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1. Induction

Induction is a proof technique that I expect that you’ve seen and grown familiar with in a course on introduction to proofs. We will review it here and discuss some different ways to use it.

1.1. Weak induction. Induction is used when we have a sequence of statements \( P(0), P(1), P(2), \ldots \) labeled by non-negative integers that we’d like to prove. For example, \( P(n) \) could be the statement: \( \sum_{i=0}^{n} i = n(n+1)/2 \). In order to prove that all of the statements \( P(n) \) are true using induction, we need to do 2 things:

- Prove that \( P(0) \) is true.
- Assuming that \( P(n) \) is true, use it to prove that \( P(n+1) \) is true.

Let’s see how that works for our example:

- \( P(0) \) is the statement \( \sum_{i=0}^{0} i = 0 \cdot 1/2 \). Both sides are 0, so the equality is valid.
- Now we assume that \( P(n) \) is true, i.e., that \( \sum_{i=0}^{n} i = n(n+1)/2 \). Now we want to prove that \( \sum_{i=0}^{n+1} i = (n+1)(n+2)/2 \). Add \( n+1 \) to both sides of the original identity. Then the left side becomes \( \sum_{i=0}^{n+1} i \) and the right side becomes \( n(n+1)/2 + n + 1 = (n+1)(n/2+1) = (n+1)(n+2)/2 \), so the new identity we want is valid.

Since we’ve completed the two required steps, we have proven that the summation identity holds for all \( n \).

**Remark 1.1.** We have labeled the statements starting from 0, but sometimes it’s more natural to start counting from 1 instead, or even some larger integer. The same reasoning as above will apply for these variations. The first step “Prove that \( P(0) \) is true” is then replaced by “Prove that \( P(1) \) is true” or wherever the start of your indexing occurs. \( \square \)

**Theorem 1.2.** There are \( 2^n \) subsets of a set of size \( n \).

For example, if \( S = \{1, \ast, U\} \), then there are \( 2^3 = 8 \) subsets, and we can list them: \( \emptyset, \{1\}, \{\ast\}, \{U\}, \{1, \ast\}, \{1, U\}, \{U, \ast\}, \{1, \ast, U\} \).

**Proof.** Let \( P(n) \) be the statement that any set of size \( n \) has exactly \( 2^n \) subsets.

We check \( P(0) \) directly: if \( S \) has 0 elements, then \( S = \emptyset \), and the only subset is \( S \) itself, which is consistent with \( 2^0 = 1 \).

Now we assume \( P(n) \) holds and use it to show that \( P(n+1) \) is also true. Let \( S \) be a set of size \( n+1 \). Pick an element \( x \in S \) and let \( S' \) be the subset of \( S \) consisting all elements that are not equal to \( x \), i.e., \( S' = S \setminus \{x\} \). Then \( S' \) has size \( n \), so by induction the number of subsets of \( S' \) is \( 2^n \). Now, every subset of \( S \) either contains \( x \) or it does not. Those which do not contain \( x \) can be thought of as subsets of \( S' \), so there are \( 2^n \) of them. To count those that do contain \( x \), we can take any subset of \( S' \) and add \( x \) to it. This accounts for all of them exactly once, so there are also \( 2^n \) subsets that contain \( x \). All together we have \( 2^n + 2^n = 2^{n+1} \) subsets of \( S \), so \( P(n+1) \) holds. \( \square \)

Continuing with our example, if \( x = 1 \), then the subsets not containing \( x \) are \( \emptyset, \{\ast\}, \{U\}, \{\ast, U\} \), while those that do contain \( x \) are \( \{1\}, \{1, \ast\}, \{1, U\}, \{1, \ast, U\} \). There are \( 2^2 = 4 \) of each kind.

A natural followup is to determine how many subsets have a given size. In our previous example, there is 1 subset of size 0, 3 of size 1, 3 of size 2, and 1 of size 3. We’ll discuss this problem in the next section.

Some more to think about:
Show that $\sum_{i=0}^{n} i^2 = n(n+1)(2n+1)/6$ for all $n \geq 0$.
Show that $\sum_{i=0}^{n} 2^i = 2^{n+1} - 1$ for all $n \geq 0$.
Show that $4n < 2^n$ whenever $n \geq 5$.

What happens with $\sum_{i=0}^{n} i^3$ or $\sum_{i=0}^{n} i^4$, or...? In the first two cases, we got polynomials in $n$ on the right side. You’ll show on homework that this always happens.

1.2. Strong induction. The version of induction we just described is sometimes called “weak induction”. Here’s a variant sometimes called “strong induction”. We have the same setup: we want to prove that a sequence of statements $P(0), P(1), P(2), \ldots$ are true. Then strong induction works by completing the following 2 steps:

- Prove that $P(0)$ is true.
- Assuming that $P(0), P(1), \ldots, P(n)$ are all true, use them to prove that $P(n+1)$ is true.

