## NOTES FOR MATH 188

STEVEN V SAM

## Contents

1. Linear recurrence relations ..... 2
1.1. Setup ..... 2
1.2. Proofs of Theorem 1.4 ..... 4
1.3. Generalizations ..... 5
2. Formal power series ..... 9
2.1. Definitions ..... 9
2.2. Binomial theorem ..... 13
2.3. Choice problems ..... 16
2.4. Rational generating functions ..... 18
2.5. Catalan numbers ..... 20
3. Fundamental counting problems ..... 22
3.1. 12-fold way, introduction ..... 22
3.2. Compositions ..... 22
3.3. Words ..... 23
3.4. Set partitions ..... 25
3.5. Integer partitions ..... 28
3.6. 12-fold way, summary ..... 31
3.7. Cycles in permutations ..... 32
3.8. Counting subspaces ..... 34
4. Walks in graphs ..... 37
4.1. Adjacency matrix ..... 37
4.2. Transfer matrix method ..... 39
5. Exponential generating functions ..... 42
5.1. Products of exponential generating functions ..... 42
5.2. Compositions of exponential generating functions ..... 45
5.3. Cayley's enumeration of labeled trees and Lagrange inversion ..... 46
6. Sieving methods ..... 49
6.1. Möbius inversion ..... 49
6.2. Boolean poset and inclusion-exclusion ..... 52
6.3. Divisor poset and classical Möbius inversion ..... 55
7. Group actions ..... 57
7.1. Terminology ..... 57
7.2. Burnside's lemma ..... 59
7.3. Redfield-Pólya theory ..... 61
7.4. Proving congruences ..... 63

## 1. Linear Recurrence Relations

1.1. Setup. A sequence of numbers $\left(a_{n}\right)_{n \geq 0}$ is said to satisfy a (homogeneous) linear recurrence relation of order $d$ if there are scalars $c_{1}, \ldots, c_{d}$ such that $c_{d} \neq 0$, and for all $n \geq d$, we have

$$
a_{n}=c_{1} a_{n-1}+c_{2} a_{n-2}+\cdots+c_{d} a_{n-d} .
$$

Example 1.1. The Fibonacci numbers $f_{n}$ are given by the sequence $0,1,1,2,3,5,8,13,21, \ldots$. This isn't really telling you what the general $f_{n}$ is, so instead let me say that for all $n \geq 2$, we have

$$
f_{n}=f_{n-1}+f_{n-2}
$$

Together with the initial conditions $f_{0}=0, f_{1}=1$, this is enough information to calculate any $f_{n}$. So (by definition), the Fibonacci numbers satisfy a linear recurrence relation of order 2.

In general, if we want to define a sequence using a linear recurrence relation of order $d$, we need to specify the first $d$ initial values $a_{0}, a_{1}, \ldots, a_{d-1}$ to allow us to calculate all of the terms.

Our goal here is to get closed formulas for sequences that satisfy linear recurrence relations.
Example 1.2. When $d=1$, this is easy to do:

$$
a_{n}=c_{1} a_{n-1}=c_{1}^{2} a_{n-2}=c_{1}^{3} a_{n-3}=\cdots=c_{1}^{n} a_{0} .
$$

So now we'll focus on the case $d=2$. So we have a sequence of numbers $a_{0}, a_{1}, a_{2}, \ldots$ that satisfies a recurrence relation of the form

$$
a_{n}=c_{1} a_{n-1}+c_{2} a_{n-2}
$$

whenever $n \geq 2$ (here $c_{1}, c_{2}$ are some constants and $c_{2} \neq 0$ ). We want to find a closed formula for $a_{n}$.

The characteristic polynomial of this recurrence relation is defined to be

$$
t^{2}-c_{1} t-c_{2}
$$

The roots of this polynomial are $\frac{c_{1} \pm \sqrt{c_{1}^{2}+4 c_{2}}}{2}$. Call them $r_{1}$ and $r_{2} .{ }^{1}$ So we can factor the characteristic polynomial as

$$
\begin{equation*}
t^{2}-c_{1} t-c_{2}=\left(t-r_{1}\right)\left(t-r_{2}\right) . \tag{1.3}
\end{equation*}
$$

Comparing constant terms, we get $r_{1} r_{2}=-c_{2}$, so $r_{1} \neq 0$ and $r_{2} \neq 0$ because we assumed that $c_{2} \neq 0$.

Here is the first statement:
Theorem 1.4. If $r_{1} \neq r_{2}$, then there are constants $\alpha_{1}$ and $\alpha_{2}$ such that

$$
a_{n}=\alpha_{1} r_{1}^{n}+\alpha_{2} r_{2}^{n}
$$

for all $n \geq 0$.

[^0]To solve for the coefficients, plug in $n=0$ and $n=1$ to get

$$
\begin{aligned}
& a_{0}=\alpha_{1}+\alpha_{2} \\
& a_{1}=r_{1} \alpha_{1}+r_{2} \alpha_{2} .
\end{aligned}
$$

Then you have to solve for $\alpha_{1}, \alpha_{2}\left(a_{0}, a_{1}\right.$ are part of the original sequence, so are given to you).

Example 1.5. Let's finish with the example of the Fibonacci numbers $f_{n}$. These are defined by

$$
\begin{aligned}
& f_{0}=0 \\
& f_{1}=1 \\
& f_{n}=f_{n-1}+f_{n-2} \quad \text { for } n \geq 2 .
\end{aligned}
$$

So the characteristic polynomial is $t^{2}-t-1$. Its roots are $\frac{1 \pm \sqrt{5}}{2}$. Set $r_{1}=(1+\sqrt{5}) / 2$ and $r_{2}=(1-\sqrt{5}) / 2$. So we have

$$
f_{n}=\alpha_{1} r_{1}^{n}+\alpha_{2} r_{2}^{n}
$$

and we have to solve for $\alpha_{1}$ and $\alpha_{2}$. Plug in $n=0,1$ to get:

$$
\begin{aligned}
& 0=\alpha_{1}+\alpha_{2} \\
& 1=\alpha_{1} r_{1}+\alpha_{2} r_{2} .
\end{aligned}
$$

So $\alpha_{1}=-\alpha_{2}$; plug this into the second formula to get $1=\alpha_{1}\left(r_{1}-r_{2}\right)=\alpha_{1} \sqrt{5}$. So $\alpha_{1}=1 / \sqrt{5}$ and $\alpha_{2}=-1 / \sqrt{5}$. In conclusion:

$$
f_{n}=\frac{1}{\sqrt{5}}\left(\frac{1+\sqrt{5}}{2}\right)^{n}-\frac{1}{\sqrt{5}}\left(\frac{1-\sqrt{5}}{2}\right)^{n} .
$$

Example 1.6. Consider a periodic sequence $x, y, x, y, \ldots$ (starting at $a_{0}$ ). This satisfies the linear recurrence relation $a_{n}=a_{n-2}$ for $n \geq 2$ so the characteristic polynomial is $t^{2}-1=$ $(t-1)(t+1)$. Hence we have $a_{n}=\alpha_{1}+\alpha_{2}(-1)^{n}$ for $\alpha_{1}, \alpha_{2}$ that satisfy $\alpha_{1}+\alpha_{2}=x$ and $\alpha_{1}-\alpha_{2}=y$, i.e., $\alpha_{1}=(x+y) / 2$ and $\alpha_{2}=(x-y) / 2$.

Remark 1.7. Here is another way to think about the theorem. Pick scalars $c_{1}, c_{2}$ such that $t^{2}-c_{1} t-c_{2}$ has distinct roots. If $\left(a_{n}\right)$ is a solution, i.e., $a_{n}=c_{1} a_{n-1}+c_{2} a_{n-2}$ and ( $a_{n}^{\prime}$ ) is also a solution, then so is any linear combination $\left(\gamma a_{n}+\delta a_{n}^{\prime}\right)$. In other words, the set of solutions to a linear recurrence relation forms a vector space. The theorem then says that $\left(r_{1}^{n}\right)$ and $\left(r_{2}^{n}\right)$ span this vector space.

Checking that they are solutions is straightforward:

$$
r_{1}^{n}-c_{1} r_{1}^{n-1}-c_{2} r_{1}^{n-2}=r_{1}^{n-2}\left(r_{1}^{2}-c_{1} r_{1}-c_{2}\right)=0
$$

In fact, since $r_{1} \neq r_{2}$, it is clear that the sequences $\left(r_{1}^{n}\right)$ and $\left(r_{2}^{n}\right)$ are linearly independent, so in fact they form a basis. This tells us that the solution space is 2-dimensional. This last fact isn't too hard to guess: in finding a solution, we can arbitrarily specify the first two terms and everything is uniquely determined. Moreover, this 2-dimensionality still holds when $r_{1}=r_{2}$, so we'll need to find a solution that's not a multiple of $r_{1}^{n}$.

The above remark gives a proof of Theorem 1.4 once a few details are filled in, but it required knowing in advance what the answer should be. Let's now discuss two other derivations of Theorem 1.4 which are a bit more general (i.e., work on other problems) and don't require knowing the answer first.
1.2. Proofs of Theorem 1.4. We will give two proofs of Theorem 1.4: using formal power series and matrices.
1.2.1. Using formal power series. We will develop formal power series more carefully in $\S 2$. For now we will sacrifice rigor to motivate why they are useful.

Define the infinite sum ( $x$ is a variable)

$$
A(x)=\sum_{n \geq 0} a_{n} x^{n} .
$$

We call this the generating function of $\left(a_{n}\right)$. The recurrence relation says that we have an identity

$$
\begin{aligned}
A(x) & =a_{0}+a_{1} x+\sum_{n \geq 2}\left(c_{1} a_{n-1}+c_{2} a_{n-2}\right) x^{n} \\
& =a_{0}+a_{1} x+c_{1} \sum_{n \geq 2} a_{n-1} x^{n}+c_{2} \sum_{n \geq 2} a_{n-2} x^{n} .
\end{aligned}
$$

Remember the recurrence is only valid for $n \geq 2$, so we have to separate out the first two terms. Now comes an important point: the last two sums are almost the same as $A(x)$ if we re-index them:

$$
\begin{aligned}
& \sum_{n \geq 2} a_{n-1} x^{n}=\sum_{n \geq 1} a_{n} x^{n+1}=x \sum_{n \geq 1} a_{n} x^{n}=x\left(A(x)-a_{0}\right) \\
& \sum_{n \geq 2} a_{n-2} x^{n}=\sum_{n \geq 0} a_{n} x^{n+2}=x^{2} A(x) .
\end{aligned}
$$

In particular,

$$
A(x)=a_{0}+a_{1} x+c_{1} x A(x)-c_{1} a_{0} x+c_{2} x^{2} A(x)
$$

We can rewrite this as

$$
\begin{equation*}
A(x)=\frac{a_{0}+\left(a_{1}-c_{1} a_{0}\right) x}{1-c_{1} x-c_{2} x^{2}} . \tag{1.8}
\end{equation*}
$$

We want to factor the denominator. To do this, plug in $t \mapsto x^{-1}$ into (1.3) and multiply by $x^{2}$ to get

$$
1-c_{1} x-c_{2} x^{2}=\left(1-r_{1} x\right)\left(1-r_{2} x\right) .
$$

Now we can apply partial fraction decomposition to (1.8) to write

$$
A(x)=\frac{\alpha_{1}}{1-r_{1} x}+\frac{\alpha_{2}}{1-r_{2} x}
$$

for some constants $\alpha_{1}, \alpha_{2}$. But these terms are both geometric series, so we can further write

$$
A(x)=\alpha_{1} \sum_{n \geq 0} r_{1}^{n} x^{n}+\alpha_{2} \sum_{n \geq 0} r_{2}^{n} x^{n} .
$$

The coefficient of $x^{n}$ on the left side is $a_{n}$ and the coefficient of $x^{n}$ on the right side is $\alpha_{1} r_{1}^{n}+\alpha_{2} r_{2}^{n}$. So we have equality for all $n$.
1.2.2. Using matrices. Our recurrence relation translates to the following matrix equation:

$$
\left[\begin{array}{cc}
c_{1} & c_{2} \\
1 & 0
\end{array}\right]\left[\begin{array}{l}
a_{n-1} \\
a_{n-2}
\end{array}\right]=\left[\begin{array}{c}
a_{n} \\
a_{n-1}
\end{array}\right]
$$

for $n \geq 2$. Set $C=\left[\begin{array}{cc}c_{1} & c_{2} \\ 1 & 0\end{array}\right]$. Then our goal is to find a formula for $C^{n}\left[\begin{array}{l}a_{1} \\ a_{0}\end{array}\right]$ for all $n \geq 0$ because the second entry is $a_{n}$. Thinking back to linear algebra, we can do this if we can diagonalize $C$ : if $C=B D B^{-1}$ for some diagonal matrix $D$, then we have $C^{n}=B D^{n} B^{-1}$ and $D^{n}$ is easy to compute. The characteristic polynomial of $C$ is conveniently $t^{2}-c_{1} t-c_{2}$, which is the characteristic polynomial of the recurrence relation, so its eigenvalues are $r_{1}, r_{2}$. Since we are assuming they are distinct, $C$ is diagonalizable, so there is some matrix $B$ such that $C=B\left[\begin{array}{cc}r_{1} & 0 \\ 0 & r_{2}\end{array}\right] B^{-1}$.

Let's just name $\left[\begin{array}{l}x \\ y\end{array}\right]=B^{-1}\left[\begin{array}{l}a_{1} \\ a_{0}\end{array}\right]$. We don't need to compute $x, y$ but note that they are constants of our recurrence relation (they do not depend on $n$ ). Then

$$
\left[\begin{array}{c}
a_{n+1} \\
a_{n}
\end{array}\right]=C^{n}\left[\begin{array}{l}
a_{1} \\
a_{0}
\end{array}\right]=B\left[\begin{array}{cc}
r_{1}^{n} & 0 \\
0 & r_{2}^{n}
\end{array}\right]\left[\begin{array}{l}
x \\
y
\end{array}\right]=\left[\begin{array}{l}
b_{1,1} x r_{1}^{n}+b_{1,2} y r_{2}^{n} \\
b_{2,1} x r_{1}^{n}+b_{2,2} y r_{2}^{n}
\end{array}\right] .
$$

To finish the proof we just set $\alpha_{1}=b_{2,1} x$ and $\alpha_{2}=b_{2,2} y$.
1.3. Generalizations. Some questions we haven't answered yet:

- What if $r_{1}=r_{2}$ ?
- What about higher degree recurrence relations?
- What about non-homogeneous recurrence relations?
- What about non-linear recurrence relations?

We can answer the first two without much additional effort, but the last two are much wider open and don't have obvious answers (we'll address some instances throughout the course). There is a good analogy here with systems of differential equations-as you learn in Math 20D, the linear homogeneous case is formulaic to solve, but in general they are a mess.

### 1.3.1. Repeated roots.

Theorem 1.9. If $r_{1}=r_{2}$, then there are constants $\alpha_{1}$ and $\alpha_{2}$ such that

$$
a_{n}=\alpha_{1} r_{1}^{n}+\alpha_{2} n r_{1}^{n}
$$

for all $n \geq 0$.
Again, to solve for $\alpha_{1}, \alpha_{2}$, just plug in $n=0,1$ to get a system of equations:

$$
\begin{aligned}
& a_{0}=\alpha_{1} \\
& a_{1}=\alpha_{1} r_{1}+\alpha_{2} r_{1} .
\end{aligned}
$$

(From this we could solve the general case, but I think it's easier to remember the way I've written it.)

Example 1.10. Suppose $\left(a_{n}\right)_{n \geq 0}$ satisfies the recurrence relation

$$
a_{n}=4 a_{n-1}-4 a_{n-2} \quad(n \geq 2)
$$

with initial conditions $a_{0}=a_{1}=1$. The characteristic polynomial is $t^{2}-4 t+4=(t-2)^{2}$. So we have constants $\alpha_{1}, \alpha_{2}$ such that $a_{n}=\left(\alpha_{1}+\alpha_{2} n\right) 2^{n}$. Plug in $n=0$ to get $1=\alpha_{1}$ and plug in $n=1$ to get $1=2\left(\alpha_{1}+\alpha_{2}\right)$, which means that $\alpha_{2}=-1 / 2$. So we get the formula

$$
a_{n}=\left(1-\frac{n}{2}\right) 2^{n}=-2^{n-1}(n-2) .
$$

Proof. We can start in the same way as in the previous proof. The only difference is that we are trying to take the partial fraction decomposition of

$$
A(x)=\frac{a_{0}+\left(a_{1}-c_{1} a_{0}\right) x}{\left(1-r_{1} x\right)^{2}} .
$$

This can still be done, but now it looks like

$$
\frac{\beta_{1}}{1-r_{1} x}+\frac{\beta_{2}}{\left(1-r_{1} x\right)^{2}}
$$

for some constants $\beta_{1}, \beta_{2}$. The first is a geometric series, and the second is obtained from taking a derivative of the geometric series: $1 /(1-x)^{2}=\sum_{n \geq 0}(n+1) x^{n}$. So we get instead

$$
A(x)=\beta_{1} \sum_{n \geq 0} r_{1}^{n} x^{n}+\beta_{2} \sum_{n \geq 0}(n+1) r_{1}^{n} x^{n} .
$$

Comparing coefficients, we get

$$
a_{n}=\beta_{1} r_{1}^{n}+\beta_{2}(n+1) r_{1}^{n}=\left(\beta_{1}+\beta_{2}\right) r_{1}^{n}+\beta_{2} n r_{1}^{n} .
$$

So $\alpha_{1}=\beta_{1}+\beta_{2}$ and $\alpha_{2}=\beta_{2}$.
How does the matrix proof adapt? The matrix

$$
C=\left[\begin{array}{cc}
c_{1} & c_{2} \\
1 & 0
\end{array}\right]
$$

now has a repeated eigenvalue $r$. It cannot be diagonalizable: if it were, then it must be $r$ times the identity since we would have $C=B\left[\begin{array}{ll}r & 0 \\ 0 & r\end{array}\right] B^{-1}=\left[\begin{array}{ll}r & 0 \\ 0 & r\end{array}\right]$. However, all matrices have something called a Jordan normal form. For non-diagonalizable $2 \times 2$ matrices, this means there is an invertible $2 \times 2$ matrix $B$ so that

$$
C=B\left[\begin{array}{ll}
r & 1 \\
0 & r
\end{array}\right] B^{-1}
$$

Remark 1.11. Why is this possible? We know that there is an eigenvector $v$ for $C$ with eigenvalue $r$. Let $w$ be any vector which is not a multiple of $v$. Then $C w=a v+b w$ for some scalars, and the matrix for $C$ with respect to the basis is $\left[\begin{array}{ll}r & a \\ 0 & b\end{array}\right]$. But $\operatorname{det} C=r^{2}$, so we see that $b=r$. Next, $a \neq 0$ since otherwise $C$ would be $r$ times the identity matrix. If $w^{\prime}=w / a$, then $C w^{\prime}=r w^{\prime}+v$, so that with respect to the basis $\left\{v, w^{\prime}\right\}$, the matrix for $C$ is now $\left[\begin{array}{ll}r & 1 \\ 0 & r\end{array}\right]$. So we can take $B$ to be the matrix whose columns are $v$ and $w^{\prime}$. (The existence of Jordan normal form for larger matrices requires a lot more effort.)

Note then that $C=B\left[\begin{array}{ll}r & 1 \\ 0 & r\end{array}\right]^{n} B^{-1}$ and that $\left[\begin{array}{cc}r & 1 \\ 0 & r\end{array}\right]^{n}=\left[\begin{array}{cc}r^{n} & n r^{n-1} \\ 0 & r^{n}\end{array}\right]$. Now we finish in the same way as before. Set $\left[\begin{array}{l}x \\ y\end{array}\right]=B^{-1}\left[\begin{array}{l}a_{1} \\ a_{0}\end{array}\right]$. Then we have

$$
\left[\begin{array}{c}
a_{n+1} \\
a_{n}
\end{array}\right]=C^{n}\left[\begin{array}{l}
a_{1} \\
a_{0}
\end{array}\right]=B\left[\begin{array}{cc}
r^{n} & n r^{n-1} \\
0 & r^{n}
\end{array}\right]\left[\begin{array}{l}
x \\
y
\end{array}\right]=B\left[\begin{array}{c}
x r^{n}+y n r^{n-1} \\
y r^{n}
\end{array}\right] .
$$

It gets a little messy to expand fully, but the point is that the second component is both $a_{n}$ and also a linear combination of $r^{n}$ and $n r^{n-1}$ (to match the theorem, we can always adjust the coefficient of $n r^{n-1}$ by $r$ to make it $n r^{n}$ ).
1.3.2. Higher order relations. Higher order recurrence relations

$$
a_{n}=c_{1} a_{n-1}+\cdots+c_{d} a_{n-d}
$$

can be solved in the same way: one has to first find the roots of the characteristic polynomial $t^{d}-c_{1} t^{d-1}-c_{2} t^{d-2}-\cdots-c_{d}$ and apply partial fraction decomposition as in the proof above. The simplest case is when the roots $r_{1}, \ldots, r_{d}$ are all distinct. In this case, we can say that there exist constants $\alpha_{1}, \ldots, \alpha_{d}$ such that

$$
a_{n}=\alpha_{1} r_{1}^{n}+\cdots+\alpha_{d} r_{d}^{n}
$$

for all $n$. In order to solve for $\alpha_{1}, \ldots, \alpha_{d}$, we have to consider $n=0, \ldots, d-1$ separately to get a system of $d$ linear equations in $d$ variables. When the roots appear with multiplicities, we have to do something like we did in Theorem 1.9. For example, if $d=5$ and the roots are $r_{1}$ with multiplicity 3 and $r_{2}$ with multiplicity 2 ( and $r_{1} \neq r_{2}$ ), then we would have

$$
a_{n}=\alpha_{1} r_{1}^{n}+\alpha_{2} n r_{1}^{n}+\alpha_{3} n^{2} r_{1}^{n}+\alpha_{4} r_{2}^{n}+\alpha_{5} n r_{2}^{n}
$$

This should look familiar to you if you've ever solved a linear homogeneous differential equation with constant coefficients.

I'll leave it to you to formulate the general case, but I'll focus on a really important case when all roots equal 1.
Example 1.12. If the characteristic polynomial is $(t-1)^{d}$, then the pattern tells us that there are constants $\alpha_{1}, \ldots, \alpha_{d}$ such that

$$
a_{n}=\alpha_{1}+\alpha_{2} n+\cdots+\alpha_{d} n^{d-1}
$$

i.e., that the sequence $\left(a_{n}\right)_{n \geq 0}$ is given by a polynomial.

Here's a different perspective. Let $T$ be the translation operation on sequences that shifts down by 1, i.e., $(T a)_{n}=a_{n+1}$. This is a linear operator on the vector space of all sequences. A recurrence relation

$$
a_{n}=c_{1} a_{n-1}+\cdots+c_{d} a_{n-d}
$$

can be rewritten as a single equation $T^{d} a=c_{1} T^{d-1} a+\cdots+c_{d} a$. If $r_{1}, \ldots, r_{d}$ are the roots of the characteristic polynomial, then we can factor this as

$$
\left(T-r_{1}\right) \cdots\left(T-r_{d}\right) a=0 .
$$

Going back to previous discussion, the solution space to the recurrence relation is the same thing as the null space of the linear operator $\left(T-r_{1}\right) \cdots\left(T-r_{d}\right)$. The previous example then translates to:

Proposition 1.13. Given a sequence $\left(a_{n}\right)_{n \geq 0}$, there is a polynomial $p(n)$ of degree $\leq d-1$ such that $a_{n}=p(n)$ if and only if $(T-1)^{d} a=0$.
Proof. We've basically seen one direction already: if $(T-1)^{d} a=0$, then $a_{n}$ is given by a polynomial of degree $\leq d-1$. Let's consider the converse statement. Suppose $p(n)=$ $p_{d-1} n^{d-1}+\cdots+p_{1} n+p_{0}$ is a polynomial of degree $\leq d-1$. Then $(T-1)$ applied to the sequence $(p(n))$ is $(p(n+1)-p(n))$ and $p(n+1)$ again has degree $\leq d-1$ and the coefficient of $n^{d-1}$ is again $p_{d-1}$, so the difference has degree $\leq d-2$. So by induction, we see that $(T-1)^{d} p=0$ (the base case $d=1$ is when we have a degree 0 polynomial, which means it is constant).
Remark 1.14. The operation $T-1$ can be interpreted as a discrete analogue of the derivative. The discrete analogue of integration $S$ should be defined as $(S a)_{n}=a_{0}+\cdots+a_{n-1}$. Then we see that $(T-1) S$ is the identity and $S(T-1)$ is the identity up to a constant, which are discrete versions of the fundamental theorem of calculus.

For more information, look up finite difference calculus.

### 1.3.3. Non-homogeneous recurrence relations.

Example 1.15. Consider a non-homogeneous linear recurrence relation of degree 1:

$$
a_{n}=c a_{n-1}+d
$$

where $c, d$ are constants. We can iterate the recursion to guess a formula:

$$
\begin{aligned}
a_{n} & =c\left(c a_{n-2}+d\right)+d=c^{2} a_{n-2}+(c+1) d \\
& =c^{3} a_{n-3}+\left(c^{2}+c+1\right) d \\
& =\cdots \\
& =c^{n} a_{0}+\left(c^{n}+\cdots+c+1\right) d .
\end{aligned}
$$

If $c=1$, this simplifies to $a_{n}=a_{0}+n d$ and otherwise we can write $a_{n}=c^{n} a_{0}+\frac{1-c^{n}}{1-c} d$.
Here's a different way to approach it. If we take the difference of $a_{n}=c a_{n-1}+d$ and $a_{n-1}=c a_{n-2}+d$, we get the relation $a_{n}=(c+1) a_{n-1}-c a_{n-2}$ for $n \geq 2$. The characteristic polynomial is $t^{2}-(c+1) t+c=(t-1)(t-c)$, so if $c \neq 1$, we must have coefficients $\alpha_{1}, \alpha_{2}$ such that $a_{n}=\alpha_{1} c^{n}+\alpha_{2}$. If $c=1$, then we have coefficients $\beta_{1}, \beta_{2}$ such that $a_{n}=\beta_{1}+\beta_{2} n$.

The approach just outlined can be generalized to show that solving a non-homogeneous linear recurrence relation of degree $d$ with constant offset reduces to solving a homogeneous linear recurrence relation of degree $d+1$. We can iterate this idea to solve a non-homogeneous linear recurrence relation with polynomial offset. I'll just illustrate the linear case and leave you to think about what happens for general polynomials. Of course, we still have the question of what to do for general offsets.

Example 1.16. Consider a non-homogeneous linear recurrence relation of degree 1 with linear offset:

$$
a_{n}=c a_{n-1}+d_{1} n+d_{2} .
$$

We try the trick from before: if $n \geq 2$, then subtract from this the equation $a_{n-1}=c a_{n-2}+$ $d_{1}(n-1)+d_{2}$ to get

$$
a_{n}=(c+1) a_{n-1}-c a_{n-2}+d_{1} .
$$

Now we do the difference trick again to this new equation to get:

$$
a_{n}=(c+2) a_{n-1}-(2 c+1) a_{n-2}+c a_{n-3}
$$

which is homogeneous with characteristic polynomial $t^{3}-(c+2) t^{2}+(2 c+1) t-c=(t-c)(t-1)^{2}$. So if $c \neq 1$, then our general solution is $a_{n}=\alpha_{1} c^{n}+\alpha_{2} n+\alpha_{3}$ and otherwise if $c=1$ it is $a_{n}=\alpha_{1} n^{2}+\alpha_{2} n+\alpha_{3}$.

