} b_{t}^{2}\right)-a_{2}\left(t \overline{b_{t}^{1}}-t^{2} b_{t}^{1}\right)+a_{3} t^{2} .
\end{aligned}
$$

Equation (5.18) implies that

$$
\left[\xi_{t}^{\prime}(w)^{*} a\right]_{s}=\binom{a_{1} s+\frac{1}{2} a_{3} s b_{t}^{2}-a_{3} \int_{0}^{s} b_{r}^{2} d r}{a_{2} s-\frac{1}{2} a_{3} s b_{t}^{1}+a_{3} \int_{0}^{s} b_{r}^{1} d r}
$$

for any $a \in \mathbb{R}^{3}$, and so

$$
\left[\xi_{t}^{\prime}(b)^{*}\right]_{s}=\left(\begin{array}{ccc}
s & 0 & \frac{1}{2} s b_{t}^{2}-\overline{b_{s}^{2}} \\
0 & s & -\frac{1}{2} s b_{t}^{1}+\overline{b_{s}^{1}}
\end{array}\right)
$$

Thus, for $0<s \leq t$,

$$
\begin{aligned}
& \mathbf{A}_{s}^{1}= {\left[\xi_{t}^{\prime}(b)^{*}\left(\xi_{t}^{\prime}(b) \xi_{t}^{\prime}(b)^{*}\right)^{-1} \tilde{A}\left(\xi_{t}\right)\right]_{1} } \\
&= \Delta_{t}^{-1}\left[a_{1} s\left(t\|b\|_{2}^{2}-\left(\overline{b_{t}^{1}}\right)^{2}-\frac{3}{2} t b_{t}^{2} \overline{b^{2}}+\frac{1}{4}\left(1+t^{2}\right)\left(b_{t}^{2}\right)^{2}\right)\right. \\
&+a_{2} s\left(-\overline{b_{t}^{1} b_{t}^{2}}+\frac{1}{2} t b_{t}^{2} \overline{b_{t}^{1}}+t b_{t}^{1} \overline{\bar{b}_{t}^{2}}-\frac{1}{2} t^{2} b_{t}^{1} b_{t}^{2}\right)+a_{3} s\left(t \overline{b_{t}^{2}}-\frac{1}{2} t^{2} b_{t}^{2}\right) \\
&\left.+\left(\frac{1}{2} s b_{t}^{2}-\overline{b_{s}^{2}}\right)\left(a_{1}\left(t \overline{b_{t}^{2}}-t^{2} b_{t}^{2}\right)-a_{2}\left(t \overline{b^{1}}-t^{2} b_{t}^{1}\right)+a_{3} t^{2}\right)\right] \\
& \mathbf{A}_{s}^{2}=[ \left.\xi_{t}^{\prime}(b)^{*}\left(\xi_{t}^{\prime}(b) \xi_{t}^{\prime}(b)^{*}\right)^{-1} \tilde{A}\left(\xi_{t}\right)\right]_{2} \\
&=\Delta_{t}^{-1}\left[a_{1} s\left(-\overline{b_{1}} \overline{b_{t}^{2}}+t b_{t}^{2} \overline{b_{t}^{1}}+\frac{1}{2} t b_{t}^{1} \overline{b_{t}^{2}}-\frac{1}{2} t^{2} b_{t}^{1} b_{t}^{2}\right)\right. \\
&+a_{2} s\left(t\|b\|_{2}^{2}-\left(\overline{b_{t}^{2}}\right)^{2}-\frac{3}{2} t b_{t}^{1} \overline{b_{t}^{1}}+\frac{1}{4}\left(1+t^{2}\right)\left(b_{t}^{1}\right)^{2}\right)+a_{3} s\left(-t \overline{b_{t}^{1}}+\frac{1}{2} t^{2} b_{t}^{1}\right) \\
&\left.+\left(-\frac{1}{2} s b_{t}^{1}+\overline{b_{s}^{1}}\right)\left(a_{1}\left(t \overline{b_{t}^{2}}-t^{2} b_{t}^{2}\right)-a_{2}\left(t \overline{b^{1}}-t^{2} b_{t}^{1}\right)+a_{3} t^{2}\right)\right] .
\end{aligned}
$$