You should convince yourself that this isn’t really anything logically distinct from weak induction. However, it can sometimes be convenient to use this variation.

Example 1.3. We know that every polynomial in $x$ is a linear combination of $1, x, x^2, x^3, \ldots$. We use strong induction to prove the statement that every polynomial is a linear combination is a linear combination of $1, (x-1), (x-1)^2, (x-1)^3, \ldots$.

Let $P(n)$ be the statement that every polynomial of degree $n$ is a linear combination of powers of $x-1$.

Then $P(0)$ is true: the only polynomials of degree 0 are constants, and we can write $c = c \cdot 1$.

Now assume that $P(0), P(1), \ldots, P(n)$ are all true. We will use them to show that $P(n+1)$ is true. Let $f(x)$ be an arbitrary polynomial of degree $n+1$. Let $\alpha$ be its leading term and define $g(x) = f(x) - \alpha \cdot (x-1)^{n+1}$. Then $g(x)$ is a polynomial of degree $\leq n$ since we have cancelled off the $x^{n+1}$ terms. So by strong induction, $g(x)$ is a linear combination of powers of $x-1$. If we add $\alpha \cdot (x-1)^{n+1}$ to this linear combination, we see that $f(x)$ is also a linear combination of powers of $x-1$. Since our argument applies to any polynomial of degree $n+1$, we have proved $P(n+1)$ is true. \qed

Some examples to think about:

- There’s nothing particular about powers of $x$ or powers of $x-1$. For example, we can take powers of any linear polynomial $ax + b$ with $a \neq 0$. Adapt the argument to work for this generalization.
- Every positive integer can be written in the form $2^n m$ where $n \geq 0$ and $m$ is an odd integer.
- Define a function $f$ on the natural numbers by $f(0) = 1$, $f(1) = 2$, and $f(n+1) = f(n) + 2f(n)$ for all $n \geq 1$. Show that $f(n) \leq 3^n$ for all $n \geq 0$.
- A chocolate bar is made up of unit squares in an $n \times m$ rectangular grid. You can break up the bar into 2 pieces by breaking on either a horizontal or vertical line. Show that you need to make $nm - 1$ breaks to completely separate the bar into $1 \times 1$ squares (if you have 2 pieces already, stacking them and breaking them counts as 2 breaks).
2. Elementary counting problems

2.1. Bijectons. Given two functions $f: X \to Y$ and $g: Y \to X$, we say that they are inverses if $f \circ g$ is the identity function on $Y$, i.e., $f(g(y)) = y$ for all $y \in Y$, and if $g \circ f$ is the identity function on $X$, i.e., $g(f(x)) = x$ for all $x \in X$. In that case, the functions $f$ and $g$ are called bijections.

The following is a very important principle in counting arguments:

**Proposition 2.1.** If there exists a bijection between $X$ and $Y$, then $|X| = |Y|$.

We can think of a bijection $f$ between $X$ and $Y$ as a way of matching the elements of $X$ with the elements of $Y$. In particular, $x \in X$ gets matched with $y = f(x) \in Y$. Note that if $x' \in X$ was also matched with $y$, i.e., $f(x') = f(x)$, then the existence of the inverse $g$ shows us that $g(f(x')) = g(f(x))$, or more simply $x = x'$. In other words, $f$ is forced to be one-to-one (or injective). On the other hand, every element is matched with something, i.e., every $y \in Y$ is of the form $f(x)$ for some $x$ because we can take $x = g(y)$. In other words, $f$ is forced to be onto (or surjective).

**Remark 2.2.** Bijections tell us that two sets have the same size without having to know how many elements are actually in the set.

Here’s a small example: imagine there is a theatre filled with hundreds of people and hundreds of seats. If we wanted to know if there are the same number of people as seats, we could count both. However, it would probably be much easier to just have each person take a seat and see if there are any empty or any standing people.

We’ll see some other examples later on.

2.2. Sum and product principle. Given two sets $X$ and $Y$ without any overlap, we have $|X \cup Y| = |X| + |Y|$. We’ll just take this for granted, though you can call it the sum principle if you’d like a name for it.

The set of pairs of elements $(x, y)$ where $x \in X$ and $y \in Y$ is the Cartesian product $X \times Y$. The related product principle says that $|X \times Y| = |X| \cdot |Y|$. Again, we will take this for granted and not usually refer to it by name.

2.3. Permutations and combinations. Given a set $S$ of objects, a permutation of $S$ is a way to put all of the elements of $S$ in order. More formally, if $|S| = n$, then a permutation is a bijection $f: S \to [n]$.