## 2. Formal power series

2.1. Definitions. A formal power series (in the variable $x$ ) is an expression of the form $A(x)=\sum_{n=0}^{\infty} a_{n} x^{n}$ where the $a_{n}$ are scalars (almost exclusively rational numbers in this class $)^{2}$. Instead of writing the sum from 0 to $\infty$, we will usually just write $A(x)=\sum_{n \geq 0} a_{n} x^{n}$. If $A(x)$ is a formal power series, let $\left[x^{n}\right] A(x)$ denote the coefficient of $x^{n}$ in $A(x)$, so in this case, $\left[x^{n}\right] A(x)=a_{n}$.

The formal power series $A(x)$ is sometimes called the generating function of the sequence $\left(a_{n}\right)$. It doesn't mean anything special since every formal power series is a generating function, but it's commonly used language, so we'll use it too sometimes.

Let $B(x)=\sum_{n \geq 0} b_{n} x^{n}$ be another formal power series.
By definition, two formal power series are equal if and only if all of their coefficients match up, i.e., $A(x)=B(x)$ if and only if $a_{n}=b_{n}$ for all $n$. We can treat these like infinite degree polynomials.

The sum of two formal power series is defined by

$$
A(x)+B(x)=\sum_{n \geq 0}\left(a_{n}+b_{n}\right) x^{n} .
$$

The product is defined by

$$
A(x) B(x)=\sum_{n \geq 0} c_{n} x^{n}, \quad c_{n}=\sum_{i=0}^{n} a_{i} b_{n-i} .
$$

This is what you get if you just distribute like normal. As a special case, if $a_{i}=0$ for $i>0$, we just get

$$
a_{0} B(x)=\sum_{n \geq 0} a_{0} b_{n} x^{n} .
$$

Polynomials are special cases of formal power series: they are the ones with only finitely many nonzero coefficients. All of the above definitions are compatible with operations on polynomials as you know them.

Addition and multiplication are commutative, so $A(x)+B(x)=B(x)+A(x)$ and $A(x) B(x)=$ $B(x) A(x)$. They are also associative, so it is unambiguous how to add or multiply 3 or more power series.

Example 2.1. Let $A(x)=B(x)=\sum_{n \geq 0} x^{n}$. Then

$$
\begin{aligned}
A(x)+B(x) & =\sum_{n \geq 0} 2 x^{n} \\
A(x) B(x) & =\sum_{n \geq 0}(n+1) x^{n} .
\end{aligned}
$$

[^1]A formal power series $A(x)$ is invertible if there is a power series $B(x)$ such that $A(x) B(x)=1$. In that case, we write $B(x)=A(x)^{-1}=1 / A(x)$ and call it the inverse of $A(x)$. If it exists, then $B(x)$ is unique.

Example 2.2. Let $A(x)=\sum_{n \geq 0} x^{n}$ and $B(x)=1-x$. Then $A(x) B(x)=1$, so $B(x)$ is the inverse of $A(x)$. For that reason, we will use the expression

$$
\frac{1}{1-x}=\sum_{n \geq 0} x^{n}
$$

This is usually referred to as a geometric series.
On the other hand, the formal power series $x$ is not invertible: the constant term of $x B(x)$ is 0 no matter what $B(x)$ is, so there is no way that an inverse exists.

Theorem 2.3. A formal power series $A(x)$ is invertible if and only if its constant term is nonzero.

Proof. Write $A(x)=\sum_{n \geq 0} a_{n} x^{n}$. We want to solve $A(x) B(x)=1$ if possible. If we multiply the left side out and equate coefficients, we get the following (infinite) system of equations:

$$
\begin{aligned}
a_{0} b_{0} & =1 \\
a_{0} b_{1}+a_{1} b_{0} & =0 \\
a_{0} b_{2}+a_{1} b_{1}+a_{2} b_{0} & =0 \\
a_{0} b_{3}+a_{1} b_{2}+a_{2} b_{1}+a_{3} b_{0} & =0
\end{aligned}
$$

If $a_{0}=0$, then there is no solution to the first equation so $A(x)$ is not invertible.
If $a_{0} \neq 0$, then we can solve the equations one by one. Formally, we can prove by induction on $n$ that there exist coefficients $b_{0}, \ldots, b_{n}$ that make the first $n+1$ equations valid. For the base case $n=0$, we have $b_{0}=1 / a_{0}$. So suppose we have found the coefficients $b_{0}, \ldots, b_{n}$ already. At the next step, we will have

$$
b_{n}=-\frac{1}{a_{0}} \sum_{i=1}^{n} a_{i} b_{n-i}
$$

In the sum, we have $i>0$, so $b_{n-i}$ is a coefficient we already solved for in a previous step. Hence we get a formula for $b_{n}$ that makes the next equation valid as well.

The proof might look unsatisfactory: how can we solve an infinite system of equations simultaneously? For this, it is convenient to introduce the notion of convergence of formal power series. Let $A_{0}(x), A_{1}(x), \ldots$ be a sequence of formal power series. We say that the sequence converges to a formal power series $A(x)$ if, for all $n \geq 0$, the sequence $\left(\left[x^{n}\right] A_{i}(x)\right)_{i}$ is equal to $\left[x^{n}\right] A(x)$ for $i \gg 0$ (i.e., there exists $N_{n}$ such that if $i \geq N_{n}$, then $\left[x^{n}\right] A_{i}(x)=$ $\left[x^{n}\right] A(x)$; note that $N_{n}$ is allowed to depend on $n$ ). In that case, we write

$$
\lim _{i \rightarrow \infty} A_{i}(x)=A(x) .
$$

This is a sort of simple-minded way to define convergence: it just says that the $n$th coefficient of the sequence is eventually constant for all $n$. It is very important to separate this from the definition of convergence that you learn in calculus.

Example 2.4. The sequence $A_{i}(x)=\sum_{n=0}^{i} x^{n}$ converges to $\sum_{n \geq 0} x^{n}$ : for each $n$ the sequence $\left[x^{n}\right] A_{i}(x)$ is $n 0$ 's followed by a constant sequence of 1 's.

The sequence $A_{i}(x)=x^{i}$ converges to 0 : for each $n$ the sequence $\left[x^{n}\right] A_{i}(x)$ is a constant sequence of 0 's with a single 1 in the $n$th position.

On the other hand, the sequence $A_{i}(x)=1 /(i+1)$ does not converge in our sense: the sequence $\left[x^{0}\right] A_{i}(x)$ is not eventually constant.

The proof of the following is left for homework.
Lemma 2.5. Assume that $\lim _{i \rightarrow \infty} A_{i}(x)=A(x)$ and $\lim _{i \rightarrow \infty} B_{i}(x)=B(x)$. Then

$$
\lim _{i \rightarrow \infty}\left(A_{i}(x)+B_{i}(x)\right)=A(x)+B(x), \quad \lim _{i \rightarrow \infty}\left(A_{i}(x) B_{i}(x)\right)=A(x) B(x) .
$$

In the proof above, by solving for the coefficients $b_{n}$ one at a time, we are essentially considering the sequence of formal power series $B_{i}(x)=\sum_{n=0}^{i} b_{n} x^{n}$ and defining $B(x)=$ $\lim _{i \rightarrow \infty} B_{i}(x)$. This makes sense by the above since we're showing that

$$
1=\lim _{i \rightarrow \infty}\left(A(x) B_{i}(x)\right)=A(x) B(x)
$$

We can also define infinite sums and products. If $A_{i}(x)$ is a sequence of formal power series, then the sum $\sum_{i \geq 0} A_{i}(x)$, is defined to be the limit of the finite partial sums if it exists, i.e.,

$$
\sum_{i \geq 0} A_{i}(x)=\lim _{i \rightarrow \infty} \sum_{j=0}^{i} A_{j}(x)
$$

Similarly for infinite products:

$$
\prod_{i \geq 0} A_{i}(x)=\lim _{i \rightarrow \infty} \prod_{j=0}^{i} A_{j}(x)
$$

Since these are just special cases of limits, Lemma 2.5 applies to give the following statement.

Lemma 2.6. Assuming the relevant limits exist, we have

$$
\begin{aligned}
\sum_{i \geq 0} A_{i}(x)+\sum_{i \geq 0} B_{i}(x) & =\sum_{i \geq 0}\left(A_{i}(x)+B_{i}(x)\right) \\
\left(\prod_{i \geq 0} A_{i}(x)\right)\left(\prod_{i \geq 0} B_{i}(x)\right) & =\prod_{i \geq 0}\left(A_{i}(x) B_{i}(x)\right) .
\end{aligned}
$$

Convergence won't usually be a problem for us, so we won't be too picky about the details whenever it comes up. But do keep in mind that to be completely rigorous we always use these definitions.

We can characterize when infinite sums or products exist as follows. Given a formal power series $A(x)$, define its minimum degree, $\operatorname{mdeg}(A(x))$, to be the smallest $n$ such that $\left[x^{n}\right] A(x) \neq 0$.
Proposition 2.7. Let $A_{0}(x), A_{1}(x), \ldots$ be a sequence of formal power series.
(1) $\sum_{i \geq 0} A_{i}(x)$ exists if and only if $\lim _{i \rightarrow \infty} \operatorname{mdeg}\left(A_{i}(x)\right)=\infty$.
(2) Now assume that each $A_{i}(x)$ has no constant term. Then $\prod_{i \geq 0}\left(1+A_{i}(x)\right)$ exists if and only if $\lim _{i \rightarrow \infty} \operatorname{mdeg}\left(A_{i}(x)\right)=\infty$.
We'll omit the proof, but you can find them in Sagan's book in §3.3.
Given two formal power series $A(x)$ and $B(x)$, suppose that $A(x)$ has no constant term. Then we can define the composition by

$$
(B \circ A)(x)=B(A(x))=\sum_{n \geq 0} b_{n} A(x)^{n} .
$$

This is well-defined by the previous result: $\operatorname{mdeg}\left(b_{n} A(x)^{n}\right) \geq n$ because $A(x)$ has no constant term.

An important special case is when $A(x)=0$ is the 0 power series. In that case, $B(0)=b_{0}$ is the constant term of $B(x)$.
Example 2.8. Let $d$ be a positive integer, $A(x)=x^{d}$ and $B(x)=\sum_{n \geq 0} x^{n}$. Then $B(A(x))=$ $\sum_{n \geq 0} x^{d n}$. We can do this substitution into the identity

$$
(1-x) B(x)=1
$$

to get

$$
\left(1-x^{d}\right) \sum_{n \geq 0} x^{d n}=1
$$

from which we conclude that

$$
\frac{1}{1-x^{d}}=\sum_{n \geq 0} x^{d n}
$$

We can also take the (formal) derivative $D$ of a formal power series. We denote it by either $D A$ or $A^{\prime}$ and deifned it as follows:

$$
(D A)(x)=A^{\prime}(x)=\sum_{n \geq 0} n a_{n} x^{n-1}=\sum_{n \geq 0}(n+1) a_{n+1} x^{n} .
$$

All of the familiar properties of derivatives hold (again we'll skip the proof):

$$
\begin{aligned}
D(A+B) & =D A+D B \\
D(A \cdot B) & =(D A) \cdot B+A \cdot(D B) \\
D(B \circ A) & =(D A) \cdot(D B \circ A) \\
D(1 / A) & =-\frac{D(A)}{A^{2}} \\
D\left(A^{n}\right) & =n D(A) A^{n-1} .
\end{aligned}
$$

Example 2.9. We have $\frac{1}{1-x}=\sum_{n \geq 0} x^{n}$. Taking the derivative of the left side gives $\frac{1}{(1-x)^{2}}$. Taking the derivative of the right side gives $\sum_{n \geq 0} n x^{n-1}=\sum_{n \geq 0}(n+1) x^{n}$. We've already seen that these two expressions are equal.

How would we simplify $B(x)=\sum_{n \geq 0} n x^{n}$ ? We have a few options. First:

$$
B(x)=\sum_{n \geq 0}(n+1) x^{n}-\sum_{n \geq 0} x^{n}=\frac{1}{(1-x)^{2}}-\frac{1}{1-x}=\frac{1-(1-x)}{(1-x)^{2}}=\frac{x}{(1-x)^{2}}
$$

Or more directly:

$$
B(x)=x \sum_{n \geq 0} n x^{n-1}=x \frac{1}{(1-x)^{2}}
$$

One more important property (which should remind you of Taylor series from calculus):

$$
\left[x^{n}\right] A(x)=\frac{\left(D^{n} A\right)(0)}{n!}
$$

where recall that the factorial is defined for non-negative integers by $0!=1$ and $n!=$ $n \cdot(n-1)!=\prod_{i=1}^{n} i$ for $n>0$.
2.2. Binomial theorem. For this section, we'll assume our formal power series have complex coefficients.

Lemma 2.10. Let $A(x)$ be a formal power series with $A(0)=1$ and let d be a positive integer. There exists a unique formal power series $B(x)$ with $B(0)=1$ such that $B(x)^{d}=A(x)$.

In the notation of the lemma, we set $B(x)=A(x)^{1 / d}$.
Proof. The idea is to equate the coefficients of $B(x)^{d}$ and $A(x)$ and solve for the coefficients of $B(x)$ one by one. However, the expansion of $B(x)^{d}$ is complicated, so we want to avoid doing it in detail. However, we can say that $\left[x^{n}\right]\left(B(x)^{d}\right)=d b_{n}+f_{n, d}$ where the $f_{n, d}$ is an expression only involving $b_{1}, \ldots, b_{n-1}$. For example, $\left[x^{1}\right]\left(B(x)^{d}\right)=d b_{1}$ and $\left[x^{2}\right]\left(B(x)^{d}\right)=d b_{2}+\frac{d(d-1)}{2} b_{1}^{2}$. Hence we can proceed as in the proof of Theorem 2.3 by solving for the $b_{n}$ by induction on $n$. By assumption, we already have $b_{0}=1$ and our equation tells us that $b_{1}=a_{1} / d$. In general, we have $d b_{n}+f_{n, d}=a_{n}$. Since $f_{n, d}$ only involves $b_{1}, \ldots, b_{n-1}$, we have already solved for them by induction and we set $b_{n}=\left(a_{n}-f_{n, d}\right) / d$. The uniqueness is clear since we have no choice about how to define $b_{n}$ at each step.

Now for any integer (positive or negative) $c$, we know that $A(x)^{c}$ has constant term 1 if $A(x)$ does, so we can apply the lemma to find a formal power series $\left(A(x)^{c}\right)^{1 / d}$ with constant term 1. We can also take the $c$ th power of the formal power series $A(x)^{1 / d}$; are they the same? Note that

$$
\left(\left(A(x)^{1 / d}\right)^{c}\right)^{d}=\left(\left(A(x)^{1 / d}\right)^{d}\right)^{c}=A(x)^{c} .
$$

By uniqueness of $d$ th roots, we conclude that $\left(A(x)^{1 / d}\right)^{c}=\left(A(x)^{c}\right)^{1 / d}$, and hence we can unambiguously define $A(x)^{c / d}$ to be either of these expressions, and so we have a definition of what it means to take a rational power of a formal power series whose constant term is 1. We do also have to check that it does not depend on the $c$ and $d$, but only on their ratio, i.e., we need to know that for any other $b$, we have $A(x)^{b c / b d}=A(x)^{c / d}$; this can be done in a similar way and we'll leave it as an exercise.

Now we find an explicit formula for this when $A(x)=1+x$.
If $m$ is a rational number and $k$ is a non-negative integer, we define (generalized) binomial coefficients by

$$
\binom{m}{0}=1, \quad\binom{m}{k}=\frac{m(m-1)(m-2) \cdots(m-k+1)}{k!} \quad(k>0)
$$

Theorem 2.11 (Binomial theorem). Let $m$ be a rational number. Then

$$
(1+x)^{m}=\sum_{n \geq 0}\binom{m}{n} x^{n} .
$$

Note also that this gives a way to compute $A(x)^{m}$ whenever $A(0)=1$ : we just substitute $A(x)-1$ into $x$ in the above formula:

$$
A(x)^{m}=(1+(A(x)-1))^{m}=\sum_{n \geq 0}\binom{m}{n}(A(x)-1)^{n}
$$

This will be useful in later calculations. Let's work out a few cases.
Example 2.12. If $m$ is a non-negative integer, then $\binom{m}{k}=0$ if $k>m$ since the product has a 0 in the numerator. For $k \leq m$, we can write it in terms of factorials:

$$
\binom{m}{k}=\frac{m!}{k!(m-k)!} .
$$

The binomial theorem then says

$$
(1+x)^{m}=\sum_{n=0}^{m}\binom{m}{n} x^{n} .
$$

Example 2.13. Consider $m=-1$. We know from before that

$$
\frac{1}{1-x}=\sum_{n \geq 0} x^{n}
$$

If we substitute in $-x$ for $x$, then we get

$$
\frac{1}{1+x}=\sum_{n \geq 0}(-1)^{n} x^{n}
$$

We should also be able to get this from the binomial theorem with $m=-1$. We have

$$
\binom{-1}{n}=\frac{(-1)(-2) \cdots(-1-n+1)}{n!}=\frac{(-1)^{n} n!}{n!}=(-1)^{n}
$$

More generally, consider $m=-d$ for some positive integer $d$. Then from what we just did, we have

$$
(1+x)^{-d}=\left(\sum_{n \geq 0}(-1)^{n} x^{n}\right)^{d}
$$

The right side could be expanded, possibly by using induction on $d$, but we'd have to know a pattern before we could proceed. Instead, let's use the binomial theorem directly:

$$
\begin{aligned}
\binom{-d}{n} & =\frac{(-d)(-d-1) \cdots(-d-n+1)}{n!}=\frac{(-1)^{n}(d+n-1)(d+n-2) \cdots(d)}{n!} \\
& =(-1)^{n} \frac{(d+n-1)!}{(d-1)!n!}=(-1)^{n}\binom{d+n-1}{n} .
\end{aligned}
$$

This gives us the identities

$$
\begin{aligned}
& \frac{1}{(1+x)^{d}}=\sum_{n \geq 0}(-1)^{n}\binom{d+n-1}{n} x^{n} \\
& \frac{1}{(1-x)^{d}}=\sum_{n \geq 0}\binom{d+n-1}{n} x^{n} .
\end{aligned}
$$

Example 2.14. Consider $m=1 / 2$. Then

$$
\binom{1 / 2}{n}=\frac{(1 / 2)(-1 / 2)(-3 / 2) \cdots(1 / 2-n+1)}{n!}=\frac{(-1)^{n-1}(2 n-3)(2 n-5) \cdots 3}{2^{n} n!} .
$$

This doesn't simplify much further, so now is a good time to introduce the double factorial: if $n$ is a positive integer, we set $n!!=n(n-2)(n-4) \cdots$. In other words, if $n$ is odd, then $n!!$ is the product of all positive odd integers between 1 and $n$, and if $n$ is even, then $n!!$ is the product of all positive even integers between 2 and $n$. Keep in mind this does not mean we do the factorial twice. With our new notation, we have

$$
\binom{1 / 2}{n}=\frac{(-1)^{n-1}(2 n-3)!!}{2^{n} n!} \quad(n \geq 2)
$$

We also have $\binom{1 / 2}{0}=1$ and $\binom{1 / 2}{1}=1 / 2$. Remember that this means that

$$
\left(1+\frac{x}{2}+\sum_{n \geq 2} \frac{(-1)^{n-1}(2 n-3)!!}{2^{n} n!} x^{n}\right)^{2}=1+x
$$

To check that by hand, we could expand the left side, but it would be a lot of work.
Now let's prove the binomial theorem. We need one preparatory result.
Lemma 2.15. If $m$ is a rational number and $A(0)=1$, then $D\left(A(x)^{m}\right)=m \cdot(D A) \cdot A(x)^{m-1}$.
Proof. Write $m=p / q$ for integers $p, q$. Then

$$
p A(x)^{p-1} D A=D\left(A(x)^{p}\right)=D\left(\left(A(x)^{m}\right)^{q}\right)=q\left(A(x)^{m}\right)^{q-1} D\left(A(x)^{m}\right)
$$

Hence

$$
D\left(A(x)^{m}\right)=\frac{p A(x)^{p-1} D A}{q\left(A(x)^{m}\right)^{q-1}}=m A(x)^{m-1} D A
$$

Proof of Binomial Theorem. We have

$$
\left[x^{n}\right](1+x)^{m}=\frac{\left(D^{n}\left((1+x)^{m}\right)\right)(0)}{n!}=\frac{m(m-1) \cdots(m-n+1)}{n!}=\binom{m}{n} .
$$

In Example 2.14, we found a square root to the formal power series $1+x$. Because $(-1)^{2}=1$, if we multiplied that solution by -1 , we'd get another solution (its constant term isn't 1 , so this doesn't contradict uniqueness!). Are there more? If we were talking about numbers, then no. The same holds for formal power series too. More generally, if we're trying to solve a quadratic equation

$$
A(x) t^{2}+B(x) t+C(x)=0
$$

where $A(x), B(x), C(x)$ are formal power series, then there are at most two different solutions $t$ in formal power series (there could be only one or none). We won't prove this because it's beyond the scope of this course, but we will use this later to solve some problems.

Conveniently, the quadratic formula can be used in this situation. If we know that there is a solution to the above equation, then we have

$$
2 A(x) t=-B(x) \pm \sqrt{B(x)^{2}-4 A(x) C(x)} .
$$

(We normally would divide by $2 A(x)$, but it might be that $A(x)$ is not invertible, so we're writing it this way.) If $B(x)^{2}=4 A(x) C(x)$, then there's only one solution. In general, $B(x)^{2}-4 A(x) C(x)$ might not have a square root (for example $x$ has no square root). In the
cases that we'll use it, we will have $B(0)^{2}-4 A(0) C(0)=1$, so that we can use Lemma 2.10. We won't worry too much about the general case. If $A(x)$ is not invertible, then $A(x)$ divides either $-B(x)+\sqrt{B(x)^{2}-4 A(x) C(x)}$ or $-B(x)-\sqrt{B(x)^{2}-4 A(x) C(x)}$ (possibly both). The one case we'll see later is when $A(x)=x$. In that case, divisibility just means that one of the possibilities $-B(x) \pm \sqrt{B(x)^{2}-4 A(x) C(x)}$ does not have a constant term.
2.3. Choice problems. We will now present some applications of the binomial theorem (these have more direct combinatorial proofs, but the point is to illustrate algebraic methods to get you used to this perspective!).

Consider counting the number of subsets of $[n]=\{1, \ldots, n\}$ of size $k$. We will encode this problem algebraically as follows. First, consider the expansion of

$$
(1+x)^{n}=(1+x)(1+x) \cdots(1+x) .
$$

To multiply this out, we have to choose either 1 or $x$ at each step. This choice is the same as a subset $S \subseteq[n]$ : given such a set of choices, we put $i \in S$ if and only if $x$ was chosen in the $i$ th factor. Note then that the result is $x^{|S|}$. For example, the subset $\{1,4,5\} \subset[5]$ corresponds to the underlined terms $(1+\underline{x})(\underline{1}+x)(\underline{1}+x)(1+\underline{x})(1+\underline{x})$. The conclusion:

Proposition 2.16. The number of subsets of $[n]$ of size $k$ is $\left[x^{k}\right](1+x)^{n}=\binom{n}{k}$.
Corollary 2.17. The total number of subsets of $[n]$ is $2^{n}$.
Proof. The number we want is $\sum_{k=0}^{n}\left[x^{k}\right](1+x)^{n}$, i.e., the sum of all of the coefficients of $(1+x)^{n}$. We can get that by plugging in $x=1$ into $(1+x)^{n}$, which gives $2^{n}$.

Remark 2.18. We can also derive Pascal's identity using this idea. Note that $(1+x)^{n}=$ $(1+x)^{n-1}(1+x)$, so in particular, the coefficient of $x^{k}$ is the same on both sides. On the left it is $\binom{n}{k}$ and on the right it is $\binom{n-1}{k}+\binom{n-1}{k-1}$ (you either take the coefficient of $x^{k}$ in $(1+x)^{n-1}$ and multiply by 1 or take the coefficient of $x^{k-1}$ and multiply by $x$ ). This says that the number of subsets of size $k$ in [ $n$ ] is the number of subsets of size $k$ or $k-1$ in $[n-1]$. It's worth thinking about how to prove that last statement directly without formulas.

You have likely already seen this in a different form. The binomial theorem tells us that $(1+x)^{n}=\sum_{k=0}^{n}\binom{n}{k} x^{k}$, which is an equality of single-variable polynomials. We can make this an identity of two-variable polynomials by homogenization. Namely, introduce a new variable $y$, do the substitution $x \mapsto x / y$, and since the polynomials are degree $n$, multiply the result by $y^{n}$ to get rid of denominators:

$$
(x+y)^{n}=\sum_{k=0}^{n}\binom{n}{k} x^{k} y^{n-k}
$$

Remark 2.19. With the 2 -variable version we can do other substitutions such as $x=2$ and $y=3$ to get identities like $5^{n}=\sum_{k=0}^{n}\binom{n}{k} 2^{k} 3^{n-k}$. We can also take derivatives of the original version to get identities like

$$
n(1+x)^{n-1}=\sum_{k=0}^{n} k\binom{n}{k} x^{k-1} .
$$

Substituting values gives plenty of other identities.

Finally, it makes sense to discuss the multinomial theorem at this point. In this case we have $d$ variables $x_{1}, \ldots, x_{d}$. Given integers $k_{i} \geq 0$ such that $k_{1}+\cdots+k_{d}=n$, we define the multinomial coefficient by

$$
\binom{n}{k_{1}, \ldots, k_{d}}=\frac{n!}{k_{1}!\cdots k_{d}!} .
$$

Theorem 2.20 (Multinomial theorem). We have

$$
\left(x_{1}+\cdots+x_{d}\right)^{n}=\sum_{k_{1}+\cdots+k_{d}=n}\binom{n}{k_{1}, \ldots, k_{d}} x_{1}^{k_{1}} \cdots x_{d}^{k_{d}}
$$

where the sum is over all choices of non-negative $k_{i}$ such that $k_{1}+\cdots+k_{d}=n$.
Proof. We prove this by induction on $d$. If $d=1$, both sides are just $x_{1}^{n}$, so the identity holds. Otherwise, consider the 2-variable binomial theorem and substitute $x \mapsto x_{1}+\cdots+x_{d-1}$ and $y \mapsto x_{d}$ and use induction:

$$
\begin{aligned}
\left(x_{1}+\cdots+x_{d}\right)^{n} & =\sum_{m=0}^{n}\binom{n}{m}\left(x_{1}+\cdots+x_{d-1}\right)^{m} x_{d}^{n-m} \\
& =\sum_{m=0}^{n}\binom{n}{m} \sum_{k_{1}+\cdots+k_{d-1}=m}\binom{m}{k_{1}, \ldots, k_{d-1}} x_{1}^{k_{1}} \cdots x_{d-1}^{k_{d-1}} x_{d}^{n-m} .
\end{aligned}
$$

Now set $k_{d}=n-m$ and use that $\binom{n}{n-k_{d}}\binom{n-k_{d}}{k_{1}, \ldots, k_{d-1}}=\binom{n}{k_{1}, \ldots, k_{d}}$.
Let's see how to extract a choice problem out of this identity. Suppose that $1, \ldots, d$ represent colors and we have to assign these colors to $n$ objects that are lined up in a row. If $n=d=3$ and we colored the first and last objects with the first color and the second object with the third color, we might think of this choice as the following underlining:

$$
\left(\underline{x_{1}}+x_{2}+x_{3}\right)\left(x_{1}+x_{2}+\underline{x_{3}}\right)\left(\underline{x_{1}}+x_{2}+x_{3}\right)
$$

Multiplying those terms gives $x_{1}^{2} x_{3}$, and from the exponents we read off that two objects are color 1 and one object is color 3 . This gives the following interpretation of multinomial coefficients.