Thus, $\mathbf{X}_{s}=\left(\mathbf{X}_{s}^{1}, \mathbf{X}_{s}^{2}\right)$, where $\mathbf{X}^{1}$ and $\mathbf{X}^{2}$ are the quadratic expressions in $b$

$$
\mathbf{X}_{s}^{1}=\Delta_{t}^{-1}\left[s\left(t\|b\|_{2}^{2}-\left(\overline{b_{t}^{1}}\right)^{2}-\frac{3}{2} t b_{t}^{2} \overline{b_{t}^{2}}+\frac{1}{4}\left(1+t^{2}\right)\left(b_{t}^{2}\right)^{2}\right)+\left(\frac{1}{2} s b_{t}^{2}-\overline{b_{s}^{2}}\right)\left(t \overline{b_{t}^{2}}-t^{2} b_{t}^{2}\right)\right]
$$

and

$$
\mathbf{X}_{s}^{2}=\Delta_{t}^{-1}\left[s\left(-\overline{b_{t}^{1} b_{t}^{2}}+t b_{t}^{2} \overline{b_{t}^{1}}+\frac{1}{2} t b_{t}^{1} \overline{b^{2}}-\frac{1}{2} t^{2} b_{t}^{1} b_{t}^{2}\right)+\left(-\frac{1}{2} s b_{t}^{1}+\overline{b_{s}^{1}}\right)\left(\overline{t b_{t}^{2}}-t^{2} b_{t}^{2}\right)\right]
$$

Also, $\mathbf{Y}_{s}=\left(\mathbf{Y}_{s}^{1}, \mathbf{Y}_{s}^{2}\right)$, where

$$
\mathbf{Y}_{s}^{1}=\Delta_{t}^{-1}\left[s\left(-\overline{b_{t}^{1} b_{t}^{2}}+\frac{1}{2} t b_{t}^{2} \overline{b_{t}^{1}}+t b_{t}^{1} \bar{b}^{2}-\frac{1}{2} t^{2} b_{t}^{1} b_{t}^{2}\right)+\left(\frac{1}{2} s b_{t}^{2}-\overline{b_{s}^{2}}\right)\left(t \overline{b_{t}^{1}}-t^{2} b_{t}^{1}\right)\right]
$$

and

$$
\begin{aligned}
\mathbf{Y}_{s}^{2}=\Delta_{t}^{-1}\left[s\left(t\|b\|_{2}^{2}-\left(\overline{b_{t}^{2}}\right)^{2}-\frac{3}{2} t b_{t}^{1} \overline{b_{t}^{1}}+\frac{1}{4}\left(1+t^{2}\right)\left(b_{t}^{1}\right)^{2}\right)\right. & \\
& \left.+\left(-\frac{1}{2} s b_{t}^{1}+\overline{b_{s}^{1}}\right)\left(t \overline{b_{t}^{1}}-t^{2} b_{t}^{1}\right)\right]
\end{aligned}
$$

and $\mathbf{Z}_{s}=\left(\mathbf{Z}_{s}^{1}, \mathbf{Z}_{s}^{2}\right)$ where

$$
\mathbf{Z}_{s}^{1}=\Delta_{t}^{-1}\left[s\left(t \overline{b_{t}^{2}}-\frac{1}{2} t^{2} b_{t}^{2}\right)+\left(\frac{1}{2} s b_{t}^{2}-\overline{b_{s}^{2}}\right) t^{2}\right]
$$

and

$$
\mathbf{Z}_{s}^{2}=\Delta_{t}^{-1}\left[s\left(-t \overline{b_{t}^{1}}+\frac{1}{2} t^{2} b_{t}^{1}\right)+\left(-\frac{1}{2} s b_{t}^{1}+\overline{b_{s}^{1}}\right) t^{2}\right] .
$$

## Bibliography

[1] I. Agricola and T. Friedrich. Global analysis, volume 52 of Graduate Studies in Mathematics. American Mathematical Society, Providence, RI, 2002. Differential forms in analysis, geometry and physics, Translated from the 2001 German original by Andreas Nestke.
[2] G. Alexopoulos. An application of homogenization theory to harmonic analysis: Harnack inequalities and Riesz transforms on Lie groups of polynomial growth. Canad. J. Math., 44(4):691-727, 1992.
[3] P. Auscher, T. Coulhon, X. Duong, and S. Hofmann. Riesz transform on manifolds and heat kernel regularity. Preprint, 2004.
[4] D. Bakry. Ricci curvature and dimension for diffusion semigroups. In Stochastic processes and their applications in mathematics and physics (Bielefeld, 1985), volume 61 of Math. Appl., pages 21-31. Kluwer Acad. Publ., Dordrecht, 1990.
[5] D. Bakry and M. Émery. Hypercontractivité de semi-groupes de diffusion. C. R. Acad. Sci. Paris Sér. I Math., 299(15):775-778, 1984.
[6] D. Bakry and M. Émery. Diffusions hypercontractives. In Séminaire de probabilités, XIX, 1983/84, volume 1123 of Lecture Notes in Math., pages 177-206. Springer, Berlin, 1985.
[7] D. R. Bell. The Malliavin calculus, volume 34 of Pitman Monographs and Surveys in Pure and Applied Mathematics. Longman Scientific \& Technical, Harlow, 1987.
[8] D. R. Bell and S. E. A. Mohammed. The Malliavin calculus and stochastic delay equations. J. Funct. Anal., 99(1):75-99, 1991.
[9] A. Bellaïche. Tangent space in sub-Riemannian geometry. In Sub-Riemannian Geometry, volume 144 of Progress in Mathematics, pages 1-78. Birkhäuser, Basel, 1996.
[10] V. I. Bogachev. Gaussian measures, volume 62 of Mathematical Surveys and Monographs. American Mathematical Society, Providence, RI, 1998.
[11] T. Coulhon and X. T. Duong. Riesz transform and related inequalities on noncompact Riemannian manifolds. Comm. Pure Appl. Math., 56(12):1728-1751, 2003.
[12] G. Da Prato and J. Zabczyk. Stochastic equations in infinite dimensions, volume 44 of Encyclopedia of Mathematics and its Applications. Cambridge University Press, Cambridge, 1992.
[13] E. B. Davies. One-parameter semigroups, volume 15 of London Mathematical Society Monographs. Academic Press Inc. [Harcourt Brace Jovanovich Publishers], London, 1980.
[14] B. K. Driver. A Cameron-Martin type quasi-invariance theorem for Brownian motion on a compact Riemannian manifold. J. Funct. Anal., 110(2):272-376, 1992.
[15] B. K. Driver. Curved wiener space analysis. Preprint, to appear in "Real and Stochastic Analysis: New Perspectives," xxx.lanl.gov/list/math/0403, 2004.
[16] B. K. Driver and L. Gross. Hilbert spaces of holomorphic functions on complex Lie groups. In New trends in stochastic analysis (Charingworth, 1994), pages 76-106. World Sci. Publishing, River Edge, NJ, 1997.
[17] K. D. Elworthy. Stochastic differential equations on manifolds, volume 70 of London Mathematical Society Lecture Note Series. Cambridge University Press, Cambridge, 1982.
[18] G. B. Folland. A fundamental solution for a subelliptic operator. Bull. Amer. Math. Soc., 79:373-376, 1973.
[19] G. B. Folland and E. M. Stein. Hardy Spaces on Homogeneous Groups. Princeton University Press, Princeton, N.J., 1982.
[20] M. Fukushima. Dirichlet forms and Markov processes, volume 23 of North-Holland Mathematical Library. North-Holland Publishing Co., Amsterdam, 1980.
[21] I I. Ī. Gīhman and A. V. Skorohod. Stochastic differential equations. Springer-Verlag, New York, 1972. Translated from the Russian by Kenneth Wickwire, Ergebnisse der Mathematik und ihrer Grenzgebiete, Band 72.
[22] R. W. Goodman. Nilpotent Lie groups: structure and applications to analysis. Springer-Verlag, Berlin, 1976. Lecture Notes in Mathematics, Vol. 562.
[23] L. Gross. Logarithmic Sobolev inequalities. Amer. J. Math., 97(4):1061-1083, 1975.
[24] L. Gross. Logarithmic Sobolev inequalities and contractivity properties of semigroups. In Dirichlet forms (Varenna, 1992), volume 1563 of Lecture Notes in Math., pages 54-88. Springer, Berlin, 1993.
[25] E. Hille. Functional Analysis and Semi-Groups. American Mathematical Society Colloquium Publications, vol. 31. American Mathematical Society, New York, 1948.
[26] L. Hörmander. Hypoelliptic second order differential equations. Acta Math., 119:147-171, 1967.
[27] E. P. Hsu. Stochastic analysis on manifolds, volume 38 of Graduate Studies in Mathematics. American Mathematical Society, Providence, RI, 2002.
[28] N. Ikeda and S. Watanabe. An introduction to Malliavin's calculus. In Stochastic analysis (Katata/Kyoto, 1982), volume 32 of North-Holland Math. Library, pages 1-52. North-Holland, Amsterdam, 1984.