Proposition 2.21. Assume we have d types of objects (type could mean color). Then $\binom{n}{k_{1}, \ldots, k_{d}}$ is the number of ways to arrange $n$ objects in a row such that exactly $k_{i}$ of them have the ith type if we interpret objects of the same type as being identical.

Example 2.22. We have 10 houses in a row and we need to paint 4 of them blue, 2 of them red, 3 of them green, and 1 orange. Then the number of different ways to choose colors is $\binom{10}{4,2,3,1}=\frac{10!}{4!2!3!1!}=12600$.

Let's see a variation. Rather than pick subsets, consider the problem of picking multisets of $[n]$. This means we can choose elements with repetition. For instance, $\{1,1,1,2,2,3,5\}$ is a multiset of size 7. It is easy to adapt the above argument. We will encode a multiset of $[n]$ as a term in the expansion of $\left(\sum_{d \geq 0} x^{d}\right)^{n}$ : given such a choice of term, we pick $i$ from [ $n$ ] exactly $d$ times if our choice of term is $x^{d}$ in the $i$ th factor. For instance, the multiset $\{1,1,1,2,2,3,5\}$ corresponds to
$\left(1+x+x^{2}+\underline{x^{3}}+\cdots\right)\left(1+x+\underline{x^{2}}+x^{3}+\cdots\right)(1+\underline{x}+\cdots)\left(\underline{1}+x+x^{2}+\cdots\right)\left(1+\underline{x}+x^{2}+\cdots\right)$
since we've chosen 1 thrice, 2 twice, 3 once, 4 not at all, and 5 once. Again, the corresponding term will be $x^{|S|}$. Finally, $\sum_{d \geq 0} x^{d}=(1-x)^{-1}$, so we conclude (using the examples from last time):

Proposition 2.23. The number of multisets of $[n]$ of size $k$ is $\left[x^{k}\right](1-x)^{-n}=(-1)^{k}\binom{-n}{k}=$ $\binom{n+k-1}{k}$.

If you have not seen this before, I encourage you to think about how to get this answer more directly. In other words, why is the number of multisets of [ $n$ ] of size $k$ the same as the number of subsets of size $k$ of $[n+k-1]$ ?
2.4. Rational generating functions. Using generating functions, we can now give a characterization of sequences (with values in a field, as usual) which satisfy a linear recurrence relation. A formal power series $F(x)$ is a rational function if there exist polynomials $P(x)$ and $Q(x)$ with $Q(x) \neq 0$ such that $F(x) Q(x)=P(x)$, or in more suggestive notation, $F(x)=P(x) / Q(x)$. The choice of $P, Q$ is not unique, since we can multiply or divide them both by a common polynomial, but in fact that is all that is possible (we won't prove it, but this follows by unique factorization for polynomials). The degree of a rational function $F(x)$ is defined to be $\operatorname{deg} F(x)=\operatorname{deg} P(x)-\operatorname{deg} Q(x)$. By our comment, this is well-defined.

Theorem 2.24. Let $F(x)=\sum_{n \geq 0} a_{n} x^{n}$ and pick an integer $N \geq 0$. Let $Q(x)=1+c_{1} x+$ $\cdots+c_{r} x^{r}$ with factorization

$$
Q(x)=\left(1-\gamma_{1} x\right)^{m_{1}} \cdots\left(1-\gamma_{s} x\right)^{m_{s}}
$$

where the $\gamma_{i}$ are distinct and nonzero. Then the following are equivalent:
(a) For all $n \geq N$, we have

$$
a_{n+r}+\cdots+c_{r-1} a_{n+1}+c_{r} a_{n}=0
$$

(b) $Q(x) F(x)$ is a polynomial of degree $<N+r$,
(c) There exist polynomials $f_{1}, \ldots, f_{s}$ with $\operatorname{deg} f_{i}<m_{i}$ such that $a_{n}=\sum_{i=1}^{s} f_{i}(n) \gamma_{i}^{n}$ for all $n \geq N$.
In other words, $\left(a_{n}\right)_{n \geq N}$ satisfies a linear recurrence relation if and only if $F(x)$ is a rational function of degree $<N$ if and only if $a_{n}$ is a linear combination of powers with polynomial coefficients for $n \geq N$.

Proof. Pick $n \geq N$; the coefficient of $x^{n+r}$ of $Q(x) F(x)$ is $a_{n+r}+\cdots+c_{r-1} a_{n+1}+c_{r} a_{n}$. Hence, all of these quantities are 0 if and only if all terms $x^{n+r}$ with $n \geq N$ of $Q(x) F(x)$ have 0 coefficient, which is the same as being a polynomial of degree $<N+r$. This shows the equivalence of (a) and (b).

Assume that (b) holds, and write $F(x)=\frac{P(x)}{Q(x)}$ where $P(x)$ is a polynomial of degree $<N+r$. Using polynomial long division, we can find polynomials $g(x)$ and $P_{0}(x)$ with $\operatorname{deg} g(x)<N$ and $\operatorname{deg} P_{0}(x)<r$ such that $F(x)=g(x)+\frac{P_{0}(x)}{Q(x)}$. Now partial fraction decomposition tells us there are polynomials $p_{1}(x), \ldots, p_{s}(x)$ with $\operatorname{deg} p_{i}(x)<m_{i}$ such that

$$
\frac{P_{0}(x)}{Q(x)}=\sum_{i=1}^{s} \frac{p_{i}(x)}{\left(1-\gamma_{i} x\right)^{m_{i}}}
$$

If $d<m$, then by the binomial theorem,

$$
\frac{x^{d}}{(1-\gamma x)^{m}}=\sum_{n \geq 0}\binom{m+n-1}{n} \gamma^{n} x^{n+d}=\sum_{n \geq d}\binom{m+n-d-1}{n-d} \gamma^{n-d} x^{n}
$$

Note that, as a polynomial in $n$, we have

$$
\binom{m+n-d-1}{n-d}=\frac{1}{(m-1)!}(m+n-d-1) \cdots(n-d+2)(n-d+1)
$$

and this has roots at $n=0, \ldots, d-1$, so we can extend the sum to $n \geq 0$ without any harm. In particular, the coefficients are given by $f(n) \gamma^{n}$ where $f(n)=\gamma^{-d}\binom{m+n-d-1}{n-d}$ is a polynomial of degree $<m-1$. So $P_{0}(x) / Q(x)$ is a linear combination of the form we claimed, which proves (c) (since we only make a claim about $a_{n}$ for $n \geq N$, we can ignore $g(x)$ ).

To prove that (c) implies (b), we can reverse these steps.
The case when $s=1$ and $\gamma_{1}=1$ is especially important: in (c) it means that $a_{n}$ is a polynomial in $n$ for $n \geq N$. We record this separately.

Corollary 2.25. Let $F(x)=\sum_{n \geq 0} a_{n} x^{n}$ and pick an integer $N \geq 0$. The following are equivalent:
(a) For all $n \geq N$, we have

$$
\sum_{i=0}^{r}(-1)^{r-i}\binom{r}{i} a_{n+i}=0
$$

(b) $(1-x)^{r} F(x)$ is a polynomial of degree $<N+r$,
(c) There exists a polynomial $f$ with $\operatorname{deg} f<r$ such that $a_{n}=f(n)$ for all $n \geq N$.

Proof. We take $Q(x)=(1-x)^{r}=\sum_{i=0}^{r}\binom{r}{i}(-1)^{i} x^{i}$ in the previous result, so $c_{i}=\binom{r}{i}(-1)^{i}$.

Example 2.26. Given triples of non-negative integers $a=\left(a_{1}, a_{2}, a_{3}\right)$ and $b=\left(b_{1}, b_{2}, b_{3}\right)$, we'll write $a \leq b$ if $a_{i} \leq b_{i}$ for $i=1,2,3$. Fix $a$ and consider all $b$ such that $b \geq a$. Letting $|b|=b_{1}+b_{2}+b_{3}$, we define

$$
F_{a}(x)=\sum_{b \geq a} x^{|b|}=\frac{x^{|a|}}{(1-x)^{3}}
$$

where the second equality follows since $\left(b_{1}-a_{1}, b_{2}-a_{2}, b_{3}-a_{3}\right)$ is equivalent to a multiset of [3] where we pick $i b_{i}-a_{i}$ many times. So we see that $|\{b|b \geq a,|b|=n\} \mid$ is a polynomial in $n$ for $n \geq|a|-2$. Given another $a^{\prime}=\left(a_{1}^{\prime}, a_{2}^{\prime}, a_{3}^{\prime}\right)$, we consider $b$ such that $b \geq a$ or $b \geq a^{\prime}$. Then setting $a^{\prime \prime}=\left(\max \left(a_{1}, a_{1}^{\prime}\right), \max \left(a_{2}, a_{2}^{\prime}\right), \max \left(a_{3}, a_{3}^{\prime}\right)\right)$, we have

$$
\sum_{\substack{b \\ b \geq a \\ \text { or } b \geq a^{\prime}}} x^{|b|}=F_{a}(x)+F_{a^{\prime}}(x)-F_{a^{\prime \prime}}(x)=\frac{x^{|a|}+x^{\left|a^{\prime}\right|}-x^{\left|a^{\prime \prime}\right|}}{(1-x)^{3}},
$$

so again $\left|\left\{b\left|b \geq a, b \geq a^{\prime},|b|=n\right\} \mid\right.\right.$ is a polynomial in $n$ for $n \geq\left|a^{\prime \prime}\right|-2$.
This generalizes to $k$-tuples rather than just triples (we just wanted to keep notation simple) and we can also allow more than just two $a$ and $a^{\prime}$ (but it'll be easier once we discuss inclusion-exclusion).
2.5. Catalan numbers. The Catalan numbers are denoted $C_{n}$ and have a lot of different interpretations. One of them is the number of ways to arrange $n$ pairs of left and right parentheses so that they are balanced: meaning that every ) pairs off with some ( that comes before it. More formally, a word consisting of parentheses is balanced, if for every initial segment, the number of (is always greater than or equal to the number of ). Our convention is that $C_{0}=1$.

Example 2.27. For $n=3$, there are 5 ways to balance 3 pairs of parentheses:

$$
()()(), \quad(())(), \quad((())), \quad(()()), \quad()(()) .
$$

Some other interpretations will be given on homework. For now, we'll see how we can use generating functions to obtain a formula for $C_{n}$. Define

$$
C(x)=\sum_{n \geq 0} C_{n} x^{n}
$$

Lemma 2.28. If $n>0$, we have

$$
C_{n}=\sum_{i=0}^{n-1} C_{i} C_{n-i-1} .
$$

Proof. Every set of balanced parentheses must begin with (. Consider the ) which pairs with it. In between the two of them is another set of balanced parentheses (possibly empty) and to the right of them is another set of balanced parentheses (again, possibly empty). So the set on the inside consists of $i$ pairs, where $0 \leq i \leq n-1$, while the set on the right consists of $n-1-i$ pairs. These sets can be chosen independently, so there are $C_{i} C_{n-i-1}$ ways for this to happen. Since the cases with different $i$ don't overlap, we sum over all possibilities to get the identity above.

Note that the right side of the equation above is the coefficient of $x^{n-1}$ in $C(x)^{2}$. So we have

$$
\begin{aligned}
C(x) & =1+\sum_{n \geq 1} C_{n} x^{n}=1+\sum_{n \geq 1}\left(\sum_{i=0}^{n-1} C_{i} C_{n-i-1}\right) x^{n} \\
& =1+x \sum_{n \geq 1}\left(\sum_{i=0}^{n-1} C_{i} C_{n-i-1}\right) x^{n-1}=1+x C(x)^{2} .
\end{aligned}
$$

This means that $C(x)$ is a solution of the quadratic polynomial $x t^{2}-t+1=0$. Using the quadratic formula, we deduce that $C(x)$ is one of the solutions

$$
\frac{1 \pm \sqrt{1-4 x}}{2 x}
$$

Note that $x$ isn't invertible as a power series, so we have to be careful here. Since $C(x)$ is a power series, it must be that $x$ divides the numerator, i.e., the numerator cannot have a constant term. Which choice of sign is correct? The constant term of $\sqrt{1-4 x}$ is, by definition, 1 , so the correct choice is a negative sign, and so

$$
C(x)=\frac{1-\sqrt{1-4 x}}{2 x} .
$$

Let's use the binomial theorem now. First, we have

$$
1-(1-4 x)^{1 / 2}=-\sum_{n \geq 1}\binom{1 / 2}{n}(-4 x)^{n}
$$

Let's simplify the coefficients (assuming $n \geq 1$ ):

$$
-(-1)^{n} 4^{n}\binom{1 / 2}{n}=-(-1)^{n} 4^{n} \frac{\frac{1}{2} \frac{-1}{2} \frac{-3}{2} \cdots \frac{-(2 n-3)}{2}}{n!}=2^{n} \frac{(2 n-3)!!}{n!} .
$$

Note that $(2 n-3)!!(2 n-2)!!=(2 n-2)!$, so we can multiply top and bottom by $(2 n-2)!!$ to get

$$
2^{n} \frac{(2 n-2)!}{n!(2 n-2)!!}=2 \frac{(2 n-2)!}{n!(n-1)!}=\frac{2}{n}\binom{2 n-2}{n-1}
$$

So:

$$
\frac{1-\sqrt{1-4 x}}{2 x}=\frac{\sum_{n \geq 1} \frac{2}{n}\binom{2 n-2}{n-1} x^{n}}{2 x}=\sum_{n \geq 1} \frac{1}{n}\binom{2 n-2}{n-1} x^{n-1}=\sum_{n \geq 0} \frac{1}{n+1}\binom{2 n}{n} x^{n} .
$$

This gives us the following formula:
Theorem 2.29. $C_{n}=\frac{1}{n+1}\binom{2 n}{n}$.
The values for $n=0, \ldots, 9$ are $1,1,2,5,14,42,132,429,1430,4862$.
Remark 2.30. $C(x)$ is not a rational function, so the $C_{n}$ do not satisfy a linear recurrence relation of any order. Since $C(x)$ is the root of a polynomial whose coefficients are polynomials in $x$, it is an example of an algebraic function.

Here are a few other things that are counted by the Catalan numbers together with the 5 instances for $n=3$ :

- The number of ways to apply a binary operation $*$ to $n+1$ elements:
$a *(b *(c * d)), \quad a *((b * c) * d), \quad(a * b) *(c * d), \quad((a * b) * c) * d, \quad(a *(b * c)) * d$.
- The number of rooted binary trees with $n+1$ leaves:

- The number of paths from $(0,0)$ to $(n, n)$ which never go above the diagonal $x=y$ and are made up of steps either moving in the direction $(0,1)$ or $(1,0)$.


It turns out that the Catalan recursion shows up a lot. There are more than 200 other known interpretations for the Catalan numbers.

## 3. Fundamental counting problems

3.1. 12-fold way, introduction. We have $k$ balls and $n$ boxes. We want to count ways to put the balls into the boxes. We can think of each assignment as a function from the set of balls to the set of boxes. Phrased this way, we will be examining how many ways to do this if we require $f$ to be injective (one-to-one), or surjective (onto), or completely arbitrary. Are the boxes supposed to be considered different or interchangeable (we also use the terminology distinguishable and indistinguishable)? And same with the balls, are they considered different or interchangeable?

Formally, we have two symmetric groups $\mathfrak{S}_{k}$ and $\mathfrak{S}_{n}$ which act on the set of functions, and "indistinguishable" corresponds to counting orbits with respect to one or both of these groups rather than counting the functions themselves. For example, if the balls are considered indistinguishable but the boxes are not, then we are considering two functions to be the same if there is a way to relabel the balls (i.e., apply a permutation in $\mathfrak{S}_{k}$ ) that turns one into the other.

All in all, this will give us 12 different problems to consider, which means we want to understand the following table:

| balls/boxes | $f$ arbitrary | $f$ injective | $f$ surjective |
| :--- | :--- | :--- | :--- |
| dist/dist |  |  |  |
| indist/dist |  |  |  |
| dist/indist |  | $\begin{cases}1 & \text { if } n \geq k \\ 0 & \text { if } n<k\end{cases}$ |  |
| indist/indist |  | $\begin{cases}1 & \text { if } n \geq k \\ 0 & \text { if } n<k\end{cases}$ |  |

Two situations have already been filled in and won't be considered interesting (so technically we only have 10 problems, not 12). Also, we don't have a column for when $f$ is bijective: that's a special case of either of the last two columns when $n=k$.
3.2. Compositions. Below, $n$ and $k$ are positive integers.

Definition 3.1. A sequence of non-negative integers $\left(a_{1}, \ldots, a_{n}\right)$ is a weak composition of $k$ if $a_{1}+\cdots+a_{n}=k$. If all of the $a_{i}$ are positive, then it is a composition. We call $n$ the number of parts of the (weak) composition.

This addresses the following counting problem: we have $k$ objects which are placed into $n$ boxes (which we can think of as a function $f$ from the $k$ objects to the $n$ boxes. The boxes are labeled $1, \ldots, n$ but the objects themselves are all identical. So the only information we can ask is how many objects are placed in each box, i.e., we get a weak composition $\left(a_{1}, \ldots, a_{n}\right)$ by letting $a_{i}$ be the number of objects placed in the $i$ th box. A composition deals with the situation where the assignment $f$ must be surjective.
Theorem 3.2. The number of weak compositions of $k$ with $n$ parts is $\binom{n+k-1}{k}=\binom{n+k-1}{n-1}$.
Proof. We've essentially done this already. Consider the expansion

$$
\frac{1}{(1-x)^{n}}=\left(\sum_{a \geq 0} x^{a}\right)^{n}=\left(\sum_{a_{1} \geq 0} x^{a_{1}}\right) \cdots\left(\sum_{a_{n} \geq 0} x^{a_{n}}\right)=\sum_{\left(a_{1}, \ldots, a_{n}\right) \in \mathbf{Z}_{\geq 0}^{n}} x^{a_{1}+\cdots+a_{n}} .
$$

So the number of weak compositions of $k$ with $n$ parts is $\left[x^{k}\right](1-x)^{-n}=\binom{n+k-1}{k}$.
Recall from Proposition 2.23 that this is the same as the number of multisets of [ $n$ ] of size $k$, so there should be a bijection between weak compositions and multisets. In fact, there is: given a weak composition $\left(a_{1}, \ldots, a_{n}\right)$ of $k$, we can associate a multiset of $[n]$ of size $k$ by picking the number $i$ exactly $a_{i}$ times.

Example 3.3. We want to distribute 20 pieces of candy (all identical) to 4 children. How many ways can we do this? If we order the children and let $a_{i}$ be the number of pieces of candy that the $i$ th child receives, then $\left(a_{1}, a_{2}, a_{3}, a_{4}\right)$ is just a weak composition of 20 into 4 parts, so we can identify all ways with the set of all weak compositions. So we know that the number of ways is $\binom{20+4-1}{20}=\binom{23}{20}$.

What if we want to ensure that each child receives at least one piece of candy? First, hand each child 1 piece of candy. We have 16 pieces left, and we can distribute them as we like, so we're counting weak compositions of 16 into 4 parts, or $\binom{19}{16}$.
Corollary 3.4. The number of compositions of $k$ into $n$ parts is $\binom{k-1}{n-1}$.
Proof. If we generalize the argument in the last example, we see that compositions of $k$ into $n$ parts are in bijection with weak compositions of $k-n$ into $n$ parts.

We can also adapt the generating function argument in the last proof (it's not much different): we have

$$
\frac{x^{n}}{(1-x)^{n}}=\left(\sum_{a_{1} \geq 1} x^{a_{1}}\right) \cdots\left(\sum_{a_{n} \geq 1} x^{a_{n}}\right)=\sum_{\left(a_{1}, \ldots, a_{n}\right) \in \mathbf{Z}_{\geq 1}^{n}} x^{a_{1}+\cdots+a_{n}}
$$

so the number of compositions of $k$ into $n$ parts is

$$
\left[x^{k}\right] \frac{x^{n}}{(1-x)^{n}}=\left[x^{k-n}\right] \frac{1}{(1-x)^{n}}=\binom{k-1}{n-1} .
$$

Corollary 3.5. The total number of compositions of $k$ (into any number of parts) is $2^{k-1}$.
Proof. The possible number of parts of a composition of $k$ is anywhere between $n=1$ to $n=k$. So the total number of compositions possible is

$$
\sum_{n=1}^{k}\binom{k-1}{n-1}=\sum_{n=0}^{k-1}\binom{k-1}{n}=2^{k-1}
$$

The answer suggests that we should be able to find a bijection between compositions of $k$ and subsets of $[k-1]$. Here's one: given a composition $\left(a_{1}, \ldots, a_{r}\right)$ of $k$, consider the subset of partial sums $\left\{a_{1}, a_{1}+a_{2}, \ldots, a_{1}+\cdots+a_{r-1}\right\}$. Since $a_{r}>0$, this is a subset of [ $k-1$ ] (we don't include the whole sum since it's always $k$ and hence redundant information). Conversely, given a subset $\left\{s_{1}, \ldots, s_{m}\right\}$ written in increasing order $s_{1}<\cdots<s_{m}$, we get a composition of $k$ by taking successive differences: $\left(s_{1}, s_{2}-s_{1}, \ldots, s_{m}-s_{m-1}, k-s_{m}\right)$.
3.3. Words. A word is a finite ordered sequence whose entries are drawn from some set $A$ (which we call the alphabet). The length of the word is the number of entries it has. Entries may repeat, there is no restriction on that. Also, the empty sequence $\varnothing$ is considered a word of length 0 .

If $A$ represents a set of distinguishable boxes and there are $k$ distinguishable balls, then a word of length $k$ is the same thing as an assignment of balls to boxes: the $i$ th entry records which box the $i$ th ball gets sent to.

Example 3.6. Say our alphabet is $A=\{a, b\}$. The words of length $\leq 2$ are:

$$
\varnothing, \quad a, \quad b, \quad a a, \quad a b, \quad b a, \quad b b .
$$

Theorem 3.7. If $|A|=n$, then the number of words in $A$ of length $k$ is $n^{k}$.
Proof. A sequence of length $k$ with entries in $A$ is an element in the product set $A^{k}=$ $A \times A \times \cdots \times A$ and $\left|A^{k}\right|=|A|^{k}$.

Alternatively, we can think of this as follows. To specify a word, we pick each of its entries, but these can be done independently of the other choices. So for each of the $k$ positions, we are choosing one of $n$ different possibilities, which leads us to $n \cdot n \cdots n=n^{k}$ different choices for words.

So that we get used to thinking about it, we record the corresponding generating function. Actually, there are two parameters, so we can either fix $|A|$ and sum over $k$ :

$$
\sum_{k \geq 0} \left\lvert\,\left.\{\text { words of length } k \text { in } A\}\left|x^{k}=\sum_{k \geq 0}\right| A\right|^{k} x^{k}=\frac{1}{1-|A| x}\right.,
$$

or we could fix $k$ and sum over $n=|A|$ :

$$
\sum_{n \geq 0} n^{k} x^{n}
$$

We've seen this already when $k=0,1$. We know that this is a rational function of the form $\frac{A_{k}(x)}{(1-x)^{k+1}}$ where $A_{k}(x)$ is a polynomial with $\operatorname{deg} A_{k}(x) \leq k$. We have $A_{0}(x)=1$ and $A_{1}(x)=x$. These are the Eulerian polynomials and have a lot of interesting properties but we won't elaborate any more here.

Example 3.8. A small city has 10 intersections. Each one could have a traffic light or gas station (or both or neither). How many different configurations could this city have?

Let $A$ be the alphabet of size 4 with symbols representing the possibilities of whether a traffic light or gas station is there. Then a configuration is just a word of length 10 , so there are $4^{10}$ possibilities.

Example 3.9. Given a subset $S \subseteq[n]$, we define a word $w_{S}$ of length $n$ in the alphabet $\{0,1\}$ as follows. If $i \in S$, then the $i$ th entry of $w_{S}$ is 1 , and otherwise the entry is 0 . This defines a function

$$
f:\{\text { subsets of }[n]\} \rightarrow\{\text { words of length } n \text { on }\{0,1\}\}
$$

which we've seen with a different description before. We can also define an inverse function: given such a word $w$, we send it to the subset of positions where there is a 1 in $w$. We omit the check that these two functions are inverse to one another. So $f$ is a bijection, and the previous result tells us that there are $2^{n}$ words of length $n$ on $\{0,1\}$.

Example 3.10. How many pairs of subsets $S, T \subseteq[n]$ satisfy $S \subseteq T$ ? We can also encode this problem as a problem about words. Let $A$ be the alphabet of size 3 whose elements are: "in $S$ and $T$ ", "in $T$ but not $S$ " and "not in $T$ or $S$ ". Then each pair $S \subseteq T$ gives a word
of length $n$ in $A$ : the $i$ th entry of the word is the element which describes the position of $i$. So there are $3^{n}$ such pairs.

A less elegant way to compute this is to use the binomial theorem:

$$
\sum_{T \subseteq[n]} \sum_{S \subseteq T} 1=\sum_{T \subseteq[n]} \mid\{\text { subsets of } T\} \left\lvert\,=\sum_{T \subseteq[n]} 2^{|T|}=\sum_{k=0}^{n}\binom{n}{k} 2^{k}=(2+1)^{n}=3^{n}\right.
$$

How about words without repeating entries? We will call these injective words. In our balls/boxes problem, this corresponds to cases where boxes don't get used more than once, i.e., injective assignments. Given $n \geq k$, define the falling factorial by

$$
(n)_{k}:=n(n-1)(n-2) \cdots(n-k+1) .
$$

There are $k$ numbers being multiplied in the above definition. When $n=k$, we have $(n)_{n}=n$ !, so this generalizes the factorial function.
Theorem 3.11. If $|A|=n$ and $n \geq k$, then there are $(n)_{k}$ different words of length $k$ in $A$ which do not have any repeating entries.