[29] N. Ikeda and S. Watanabe. Malliavin calculus of Wiener functionals and its applications. In From local times to global geometry, control and physics (Coventry, 1984/85), volume 150 of Pitman Res. Notes Math. Ser., pages 132-178. Longman Sci. Tech., Harlow, 1986.
[30] N. Ikeda and S. Watanabe. Stochastic differential equations and diffusion processes, volume 24 of North-Holland Mathematical Library. North-Holland Publishing Co., Amsterdam, second edition, 1989.
[31] K. Itô, editor. Proceedings of the International Symposium on Stochastic Differential Equations. John Wiley \& Sons, New York-Chichester-Brisbane, 1978. Held at the Research Institute for Mathematical Sciences, Kyoto Univ., Kyoto, July 9-14, 1976, A Wiley-Interscience Publication.
[32] S. Kusuoka and D. Stroock. Applications of the Malliavin calculus. I. In Stochastic analysis (Katata/Kyoto, 1982), volume 32 of North-Holland Math. Library, pages 271-306. North-Holland, Amsterdam, 1984.
[33] S. Kusuoka and D. Stroock. Applications of the Malliavin calculus. II. J. Fac. Sci. Univ. Tokyo Sect. IA Math., 32(1):1-76, 1985.
[34] S. Kusuoka and D. Stroock. Applications of the Malliavin calculus. III. J. Fac. Sci. Univ. Tokyo Sect. IA Math., 34(2):391-442, 1987.
[35] M. Ledoux. The geometry of Markov diffusion generators. Ann. Fac. Sci. Toulouse Math. (6), 9(2):305-366, 2000. Probability theory.
[36] P. Malliavin. $C^{k}$-hypoellipticity with degeneracy. In Stochastic analysis (Proc. Internat. Conf., Northwestern Univ., Evanston, Ill., 1978), pages 199-214. Academic Press, New York, 1978.
[37] P. Malliavin. $C^{k}$-hypoellipticity with degeneracy. II. In Stochastic analysis (Proc. Internat. Conf., Northwestern Univ., Evanston, Ill., 1978), pages 327-340. Academic Press, New York, 1978.
[38] P. Malliavin. Stochastic calculus of variation and hypoelliptic operators. In Proceedings of the International Symposium on Stochastic Differential Equations (Res. Inst. Math. Sci., Kyoto Univ., Kyoto, 1976), pages 195-263, New York, 1978. Wiley.
[39] J. Norris. Simplified Malliavin calculus. In Séminaire de Probabilités, XX, 1984/85, volume 1204 of Lecture Notes in Math., pages 101-130. Springer, Berlin, 1986.
[40] D. Nualart. The Malliavin calculus and related topics. Probability and its Applications (New York). Springer-Verlag, New York, 1995.
[41] A. Pazy. Semigroups of linear operators and applications to partial differential equations, volume 44 of Applied Mathematical Sciences. Springer-Verlag, New York, 1983.
[42] J. Picard. Gradient estimates for some diffusion semigroups. Probab. Theory Related Fields, 122(4):593-612, 2002.
[43] L. P. Rothschild and E. M. Stein. Hypoelliptic differential operators and nilpotent groups. Acta Math., 137(3-4):247-320, 1976.
[44] I. Shigekawa. Derivatives of Wiener functionals and absolute continuity of induced measures. J. Math. Kyoto Univ., 20(2):263-289, 1980.
[45] D. W. Stroock. The Malliavin calculus and its application to second order parabolic differential equations. I. Math. Systems Theory, 14(1):25-65, 1981.
[46] D. W. Stroock. The Malliavin calculus and its application to second order parabolic differential equations. II. Math. Systems Theory, 14(2):141-171, 1981.
[47] S. Taniguchi. Malliavin's stochastic calculus of variations for manifold-valued Wiener functionals and its applications. Z. Wahrsch. Verw. Gebiete, 65(2):269290, 1983.
[48] S. Thangavelu. Harmonic analysis on the Heisenberg group, volume 159 of Progress in Mathematics. Birkhäuser Boston Inc., Boston, MA, 1998.
[49] V. S. Varadarajan. Lie groups, Lie algebras, and their representations, volume 102 of Graduate Texts in Mathematics. Springer-Verlag, New York, 1984. Reprint of the 1974 edition.
[50] N. T. Varopoulos. Analysis on nilpotent groups. J. Funct. Anal., 66(3):406-431, 1986.
[51] N. T. Varopoulos, L. Saloff-Coste, and T. Coulhon. Analysis and geometry on groups, volume 100 of Cambridge Tracts in Mathematics. Cambridge University Press, Cambridge, 1992.
[52] S. Watanabe. Lectures on stochastic differential equations and Malliavin calculus, volume 73 of Tata Institute of Fundamental Research Lectures on Mathematics and Physics. Published for the Tata Institute of Fundamental Research, Bombay, 1984. Notes by M. Gopalan Nair and B. Rajeev.
[53] S. Watanabe. Analysis of Wiener functionals (Malliavin calculus) and its applications to heat kernels. Ann. Probab., 15(1):1-39, 1987.
[54] S. Watanabe. Stochastic analysis on Wiener space. In Proceedings of the Norbert Wiener Centenary Congress, 1994 (East Lansing, MI, 1994), volume 52 of Proc. Sympos. Appl. Math., pages 359-370, Providence, RI, 1997. Amer. Math. Soc.
[55] K. Yosida. Functional analysis, volume 123 of Grundlehren der Mathematischen Wissenschaften [Fundamental Principles of Mathematical Sciences]. SpringerVerlag, Berlin, sixth edition, 1980.