Proof. Start with a permutation of $A$. The first $k$ elements in that permutation give us a word of length $k$ with no repeating entries. But we've overcounted because we don't care how the remaining $n-k$ things we threw away are ordered. In particular, this process returns each word exactly $(n-k)$ ! many times, so our desired quantity is

$$
\frac{n!}{(n-k)!}=(n)_{k}
$$

Let's record the generating function. Should we fix $k$ or $n$ ? Since $(n)_{k}$ is only nonzero if $n \geq k$, let's fix $k$ and sum over $n$ (otherwise it's a finite sum):

$$
\sum_{n \geq k}(n)_{k} x^{n}=k!\sum_{n \geq k}\binom{n}{k} x^{n} .
$$

This is close to a binomial expansion. More precisely, we have

$$
\frac{k!x^{k}}{(1-x)^{k+1}}=k!\sum_{m \geq 0}\binom{m+k}{k} x^{m+k}=k!\sum_{n \geq k}\binom{n}{k} x^{n} .
$$

Even though $\left\{n^{k} \mid k \geq 0\right\}$ is a natural basis for the space of polynomials in $n$, their generating functions are complicated when compared to $\left\{(n)_{k} \mid k \geq 0\right\}$, so for some purposes it makes sense to write polynomials as linear combinations of $(n)_{k}$ rather than $n^{k}$.
3.4. Set partitions. (Weak) compositions were about indistinguishable objects into distinguishable boxes. Now we reverse the roles and consider distinguishable objects into indistinguishable boxes.

Definition 3.12. Let $X$ be a set. A partition of $X$ is an unordered collection of nonempty subsets $S_{1}, \ldots, S_{k}$ of $X$ such that every element of $X$ belongs to exactly one of the $S_{i}$. An ordered partition of $X$ is the same, except the subsets are ordered. The $S_{i}$ are the blocks of the partition. Partitions of sets are also called set partitions to distinguish from integer partitions, which will be discussed next.

These are another way of thinking about surjective assignments of $n$ distinguishable balls to $k$ boxes (distinguishable in the ordered case and indistinguishable otherwise). We have swapped $n$ and $k$ here, to make the notation for set partitions a little more standard. We think of each block $S_{i}$ as describing the contents of the $i$ th box.

Example 3.13. Let $X=\{1,2,3\}$. There are 5 partitions of $X$ :

$$
\{\{1,2,3\}\}, \quad\{\{1,2\},\{3\}\}, \quad\{\{1,3\},\{2\}\}, \quad\{\{2,3\},\{1\}\}, \quad\{\{1\},\{2\},\{3\}\} .
$$

When we say unordered collection of subsets, we mean that $\{\{1,2\},\{3\}\}$ and $\{\{3\},\{1,2\}\}$ are to be considered the same partition.

The notation above is a little cumbersome, so we can also write the above partitions as follows:

$$
123, \quad 12|3, \quad 13| 2, \quad 23|1, \quad 1| 2 \mid 3 .
$$

The number of partitions of $X$ with $k$ blocks only depends on the number of elements of $X$. So for concreteness, we will usually assume that $X=[n]$.

Example 3.14. If we continue with our previous example of candy and children: imagine the 20 pieces of candy are now labeled 1 through 20 and that the 4 children are all identical clones. The number of ways to distribute candy to them so that each gets at least 1 piece of candy is then the number of partitions of [20] into 4 blocks.

Definition 3.15. We let $S(n, k)$ be the number of partitions of a set of size $n$ into $k$ blocks. These are called the Stirling numbers of the second kind. By convention, we define $S(0,0)=1$. Note that $S(n, k)=0$ if $k>n$ or if $k=0$ and $n>0$.

The number of ordered partitions of a set of size $n$ into $k$ blocks is $k!S(n, k)$ : the extra data we need is a way to order the blocks.

At the moment, it will be hard to get nice, exact formulas for $S(n, k)$, but we can do some special cases:

Example 3.16. For $n \geq 1, S(n, 1)=S(n, n)=1$. For $n \geq 2, S(n, 2)=2^{n-1}-1$ and $S(n, n-1)=\binom{n}{2}$. Can you see why?

We also have the following recursive formula:
Theorem 3.17. If $n \geq k \geq 1$, then

$$
S(n, k)=S(n-1, k-1)+k S(n-1, k) .
$$

Proof. Consider two kinds of partitions of $[n]$. The first kind is when $n$ is in its own block. In that case, if we remove this block, then we obtain a partition of $[n-1]$ into $k-1$ blocks. To reconstruct the original partition, we just add a block containing $n$ by itself. So the number of such partitions is $S(n-1, k-1)$.

The second kind is when $n$ is not in its own block. This time, if we remove $n$, we get a partition of $n-1$ into $k$ blocks. However, it's not possible to reconstruct the original block because we can't remember which block it belonged to. So in fact, there are $k$ different ways to try to reconstruct the original partition. This means that the number of such partitions is $k S(n-1, k)$.

If we add both answers, we account for all possible partitions of $[n]$, so we get the identity we want.

Here's a table of small values of $S(n, k)$ :

| $n \backslash k$ | 1 | 2 | 3 | 4 | 5 |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | 1 | 0 | 0 | 0 | 0 |
| 2 | 1 | 1 | 0 | 0 | 0 |
| 3 | 1 | 3 | 1 | 0 | 0 |
| 4 | 1 | 7 | 6 | 1 | 0 |
| 5 | 1 | 15 | 25 | 10 | 1 |

Let's consider the generating function of $S(n, k)$. Since this is nonzero only for $n \geq k$, let's fix $k$ and let $n$ vary. For $k \geq 0$, define

$$
F_{k}(x)=\sum_{n \geq k} S(n, k) x^{n}
$$

For $k \geq 1$, the recursion above becomes

$$
\sum_{n \geq k} S(n, k) x^{n}=x \sum_{n \geq k} S(n-1, k-1) x^{n-1}+k x \sum_{n \geq k} S(n-1, k) x^{n-1}
$$

which translates to

$$
F_{k}(x)=x F_{k-1}(x)+k x F_{k}(x),
$$

or more simply $F_{k}(x)=F_{k-1}(x) \frac{x}{1-k x}$. Since $F_{0}(x)=1$, we conclude that

$$
F_{k}(x)=\frac{x^{k}}{(1-x)(1-2 x) \cdots(1-k x)} .
$$

So $F_{k}(x)$ is rational and Theorem 2.24 tells us that there are constants $\alpha_{i, k}$ for $i=1, \ldots, k$ so that $S(n, k)=\sum_{i=1}^{k} \alpha_{i, k} i^{n}$ for $n \geq 1$. We'll see how to figure out these constants when we talk about inclusion-exclusion.

We define $B(n)$ to be the number of partitions of $[n]$ into any number of blocks. This is the $n$th Bell number. By definition,

$$
B(n)=\sum_{k=0}^{n} S(n, k) .
$$

We can also consider its generating function

$$
\sum_{n \geq 0} B(n)=\sum_{k \geq 0} F_{k}(x)=\sum_{k \geq 0} \frac{x^{k}}{(1-x) \cdots(1-k x)}
$$

Here we see an infinite sum of formal power series. It converges: for each $n$, we have $\left[x^{n}\right] F_{k}(x)=0$ if $k>n$, so the partial sums $\left[x^{n}\right] \sum_{k=0}^{m} F_{k}(x)$ are constant for $m>n$. However, there doesn't seem to be any way to "simplify" this expression.

We have the following recursion:
Theorem 3.18. $B(n+1)=\sum_{i=0}^{n}\binom{n}{i} B(i)$.
Proof. We separate all of the set partitions of $[n+1]$ based on the number of elements in the block that contains $n+1$. Consider those where the size is $j$. To count the number of these, we need to first choose the other elements to occupy the same block as $n+1$. These numbers come from $[n]$ and there are $j-1$ to be chosen, so there are $\binom{n}{j-1}$ ways to do this. We have to then choose a set partition of the remaining $n+1-j$ elements, and there are
$B(n+1-j)$ many of these. So the number of such partitions is $\binom{n}{j-1} B(n+1-j)$. The possible values for $j$ are between 1 and $n+1$, so we get the identity

$$
B(n+1)=\sum_{j=1}^{n+1}\binom{n}{j-1} B(n+1-j)
$$

Re-index the sum by setting $i=n+1-j$ and use the identity $\binom{n}{n-i}=\binom{n}{i}$ to get the desired identity.
3.5. Integer partitions. Now we come to the situation where both balls and boxes are indistinguishable. In this case, the only relevant information is how many boxes are empty, how many contain exactly 1 ball, how many contain exactly 2 balls, etc. We use the following structure:

Definition 3.19. An partition of an integer $n$ is a sequence of non-negative integers $\lambda=$ $\left(\lambda_{1}, \lambda_{2}, \ldots, \lambda_{k}\right)$ so that $\lambda_{1}+\lambda_{2}+\cdots+\lambda_{k}=n$ and so that $\lambda_{1} \geq \lambda_{2} \geq \cdots \geq \lambda_{k}$. The $\lambda_{i}$ are the parts of $\lambda$. We use the notation $|\lambda|=n$ (size of the partition) and $\ell(\lambda)$ (length of the partition) is the number of $\lambda_{i}$ which are positive. These are also called integer partitions to distinguish from set partitions.

We will consider two partitions the same if they are equal after removing all of the parts equal to 0 .

The number of partitions of $n$ is denoted $p(n)$, and the number of partitions of $n$ with $k$ parts is denoted $p_{k}(n)$.

We've reversed the roles of $n$ and $k$, but the partition $\left(\lambda_{1}, \ldots, \lambda_{k}\right)$ encodes an assignment of $n$ balls to $k$ boxes where some box has $\lambda_{1}$ balls, another box has $\lambda_{2}$ balls, etc. Remember we don't distinguish the boxes, so we can list the $\lambda_{i}$ in any order and we'd get an equivalent assignment. But our convention will be that the $\lambda_{i}$ are listed in weakly decreasing order.

Example 3.20. $p(5)=7$ since there are 7 partitions of 5 :

$$
(5), \quad(4,1), \quad(3,2), \quad(3,1,1), \quad(2,2,1), \quad(2,1,1,1), \quad(1,1,1,1,1)
$$

We can visualize partitions using Young diagrams. To illustrate, the Young diagram of $(4,2,1)$ is

$$
Y(\lambda)=\begin{array}{|l|l|l}
\hline & & \\
\hline & \\
\hline &
\end{array}
$$

In general, it is a left-justified collection of boxes with $\lambda_{i}$ boxes in the $i$ th row (counting from top to bottom).

It's not practical to find a closed formula for $p(n)$. Instead, we'll study relationships between different classes of partitions. Also, it turns out the generating functions have relatively simple-looking formulas even though the coefficients do not.

The transpose (or conjugate) of a partition $\lambda$ is the partition whose Young diagram is obtained by flipping $Y(\lambda)$ across the main diagonal. For example, the transpose of $(4,2,1)$ is $(3,2,1,1)$ :


Note that we get the parts of a partition from a Young diagram by reading off the row lengths. The transpose is obtained by instead reading off the column lengths. The notation is $\lambda^{T}$. If we want a formula: $\lambda_{i}^{T}=\left|\left\{j \mid \lambda_{j} \geq i\right\}\right|$.

Note that $\left(\lambda^{T}\right)^{T}=\lambda$. A partition $\lambda$ is self-conjugate if $\lambda=\lambda^{T}$.
Example 3.21. Some self-conjugate partitions: $(4,3,2,1),(5,1,1,1,1),(4,2,1,1)$ :


Theorem 3.22. The number of partitions $\lambda$ of $n$ with $\ell(\lambda) \leq k$ is the same as the number of partitions $\mu$ of $n$ such that all $\mu_{i} \leq k$.

Proof. We get a bijection between the two sets by taking transpose. Details omitted.
Let's write down generating functions for various classes of partitions to get a feel for how it works.

Example 3.23. Let $p_{\leq k}(n)$ be the number of integer partitions of $n$ with at most $k$ parts $\left(p_{\leq k}(0)=1\right)$. Using the transpose of partitions, this is also the number of integer partitions of $n$ using only the numbers $1, \ldots, k$, and we will instead use this interpretation. We want a simple expression for $\sum_{n>0} p_{\leq k}(n) x^{n}$. When $k=1$, we get $p_{\leq 1}(n)=1$ for all $n$, so $\sum_{n \geq 0} p_{\leq 1}(n) x^{n}=\frac{1}{1-x}$.

I claim that

$$
\sum_{n \geq 0} p_{\leq k}(n) x^{n}=\prod_{i=1}^{k} \frac{1}{1-x^{i}}=\frac{1}{(1-x)\left(1-x^{2}\right) \cdots\left(1-x^{k}\right)}
$$

Why? We have $\left(1-x^{i}\right)^{-1}=1+x^{i}+x^{2 i}+\cdots$, and when we multiply out $\prod_{i=1}^{k}\left(1-x^{i}\right)^{-1}$, we make a choice of non-negative integers $m_{1}, \ldots, m_{k}$ and select the term $x^{i m_{i}}$ from the $i$ th factor. This corresponds to the partition where $i$ gets used $m_{i}$ times. For example, for the partition $(4,3,3,1,1,1)$ and $k=4$, we underline the following terms:

$$
\left(1+x+x^{2}+\underline{x^{3}}+\cdots\right)\left(\underline{1}+x^{2}+x^{4}+\cdots\right)\left(1+x^{3}+\underline{x^{6}}+\cdots\right)\left(1+\underline{x^{4}}+x^{8}+\cdots\right) .
$$

Note that $p_{\leq k}(n)=p(n)$ if $k \geq n$. Hence we can take $k \rightarrow \infty$ to get Euler's formula

$$
\sum_{n \geq 0} p(n) x^{n}=\lim _{k \rightarrow \infty} \sum_{n \geq 0} p_{\leq k}(n) x^{n}=\prod_{i \geq 1} \frac{1}{1-x^{i}}
$$

More generally, for any subset $S$ of the positive integers, the generating function for the number of partitions that only use parts from $S$ is $\prod_{i \in S} \frac{1}{1-x^{i}}$.

Example 3.24. We can write down the generating function for $p(n)$ in a different way. Given a partition $\lambda$, its Durfee square is the largest $r \times r$ block sitting in the top-left of the Young diagram of $\lambda$. In formulas, $r=\max \left\{i \mid \lambda_{i} \geq i\right\}$. Every Young diagram can be built as follows: choose the size $r$ of the Durfee square, choose a partition with at most $r$ parts to
put to the right of the Durfee square, and choose a partition whose parts are at most $r$ to but below the Durfee square. These 2 partitions can be chosen independently, so we have

$$
\sum_{n \geq 0} p(n) x^{n}=\sum_{r \geq 0} x^{r^{2}}\left(\sum_{n \geq 0} p_{\leq r}(n) x^{n}\right)^{2}=\sum_{r \geq 0} \frac{x^{r^{2}}}{(1-x)^{2} \cdots\left(1-x^{r}\right)^{2}}
$$

Let $p_{\text {odd }}(n)$ be the number of partitions of $n$ such that all parts are odd. Let $p_{\text {dist }}(n)$ be the number of partitions of $n$ such that all parts are distinct.

Theorem 3.25 (Euler). $p_{\text {odd }}(n)=p_{\text {dist }}(n)$.
For example, when $n=5$, both quantities are 3 since we have $(5),(3,1,1),(1,1,1,1,1)$ for $p_{\text {odd }}(5)$ and $(5),(4,1),(3,2)$ for $p_{\text {dist }}(5)$.

Proof. There are ways to build bijections, but we'll prove this by showing that they have the same generating function.

By the last example, we have

$$
\sum_{n \geq 0} p_{\text {odd }}(n) x^{n}=\prod_{i \geq 0} \frac{1}{1-x^{2 i+1}}=\frac{1}{(1-x)\left(1-x^{3}\right)\left(1-x^{5}\right)\left(1-x^{7}\right) \cdots}
$$

How about for $p_{\text {dist }}(n)$ ? I claim that

$$
\sum_{n \geq 0} p_{\text {dist }}(n) x^{n}=\prod_{i \geq 1}\left(1+x^{i}\right)=(1+x)\left(1+x^{2}\right)\left(1+x^{3}\right)\left(1+x^{4}\right) \cdots
$$

To multiply out the right side, we either choose 1 or $x^{i}$ from the $i$ th term, and we can only avoid choosing 1 finitely many times (due to how infinite products are defined). What we get then is $x^{N}$ where $N$ is the sum of the $i$ where we chose $x^{i}$. But we get $x^{N}$ one time for every partition of $N$ into distinct parts, so the coefficient is $p_{\text {dist }}(N)$.

Now we observe that $1+x^{i}=\frac{1-x^{2 i}}{1-x^{i}}$, so we can rewrite it as

$$
\sum_{n \geq 0} p_{\text {dist }}(n) x^{n}=\frac{1-x^{2}}{1-x} \cdot \frac{1-x^{4}}{1-x^{2}} \cdot \frac{1-x^{6}}{1-x^{3}} \cdot \frac{1-x^{8}}{1-x^{4}} \cdot \frac{1-x^{10}}{1-x^{5}} \cdots
$$

We can start cancelling: each $1-x^{2 i}$ on the top cancels with the corresponding $1-x^{2 i}$ on the bottom. What we're left with is $\prod_{i \geq 0} \frac{1}{1-x^{2 i+1}}=\sum_{n \geq 0} p_{\text {odd }}(n) x^{n}$.

Let's end with an example where we can construct a bijection directly.
Theorem 3.26. The number of self-conjugate partitions of $n$ is equal to the number of partitions of $n$ using only distinct odd parts.

Proof. Given a self-conjugate partition, take all of the boxes in the first row and column of its Young diagram. Since it's self-conjugate, there are an odd number of boxes. Use this as the first part of a new partition. Now remove those boxes and repeat. For example, starting
with
 we get


In formulas, if $\lambda$ is self-conjugate, then $\mu_{i}=\lambda_{i}-(i-1)+\lambda_{i}^{T}-(i-1)-1=2 \lambda_{i}-2 i+1$ and so $\mu_{1}>\mu_{2}>\cdots$.

This process is reversible: let $\mu$ be a partition with distinct odd parts. Each part $\mu_{i}$ can be turned into a shape with a single row and column, both of length $\left(\mu_{i}+1\right) / 2$. Since the $\mu_{i}$ are distinct, these shapes can be nested into one another to form the partition $\lambda$ (this is easiest to understand by studying the two examples above).

Example 3.27. What about their generating functions? Following our previous examples, the generating function for the number of partitions using only distinct odd parts is given by

$$
\prod_{i \geq 0}\left(1+x^{2 i+1}\right) .
$$

How about for self-conjugate partitions? It has the same generating function by the previous result, but let's think differently. Going back to the Durfee square decomposition of a partition $\lambda$, the partition we choose to put to the right of the Durfee square is the transpose of the partition we put below, so they determine each other. That gives the following generating function for self-conjugate partitions:

$$
\sum_{r \geq 0} x^{r^{2}} \sum_{n \geq 0} p_{\leq r}(n) x^{2 n}=\sum_{r \geq 0} \frac{x^{r^{2}}}{\left(1-x^{2}\right)\left(1-x^{4}\right) \cdots\left(1-x^{2 r}\right)}
$$

So the result above says this is the same as $\prod_{i \geq 0}\left(1+x^{2 i+1}\right)$ which doesn't seem so obvious.
3.6. 12-fold way, summary. We have $k$ balls and $n$ boxes. We want to count the number of assignments $f$ of balls to boxes. We considered 3 conditions on $f$ : arbitrary (no conditions at all), injective (no box receives more than one ball), surjective (every box has to receive at least one ball). We also considered conditions on the balls: indistinguishable (we can't tell the balls apart) and distinguishable (we can tell the balls apart) and similarly for the boxes: they can be distinguishable or indistinguishable.

| balls/boxes | $f$ arbitrary | $f$ injective | $f$ surjective |
| :--- | :--- | :--- | :--- |
| dist/dist | $n^{k}$, see (1) | $(n)_{k}$, see (2) | $n!S(k, n)$, see (3) |
| indist/dist | $\binom{n+k-1}{k}$, see (4) | $\binom{n}{k}$, see (5) | $\binom{k-1}{n-1}$, see (6) |
| dist/indist | $\sum_{i=1}^{n} S(k, i)$, see (7) | $\left\{\begin{array}{ll}1 & \text { if } n \geq k \\ 0 & \text { if } n<k\end{array}\right.$, see (8) | $S(k, n)$, see (9) |
| indist/indist | $p_{\leq n}(k)$, see (10) | $\left\{\begin{array}{ll}1 & \text { if } n \geq k \\ 0 & \text { if } n<k\end{array}\right.$, see (11) | $p_{n}(k)$, see (12) |

(1) These are words of length $k$ in an alphabet of size $n$.
(2) These are words of length $k$ without repetitions in an alphabet of size $n$. Recall that

$$
(n)_{k}=n(n-1)(n-2) \cdots(n-k+1)
$$

(3) These are ordered (set) partitions of [k] into $n$ blocks. Recall that $S(k, n)$ is the Stirling number of the second kind, i.e., the number of partitions of $[k]$ into $n$ blocks.
(4) These are multisets of $[n]$ of size $k$; equivalently, weak compositions of $k$ into $n$ parts.
(5) These are subsets of $[n]$ of size $k$.
(6) These are compositions of $k$ into $n$ parts.
(7) These are set partitions of $[k]$ where the number of blocks is $\leq n$.
(8) If $n<k$, then we can't assign $k$ balls to $n$ boxes without some box receiving more than one ball (pigeonhole principle), so the answer is 0 in that case. If $n \geq k$, then there is certainly a way to make an assignment, but they're all the same: we can't tell the boxes apart, so it doesn't matter where the balls go.
(9) These are set partitions of $[k]$ into $n$ blocks.
(10) These are the number of integer partitions of $k$ where the number of parts is $\leq n$.
(11) The reasoning here is the same as (8).
(12) These are the number of integer partitions of $k$ into $n$ parts.
3.7. Cycles in permutations. Recall that $(x)_{k}=x(x-1) \cdots(x-k+1)$.

Proposition 3.28. We have an equality of polynomials in $x$ :

$$
x^{n}=\sum_{k=0}^{n} S(n, k)(x)_{k}
$$

Proof. Pick a positive integer $d \geq n$. Then $d^{n}$ is the number of functions $[n] \rightarrow[d]$. We separate functions by their image, which is some subset $S \subseteq[d]$. These are equivalently surjective functions $[n] \rightarrow S$, and if $|S|=k$, then there are $S(n, k) k$ ! many such functions. This only depends on $|S|$, so we have

$$
d^{n}=\sum_{k=0}^{d}\binom{d}{k} k!S(n, k)=\sum_{k=0}^{d} S(n, k)(d)_{k} .
$$

Since $S(n, k)=0$ if $k>n$, for the last sum, we only have to go up to $n$ instead of $d$. This means $d$ is a root of $x^{n}-\sum_{k=0}^{n} S(n, k)(x)_{k}$. Since nonzero polynomials only have finitely many roots, this difference must be identically 0 .

We can ask the inverse problem: what is the coefficient of $x^{k}$ in $(x)_{n}$ ? It turns out to have a nice interpretation in terms of permutations.

Recall the cycle decomposition of a permutation $\sigma \in \mathfrak{S}_{n}$ : starting with any $1 \leq i \leq n$, we consider the sequence $i, \sigma(i), \sigma^{2}(i), \ldots, \sigma^{k-1}(i)$ where $\sigma^{k}(i)=i$ (there is guaranteed to be such a $k$ since $\sigma$ has finite order). We write the cycle as $i \rightarrow \sigma(i) \rightarrow \cdots \rightarrow \sigma^{k-1}(i) \rightarrow i$. Note that $k$ could be 1 , in which case the cycle has length 1 and also that there isn't a unique beginning (we could have started and ended with $\sigma(i)$ instead of $i$ ). For example, if in 1-line notation, $\sigma=135624$, then the cycle decomposition is $1 \rightarrow 1,2 \rightarrow 3 \rightarrow 5 \rightarrow 2,4 \rightarrow 6 \rightarrow 4$. The cycles form a set partition of $[n]$ but they encode more information.

Let $c(n, k)$ be the number of permutations in $\mathfrak{S}_{n}$ with exactly $k$ different cycles. We use the convention that $c(0,0)=1$. Note that $c(n, 0)=0$ if $n>0$. These are the (signless) Stirling numbers of the first kind.

Proposition 3.29. If $n \geq k \geq 1$, we have

$$
c(n, k)=c(n-1, k-1)+(n-1) c(n-1, k) .
$$

Proof. We break up the permutations with $k$ cycles into 2 types.
The first type consists of permutations such that $n$ is its own cycle. Removing this cycle gives a bijection between such permutations and permutations of $\mathfrak{S}_{n-1}$ with $k-1$ cycles, so the total number is $c(n-1, k-1)$.

The second type consists of permutations such that $n$ is not in its own cycle. In that case, consider its cycle, more specifically, the portion $i \rightarrow n \rightarrow j$. Since $n$ is not in its own
cycle, we know that $i \neq n$ and $j \neq n$. We define a permutation $\tau \in \mathfrak{S}_{n-1}$ by $\tau(i)=j$ and $\tau(k)=\sigma(k)$ for all $k \neq i$. If we remember $i$, then we can reconstruct $\sigma$ uniquely, so we get a bijection between the second type of permutations and pairs $(i, \tau)$ where $1 \leq i \leq n-1$ and $\tau \in \mathfrak{S}_{n-1}$ has $k$ cycles. So there are $(n-1) c(n-1, k)$ many of the second type of permutations.

Corollary 3.30. For $n \geq 0$, we have

$$
\sum_{k=0}^{n} c(n, k) x^{k}=x(x+1) \cdots(x+n-1)
$$

where the right side is 1 if $n=0$, and in particular,

$$
\sum_{k=0}^{n}(-1)^{n-k} c(n, k) x^{k}=(x)_{n}
$$

Proof. We prove the first identity by induction on $n$. For $n=0$, both sides are 1 . Similarly, if $n=1$, both sides are $x$. Now suppose $n \geq 2$. Then $c(n, 0)=c(n-1,0)=0$ and

$$
\begin{aligned}
\sum_{k=1}^{n} c(n, k) x^{k} & =x \sum_{k=1}^{n} c(n-1, k-1) x^{k-1}+(n-1) \sum_{k=1}^{n} c(n-1, k) x^{k} \\
& =x \sum_{k=0}^{n-1} c(n-1, k) x^{k}+(n-1) \sum_{k=0}^{n-1} c(n-1, k) x^{k} \\
& =(x+n-1) \sum_{k=0}^{n-1} c(n-1, k) x^{k} \\
& =(x+n-1) \cdot x(x+1) \cdots(x+n-2)
\end{aligned}
$$

where the last equality is by induction, and this proves what we claimed.
The second identity follows by doing the substitution $x \mapsto-x$ and multiplying by $(-1)^{n}$.

The coefficients $(-1)^{n-k} c(n, k)$ are the (actual) Stirling numbers of the first kind, and are usually denoted $s(n, k)$.

Finally, we have the following "inversion formula":
Corollary 3.31. For $n, \ell \geq 0$, we have

$$
\sum_{k=0}^{n} S(n, k) s(k, \ell)=\delta_{n, \ell}=\sum_{k=0}^{n} s(n, k) S(k, \ell)
$$

Proof. We have

$$
x^{n}=\sum_{k=0}^{n} S(n, k)(x)_{k}=\sum_{k=0}^{n} S(n, k) \sum_{\ell=0}^{k} s(k, \ell) x^{\ell} .
$$

The first sum is the coefficient of $x^{\ell}$ in $x^{n}$, i.e., $\delta_{n, \ell}$ since the $x^{\ell}$ are linearly independent. The second sum is similar keeping in mind that the $(x)_{\ell}$ are also linearly independent.
3.8. Counting subspaces. It turns out that a lot of what we've done has a certain " $q$ enhancement", meaning that we can replace the numbers in our formulas by polynomials in $q$ such that plugging in $q=1$ gives the original result. There's a lot one can do here, so we'll just focus on one important example coming from linear algebra over finite fields.

We've been somewhat ambivalent about our scalars up until now, but mostly just using that they are fields: this is just a structure consisting of "numbers" where you can add, subtract, multiply, and divide (by nonzero elements). The familiar examples are the field of rational numbers, real numbers, and complex numbers. We'll focus here on cases where our field of scalars is finite.

The most important example to keep in mind is $\mathbf{Z} / p$ with $p$ a prime number. You've already seen in your algebra class that we can add, subtract, and multiply. You've probably also seen that we can divide. ${ }^{3}$ For example, if $p=7$, then in $\mathbf{Z} / 7$, we have

$$
\frac{1}{1}=1, \quad \frac{1}{2}=4, \quad \frac{1}{3}=5, \quad \frac{1}{4}=2, \quad \frac{1}{5}=3, \quad \frac{1}{6}=6 .
$$

We won't actually be doing this explicitly, we'll just need to know that it can be done. If $q=p^{k}$ is a positive power of a prime, then there is a unique field (up to isomorphism), denoted $\mathbf{F}_{q}$ of size $q$. If $q=p$, then $\mathbf{F}_{p} \cong \mathbf{Z} / p$. The main thing you need to know is that we can add, subtract, multiply, and divide by nonzero numbers, and we can also do linear algebra with the scalars coming from $\mathbf{F}_{q}$ (if it makes it easier, you can just think of the case $q=p$ and you won't really lose much). I won't reprove everything from scratch in the interest of time. The main point of doing this is that we can count things since everything is a finite set.

Lemma 3.32. If $V$ is an $n$-dimensional $\mathbf{F}_{q}$-vector space, then $|V|=q^{n}$.
Proof. Pick a basis $v_{1}, \ldots, v_{n}$ for $V$. Then every element of $V$ is uniquely of the form $c_{1} v_{1}+\cdots+c_{n} v_{n}$ for $c_{i} \in \mathbf{F}_{q}$ which can be chosen arbitrarily. So vectors are the same as words of length $n$ in $\mathbf{F}_{q}$, and there are $q^{n}$ many of them.

Example 3.33. How many invertible $2 \times 2$ matrices are there with entries in $\mathbf{F}_{q}$ ? This is the set of matrices $\left[\begin{array}{ll}a & b \\ c & d\end{array}\right]$ with $a, b, c, d \in \mathbf{F}_{q}$ such that $a d-b c \neq 0$. Equivalently, the column vectors $(a, c)$ and $(b, d)$ form a basis for $\mathbf{F}_{q}^{2}$, which is also equivalent to saying that $(a, c) \neq(0,0)$ and $(b, d)$ is not a scalar multiple of $(a, c)$. We use this as follows: the number of ways to pick a nonzero vector is $\left|\mathbf{F}_{q}^{2}\right|-1=q^{2}-1$. How many scalar multiples does it have? It spans a 1 -dimensional subspace, so there are $q$ multiples. Hence there are $q^{2}-q$ ways to choose $(b, d)$, and our answer is $\left(q^{2}-1\right)\left(q^{2}-q\right)$.

We'll generalize this to size $n$. The set of invertible $n \times n$ matrices is denoted $\mathbf{G L} \mathbf{L}_{n}\left(\mathbf{F}_{q}\right)$.
Proposition 3.34. The number of $k$-tuples of linearly independent vectors $\left(v_{1}, \ldots, v_{k}\right)$ in an n-dimensional $\mathbf{F}_{q}$-vector space is

$$
\prod_{i=0}^{k-1}\left(q^{n}-q^{i}\right)=\left(q^{n}-1\right)\left(q^{n}-q\right) \cdots\left(q^{n}-q^{k-1}\right)
$$

[^2]In particular, the number of invertible $n \times n$ matrices is

$$
\left|\mathbf{G L}_{n}\left(\mathbf{F}_{q}\right)\right|=\prod_{i=0}^{n-1}\left(q^{n}-q^{i}\right)=\left(q^{n}-1\right)\left(q^{n}-q\right) \cdots\left(q^{n}-q^{n-1}\right)
$$

Proof. These can be chosen as follows: $v_{1}$ is nonzero, $v_{2}$ is not a multiple of $v_{1}, v_{3}$ is not in the linear span of $\left\{v_{1}, v_{2}\right\}$, and in general, the $v_{i}$ th vector is not in the linear span of the first $\left\{v_{1}, \ldots, v_{i-1}\right\}$. The linear span of $\left\{v_{1}, \ldots, v_{i-1}\right\}$ vectors has dimension $i-1$ since they are linearly independent, and so there are $q^{n}-q^{i-1}$ choices for $v_{i}$. This gives the product formula above.

For the second statement, we use that a square matrix is invertible if and only if its columns form a basis.

Finally, we come to our " $q$-analogue". Let $\mathbf{G r}_{k}\left(\mathbf{F}_{q}^{n}\right)$ be the set of $k$-dimensional subspaces of $\mathbf{F}_{q}^{n}$. This is called the Grassmannian.
Theorem 3.35. The number of $k$-dimensional subspaces is

$$
\mathbf{G r}_{k}\left(\mathbf{F}_{q}^{n}\right)=\frac{\left(q^{n}-1\right)\left(q^{n}-q\right) \cdots\left(q^{n}-q^{k-1}\right)}{\left(q^{k}-1\right)\left(q^{k}-q\right) \cdots\left(q^{k}-q^{k-1}\right)}
$$

Proof. Consider pairs $(W, \mathbf{v})$ where $W$ is a $k$-dimensional subspace and $\mathbf{v}$ is an ordered basis for $W$. Note that $\mathbf{v}$ determines $W$, and the choice of $\mathbf{v}$ is just a choice of $k$ ordered linearly independent vectors. By Proposition 3.34, there are $\left(q^{n}-1\right) \cdots\left(q^{n}-q^{k-1}\right)$ many of these. The number of ordered bases of each $W$ is the same, and again by Proposition 3.34 is given by $\left(q^{k}-1\right) \cdots\left(q^{k}-q^{k-1}\right)$. Hence the number of $k$-dimensional spaces is the ratio as claimed.

Let's simplify this ratio as follows. Define the $q$-number and $q$-factorial by

$$
[n]_{q}=\frac{q^{n}-1}{q-1}=1+q+\cdots+q^{n-1}, \quad[n]_{q}!=[n]_{q} \cdots[2]_{q}[1]_{q}
$$

By pulling out common powers of $q$, we have

$$
\operatorname{Gr}_{k}\left(\mathbf{F}_{q}^{n}\right)=\frac{[n]_{q}[n-1]_{q} \cdots[n-k+1]_{q}}{[k]_{q}!}=\frac{[n]_{q}!}{[k]_{q}![n-k]_{q}!},
$$

which we will call the $q$-binomial coefficient, and denote by $\left[\begin{array}{l}n \\ k\end{array}\right]_{q}$. It's not obvious from the formula, but this turns out to be a polynomial in $q$. Note that evaluating $[n]_{q}$ at $q=1$ gives $n$, so evaluating the $q$-binomial coefficient at $q=1$ gives $\binom{n}{k}$.

Example 3.36. For $n=4$ and $k=2$, we have

$$
\left[\begin{array}{l}
4 \\
2
\end{array}\right]_{q}=\frac{\left(q^{4}-1\right)\left(q^{4}-q\right)}{\left(q^{2}-1\right)\left(q^{2}-q\right)}=1+q+2 q^{2}+q^{3}+q^{4}
$$

which is visibly a polynomial that evaluates to $6=\binom{4}{2}$ at $q=1$.
We can see directly that $\left|\mathbf{G r}_{k}\left(\mathbf{F}_{q}^{n}\right)\right|$ is a polynomial in $q$ using reduced-row echeleon form. We can represent $k$-dimensional subspaces of $\mathbf{F}_{q}^{n}$ as the row space of a $k \times n$ matrix. This is not a unique way to represent subspaces, but if we put them into reduced row-echelon form (row operations don't change the row space!), then we do get a unique representation.

Example 3.37. For $n=4$ and $k=2$, every full rank $2 \times 4$ matrix in reduced row-echelon form has exactly 1 of the 6 following types:

$$
\begin{array}{lll}
1 & 0 & *
\end{array} *+1\left[\begin{array}{cccc}
1 & * & 0 & * \\
0 & 1 & * & *
\end{array}\right], \quad\left[\begin{array}{llll}
1 & * & * & 0 \\
0 & 0 & 0 & 1
\end{array}\right],
$$

where the entries in the $*$ are some elements of $\mathbf{F}_{q}$. Furthermore, every choice of scalars gives the reduced row-echelon form for some subspace, so the number of subspaces with reduced row-echelon form of each of the 6 types is given by a power of $q$ (the power being the number of $*$ ). This gives a manifestly polynomial expression $\left|\mathbf{G r}_{2}\left(\mathbf{F}_{q}^{4}\right)\right|=1+q+2 q^{2}+q^{3}+q^{4}$.

Note one thing: if we delete the columns with pivots and reverse the remaining columns, the $*$ form a Young diagram. We actually get every single one that fits into a $2 \times 2$ square.

This is the Schubert decomposition of $\mathbf{G r}_{2}\left(\mathbf{F}_{q}^{4}\right)$. Here's how it generalizes for $\mathbf{G r}_{k}\left(\mathbf{F}_{q}^{n}\right)$ :

- The reduced-row echelon form of a full rank $k \times n$ matrix splits into one of $\binom{n}{k}$ types: these are indexed a subset of $k$ columns (the ones that contain the pivots).
- Given a subset $S=\left\{s_{1}<s_{2}<\cdots<s_{k}\right\}$ of the indices of the $k$ columns, the number of $*$ in row $i$ is $n-s_{i}-(k-i)$. This gives us a bijection between the subsets and Young diagrams fitting into the $k \times(n-k)$ box.
- Hence we get $\left|\mathbf{G r}_{k}\left(\mathbf{F}_{q}^{n}\right)\right|=\sum_{\lambda \subseteq k \times(n-k)} q^{|\lambda|}$.

Remark 3.38. There is no field of size 1 , but if it were to exist, the formula

$$
\lim _{q \rightarrow 1}\left|\mathbf{G r}_{k}\left(\mathbf{F}_{q}^{n}\right)\right|=\binom{n}{k}
$$

suggests that finite sets should play the role of $n$-dimensional vector spaces and $k$-element subsets are then $k$-dimensional subspaces. As stated, this is not rigorous, and is part of a general heuristic that the combinatorics of finite sets should be interpreted as linear algebra over this pretend field. For more information, look up the phrase "field with one element".

Example 3.39. Now fix two integers $k_{1}<k_{2}$ such that $1 \leq k_{1}<k_{2} \leq n$. How many pairs of subspaces $W_{1}$ and $W_{2}$ are there such that $\operatorname{dim} W_{i}=k_{i}$ and $W_{1} \subseteq W_{2}$ ? First we can choose $W_{2}$ in $\left[\begin{array}{c}n \\ k_{2}\end{array}\right]_{q}$ many ways. Now $W_{1}$ is just a choice of $k_{1}$-dimensional subspace of $W_{2}$, and $W_{2}$ is abstractly isomorphic to $\mathbf{F}_{q}^{k_{2}}$. Hence once $W_{2}$ is chosen, there are always $\left[\begin{array}{l}k_{2} \\ k_{1}\end{array}\right]_{q}$ many choices for $W_{1}$. In particular, the total number is

$$
\left[\begin{array}{c}
n \\
k_{2}
\end{array}\right]_{q}\left[\begin{array}{c}
k_{2} \\
k_{1}
\end{array}\right]_{q}=\frac{[n]_{q}!}{\left[k_{1}\right]_{q}!\left[k_{2}-k_{1}\right]_{q}!\left[n-k_{2}\right]_{q}!} .
$$

It's sensible to denote the latter quantity by $\left[\begin{array}{c}n \\ k_{1}, k_{2}-k_{1}, n-k_{2}\end{array}\right]_{q}$ and call it a $q$-multinomial coefficient.

What if we chose $W_{1}$ first? How do we count subspaces $W_{2}$ that contain $W_{1}$ ? The key is to consider the quotient vector space $\mathbf{F}_{q}^{n} / W_{1}$ together with the quotient map $\mathbf{F}_{q}^{n} \rightarrow \mathbf{F}_{q}^{n} / W_{1}$.

The image of $W_{2}$ will always have dimension $k_{2}-k_{1}$. Similarly, for any $\left(k_{2}-k_{1}\right)$-dimensional subspace of $\mathbf{F}_{q}^{n} / W_{1}$, its preimage is a $k_{2}$-dimensional subspace that contains $W_{1}$, so we get a bijection between choices of $W_{2}$ and $\left(k_{2}-k_{1}\right)$-dimensional subspaces of a vector space of dimension $n-k_{1}$, and there are $\left[\begin{array}{c}n-k_{1} \\ k_{2}-k_{1}\end{array}\right]_{q}$ many of them.

We can generalize to more containments of subspaces, but I'll leave it as an exercise.
What else? For example, invertible matrices seem like the right generalization of the symmetric group (they're permutations of $\mathbf{F}_{q}^{n}$ that preserve the linear structure). However, plugging in $q=1$ into the formula for $\left|\mathbf{G L}_{n}\left(\mathbf{F}_{q}\right)\right|$ gives 0 , so it's maybe not a good $q$-analogue of $n!$. But if we divide by all of the powers of $q-1$ and evaluate, we do get $n!$. It turns out that a more direct $q$-analogue is to count complete flags: choices of increasing sequences of subspaces $W_{1} \subset W_{2} \subset \cdots \subset W_{n-1} \subset \mathbf{F}_{q}^{n}$ such that $\operatorname{dim} W_{i}=i$. Generalizing the previous example leads to the conclusion that there are $[n]_{q}$ ! many of these and plugging in $q=1$ gives $n!$.

## 4. Walks in graphs

### 4.1. Adjacency matrix.

Definition 4.1. A graph $G$ is a pair of sets $(V, E)$ where $V$ is the set of vertices, and $E$ is a multiset from $V \cup\binom{V}{2}$, called the edges. The edges in $V$ are called loops. Given an edge, the vertices that it uses are called its endpoints (they could be the same in the case of a loop). If there are no loops and each pair of vertices has at most one edge between them, then the graph is called simple.

We think of elements of $V$ as representing nodes, and the elements of $E$ in $\binom{V}{2}$ tell us which nodes are connected to each other (and how many times). We can think of the elements of $E$ in $V$ as representing self-connections, so we can draw them as loops beginning and ending at the same node. Note there is nothing about locations or lengths in this definition. So while we can draw a graph, such a pictorial representation is not unique.
Definition 4.2. A directed graph is a graph where every edge has a direction. Alternatively, the set of edges is now a multiset from $V \times V$, where $\left(v_{1}, v_{2}\right)$ means an edge that is directed from $v_{1}$ towards $v_{2}$.

If $G$ is a graph with $n$ vertices, then the adjacency matrix is an $n \times n$ matrix which encodes $G$. There are 2 versions, depending on if $G$ is a directed graph or a plain graph.

If $G$ is a graph with vertices $v_{1}, \ldots, v_{n}$, then set $a_{i j}$ to be the number of edges between $v_{i}$ and $v_{j}$. Define its adjacency matrix $A_{G}$ to be the $n \times n$ matrix whose $(i, j)$ entry is $a_{i j}$. Note that $A_{G}$ is a symmetric matrix since $a_{i j}=a_{j i}$. Writing down $A_{G}$ depends on how we order the vertices, but things we calculate from $A_{G}$ won't usually depend on the ordering.

Here is a graph and its adjacency matrix:


If $G$ is a directed graph with vertices $v_{1}, \ldots, v_{n}$, then set $a_{i j}$ to be the number of edges between $v_{i}$ and $v_{j}$ that are going from $v_{i}$ to $v_{j}$ and again $A_{G}$ is the $n \times n$ matrix whose $(i, j)$ entry is $a_{i j}$. In general, $A_{G}$ need not be symmetric.

For example, here's a directed graph and its adjacency matrix:


Definition 4.3. A walk in a graph $G$ is a sequence $v_{0}, e_{1}, v_{1}, e_{2}, v_{2}, \ldots, e_{k}, v_{k}$ which alternates between vertices and edges such that for all $i=1, \ldots, k, v_{i-1}$ and $v_{i}$ are the endpoints of $e_{i}$ (so they must be different unless $e_{i}$ is a loop). The beginning of the walk is $v_{0}$ and the ending is $v_{k}$. If $G$ is directed, we also require that $e_{i}=\left(v_{i-1}, v_{i}\right)$, i.e., that it is oriented from $v_{i-1}$ towards $v_{i}$. A walk is closed if $v_{0}=v_{k}$. The length of the walk is $k$.

Note that a walk of length 0 is just a choice of vertex, so there is always 1 walk of length 0 from a vertex to itself. Also, $A^{0}$ is the identity matrix.

Theorem 4.4. Let $G$ be a graph (directed or not) with vertices $v_{1}, \ldots, v_{n}$. Let $A=A_{G}$ be its adjacency matrix. Then for all integers $k \geq 0$, the number of walks of length $k$ starting at $v_{i}$ and ending at $v_{j}$ is $\left(A^{k}\right)_{i, j}$.

Proof. We prove this by induction on $k$. When $k=0, A^{0}$ is the identity and there is exactly 1 walk of length 0 between a vertex and itself and 0 for different endpoints.

Now suppose the result is known for $k$, we need to prove it for $k+1$. Let $B=A^{k}$. Then

$$
\left(A^{k+1}\right)_{i, j}=(B A)_{i, j}=\sum_{\ell=1}^{n} B_{i, \ell} A_{\ell, j}
$$

by definition of matrix multiplication. The term $B_{i, \ell} A_{\ell, j}$ counts the following: the number of pairs of walks from $v_{i}$ to $v_{\ell}$ of length $k$ and the number of walks from $v_{\ell}$ to $v_{j}$ of length 1 . For each such pair, we can concatenate the walks to get a walk from $v_{i}$ to $v_{j}$ of length $k+1$. Every such walk is accounted for if we sum over all $\ell$ since the last step before reaching $v_{j}$ is some $v_{\ell}$ (and we're including all of them). In particular, the sum counts the number of walks of length $k+1$ from $v_{i}$ to $v_{j}$, so we've proven the statement for $k+1$.

Remark 4.5. A general fact from linear algebra (probably not discussed in Math 18) is that an $n \times n$ symmetric matrix whose entries are all real numbers is always diagonalizable (this is the spectral theorem, and you can use it in this course if needed). In particular, this applies to $A_{G}$ in the undirected case. So we can write $A_{G}=B D B^{-1}$ where $D$ is a diagonal matrix whose entries are the eigenvalues of $A_{G}$. In particular, $A_{G}^{k}=B D^{k} B^{-1}$. So if we want general formulas for the number of walks of length $k$ as $k$ varies, it's enough to diagonalize $A_{G}$. So we see that the eigenvalues of the adjacency matrix are relevant for counting walks, which is surprising.

In the directed case, $A_{G}$ need not be diagonalizable in general.
Let $\operatorname{tr}$ denote the trace of a square matrix. This is the sum of the diagonal entries, and also the sum of the eigenvalues (with multiplicity).

Proposition 4.6. If $\lambda_{1}, \ldots, \lambda_{n}$ are the eigenvalues of $A$, then $\lambda_{1}^{k}, \ldots, \lambda_{n}^{k}$ are the eigenvalues of $A^{k}$, and hence

$$
\sum_{i=1}^{n}\left(A^{k}\right)_{i, i}=\operatorname{tr}\left(A^{k}\right)=\lambda_{1}^{k}+\cdots+\lambda_{n}^{k}
$$

In particular, if $A=A_{G}$, then $\operatorname{tr}\left(A_{G}^{k}\right)$ is the number of closed walks in $G$ of length $k$.
Proof. The trace of a matrix is the sum of its eigenvalues, so the main content is that the eigenvalues of $A^{k}$ are $\lambda_{1}^{k}, \ldots, \lambda_{n}^{k}$ (with these multiplicities). I'll only prove this when $A$ is diagonalizable. Let $v_{1}, \ldots, v_{n}$ be an eigenbasis with eigenvalues $\lambda_{1}, \ldots, \lambda_{n}$. Then $A^{k} v_{i}=$ $\lambda_{i}^{k} v_{i}$, so $v_{1}, \ldots, v_{n}$ is also an eigenbasis for $A^{k}$.

Remark 4.7. If we don't assume that $A$ is diagonalizable, we can prove the result using the Jordan canonical form.

Example 4.8. Fix an integer $n \geq 1$. Let $f(k)$ be the number of words $a_{1} \cdots a_{k}$ of length $k$ in $[n]$ such that (1) $a_{i} \neq a_{i+1}$ for $i=1, \ldots, k-1$ and (2) $a_{k} \neq a_{1}$. If we only impose (1) and don't care about (2), the number of such words is just $n(n-1)^{k-1}$, but how do we control the relation between $a_{k}$ and $a_{1}$ ?

We can interpret a word as a walk of length $k-1$ in a graph with $n$ vertices $v_{1}, \ldots, v_{n}$ so that there is an edge between $v_{i}$ and $v_{j}$ if and only if $i \neq j$. This is the complete graph, denoted $K_{n}$. Letting $A$ be its adjacency matrix, we're interested in $\sum_{i \neq j}\left(A^{k-1}\right)_{i, j}$.

From what we just said, $n(n-1)^{k-1}$ is the total number of walks of length $k-1$ in $K_{n}$ (with any starting or ending point). The number of words satisfying (1) but not (2) is the number of closed walks, so

$$
n(n-1)^{k-1}=f(k)+\operatorname{tr}\left(A^{k-1}\right)
$$

So we want to know the eigenvalues of $A$. Letting $e_{1}, \ldots, e_{n}$ be the standard basis vectors, we have $A e_{i}=\left(\sum_{j=1}^{n} e_{j}\right)-e_{i}$. So a quick calculation shows that

$$
e_{1}+\cdots+e_{n}, e_{1}-e_{2}, e_{2}-e_{3}, \ldots, e_{n-1}-e_{n}
$$

are eigenvectors of $A$ with eigenvalues $n-1,-1,-1, \ldots,-1$. Further, they are linearly independent and hence an eigenbasis. So $\operatorname{tr}\left(A^{k-1}\right)=(n-1)^{k-1}+(n-1)(-1)^{k-1}$, and we solve:

$$
\begin{aligned}
f(k) & =n(n-1)^{k-1}-\operatorname{tr}\left(A^{k-1}\right) \\
& =n(n-1)^{k-1}-(n-1)^{k-1}-(n-1)(-1)^{k-1} \\
& =(n-1)^{k}+(n-1)(-1)^{k}
\end{aligned}
$$

4.2. Transfer matrix method. As we just saw, we can set up some counting problems as counting paths in a graph. So we now study generating functions of the form

$$
F_{A ; i, j}(x)=\sum_{k \geq 0}\left(A^{k}\right)_{i, j} x^{k}
$$

where $A$ is an $n \times n$ matrix and $1 \leq i, j \leq n$.
Below, if $B$ is any $n \times n$ matrix, then $(B ; j, i)$ is the $(n-1) \times(n-1)$ submatrix obtained by deleting row $j$ and column $i$ from it.

Theorem 4.9. We have

$$
F_{A ; i, j}(x)=(-1)^{i+j} \frac{\operatorname{det}\left(\left(\mathrm{id}_{n}-x A\right) ; j, i\right)}{\operatorname{det}\left(\operatorname{id}_{n}-x A\right)}
$$

so that $F_{A ; i, j}(x)$ is a rational generating function.
Proof. Consider $\sum_{k \geq 0} A^{k} x^{k}$, which we think of as an $n \times n$ matrix whose entries are formal power series in $x$. A calculation similar to Example 2.2 shows that this is the inverse of the matrix $\mathrm{id}_{n}-x A$. So we want the $(i, j)$-entry of $\left(\mathrm{id}_{n}-x A\right)^{-1}$. So the formula that we claim follows from Cramer's rule, which computes this quantity.

As for the last claim, both determinants give us polynomials in $x$.
Let $\lambda_{1}, \ldots, \lambda_{n}$ be the eigenvalues of $A$. These are the roots of the characteristic polynomial

$$
\left(t-\lambda_{1}\right) \cdots\left(t-\lambda_{n}\right)=\operatorname{det}\left(\operatorname{tid}_{n}-A\right)
$$

By doing the substitution $t \mapsto 1 / x$ and multiplying both sides by $x^{n}$, we get

$$
\left(1-\lambda_{1} x\right) \cdots\left(1-\lambda_{n} x\right)=\operatorname{det}\left(\mathrm{id}_{n}-x A\right)
$$

Hence from $\S 2.4$, we know that $\left(A^{k}\right)_{i, j}$ has a formula in terms of the eigenvalues $\lambda_{1}, \ldots, \lambda_{n}$ of $A$. In the simplest case when all of the eigenvalues are distinct, we know that there exist constants $c_{1}, \ldots, c_{n}$ (they depend on $i, j$ ) such that

$$
\left(A^{k}\right)_{i, j}=c_{1} \lambda_{1}^{k}+\cdots+c_{n} \lambda_{n}^{k}
$$

for $k \geq \max \left(0, \operatorname{deg} F_{A ; i, j}(x)+1\right)$, and in all cases, $n-1 \geq \operatorname{deg} F_{A ; i, j}(x)$, so we can say it holds for $k \geq n$ without any specific knowledge.

In particular, without doing any calculations, we see that the number of paths in a graph between $v_{i}$ and $v_{j}$ of length $k$ satisfies a linear recurrence relation with respect to $k$.

Example 4.10. Consider length $n$ words in [3] such that 11 and 23 do not appear consecutively. These can be encoded as walks of length $n-1$ in the following graph $G$
with adjacency matrix


$$
A_{G}=\left[\begin{array}{lll}
0 & 1 & 1 \\
1 & 1 & 0 \\
1 & 1 & 1
\end{array}\right]
$$

Then $\operatorname{det}\left(\mathrm{id}_{3}-x A_{G}\right)=1-2 x-x^{2}+x^{3}$. For simplicity, let's just compute the number of words starting with 1 and ending at 3 . Then $\operatorname{det}\left(\operatorname{id}_{3}-x A_{G} ; 3,1\right)=x-x^{2}$, so the generating function is

$$
\frac{x-x^{2}}{1-2 x-x^{2}+x^{3}},
$$

but remember that the coefficient of $x^{n}$ is the number words of length $n+1$.
Example 4.11. Consider the problem of covering a $n \times k$ chessboard with $1 \times 2$ size dominoes (placed horizontally or vertically) so that no two overlap. We want to know how many ways $f_{n}(k)$ this can be done.
(1) If $n=1$, there is exactly 1 way if $k$ is even and no way to do it if $k$ is odd.
(2) If $n=2$, consider the right most portion of the board. Either the last column has a vertical domino, or the last two columns have 2 horizontal dominoes. This tells us there is a recurrence $f_{2}(k)=f_{2}(k-1)+f_{2}(k-2)$ for $k \geq 3$, while $f_{2}(1)=1$ and $f_{2}(2)=2$. So these are the Fibonacci numbers.
(3) We can try the same analysis for $n=3$, though it seems harder to guess what the different cases should be, but you could squeeze out a recurrence relation with enough effort.

For general $n$, we can encode tilings as walks in a graph. First, given a valid tiling, consider a single column. In each square, we either have a horizontal domino or vertical domino. If we remember which squares are the left portion of a horizontal domino, we can reconstruct the rest of the dominoes uniquely. Hence we can encode each column by a word of length $n$ in the alphabet $\{0,1\}$ ( 1 denotes the squares which are the left portion of a horizontal domino and 0 denotes the others).

Let $G_{n}$ be the directed graph whose vertices are length $n$ words in $\{0,1\}$. We have an edge from one word $w$ to another word $w^{\prime}$ if it's possible to place dominoes so that two consecutive columns are labeled by $w$ and $w^{\prime}$ (don't worry about whether it can be completed to a tiling of the whole chessboard). The benefit is that determining the edges is a local problem: we only consider 2 columns at a time. For $n=2, G_{2}$ looks like


Finally we have boundary conditions. The rightmost column always has to be all 0's. On the other hand, there are restrictions on which words can appear as the leftmost column. For $n=2$, only 00 and 11 are the only valid leftmost columns. Note that 10 and 01 appearing in their own component just says you never see these in a tiling of a chessboard. Without figuring out what all of the conditions are, we can say that the number of valid tilings is the number of walks of length $k-1$ starting where the first vertex is a valid leftmost column and the last vertex is all 0's. This is a finite sum of generating functions of the form $F_{A_{G_{n}} ; i, j}(x)$, so it is rational by Theorem 4.9.

In particular, there exists a linear recurrence relation for $f_{n}(k)$ with respect to $k$ whenever $n$ is fixed. We won't pursue it here, but it is possible to get a closed formula for all $n$.

Another benefit of this approach is that we can change the boundary conditions to address related problems. For example, we could glue together the left and right ends of the chessboard and ask about tiling the resulting circular band. Then we'd be asking about closed walks in $G_{n}$.

Remark 4.12 (If you know some complexity theory). If we fix a set of starting vertices and ending vertices, then a directed graph can also be interpreted as a deterministic finite-state automaton (DFA). Hence in the above problems, the words we are counting are the words in a regular language. In fact, it's no more general, so that counting words in a regular language can be done with the transfer-matrix method. Hence, if $L$ is a regular language with respect to some finite alphabet and $a_{k}$ is the number of words of length $k$ in $L$, then $\sum_{k \geq 0} a_{k} x^{k}$ is a rational function.

There is a natural hierarchy of languages. For example, the languages described by an unambiguous context-free grammar have an algebraic generating function, i.e., they are the
solution to a polynomial equation (in $t$ ) whose coefficients are themselves polynomials (in $x)$.

## 5. Exponential generating functions

5.1. Products of exponential generating functions. Let $a_{0}, a_{1}, a_{2}, \ldots$ be a sequence of numbers. The associated exponential generating function (EGF) is the formal power series

$$
A(x)=\sum_{n \geq 0} a_{n} \frac{x^{n}}{n!}
$$

where recall that $n!=n(n-1)(n-2) \cdots 2 \cdot 1$ and $0!=1$. When $a_{n}=1$ for all $n$, we use the notation

$$
e^{x}=\exp (x)=\sum_{n \geq 0} \frac{x^{n}}{n!}
$$

Lemma 5.1. If $A(x)=\sum_{n \geq 0} a_{n} \frac{x^{n}}{n!}$ and $B(x)=\sum_{n \geq 0} b_{n} \frac{x^{n}}{n!}$, then $A(x) B(x)=\sum_{n \geq 0} c_{n} \frac{x^{n}}{n!}$ where $c_{n}=\sum_{i=0}^{n}\binom{n}{i} a_{i} b_{n-i}$.

Proof. The coefficient of $x^{n}$ in $A(x) B(x)$ is $\sum_{i=0}^{n} \frac{a_{i}}{i!} \frac{b_{n-i}}{(n-i)!}$. By definition it is also $c_{n} / n$ !, so $c_{n}=\sum_{i=0}^{n}\binom{n}{i} a_{i} b_{n-i}$.

It will be convenient to think of the coefficients of an EGF as counting the number of structures on a set. Formally, a structure is a function $\alpha$ that takes as input a finite set $S$ (including $S=\varnothing$ ) and outputs another finite set $\alpha(S)$, with the key property that if $|S|=|T|$, then $|\alpha(S)|=|\alpha(T)|$. We've been dealing with many of these. Some examples are $\alpha(S)$ is the set of 2-element subsets of $S$, or the set of set partitions of $S$, or the set of bijections from $S$ to itself, etc. We will say that elements of $\alpha(S)$ are structures of type $\alpha$, and the associated exponential generating function is

$$
E_{\alpha}(x)=\sum_{n \geq 0}|\alpha([n])| \frac{x^{n}}{n!}
$$

Let $\alpha, \beta$ be structures. We can add and multiply structures:

$$
\begin{aligned}
(\alpha+\beta)(S) & =\alpha(S) \amalg \beta(S) \\
(\alpha \cdot \beta)(S) & =\coprod_{T \subseteq S} \alpha(T) \times \beta(S \backslash T)
\end{aligned}
$$

The sum is just taking disjoint union. The product requires more explanation: we are taking the disjoint union over all subsets $T$ in $S$, picking an $\alpha$-structure on $T$ and a $\beta$-structure on its complement. We'll see in examples why this is a sensible thing to do, but first, we show that these operations behave well with respect to EGFs:

Proposition 5.2. We have

$$
E_{\alpha+\beta}(x)=E_{\alpha}(x)+E_{\beta}(x), \quad E_{\alpha \cdot \beta}(x)=E_{\alpha}(x) E_{\beta}(x) .
$$

Proof. For the sum, we have $|(\alpha+\beta)([n])|=|\alpha([n])|+|\beta([n])|$ since we're taking a disjoint union.

For the product, we have

$$
|(\alpha \cdot \beta)([n])|=\sum_{T \subseteq[n]}|\alpha(T)| \cdot|\beta([n] \backslash T)|
$$

Since the size of $\alpha(T)$ only depends on $|T|$ and similarly for $\beta([n] \backslash T)$, we can just sum over possible sizes of $T$ :

$$
\sum_{i=0}^{n}\binom{n}{i}|\alpha([i])| \cdot|\beta([n-i])|
$$

which is the coefficient of $E_{\alpha}(x) E_{\beta}(x)$ by Lemma 5.1.
Example 5.3. Consider a set of $n$ football players. We want to split them up into two groups. Both groups needs to be assigned an ordering and the second group additionally needs to choose one of 3 colors for their uniform. Let $c_{n}$ be the number of ways to do this.

This scenario calls for a product of structures:

- Let $\alpha(S)$ be the set of orderings of $S$, so $|\alpha(S)|=|S|$ !. We have

$$
E_{\alpha}(x)=\sum_{n \geq 0} n!\frac{x^{n}}{n!}=\frac{1}{1-x}
$$

- Let $\beta(S)$ be the set of pairs $(\sigma, f)$ where $\sigma$ is an ordering of $S$ and $f: S \rightarrow[3]$ is an assignment of the 3 colors to each element. So $|\beta(S)|=|S|!3^{|S|}$. We have

$$
E_{\beta}(x)=\sum_{n \geq 0} n!3^{n} \frac{x^{n}}{n!}=\frac{1}{1-3 x}
$$

Then $(\alpha \cdot \beta)([n])$ is the set of things we're asking about (I glossed over it, but it's important that the definitions above make sense and give the correct thing when $S=\varnothing$, otherwise our product interpretation will be incorrect when $T=\varnothing$, for example), so its EGF is

$$
E_{\alpha \cdot \beta}(x)=\frac{1}{(1-x)(1-3 x)} .
$$

In particular,

$$
c_{n} / n!=\left[x^{n}\right]_{\frac{1}{(1-x)(1-3 x)}}=\left[x^{n}\right]\left(\frac{3 / 2}{1-3 x}-\frac{1 / 2}{1-x}\right)=\frac{3}{2} 3^{n}-\frac{1}{2},
$$

and hence

$$
c_{n}=n!\left(\frac{3}{2} 3^{n}-\frac{1}{2}\right)=\frac{n!}{2}\left(3^{n+1}-1\right) .
$$

Example 5.4. We have $n$ distinguishable telephone polls which are to be painted either red or blue. The number which are blue must be even. Let $c_{n}$ be the number of ways to do this.

Again we want to interpret this as counting the product of two structures (we'll think of the elements of sets as telephone polls):

- Let $\alpha(S)$ be the set of ways to paint the polls red according to our rules, so $|\alpha(S)|=1$ for all $S$ (even $S=\varnothing$ ) and $E_{\alpha}(x)=e^{x}$.
- Let $\beta(S)$ be the set of ways to paint the polls blue according to our rules, so $|\alpha(S)|=1$ if $|S|$ is even and $|\alpha(S)|=0$ if $|S|$ is odd. Hence

$$
E_{\beta}(x)=\sum_{n \geq 0} \frac{x^{2 n}}{(2 n)!}
$$

Here we are deleting all of the odd powers of $x$ from $e^{x}$. To get a nice expression, note that this is the same as $\left(e^{x}+e^{-x}\right) / 2$. (How about if we wanted to delete the even terms instead?)
Hence we get (I leave it as an exercise to check that $e^{A(x)} e^{B(x)}=e^{A(x)+B(x)}$ for any formal power series $A, B$ with no constant term):

$$
E_{\alpha \cdot \beta}(x)=\frac{1}{2} e^{x}\left(e^{x}+e^{-x}\right)=\frac{1}{2}\left(e^{2 x}+1\right)=\frac{1}{2} \sum_{n \geq 0} \frac{2^{n} x^{n}}{n!}+\frac{1}{2}
$$

So $c_{n}=2^{n-1}$ if $n>0$ and $c_{0}=1$.
Actually we could have derived this formula using earlier stuff: we're just trying to pick a subset of even size to be painted blue. We know that half of the subsets of $[n]$ have even size and half have odd size, so we can also see $2^{n-1}$. However, the approach given here generalizes more easily if we introduce more colors, for example.

We can multiply more than 2 structures at once. By iterating the case of 2 structures, we come to the following definition and result. Let $\alpha_{1}, \ldots, \alpha_{k}$ be structures. Then their product is

$$
\left(\alpha_{1} \cdots \alpha_{k}\right)(S)=\coprod_{\substack{\left(T_{1}, \ldots, T_{k}\right) \\ T_{1}, \cdots T_{k}=S \\ T_{i} \cap T_{j}=\varnothing \text { for } i \neq j}} \alpha_{1}\left(T_{1}\right) \times \cdots \times \alpha_{k}\left(T_{k}\right)
$$

where the disjoint union is over all ways to write $S$ as a disjoint union of $k$ subsets. This is almost like an ordered set partition, except that the $T_{i}$ are allowed to be empty. Then

$$
E_{\alpha_{1} \cdots \alpha_{k}}(x)=E_{\alpha_{1}}(x) \cdots E_{\alpha_{k}}(x)
$$

Example 5.5. Continuing from the previous example, suppose we can also color some telephone polls green and there are no restrictions on how many are green. This introduces a third structure: let $\gamma(S)$ be the ways to paint the polls green, so $|\gamma(S)|=1$ for all $S$. Our new EGF is

$$
E_{\alpha \cdot \beta \cdot \gamma}(x)=\frac{1}{2} e^{x}\left(e^{x}+e^{-x}\right) e^{x}=\frac{1}{2}\left(e^{3 x}+e^{x}\right)=\frac{1}{2}\left(\sum_{n \geq 0} \frac{(3 x)^{n}}{n!}+\sum_{n \geq 0} \frac{x^{n}}{n!}\right),
$$

so the answer we want is $\frac{1}{2}\left(3^{n}+1\right)$.
Example 5.6. Consider the following structure:

$$
\alpha(S)= \begin{cases}\{1\} & \text { if }|S|>0 \\ \varnothing & \text { if }|S|=0\end{cases}
$$

We can think of this as a "selection structure" which picks out nonempty subsets. In particular, $(\alpha \cdot \alpha)(S)$ is the number of nonempty subsets $T \subseteq S$ such that $S \backslash T$ is also nonempty. In other words, it is an ordered set partition with 2 blocks. More generally, $\alpha^{k}(S)$ is the set of ordered set partitions of $S$ with $k$ blocks. Hence

$$
\sum_{n \geq 0} k!S(n, k) \frac{x^{n}}{n!}=E_{\alpha^{k}}(x)=E_{\alpha}(x)^{k}=\left(e^{x}-1\right)^{k}
$$

and also

$$
\sum_{n \geq 0} S(n, k) \frac{x^{n}}{n!}=\frac{\left(e^{x}-1\right)^{k}}{k!}
$$

By modifying the definition of $\alpha$ we can get formulas for set partitions with different conditions on the sizes of the blocks (or even using $k$ different modifications).
5.2. Compositions of exponential generating functions. Now we consider a structure $\alpha$ such that $\alpha(\varnothing)=\varnothing$. For a finite set $S$, let $\Pi_{S}$ be the set of set partitions of $S$. We define $e^{\alpha}$ to be the following structure:

$$
e^{\alpha}(S)=\coprod_{\left\{S_{1}, \ldots, S_{k}\right\} \in \Pi_{S}} \alpha\left(S_{1}\right) \times \cdots \times \alpha\left(S_{k}\right)
$$

In other words, we consider all set partitions of $S$, and put structures of type $\alpha$ on each block. There is some ambiguity about the order to take the product since the $S_{i}$ are not ordered, but this choice won't matter much since we only care about the size of $e^{\alpha}(S)$. To make this well-defined, we could pick an ordering of $S$ and can take the convention that we list blocks so that $\min \left(S_{1}\right)<\min \left(S_{2}\right)<\cdots<\min \left(S_{k}\right)$.
Theorem 5.7. We have

$$
E_{e^{\alpha}}(x)=\exp \left(E_{\alpha}(x)\right)
$$

Proof. Since $|\alpha(\varnothing)|=0$, we have $\left[x^{n}\right] E_{\alpha}(x)^{k}=0$ if $k>n$. So

$$
\left[x^{n}\right] \exp \left(E_{\alpha}(x)\right)=\left[x^{n}\right] \sum_{k \geq 0} \frac{E_{\alpha}(x)^{k}}{k!}=\left[x^{n}\right] \sum_{k=0}^{n} \frac{E_{\alpha}(x)^{k}}{k!} .
$$

From our discussion on products of EGFs, $\left[x^{n}\right] E_{\alpha}(x)^{k}$ is the number of ways to pick an ordered set partition of $[n]$ into $k$ blocks and put structures of type $\alpha$ on each block; if we divide by $k$ ! we just remove the ordering. Hence the coefficient of $x^{n}$ above is exactly the size of $e^{\alpha}([n])$.

If an EGF has the form $\exp (A(x))$, we can try to use the "logarithmic derivative" to get recurrence relations on its coefficients. I'll explain the setup in a homework problem and instead focus on examples of the above theorem.

Example 5.8. We continue with Example 5.6 and define the structure

$$
\alpha(S)= \begin{cases}\{1\} & \text { if }|S|>0 \\ \varnothing & \text { if }|S|=0\end{cases}
$$

Then $\left|e^{\alpha}(S)\right|$ is the number of set partitions of $S$, so we get the EGF for Bell numbers:

$$
\sum_{n \geq 0} B(n) \frac{x^{n}}{n!}=E_{e^{\alpha}}(x)=\exp \left(E_{\alpha}(x)\right)=\exp \left(e^{x}-1\right)
$$

There is a general heuristic that if $\alpha(S)$ consists of the set of "connected" structures, then $e^{\alpha}(S)$ is the set of all structures (not necessarily connected) since generally speaking, arbitrary structures are disjoint unions of connected ones. Rather than make this precise, let's illustrate with some examples.

Example 5.9. An undirected graph is connected if, for any two vertices $x$ and $y$, there exists a walk from $x$ to $y$. For a general undirected graph $G$, we put an equivalence relation on the vertices by $x \sim y$ if there is a walk from $x$ to $y$, and the equivalence classes are connected graphs, which are called the connected components of $G$. Intuitively, this is just saying that every graph is a disjoint union of connected ones.

For a finite set $S$, let $\beta(S)$ be the set of all simple graphs whose vertex set is $S$. This is just the set of all subsets of the 2-element subsets of $S$, so $|\beta(S)|=2^{d}$ with $d=\binom{|S|}{2}$, and $\beta$ is a structure. If $|S|>0$, let $\alpha(S)$ be the set of all simple connected graphs whose vertex set is $S$ and let $\alpha(\varnothing)=\varnothing$. It's harder to get a formula for $|\alpha(S)|$, but it's not hard to see it only depends on $|S|$. But we can say that $e^{\alpha}=\beta$ and $\exp \left(E_{\alpha}(x)\right)=E_{\beta}(x)$.

We might try to use this since we have a formula for $|\beta([n])|$, but the EGF doesn't seem to have a nice form. We can define a formal version of the logarithm function to write $E_{\alpha}(x)$ in terms of $E_{\beta}(x)$, but we'll omit this discussion. See Homework 2, Problem 8 for details and various properties.

Example 5.10. Let's do something more abstract. Recall that every permutation has a cycle decomposition. So we can think of cycles as being "connected" permutations. For $S \neq \varnothing$, let $\alpha(S)$ be the set of ways to cyclically order the elements of $S$, so that $|\alpha(S)|=(|S|-1)$ !, and let $\alpha(\varnothing)=\varnothing$. Then $e^{\alpha}(S)$ can be interpreted as the set of permutations of $S$, so that $\left|e^{\alpha}(S)\right|=|S|$ !. Hence we have

$$
\exp \left(\sum_{n \geq 1} \frac{x^{n}}{n}\right)=\exp \left(E_{\alpha}(x)\right)=E_{e^{\alpha}}(x)=\sum_{n \geq 0} x^{n}=\frac{1}{1-x},
$$

which gives us a standard identity $\sum_{n \geq 1} \frac{x^{n}}{n}=\log \left((1-x)^{-1}\right)$.
Example 5.11. A bijection $f:[n] \rightarrow[n]$ is an involution if $f \circ f$ is the identity function. This is just a permutation whose cycles all have length 1 or 2 . So the connected involutions are cycles of length 1 or 2 . Hence let $\alpha(S)$ be the set of cyclic orderings of $S$ if $1 \leq|S| \leq 2$ and $\alpha(S)=\varnothing$ otherwise. Then $E_{\alpha}(x)=x+x^{2} / 2$. Also, $e^{\alpha}(S)$ can be interpreted as the set of involutions on $S$. So we have

$$
E_{e^{\alpha}}(x)=\exp \left(x+\frac{x^{2}}{2}\right)
$$

Finally, we can interpret general compositions as follows. Let $\alpha$ be a structure such that $\alpha(\varnothing)=\varnothing$ and let $\beta$ be a general structure. We define the composition by

$$
(\beta \circ \alpha)(S)=\coprod_{\left\{S_{1}, \ldots, S_{k}\right\} \in \Pi_{S}} \beta([k]) \times \alpha\left(S_{1}\right) \times \cdots \times \alpha\left(S_{k}\right)
$$

Again, the blocks aren't ordered, see the discussion above on how to deal with that issue. We can think of this as picking a set partition and putting an $\alpha$ structure on each block as before, but we also additionally put a $\beta$ structure on the set of blocks. The proof of the following is pretty close to the proof of the exponential formula, so we'll skip it.

Theorem 5.12 (Composition formula). With the notation above,

$$
E_{\beta \circ \alpha}(x)=E_{\beta}\left(E_{\alpha}(x)\right) .
$$

5.3. Cayley's enumeration of labeled trees and Lagrange inversion. As discussed in Example 5.9, a labeled (simple) graph on a set $S$ is a collection of 2-element subsets of $S$. A cycle is a closed walk which does not repeat any edges. If the graph has no cycles, it is called a labeled forest. If, in addition, it is connected, then is a labeled tree. Let $t_{n}$ be the number of labeled trees on $[n]$. Our goal is the following formula for $t_{n}$ (as discussed before, the number of labeled graphs is $2^{\binom{n}{2}}$ so there isn't much to discuss):

Theorem 5.13 (Cayley). $t_{n}=n^{n-2}$.
There are a lot of different ways to get this, but we will focus on using EGF.
Example 5.14. When $n=1$ or $n=2$, we get 1 labeled tree. When $n=3$, we get 3 , corresponding to the following pictures:


When $n=4$, there are 2 types of unlabeled trees:


There are 4 labelings of the first kind since it only matters what goes in the middle, and the second has $12=4!/ 2$ labelings since a labeling can be thought of as a permutation of size 4 , except that reversing the order gives the same tree.

We need one more definition: a rooted labeled tree is a labeled tree where one of the points has been designated as the "root". The number of rooted labeled trees is then $n t_{n}$. Similarly, we define a planted labeled forest to be a labeled forest in which each connected component is a rooted labeled tree. Let $f_{n}$ be the number of planted labeled forests. Define EGFs

$$
F(x)=\sum_{n \geq 0} f_{n} \frac{x^{n}}{n!}, \quad R(x)=\sum_{n \geq 0} n t_{n} \frac{x^{n}}{n!}
$$

Lemma 5.15. $F(x)=e^{R(x)}$.
Proof. Every planted labeled forest is a disjoint union of rooted labeled trees, so this follows from the exponential formula.

Lemma 5.16. $R(x)=x F(x)$.
Proof. For $n \geq 1$, we can construct all rooted labeled trees on $n n]$ uniquely in the following way. First, pick some element $i$ to be the root. Second, put the structure of a planted labeled forest on $[n] \backslash\{i\}$. Given this information, we join $i$ to each of the roots of the trees that make up our forest and then forget that they are roots.

Conversely, given a rooted labeled tree, if we delete the root, then we are left with a labeled forest. Each point that was previously connected to the root is now in a separate component (if they were still connected, then the original graph had a cycle because we could go through the root and then through go through whatever path remains), so we can declare all of them to be the roots of their respective components.

In conclusion, we see that $R(x)$ is the EGF for first picking an element of $[n]$ and then putting a planted labeled forest on the remaining elements. Hence

$$
n!\left[x^{n}\right] R(x)=n t_{n}=n f_{n-1}=n!\left[x^{n-1}\right] F(x)=n!\left[x^{n}\right] x F(x)
$$

for all $n \geq 1$ (and the constant terms of $R(x)$ and $x F(x)$ are 0 ), and so $R(x)=x F(x)$.
In particular, we have the recursion

$$
R(x)=x e^{R(x)}
$$

We can try to solve this coefficient by coefficient: say that $R(x)=\sum_{n \geq 0} r_{n} x^{n}$ and we are trying to solve for the $r_{i}$. The left hand side has no constant term, so we must have $r_{0}=0$. This tells us that $R(x)^{n}$ starts with the term $x^{n}$. Expanding the equation, we get

$$
R(x)=x\left(1+R(x)+\frac{R(x)^{2}}{2!}+\cdots\right) .
$$

So if we want to solve for $r_{n}$ we just need to consider $x\left(1+R(x)+\cdots+\frac{R(x)^{n-1}}{(n-1)!}\right)$ since all other terms don't have a $x^{n}$ term. In particular,

$$
\begin{aligned}
& r_{1}=\left[x^{1}\right] R(x)=\left[x^{1}\right] x=1 \\
& r_{2}=\left[x^{2}\right] R(x)=\left[x^{2}\right] x(1+R(x))=0+r_{1}=1, \\
& r_{3}=\left[x^{3}\right] R(x)=\left[x^{3}\right] x\left(1+R(x)+\frac{R(x)^{2}}{2}\right)=0+r_{2}+\frac{r_{0} r_{2}+r_{1}^{2}+r_{2} r_{0}}{2}=\frac{3}{2}, \\
& \vdots
\end{aligned}
$$

We can continue like this, but it would be nice to have a closed formula without having to guess one. This can be done with the Lagrange inversion formula (due to time constraints, we'll skip the proof):

Theorem 5.17 (Lagrange inversion formula). Let $G(x)$ be a formal power series whose constant term is nonzero. There is a unique formal power series $A(x)$ such that $A(0)=0$ and

$$
A(x)=x G(A(x))
$$

For $k, n \geq 0$, we have

$$
n\left[x^{n}\right] A(x)^{k}=k\left[x^{n-k}\right]\left(G(x)^{n}\right) .
$$

Remark 5.18. Showing that there exists a unique formal power series $A(x)$ satisfying the equation $A(x)=x G(A(x))$ when $G(0) \neq 0$ follows the example above since the coefficients of $A(x)$ can be determined one at a time. The real content (whose proof we're skipping) is the nice formula for these coefficients in terms of $G(x)$.
Proof of Theorem 5.13. We take $A(x)=R(x)$ and $G(x)=e^{x}$ and for the moment, let's take $k$ general. For $n>0$, the Lagrange inversion formula tells us that

$$
\left[x^{n}\right] R(x)^{k}=\frac{k}{n}\left[x^{n-k}\right] e^{n x}=\frac{k}{n}\left[x^{n-k}\right] \sum_{d \geq 0} \frac{n^{d}}{d!} x^{d}=\frac{k}{n} \frac{n^{n-k}}{(n-k)!}
$$

We're interested in $k=1$ which simplifies to $n^{n-1} / n!$. Remember that $\left[x^{n}\right] R(x)=n t_{n} / n$ !, so we conclude that $t_{n}=n^{n-2}$.
Corollary 5.19. The number of planted labeled forests on $n$ vertices with $k$ connected components is

$$
\binom{n-1}{k-1} n^{n-k}
$$

Proof. Using the product formula for EGF, $n!\left[x^{n}\right] R(x)^{k}$ is the number of planted labeled forests with $k$ connected components, together with an ordering of the components, and

$$
\frac{n!}{k!}\left[x^{n}\right] R(x)^{k}=\frac{n!}{k!} \frac{k}{n} \frac{n^{n-k}}{(n-k)!}=\binom{n-1}{k-1} n^{n-k} .
$$

There doesn't seem to be a simple way to get a nice formula for the number of labeled forests with $k$ connected components (for $k=1$ we just divide by $n$ because we know the size of the single component, but in general, the sizes can vary).

We'll finish with another example of Lagrange inversion which generalizes the formula for Catalan numbers.

Example 5.20. Recall that we discussed why Catalan numbers count the number of rooted binary trees with $n+1$ leaves. Equivalently, this is the number of rooted binary trees with $n$ internal vertices. More generally, we can fix $k$ and consider rooted $k$-ary trees with $n$ internal vertices. We'll leave $k$ out of the notation for simplicity, and let $c_{n}$ be the number of rooted $k$-ary trees with $n$ internal vertices. To build one when $n>0$, we start with a single node for our root, and then attach $k$ rooted $k$-ary trees below it. This gives us the relation

$$
c_{n}=\sum_{\substack{\left(i_{1}, i_{2}, \ldots, i_{k}\right) \\ i_{1}+\cdots+i_{k}=n-1}} c_{i_{1}} c_{i_{2}} \cdots c_{i_{k}} \quad \text { for } n>0
$$

The sum is over all weak compositions of $n-1$ with $k$ parts. Here $i_{j}$ represents the number of internal vertices that are in the $j$ th tree connected to our original root. As before, if $C(x)=\sum_{n \geq 0} c_{n} x^{n}$, this leads to the relation

$$
C(x)=1+x C(x)^{k} .
$$

Now we don't have a general method of solving this polynomial equation, but we can use Lagrange inversion. We set $A(x)=C(x)-1$ to convert the relation into

$$
A(x)=x(A(x)+1)^{k} .
$$

So we take $G(x)=(x+1)^{k}$ and we conclude that

$$
\left[x^{n}\right] A(x)=\frac{1}{n}\left[x^{n-1}\right](x+1)^{k n}=\frac{1}{n}\binom{k n}{n-1}=\frac{1}{(k-1) n+1}\binom{k n}{n} .
$$

## 6. Sieving methods

6.1. Möbius inversion. A partially ordered set (poset for short) is an abstraction for systems where some things can be compared and some things might not be comparable. First, we give the formal definition. Recall that a relation $R$ on a set $S$ is a collection of ordered pairs of elements. If $(x, y)$ is in the relation, we usually just write $x R y$. Below, our relation will be written as $\leq$ to be suggestive that it is a comparison.

Definition 6.1. Let $P$ be a set. A relation $\leq$ on $P$ is a partial ordering if it satisfies the following 3 conditions:
(1) (Reflexive) For all $x \in P, x \leq x$.
(2) (Transitive) If $x \leq y$ and $y \leq z$, then $x \leq z$.
(3) (Anti-symmetric) If $x \leq y$ and $y \leq x$, then $x=y$.

The pair $(P, \leq)$ is a partially ordered set (poset). Given two elements $x, y \in P$, they are comparable if either $x \leq y$ or $y \leq x$, and otherwise they are incomparable.

The notation $x<y$ is always shorthand to mean that $x \leq y$ and $x \neq y$.
We will assume $P$ is a finite set (much weaker assumptions can be made, but it won't benefit us immediately, so let's keep things simple). If all pairs of elements are comparable,
then $\leq$ is called a total ordering, which are perhaps more familiar to you. Most of the examples we deal with are not total orderings.

Example 6.2. Let $P=[n]$ and write $x \leq y$ if $x$ is smaller than $y$ in the usual sense. If we write $[n]$ for a poset, we will mean this one. This is a total ordering.

Example 6.3. Let $S$ be a set and let $P$ be the set of all subsets of $S$. Given $x, y \in P$, we define $x \leq y$ to mean that $x$ is a subset of $y$. Then $(P, \leq)$ is a poset, called the Boolean poset of $S$. When $S=[n]$, we will use the notation $B_{n}$ for $P$. When $n \geq 2$ this is not a total ordering.

Example 6.4. Let $P$ be the set of positive integers. Given $x, y \in P$, we define $x \leq y$ if $x$ divides $y$. Since it can be confusing, we will usually write | instead of $\leq$, so that the notation is $x \mid y$. We will use the notation $\left(\mathbf{Z}_{>0}, \mid\right)$ for this poset. Related to that, for any positive integer $n$, let $D_{n}$ be the set of positive integers that divide $n$. We put the divisibility relation on $D_{n}$. If $n$ is not a prime power, then this is not a total ordering.
Example 6.5. Let $P$ be the set of set partitions of $[n]$. Given two set partitions $x$ and $y$, we say that $x$ refines $y$ if every block of $x$ is a subset of some block of $y$. For example, $12|34| 5$ refines $125 \mid 34$. We write $x \leq y$ if $x$ refines $y$. Then $(P, \leq)$ is a poset, which we will denote by $\Pi_{n}$. This is not a total ordering when $n \geq 3$.

We can draw posets using Hasse diagrams. Let $(P, \leq)$ be a poset. First, if $x \leq y$ and $x \neq y$, then we will write $x<y$. We say $y$ covers $x$ if there does not exist an element $z$ such that $x<z$ and $z<y$. The Hasse diagram of $P$ is a picture with the elements of $P$ as nodes, and an edge drawn from $x$ up to $y$ whenever $y$ covers $x$.

Example 6.6. Here is the Hasse diagram of $B_{3}$, the poset of subsets of [3]:


We define various kinds of intervals

$$
\begin{aligned}
& {[x, y]=\{z \in P \mid x \leq z \text { and } z \leq y\}} \\
& {[x, y)=\{z \in P \mid x \leq z \text { and } z<y\}} \\
& (x, y]=\{z \in P \mid x<z \text { and } z \leq y\} \\
& (x, y)=\{z \in P \mid x<z \text { and } z<y\}
\end{aligned}
$$

For $x \leq y$, we define the Möbius function $\mu(x, y)$ recursively as follows:

$$
\begin{aligned}
& \mu(x, x)=1 \text { for all } x \in P \\
& \mu(x, y)=-\sum_{z \in[x, y)} \mu(x, z) \text { for all } x<y
\end{aligned}
$$

This is equivalent to the more compact formula

$$
\sum_{z \in[x, y]} \mu(x, z)=\delta_{x, y} .
$$

where $\delta$ is the Kronecker delta.
Lemma 6.7. We have

$$
\sum_{z \in[x, y]} \mu(z, y)=\delta_{x, y} .
$$

Proof. Define a function $\mu^{\prime}(x, y)$ for $x \leq y$ by $\mu^{\prime}(x, x)=1$ for all $x \in P$ and $\mu^{\prime}(x, y)=$ $-\sum_{z \in(x, y]} \mu^{\prime}(z, y)$ for $x<y$. Then we have

$$
\begin{aligned}
\mu(x, y) & =\sum_{w \in[x, y]} \mu(x, w) \delta_{w, y} \\
& =\sum_{w \in[x, y]} \mu(x, w) \sum_{z \in[w, y]} \mu^{\prime}(z, y) \\
& =\sum_{z \in[x, y]} \mu^{\prime}(z, y) \sum_{w \in[x, z]} \mu(x, w) \\
& =\sum_{z \in[x, y]} \mu^{\prime}(z, y) \delta_{x, z}=\mu^{\prime}(x, y) .
\end{aligned}
$$

Hence $\mu$ satisfies the desired relation by definition of $\mu^{\prime}$.
Remark 6.8. The last lemma is more transparent if we introduce the incidence algebra, but I've avoided that to keep notation to a minimum.

Example 6.9. Suppose the following is the Hasse diagram of our poset $P$ :


First, $\mu(a, b)=-\mu(a, a)=-1$ and similarly, $\mu(a, c)=-1=\mu(a, d)$. Say we want to compute $\mu(a, e)$. Then we use the recursive formula:

$$
\mu(a, e)=-(\mu(a, a)+\mu(a, b)+\mu(a, c))=-(1-1-1)=1 .
$$

In the same way, $\mu(a, f)=1$. Now to compute $\mu(a, g)$ :

$$
\begin{aligned}
\mu(a, g) & =-(\mu(a, a)+\mu(a, b)+\mu(a, c)+\mu(a, d)+\mu(a, e)+\mu(a, f)) \\
& =-(1-1-1-1+1+1)=0
\end{aligned}
$$

The purpose of the Möbius function is the inversion formula and its dual form. Before stating it, we first note that $\mu(x, y)$ is always integer-valued. Next, if $a$ is an element in an abelian group (whose operation is written as addition) and $n$ is an integer, we can make sense of $n a$ : it's either $a+\cdots+a(n$ copies of $a)$ if $n \geq 0$, or else it is $-(a+\cdots+a)(-n$ copies of $a$ ) if $n<0$.

Theorem 6.10 (Möbius inversion formula). Let $P$ be a poset and let $f, g$ be functions from $P$ to some abelian group.
(a) We have

$$
g(y)=\sum_{x \leq y} f(x) \text { for all } y \in P
$$

if and only if

$$
f(y)=\sum_{x \leq y} g(x) \mu(x, y) \text { for all } y \in P
$$

(b) (Dual version) We have

$$
g(y)=\sum_{x \geq y} f(x) \text { for all } y \in P
$$

if and only if

$$
f(y)=\sum_{x \geq y} \mu(y, x) g(x) \text { for all } y \in P
$$

Proof. (a) Suppose the first equality holds for all $y \in P$. Then for any $y \in P$, we have

$$
\sum_{x \leq y} g(x) \mu(x, y)=\sum_{x \leq y} \mu(x, y) \sum_{z \leq x} f(z)=\sum_{z \leq y} f(z) \sum_{x \in[z, y]} \mu(x, y)=\sum_{z \leq y} f(z) \delta_{z, y}=f(y)
$$

Now suppose the second equality holds for all $y \in P$. Then for any $y \in P$, we have

$$
\sum_{x \leq y} f(x)=\sum_{x \leq y} \sum_{z \leq x} g(z) \mu(z, x)=\sum_{z \leq y} g(z) \sum_{x \in[z, y]} \mu(z, x)=\sum_{z \leq y} g(z) \delta_{z, y}=g(y) .
$$

(b) is similar.
6.2. Boolean poset and inclusion-exclusion. Inclusion-exclusion is a formula for the size of a union of sets in terms of the sizes of their intersections.

Example 6.11. This is likely familiar when we have 2 or 3 sets. Let's draw Venn diagrams to visualize. For 3 sets we get

which says

$$
|A \cup B|=|A|+|B|-|A \cap B| .
$$

For 3 sets we get

from which we can verify the following formula:

$$
|A \cup B \cup C|=|A|+|B|+|C|-|A \cap B|-|A \cap C|-|B \cap C|+|A \cap B \cap C| .
$$

To see this, the total diagram has 7 regions and we need to make sure that each region get counted exactly once in the right side expression. For example, the elements that belong to $A$ and $B$ but not $C$ appear in $A, B, A \cap B$, and the coefficients are $1+1-1=1$.

The examples above have a generalization to $n$ sets.
Theorem 6.12 (Inclusion-Exclusion). Let $A_{1}, \ldots, A_{n}$ be finite sets. Then

$$
\left|A_{1} \cup \cdots \cup A_{n}\right|=\sum_{j=1}^{n}(-1)^{j-1} \sum_{1 \leq i_{1}<i_{2}<\cdots<i_{j} \leq n}\left|A_{i_{1}} \cap A_{i_{2}} \cap \cdots \cap A_{i_{j}}\right| .
$$

Let's interpret this as a special case of Möbius inversion on the Boolean poset. First, we determine the Möbius function.

Lemma 6.13. If $S \subseteq T$ are subsets of $[n]$, then

$$
\mu(S, T)=(-1)^{|T|-|S|}
$$

Proof. For $S \subseteq T$, it suffices to show that

$$
\delta_{S, T}=\sum_{U \in[S, T]}(-1)^{|U|-|S|} .
$$

If $|S|=s$ and $|T|=t$, then the number of subsets of size $k$ in the interval $[S, T]$ is $\binom{t-s}{k-s}$, so the latter sum becomes

$$
\sum_{k=s}^{t}\binom{t-s}{k-s}(-1)^{k-s}=\sum_{i=0}^{t-s}\binom{t-s}{i}(-1)^{i}=\delta_{s, t}=\delta_{S, T}
$$

Proof of Inclusion-Exclusion. Set $A=A_{1} \cup \cdots \cup A_{n}$. For a subset $S \subseteq[n]$, define

$$
\begin{aligned}
& f(S)=\mid\left\{x \in A \mid x \in A_{i} \text { if and only if } i \in S\right\} \mid, \\
& g(S)=\mid\left\{x \in A \mid x \in A_{i} \text { if } i \in S\right\}\left|=\left|\bigcap_{i \in S} A_{i}\right| .\right.
\end{aligned}
$$

Note that $g(\varnothing)=|A|$ and $f(\varnothing)=0$. By definition,

$$
g(T)=\sum_{S \supseteq T} f(S),
$$

for all $T \subseteq[n]$ and hence by the dual version of Möbius inversion, we have

$$
0=f(\varnothing)=\sum_{S \subseteq[n]}(-1)^{|S|} g(S)=\sum_{S \subseteq[n]}(-1)^{|S|}\left|\bigcap_{i \in S} A_{i}\right| .
$$

Subtract $g(\varnothing)$ and multiply by -1 to get

$$
|A|=g(\varnothing)=\sum_{\substack{S \subseteq[n] \\ S \neq \varnothing}}(-1)^{|S|-1}\left|\bigcap_{i \in S} A_{i}\right|
$$

which is a reformulation of what we wanted to prove.
We use this to address two counting problems.
First, we can think of a permutation of $[n]$ as the same thing as a bijection $f:[n] \rightarrow[n]$ (given the bijection, $f(i)$ is the position in the permutation where $i$ is supposed to appear). A derangement of size $n$ is a permutation such that for all $i, i$ does not appear in position $i$. Equivalently, it is a bijection $f$ such that $f(i) \neq i$ for all $i$.

Theorem 6.14. The number of derangements of size $n$ is

$$
\sum_{i=0}^{n}(-1)^{i} \frac{n!}{i!}
$$

Proof. For a subset $S \subseteq[n]$, define

$$
\begin{aligned}
f(S) & =\mid\{\text { permutations } \sigma \text { of }[n] \text { such that } \sigma(i)=i \text { if and only if } i \in S\} \mid, \\
g(S) & =\mid\{\text { permutations } \sigma \text { of }[n] \text { such that } \sigma(i)=i \text { if } i \in S\} \mid=(n-|S|)!.
\end{aligned}
$$

The number of derangements is $f(\varnothing)$. Also, $g(S)=\sum_{S \subseteq T} f(T)$, so by the dual version of inclusion-exclusion, we have

$$
f(\varnothing)=\sum_{S \subseteq[n]}(-1)^{|S|} g(S)=\sum_{S \subseteq[n]}(-1)^{|S|}(n-|S|)!=\sum_{i=0}^{n}(-1)^{i}\binom{n}{i}(n-i)!=\sum_{i=0}^{n}(-1)^{i} \frac{n!}{i!} .
$$

Remark 6.15. A typical way to phrase the problem of counting derangements is to ask for the percentage of permutations that are derangements (the two quantities differ by dividing/multiplying by $n!$ ). This can be thought as follows: $n$ people put their hat in a bin and they are distributed at random. What is the chance that no one receives their original hat?

We temporarily switch gears to calculus and use Taylor series instead of formal power series. We have a formula for the exponential function

$$
e^{x}=\sum_{i=0}^{\infty} \frac{x^{i}}{i!} .
$$

This converges everywhere, so we can in particular plug in $x=-1$. If we only take the terms up to $i=n$, then we get the number of derangements divided by $n$ !, i.e., the percentage of permutations that are derangements. By Taylor's theorem, the limit of the first $n$ terms of the Taylor expansion of a "well-behaved" (= analytic) function like $e^{x}$ converges to the function. So for $n \rightarrow \infty$, the proportion of permutations that are derangements is $e^{-1} \approx .368$, or roughly $36.8 \%$, which may be surprising depending on your own intuition. Here are the values of this truncated sum for $n=0, \ldots, 10$ up to 6 digits:

$$
1,0, .5, .333333, .375, .366667, .368056, .367857, .367882, .367879, .367879
$$

If you use Taylor's theorem more carefully, you can actually prove that the number of derangements is exactly round $(n!/ e)$ where round means take the closest integer. I'll put the details as an optional homework problem.

We can also use inclusion-exclusion to get an alternating sum formula for Stirling numbers.
Theorem 6.16. For all $n \geq k>0$,

$$
S(n, k)=\frac{1}{k!} \sum_{i=0}^{k}(-1)^{k-i}\binom{k}{i} i^{n}
$$

Proof. As we discussed before, $k!S(n, k)$ counts the number of surjective functions $f:[n] \rightarrow$ [ $k$ ]. So we will count this quantity. For a subset $S \subseteq[k]$, define

$$
\begin{aligned}
f(S) & =\mid\{\text { surjective functions }[n] \rightarrow S\} \mid, \\
g(S) & =\mid\{\text { functions }[n] \rightarrow S\}\left|=|S|^{n} .\right.
\end{aligned}
$$

Then $g(T)=\sum_{S \subseteq T} f(S)$ since by definition every function to $T$ is surjective onto its image (which is some subset of $T$ ). By inclusion-exclusion,

$$
k!S(n, k)=f([k])=\sum_{S \subseteq[k]}(-1)^{k-|S|} g(S)=\sum_{S \subseteq[k]}(-1)^{k-|S|}|S|^{n}=\sum_{i=0}^{k}(-1)^{k-i}\binom{k}{i} i^{n} .
$$

Now divide both sides by $k$ !.
Remark 6.17. We know from the generating function for $S(n, k)$ ( $k$ fixed) that $S(n, k)$ is a linear combination of the powers $1^{n}, 2^{n}, \ldots, k^{n}$. This formula tells us that the coefficient of $i^{n}$ is $\frac{(-1)^{k-i}}{i!(k-i)!}$.
6.3. Divisor poset and classical Möbius inversion. One of the original sources for Möbius inversion is number theory. Here we pick a positive integer $n$ and consider the poset $D_{n}$ of positive integers dividing $n$ (with divisibility as the partial order).

Proposition 6.18. Suppose $x$ divides $y$ (and both divide $n$ ). Then

$$
\mu(x, y)= \begin{cases}0 & \text { if } y / x \text { is divisible by the square of a prime number } \\ (-1)^{k} & \text { if } y / x \text { is a product of } k \text { distinct prime numbers }\end{cases}
$$

If $x=y$, then it falls into the second case with $k=0$.
This is the "classical" Möbius function.
In other words, if any prime divides $y / x$ more than once, then $\mu(x, y)=0$. Otherwise, we count how many different prime numbers divide $y / x ; \mu(x, y)=1$ if that number is even and $\mu(x, y)=-1$ if that number is odd.

Proof. If $x$ divides $y$, let $\mu^{\prime}(x, y)$ be the proposed formula; we need to show that $\sum_{z \in[x, y]} \mu^{\prime}(x, z)=$ $\delta_{x, y}$.

Let $y=p_{1}^{a_{1}} \cdots p_{r}^{a_{r}}$ and $x=p_{1}^{b_{1}} \cdots p_{r}^{b_{r}}$ be prime factorizations. The sum can be rewritten

$$
\sum_{z \in[x, y]} \mu^{\prime}(x, z)=\sum_{\substack{b_{1} \leq e_{1} \leq a_{1} \\ b_{2} \leq e_{2} \leq a_{2}}} \mu^{\prime}\left(x, p_{1}^{e_{1}} \cdots p_{r}^{e_{r}}\right)=\sum_{\substack{b_{1} \leq e_{1} \leq \min \left(b_{1}+1, a_{1}\right) \\ b_{2} \leq e_{2} \leq \min \left(b_{2}+1, a_{2}\right) \\ b_{r} \leq e_{r} \leq a_{r}}} \mu^{\prime}\left(x, p_{1}^{e_{1}} \cdots p_{r}^{e_{r}}\right) .
$$

The second equality holds because if any $e_{i} \geq b_{i}+2$ then $p_{1}^{e_{1}} \cdots p_{r}^{e_{r}} / x$ is divisible by $p_{i}^{2}$. The last sum is a sum over all products of subsets of the primes $Q=\left\{p_{i} \mid b_{i}<a_{i}\right\}$, so we get

$$
\sum_{S \subseteq Q} \mu^{\prime}\left(x, x \cdot \prod_{p \in S} p\right)=\sum_{S \subseteq Q}(-1)^{|S|}=\sum_{k=0}^{|Q|}(-1)^{k}\binom{|Q|}{k}=\delta_{0,|Q|},
$$

and finally, $|Q|=0$ if and only if $y=x$.
Let's consider an enumerative application with a number-theoretic flavor.
Let $A$ be an alphabet of size $k$. We want to count the number of words of length $n$ in $A$ up to cyclic symmetry. This means that two words are considered the same if one is a cyclic shift of another. For example, for words of length 4 , the following 4 words are all the same:

$$
a_{1} a_{2} a_{3} a_{4}, \quad a_{2} a_{3} a_{4} a_{1}, \quad a_{3} a_{4} a_{1} a_{2}, \quad a_{4} a_{1} a_{2} a_{3}
$$

We call these necklaces: the elements of $A$ might be different beads we can put on the necklace, but we would consider two to be the same if we can rotate one to get the other. Naively, we might say that the number of necklaces of length $n$ is $k^{n} / n$ since we have $n$ rotations for each necklace. However, there is a problem: the $n$ rotations might not all be the same. For example there are only 2 different rotations of 0101.

We have to separate necklaces into different groups based on their period: this is the smallest $d$ such that rotating $d$ times gives the same thing. So for $n=4$, we can have necklaces of periods 1,2 , or 4 , examples being $0000,0101,0001$. There aren't any of period 3: the period must divide the length (this can be translated into a group theory fact about the order of a subgroup dividing the order of a group). Let $\omega(d, k)$ be the number of words of period $d$. Hence for necklaces of length 4 , we get the following formula:

$$
\omega(1, k)+\frac{1}{2} \omega(2, k)+\frac{1}{4} \omega(4, k) .
$$

For general $n$, we would have

$$
\mid \text { necklaces of length } n \left\lvert\,=\sum_{d \mid n} \frac{1}{d} \omega(d, k)\right.
$$

We'll see how to get a different formula once we develop group actions in the next section, but let's connect $\omega(d, k)$ with the classical Möbius function.

Theorem 6.19 (Witt's formula). For any positive integer d, we have

$$
\omega(d, k)=\sum_{e \mid d} \mu(e, d) k^{e}
$$

where the sum is over all positive integers e that divide $d$.
Proof. Let $g(e)$ be the number of words of length $d$ whose period divides $e$ and let $f(e)$ be the number of words of length $d$ whose period is exactly $e$. Note that a word that has period dividing $e$ is determined by its first $e$ letters, which can be anything, so $g(e)=k^{e}$.

Also, $g(e)=\sum_{e^{\prime} \mid e} f\left(e^{\prime}\right)$, so by Möbius inversion, we get $f(d)=\sum_{e \mid d} \mu(e, d) k^{e}$.

Example 6.20. Let's apply this to the case $n=4$. Then we have the following formulas:

$$
\begin{aligned}
& \omega(1, k)=\mu(1,1) k=k \\
& \omega(2, k)=\mu(1,2) k+\mu(2,2) k^{2}=-k+k^{2} \\
& \omega(4, k)=\mu(1,4) k+\mu(2,4) k^{2}+\mu(4,4) k^{4}=0-k^{2}+k^{4} .
\end{aligned}
$$

So the number of necklaces of length 4 is $k+\frac{k^{2}-k}{2}+\frac{k^{4}-k^{2}}{4}=\left(k^{4}+k^{2}+2 k\right) / 4$.
Let's also do $n=6$ :

$$
\begin{aligned}
& \omega(3, k)=\mu(1,3) k+\mu(3,3) k^{3}=-k+k^{3} \\
& \omega(6, k)=\mu(1,6) k+\mu(2,6) k^{2}+\mu(3,6) k^{3}+\mu(6,6) k^{6}=k-k^{2}-k^{3}+k^{6} .
\end{aligned}
$$

So the number of necklaces of length 6 is $k+\frac{k^{2}-k}{2}+\frac{k^{3}-k}{3}+\frac{k^{6}-k^{3}-k^{2}+k}{6}=\left(k^{6}+k^{3}+2 k^{2}+\right.$ $2 k) / 6$.
Example 6.21. For this example, $i$ is a (complex) square root of -1 . Consider the polynomial $x^{d}-1$. Let $\omega_{d}=e^{2 \pi i / d}$. By Euler's formula $e^{2 \pi i}=1$, the roots of $x^{d}-1$ are the complex numbers $\omega_{d}^{j}$ for $j=0,1, \ldots, d-1$. In particular, if $d^{\prime}$ divides $d$, then $x^{d^{\prime}}-1$ divides $x^{d}-1$. The $d$ th cyclotomic polynomial $\Phi_{d}(x)$ is the result of dividing $x^{d}-1$ by all $\left(x-\omega_{d}^{j}\right)$ where $\operatorname{gcd}(j, d) \neq 1$. Then $\Phi_{d}(x)$ and $x^{d^{\prime}}-1$ have no common factors when $d^{\prime} \mid d$ and $d^{\prime}<d$ and we get the formula $x^{d}-1=\prod_{d^{\prime} \mid d} \Phi_{d^{\prime}}(x)$.

Then define $f(d)=\Phi_{d}(x)$ and $g(d)=x^{d}-1$; then $f, g$ take values in the abelian group of nonzero rational functions in $x$ with the operation of multiplication. So we can use Möbius inversion to conclude that

$$
\Phi_{d}(x)=\prod_{d^{\prime} \mid d}\left(x^{d^{\prime}}-1\right)^{\mu\left(d^{\prime}, d\right)}
$$

For example,

$$
\left(x-\omega_{6}\right)\left(x-\omega_{6}^{5}\right)=\Phi_{6}(x)=\frac{\left(x^{6}-1\right)(x-1)}{\left(x^{3}-1\right)\left(x^{2}-1\right)}=x^{2}-x+1 .
$$

## 7. Group actions

7.1. Terminology. Hopefully you have seen group actions in your abstract algebra course. We review the key definitions and facts now.

Recall that a group is a set $G$ with an associative binary operation such that there is an identity element and every element has an inverse, and that a subgroup is a subset $H \subseteq G$ which contains the identity and is closed under the binary operation. We will only deal with finite groups, so this will be a running assumption. For the moment, we will use concatenation to denote the binary operation.

The key example for us are permutation groups. Given a set $X$, the set of invertible functions $f: X \rightarrow X$ is a group under composition, which we call $\mathfrak{S}_{X}$. If $X=[n]$, we write $\mathfrak{S}_{n}$ in place of $\mathfrak{S}_{[n]}$, and call it the symmetric group. In our examples, we consider subgroups of $\mathfrak{S}_{X}$ where $X$ is a finite set.

If $G$ is a group and $X$ is a set, then an action of $G$ on $X$ is a function $\varphi: G \times X \rightarrow X$ (we write $g x$ or $g \cdot x$ instead of $\varphi(g, x)$ ) such that
(1) $1 \cdot x=x$ for all $x \in X$ where $1 \in G$ is the identity,
(2) $g \cdot(h \cdot x)=(g h) \cdot x$ for all $g, h \in G$ and $x \in X$.

Again, we will always assume that the set $X$ is finite.
In this way, each $g \in G$ gives a function $\varphi_{g}: X \rightarrow X$ via $\varphi_{g}(x)=g \cdot x$.
Proposition 7.1. $g \mapsto \varphi_{g}$ is a group homomorphism $G \rightarrow \mathfrak{S}_{X}$. Conversely, every group homomorphism $G \rightarrow \mathfrak{S}_{X}$ is of this form.
Proof. First, $\varphi_{1}$ is the identity function because of axiom (1). Second, for any $x \in X$, and $g, h \in G$, we have $\varphi_{g}\left(\varphi_{h}(x)\right)=g \cdot(h \cdot x)=(g h) \cdot x=\varphi_{g h}(x)$ by axiom (2). So $\varphi_{g} \varphi_{h}=\varphi_{g h}$. In particular, this implies that $\varphi_{g}$ is invertible with inverse $\varphi_{g^{-1}}$, so $\varphi_{g} \in \mathfrak{S}_{X}$. Hence $g \mapsto \varphi_{g}$ is a homomorphism.

Conversely, suppose we are given a homomorphism $\psi: G \rightarrow \mathfrak{S}_{X}$. For $g \in G$ and $x \in X$, define $g \cdot x=\psi(g)(x)$. Since $\psi(1)$ is the identity function, we have $1 \cdot x=x$ for all $x$. Second, for $g, h \in G$, since $\psi(g h)=\psi(g) \psi(h)$, we have $(g h) \cdot x=\psi(g h)(x)=\psi(g)(\psi(h)(x))=g \cdot(h \cdot x)$. Hence $\psi$ comes from a group action.

Hence, given an action on a set, it makes sense to ask about the number of cycles of $g$, interpreted as a permutation. We denote it by $c_{X}(g)$ (this depends on the set and the action).

Given $x \in X$, define the orbit of $x$ by

$$
G \cdot x=\{g \cdot x \mid g \in G\}
$$

This is the subset of $X$ consisting of all elements which can be "reached" by multiplying $x$ by elements of $G$. Being in the same orbit is an equivalence relation, so any two distinct orbits are disjoint.

The set of orbits is denoted $X / G$.
We also define the stabilizer of $x$ by

$$
G_{x}=\{g \in G \mid g \cdot x=x\},
$$

which is the subset of $G$ of elements that act on $x$ as the identity. This is a subgroup.
Lemma 7.2 (Orbit-stabilizer formula). We have $|G| /\left|G_{x}\right|=|G \cdot x|$.
Proof. Let $\left\{x_{1}, \ldots, x_{r}\right\}$ be the elements of $G \cdot x$. Then there exist $g_{i} \in G$ such that $x_{i}=g_{i} \cdot x$ by definition. If $g \in G$, then $g \cdot x=x_{i}$ for some $i$, and hence $g_{i}^{-1} g \in G_{x}$. In particular, every element of $G$ can be written in the form $g_{i} h$ for $h \in G_{x}$. This is also unique: if $g_{i} h=g_{j} h^{\prime}$ are two different ways, then $g_{j}^{-1} g_{i}=h^{\prime} h^{-1} \in G_{x}$ which means that $g_{j}^{-1} g_{i} x=x$, so $x_{j}=g_{j} \cdot x=g_{i} \cdot x=x_{i}$ and so $i=j$. Next, $h=h^{\prime}$ as well by multiplying $g_{i} h=g_{i} h^{\prime}$ on the left by $g_{i}^{-1}$. This means that $|G|=r\left|G_{x}\right|$.
Example 7.3. The standard example of a group action is $G=\mathfrak{S}_{n}$ and $X=[n]$ with the natural action $\sigma \cdot i=\sigma(i)$. In that case, the orbit of any element of $X$ is all of $X$. The stabilizer of $i$ is the set of permutations $\sigma$ that satisfy $\sigma(i)=i$, so there are ( $n-1$ )! many of them. This is consistent with the orbit-stabilizer formula.

Example 7.4. For a more geometric example, we can consider the dihedral group $D_{4}$ of order 8 with its action on the set of vertices of a square. Again, the orbit of any vertex is the whole set. The orbit-stabilizer formula then tells us that the stabilizer of any vertex has size 2. In fact it consists of the identity element and the reflection with respect to the diagonal that contains that vertex.

If we label the vertices $1,2,3,4$ in clockwise order, then the non-trivial rotations are the permutations (in cycle notation) (1234), (13)(24), and (1432).

The reflection across the diagonal through 1 and 3 is $(1)(3)(24)$ and for the diagonal through 2 and 4 , the reflection is $(13)(2)(4)$.

There are 2 more elements, corresponding to reflection across the lines through opposite sides. They give the permutations (12)(34) and (14)(23).
Example 7.5. Pick $n \geq 2$. Consider $G=\mathfrak{S}_{n}$ and $X=[n] \times[n]$ with $\sigma \cdot(i, j)=(\sigma(i), \sigma(j))$. Now you can check that there are 2 orbits of pairs $(i, j)$, one orbit consisting of pairs $(i, i)$ and one consisting of pairs $(i, j)$ with $i \neq j$. The first orbit behaves just like the action of $\mathfrak{S}_{n}$ on $[n]$. For the second orbit, the stabilizer of $(i, j)$ is the set of permutations such that $\sigma(i)=i$ and $\sigma(j)=j$, so there are $(n-2)$ ! many of them.

Now consider a variation where $X$ is the set of 2 -element subsets of $[n]$. Then there is 1 orbit. This time, the stabilizer of $\{i, j\}$ is the set of permutations such that $\sigma(\{i, j\})=\{i, j\}$ which happens in two cases:

- $\sigma(i)=i$ and $\sigma(j)=j$, or
- $\sigma(i)=j$ and $\sigma(j)=i$,
so there are $2(n-2)$ ! many of them.
Example 7.6. Another important example comes when $X$ has some kind of structure and we take $G$ to be the subgroup of $\mathfrak{S}_{n}$ which preserves this structure. For example, if $X$ is the set of vertices of a simple graph $\Gamma$, then we take $G$ to be the set of permutations such that $\{i, j\}$ is an edge if and only if $\{\sigma(i), \sigma(j)\}$ is an edge for all $i, j \in X$. A quick check shows that $G$ is a subgroup. In this case, $G$ is called the automorphism group of $\Gamma$, and denoted $\operatorname{Aut}(\Gamma)$. For a specific example, think of the edges of a square as giving a simple graph. In that case, $\operatorname{Aut}(\Gamma)=D_{4}$.

For another example, take $G=\mathbf{G} \mathbf{L}_{n}\left(\mathbf{F}_{q}\right)$ and $X=\mathbf{F}_{q}^{n}$.
7.2. Burnside's lemma. Let $G$ act on $X$. For $g \in G$, define the set of fixed points of $g$ to be

$$
X^{g}=\{x \in X \mid g \cdot x=x\} .
$$

Theorem 7.7 (Burnside's lemma). The number of $G$-orbits on $X$ is

$$
|X / G|=\frac{1}{|G|} \sum_{g \in G}\left|X^{g}\right|
$$

Proof. Define

$$
S=\{(g, x) \in G \times X \mid g \cdot x=x\} .
$$

By definition, for given $g \in G$, we have $(g, x) \in S$ if and only if $x \in X^{g}$, so $|S|=\sum_{g \in G}\left|X^{g}\right|$. On the other hand, for given $x \in X$, we have $(g, x) \in S$ if and only if $g \in G_{x}$, so $|S|=$ $\sum_{x \in X}\left|G_{x}\right|$. Hence we get

$$
\frac{1}{|G|} \sum_{g \in G}\left|X^{g}\right|=\frac{|S|}{|G|}=\sum_{x \in X} \frac{\left|G_{x}\right|}{|G|}=\sum_{x \in X} \frac{1}{|G \cdot x|}
$$

The last sum is the number of orbits: if an orbit has $r$ elements, then each element of it contributes $1 / r$ to the sum, and hence the total contribution is 1 for each orbit.

Example 7.8. Consider $G=\mathfrak{S}_{n}$ and $X=[n]$ with the standard action: $\sigma \cdot x=\sigma(x)$. Then $\left|X^{\sigma}\right|$ is the number of cycles of size 1 of $\sigma$, so we can interpret $\frac{1}{n!} \sum_{\sigma \in \mathfrak{G}_{n}}\left|X^{\sigma}\right|$ as the average number of cycles of size 1 (or average number of fixed points) of a permutation. The last result tells us that this is the number of orbits of $\mathfrak{S}_{n}$ on $[n]$ which is 1 .

We will consider the following situation. Let $Y$ be a finite set and let $Y^{X}$ be the set of functions $X \rightarrow Y$. We introduce a $G$-action on $Y^{X}$ via $g \cdot f=f \circ \varphi_{g^{-1}}$ for $g \in G$ and $f \in Y^{X}$. This is indeed an action since $\varphi_{1^{-1}}$ is the identity function, and

$$
g \cdot(h \cdot f)=\left(f \circ \varphi_{h^{-1}}\right) \circ \varphi_{g^{-1}}=f \circ \varphi_{h^{-1} g^{-1}}=f \circ \varphi_{(g h)^{-1}}=(g h) \cdot f .
$$

Alternatively, for $x \in X$, we have $(g \cdot f)(x)=f\left(g^{-1} \cdot x\right)$. We think of functions $X \rightarrow Y$ as labelings of the elements of $X$ by elements of $Y$ (which we might think of as colors) and $G$-orbits of $Y^{X}$ as equivalence classes of labelings. As before, for $g \in G$, let $c_{X}(g)$ denote the number of cycles of $\varphi_{g}$ acting on $X$.

Theorem 7.9. The number of $G$-orbits on $Y^{X}$ is

$$
\left|Y^{X} / G\right|=\frac{1}{|G|} \sum_{g \in G}|Y|^{c_{X}(g)}
$$

Proof. Via Burnside's lemma, it suffices to prove that $|Y|^{c_{X}(g)}=\left|\left(Y^{X}\right)^{g}\right|$. By definition, if $f \in\left(Y^{X}\right)^{g}$, then $f\left(g^{-1} x\right)=f(x)$ for all $x \in X$. This translates into the condition that $f$ is constant on the cycles of $g$, i.e., $\left|\left(Y^{X}\right)^{g}\right|$ is the number of functions constant on the cycles of $g$. To count such functions, we can pick the values on each cycle independently, so there are $|Y|^{c_{X}(g)}$ many functions fixed by $g$.

Example 7.10. We can use this to revisit the problem of counting necklaces. Consider necklaces of length $n$ in an alphabet of size $k$. In our new setup, let $X=\mathbf{Z} / n$. Let also $G=\mathbf{Z} / n$ and let the action be given by addition, i.e., $i \cdot j=i+j$. If the elements of $X$ are placed in a circle, we can think of the action of $i$ as cyclic rotation by $i$ places. Let $Y$ be our alphabet. Then a function $X \rightarrow Y$ is a word of length $n$, and a $G$-orbit represents a word up to cyclic shift, i.e., a necklace. So necklaces are in bijection with $G$-orbits of $Y^{X}$.

Consider the case $n=4$. The elements of $G$ are the powers of the permutation (0123) (written in cycle notation), so more specifically they are (0123), (02)(13), (0321), (0)(1)(2)(3), with $1,2,1,4$ cycles respectively. So the number of necklaces is $\frac{1}{4}\left(k^{4}+k^{2}+2 k\right)$, agreeing with our previous result using Möbius inversion.

For general $n$, we need to compute the number of cycles of $(01 \cdots n-1)^{i}$ for $i=0,1, \ldots, n-$ 1. Recall (or prove as an exercise) that the order of this element is $n / \operatorname{gcd}(n, i)$ and its cycles all have equal length (we can interpret cycles as cosets of the subgroup generated by $\left.(01 \cdots n-1)^{i}\right)$, hence there are $\operatorname{gcd}(n, i)$ many cycles, so the answer is

$$
\frac{1}{n} \sum_{i=1}^{n} k^{\operatorname{gcd}(n, i)}
$$

Example 7.11. As a variation, we might consider two necklaces the same if they are reflections of one another. Then we have the same setup, but we replace the cyclic group $\mathbf{Z} / n$ with the dihedral group $D_{n}$, i.e., the symmetries of a regular $n$-gon. We'll just do a single example with $n=4$. In that case, $\left|D_{4}\right|=8$ and includes the cyclic group as the subgroup of rotations, so their cycle lengths are $1,1,2,4$ from before.

Thinking of $0,1,2,3$ as the vertices of a square in clockwise order, we can reflect across the diagonals to get $(0)(13)(2)$ and $(02)(1)(3)$ or reflect across the vertical or horizontal axis to get $(01)(23)$ and $(03)(12)$. Then the number of equivalence classes is $\frac{1}{8}\left(2 k+3 k^{2}+2 k^{3}+k^{4}\right)$.
7.3. Redfield-Pólya theory. Continuing from the previous section, we may want more detailed information: rather than ask how many colorings of $X$ there are up to $G$, we can ask how many use each color a specific number of times. To do this, it is convenient to think of $Y$ as a set of variables now. Then we define the weight of a function $f: X \rightarrow Y$ to be

$$
\mathrm{wt}(f)=\prod_{x \in X} f(x)
$$

For example, if 2 things are blue and 3 are red, the weight is $B^{2} R^{3}$.
Note that if two functions are in the same $G$-orbit, they necessarily have the same weight since we are keeping the same labels, just redistributing where they go. Hence it makes sense to define the weight of an orbit $\operatorname{wt}(G \cdot f)$ to be the weight $\mathrm{wt}(f)$ of any element in the orbit.

Let $n=|X|$. Given $g \in G$, let $c_{X, i}(g)$ be the number of cycles of length $i$, so that $c_{X}(g)=\sum_{i=1}^{n} c_{X, i}(g)$. We introduce new variables $t_{1}, \ldots, t_{n}$ and define the cycle indicator of $G$ acting on $X$ to be

$$
Z_{X}\left(G ; t_{1}, \ldots, t_{n}\right)=\frac{1}{|G|} \sum_{g \in G} t_{1}^{c_{X, 1}(g)} \cdots t_{n}^{c_{X, n}(g)}
$$

Theorem 7.12 (Redfield-Pólya). The sum of the weights of each orbit of $Y^{X}$ is given by

$$
\sum_{\alpha \in Y^{x} / G} \operatorname{wt}(\alpha)=Z_{X}\left(G ; \sum_{y \in Y} y, \sum_{y \in Y} y^{2}, \ldots, \sum_{y \in Y} y^{n}\right) .
$$

Before proving this, note that if we set $y=1$ for each $y \in Y$, then the left side is just the number of orbits and the right side is the sum $\frac{1}{|G|} \sum_{g \in G}|Y|^{c_{X}(g)}$, so this generalizes our previous result. In fact, we can follow a similar strategy to prove it.

Proof. Define

$$
S=\left\{(g, f) \in G \times Y^{X} \mid g \cdot f=f\right\}
$$

and define $\mathrm{wt}(S)=\sum_{(g, f) \in S} \mathrm{wt}(f)$. As before, we can compute this as a sum over $Y^{X}$ or over $G$. Summing over $Y^{X}$ gives

$$
\mathrm{wt}(S)=\sum_{f \in Y^{X}}\left|G_{f}\right| \mathrm{wt}(f)
$$

Now divide by $|G|$ and use the orbit-stabilizer formula to get:

$$
\frac{\mathrm{wt}(S)}{|G|}=\sum_{f \in Y^{X}} \frac{\left|G_{f}\right|}{|G|} \mathrm{wt}(f)=\sum_{f \in Y^{X}} \frac{\mathrm{wt}(f)}{|G \cdot f|}=\sum_{\alpha \in Y^{X} / G} \mathrm{wt}(\alpha)
$$

where the last sum is over all orbits, and the last equality follows as before: each $f$ contributes $\mathrm{wt}(G \cdot f) /|G \cdot f|$ and there are $|G \cdot f|$ many elements in $G \cdot f$, so the total contribution of the elements in this orbit is $\mathrm{wt}(G \cdot f)$.

If instead we sum over $G$, we get

$$
\mathrm{wt}(S)=\sum_{g \in G} \sum_{f \in\left(Y^{x}\right)^{g}} \mathrm{wt}(f) .
$$

Consider the inner sum. Every $f \in\left(Y^{X}\right)^{g}$ is constant on the cycles of $g$, call the cycles $C_{1}, \ldots, C_{r}$. So to specify $f$, we can pick a value $f\left(C_{i}\right)$ for each $i$ independently. Its weight
is then $f\left(C_{1}\right)^{\left|C_{1}\right|} \cdots f\left(C_{r}\right)^{\left|C_{r}\right|}$. Hence if we sum over all choices of $f$, we get

$$
\sum_{f \in\left(Y^{X}\right)^{g}} \mathrm{wt}(f)=\prod_{i=1}^{r}\left(\sum_{y \in Y} y^{\left|C_{i}\right|}\right)=\prod_{j=1}^{n}\left(\sum_{y \in Y} y^{j}\right)^{c_{X, j}(g)} .
$$

Comparing that to the definition of $Z_{X}\left(G ; t_{1}, \ldots, t_{n}\right)$, we see that

$$
\frac{\mathrm{wt}(S)}{|G|}=\frac{1}{|G|} \sum_{g \in G} \sum_{f \in\left(Y^{X}\right)^{g}} \mathrm{wt}(f)=Z_{X}\left(G ; \sum_{y \in Y} y, \ldots, \sum_{y \in Y} y^{n}\right) .
$$

By what we've shown, the left hand side is also $\sum_{\alpha \in Y^{X} / G} \mathrm{wt}(\alpha)$.
Example 7.13. Consider again necklaces of length 4. Let $y_{1}, \ldots, y_{k}$ be variables representing possible colors. The elements of $G$ are (0123), (02)(13), (0321), (0)(1)(2)(3), so the cycle indicator is

$$
Z_{X}\left(G ; t_{1}, t_{2}, t_{3}, t_{4}\right)=\frac{1}{4}\left(2 t_{4}+t_{2}^{2}+t_{1}^{4}\right) .
$$

Doing the substitution $t_{d} \mapsto \sum_{i=1}^{k} y_{i}^{d}$, we get

$$
\frac{1}{4}\left(2 \sum_{i} y_{i}^{4}+\left(\sum_{i} y_{i}^{2}\right)^{2}+\left(\sum_{i} y_{i}\right)^{4}\right)
$$

This is symmetric in the $y_{i}$, so the only relevant coefficients are those of $y_{1}^{4}, y_{1}^{3} y_{2}, y_{1}^{2} y_{2}^{2}, y_{1}^{2} y_{2} y_{3}$, $y_{1} y_{2} y_{3} y_{4}$, which correspond to the integer partitions of size 4 .

- The coefficient of $y_{1}^{4}$ is 1 , which tells us there is only 1 necklace of length 4 where all colors are 1 (of course) and since there are $k$ different ways to choose the subscript 1 , there are $k$ necklaces where all colors are the same.
- The coefficient of $y_{1}^{3} y_{2}$ is 1 , which tells us there is only 1 necklace that uses 1 exactly 3 times and 2 exactly once. In general, if we want to know many use some color 3 times and a different color once, then we consider the coefficients $y_{i}^{3} y_{j}$ over all choices of $i \neq j$, of which there are $k(k-1)$.
- The coefficient of $y_{1}^{2} y_{2}^{2}$ is 2 , which tells us there are 2 necklaces that use 1 exactly twice and 2 exactly twice. In general, there are $\binom{k}{2}$ many monomials $y_{i}^{2} y_{j}^{2}$ with $i \neq j$, so there are $2\binom{k}{2}$ necklaces that use exactly 2 colors, each used twice.
- The coefficient of $y_{1}^{2} y_{2} y_{3}$ is 3 , which tells us there are 3 necklaces that use 1 exactly twice, 2 exactly once, and 3 exactly once. In general there are $k\binom{k-1}{2}$ many monomials $y_{i}^{2} y_{j} y_{k}$ with $i, j, k$ distinct, so there are a total of $3 k\binom{k-1}{2}$ many necklaces that use 3 different colors, one of which is used twice.
- The coefficient of $y_{1} y_{2} y_{3} y_{4}$ is 6 , which tells us there are 6 necklaces of length 4 where all colors are different and use $1,2,3,4$. Finally, there are $\binom{k}{4}$ many ways to choose 4 different colors, so we see that there are $6\binom{k}{4}$ many necklaces where all colors are different.

Example 7.14. Now let's consider the case of necklaces of size 4 up to reflection, so that $G$ in the previous example is now the dihedral group $D_{4}$. The new elements that we get are $(0)(13)(2),(0)(12)(3),(01)(23)$, and $(03)(12)$, so the cycle indicator is

$$
Z_{X}\left(D_{4} ; t_{1}, t_{2}, t_{3}, t_{4}\right)=\frac{1}{8}\left(2 t_{4}+3 t_{2}^{2}+2 t_{1}^{2} t_{2}+t_{1}^{4}\right)
$$

Doing the substitution $t_{d} \mapsto \sum_{i=1}^{k} y_{i}^{d}$ gives

$$
\frac{1}{8}\left(2 \sum_{i} y_{i}^{4}+3\left(\sum_{i} y_{i}^{2}\right)^{2}+2\left(\sum_{i} y_{i}\right)^{2}\left(\sum_{i} y_{i}^{2}\right)+\left(\sum_{i} y_{i}\right)^{4}\right)
$$

Again let's compute the coefficients of the different types of monomials.

- The coefficient of $y_{1}^{4}$ is 1 .
- The coefficient of $y_{1}^{3} y_{2}$ is 1 .
- The coefficient of $y_{1}^{2} y_{2}^{2}$ is 2 .
- The coefficient of $y_{1}^{2} y_{2} y_{3}$ is 2 .
- The coefficient of $y_{1} y_{2} y_{3} y_{4}$ is 3 .

We see that in all cases, the coefficient is at most the coefficient in the previous case (which it should be, since we're only making more things equivalent).

There is a whole developed theory for working with multivariate polynomials which are symmetric in their variables, known otherwise as symmetric polynomials. There isn't enough time to develop them adequately in this course, but you can see my notes for Math 202B if you're interested in seeing the basic development.
7.4. Proving congruences. In this last section, we will focus on proving congruences modulo a prime $p$. Recall that for any integers $a, b, n$, we write $a \equiv b(\bmod n)$ to mean that $a-b$ is divisible by $n$. Everything will be a consequence of choosing good examples in the following lemma.
Lemma 7.15. Let $p$ be a prime integer. Let $G$ be a group of order $p$ acting on a finite set $X$. Let $g \in G$ be a generator of $G$. Then

$$
|X| \equiv\left|X^{g}\right| \quad(\bmod p)
$$

Proof. From the orbit-stabilizer formula, every orbit of $G$ on $X$ has size dividing $|G|=p$, and hence is either $p$ or 1 . Since the orbits form a set partition of $X$, this implies that $|X|$ is the same as the number of orbits of size 1 modulo $p$. An orbit of size 1 is an element $x \in X$ such that $h x=x$ for all $h \in G$, but $h$ is a power of $g$, so it is the same to know that $g x=x$. In particular, $X^{g}$ is the union of the orbits of size 1.
Theorem 7.16 (Fermat's little theorem). Pick $a, p \in \mathbf{Z}$ with $p$ a prime. Then

$$
a^{p} \equiv a \quad(\bmod p)
$$

Proof. If $a \equiv b(\bmod p)$, then $a^{p} \equiv b^{p}(\bmod p)$, so we only have to prove this for one representative of each congruence class, i.e., we can assume that $1 \leq a \leq p$. Let $X$ be the set of functions $\mathbf{Z} / p \rightarrow[a]$. So $|X|=a^{p}$. Let $g \in \mathfrak{S}_{X}$ be given by $(g f)(i)=f(i+1)$. Then $g$ generates a cyclic group of order $p$. If $g f=f$, then $f$ is a constant function, and there are $a$ of those, so $\left|X^{g}\right|=a$.
Theorem 7.17 (Wilson). If $p$ is a prime, then

$$
(p-1)!\equiv-1 \quad(\bmod p)
$$

Proof. Let $X$ be the set of ways of placing the elements of $\mathbf{Z} / p$ around a circular table (each element used exactly once). In particular, we consider two placements the same if they differ by rotation. So $|X|=(p-1)$ !. Let $g \in \mathfrak{S}_{X}$ have the effect of adding 1 to each entry of a placement. Then $g$ generates a cyclic group $G$ of order $p$.

Suppose that $x \in X^{g}$. List the elements in clockwise order $a_{0}, \ldots, a_{p-1}$ so that $a_{0}=0$ (this is all considered up to cyclic shift). There is a unique $i$ such that $a_{i}=1$ and $1 \leq i \leq$ $p-1$. Applying $g$ gives the sequence $a_{0}+1, \ldots, a_{p-1}+1$, and the first element is $a_{i}$, so $a_{2 i}=a_{i}+1=2$. Iterating this, we deduce that $a_{i j}=j$ for all $j$ (thinking of the indices modulo $p$ ). Since $i \in \mathbf{Z} / p$ has order $p$, all of the elements $\{i, 2 i, \ldots,(p-1) i\}$ are distinct elements modulo $p$, so we see that knowing the distance $i$ between 0 and 1 determines the whole placement. Since there are $p-1$ possibilities for this distance, $\left|X^{g}\right|=p-1$.
Lemma 7.18. Let $p$ be a prime and $n \geq p$. Then

$$
\binom{n}{k} \equiv\binom{n-p}{k-p}+\binom{n-p}{k} \quad(\bmod p) .
$$

Proof. Let $X$ be the set of $k$-element subsets of $[n$ ]. Let $\sigma$ be the permutation which is the $p$-cycle $(12 \cdots p)$. Define $g \in \mathfrak{S}_{X}$ as follows: if $S=\left\{s_{1}, \ldots, s_{k}\right\}$, then $g(S)=$ $\left\{\sigma\left(s_{1}\right), \ldots, \sigma\left(s_{k}\right)\right\}$. Then $g$ generates a cyclic group of order $p$. Now we describe $X^{g}$; suppose $S \in X^{g}$. If $S \cap[p]=\varnothing$, then $S \in X^{g}$ since $g$ does nothing to its elements. There are $\binom{n-p}{k}$ many such subsets. Otherwise, suppose $S$ contains some $i \leq p$. If $g(S)=S$, then $S$ must also contain $\sigma(i)$. Iterating that argument, it must also contain $\sigma^{2}(i)$, and actually all $\sigma^{k}(i)$, i.e., $[p] \subseteq S$. The number of such subsets is $\binom{n-p}{k-p}$ since we can freely choose the $k-p$ elements of $S \backslash[p]$ from $\{p+1, \ldots, n\}$.


[^0]:    ${ }^{1}$ We haven't said what kind of numbers we're dealing with here; for simplicity, you may assume they are complex numbers, but they can be taken from any field, even finite fields $\mathbf{Z} / p$, though the quadratic formula makes no sense for $\mathbf{Z} / 2$ of course.

[^1]:    ${ }^{2}$ You could use any field with basically no changes below, and in fact you could even use a commutative ring with some appropriate adjustments to the results below.

[^2]:    ${ }^{3}$ If not, note that multiplication by any nonzero $a \in \mathbf{Z} / p$ is injective because $a b=0$ implies that $p$ divides $a b$; since it does not divide $a$, it has to divide $b$. Since $\mathbf{Z} / p$ is finite, injective implies surjective, so there is a $b$ such that $a b=1$. Or if you know Fermat's little theorem, you can just take $b=a^{p-2}$.