**Example 2.3.** There are 6 permutations of $\{1, 2, 3\}$ which we list:

\[123, \ 132, \ 213, \ 231, \ 312, \ 321.\]

To count permutations in general, we define the factorial as follows: $0! = 1$ and if $n$ is a positive integer, then $n! = n \cdot (n - 1)!$. Here are the first few values:

\[0! = 1, \ 1! = 1, \ 2! = 2, \ 3! = 6, \ 4! = 24, \ 5! = 120, \ 6! = 720.\]

In the previous example, we had 6 permutations of 3 elements, and $6 = 3!$. This holds more generally:

**Theorem 2.4.** If $S$ has $n$ elements and $n > 0$, then there are $n!$ different permutations of $S$. 

Proof. We do this by induction on \( n \). Let \( P(n) \) be the statement that a set of size \( n \) has exactly \( n! \) elements. The statement \( P(1) \) follows from the definition: there is exactly 1 way to order a single element, and \( 1! = 1 \). Now assume for our induction hypothesis that \( P(n) \) has been proven. Let \( S \) be a set of size \( n + 1 \). To order the elements, we can first pick any element to be first, and then we have to order the remaining \( n \) elements. There are \( n + 1 \) different elements that can be first, and for each such choice, there are \( n! \) ways to order the remaining elements by our induction hypothesis. So all together, we have \((n+1) \cdot n! = (n+1)!\) different ways to order all of them, which proves \( P(n+1) \).

We can use factorials to answer related questions. For example, suppose that some of the objects in our set can’t be distinguished from one another, so that some of the orderings end up being the same.

**Example 2.5.** (1) Suppose we are given 2 red flowers and 1 yellow flower. Aside from their color, the flowers look identical. We want to count how many ways we can display them in a single row. There are 3 objects total, so we might say there are \( 3! = 6 \) such ways. But consider what the 6 different ways look like:

\[
RRY, \quad RRY, \quad RYR, \quad YRR, \quad YRR.
\]

Since the two red flowers look identical, we don’t actually care which one comes first. So there are really only 3 different ways to do this – the answer \( 3! \) has included each different way twice, but we only wanted to count them a single time.

(2) Consider a larger problem: 10 red flowers and 5 yellow flowers. There are too many to list, so we consider a different approach. As above, if we naively count, then we would get \( 15! \) permutations of the flowers. But note that for any given arrangement, the 10 red flowers can be reordered in any way to get an identical arrangement, and same with the yellow flowers. So in the list of \( 15! \) permutations, each arrangement is being counted \( 10! \cdot 5! \) times. The number of distinct arrangements is then \( \frac{15!}{10! \cdot 5!} \).

(3) The same reasoning allows us to generalize. If we have \( r \) red flowers and \( y \) yellow flowers, then the number of different ways to arrange them is \( \frac{(r+y)!}{r! \cdot y!} \).

(4) How about more than 2 colors of flowers? If we threw in \( b \) blue flowers, then again the same reasoning gives us \( \frac{(r+y+b)!}{r! \cdot y! \cdot b!} \) different arrangements. 

Now we state a general formula, which again can be derived by the same reasoning as in (2) above. Suppose we are given \( n \) objects, which have one of \( k \) different types (for example, our objects could be flowers and the types are colors). Also, objects of the same type are considered identical. For convenience, we will label the “types” with numbers \( 1, 2, \ldots, k \) and let \( a_i \) be the number of objects of type \( i \) (so \( a_1 + a_2 + \cdots + a_k = n \)).

**Theorem 2.6.** The number of ways to arrange the \( n \) objects in the above situation is

\[
\frac{n!}{a_1! \cdot a_2! \cdots a_k!}.
\]

As an exercise, you should adapt the reasoning in (2) to give a proof of this theorem. The quantity above will be used a lot, so we give it a symbol, called the * multinomial coefficient *:

\[
\binom{n}{a_1, a_2, \ldots, a_k} := \frac{n!}{a_1! \cdot a_2! \cdots a_k!}.
\]
In the case when \( k = 2 \) (a very important case), it is called the **binomial coefficient**. Note that in this case, \( a_2 = n - a_1 \), so for shorthand, one often just writes \( \binom{n}{a_1} \) instead of \( \binom{n}{a_1, a_2} \). For similar reasons, \( \binom{n}{a_2} \) is also used as a shorthand.

2.4. **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 \( \emptyset \) is considered a word of length 0.

**Example 2.7.** Say our alphabet is \( A = \{a, b\} \). The words of length \( \leq 2 \) are:

\[
\emptyset, \ a, \ b, \ aa, \ ab, \ ba, \ bb.
\]

**Theorem 2.8.** 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 \( |A^k| = |A|^k \).

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 \( n \cdot n \cdot \cdots n = n^k \) different choices for words.

For a positive integer \( n \), let \([n]\) denote the set \( \{1, \ldots, n\} \).

**Example 2.9.** We use words to show that the number of subsets of \([n]\) is \( 2^n \) (we’ve already seen this result, so now we’re using a different proof method).

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

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

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 2.10.** 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 \( S \)”, “in \( T \) but not \( S \)” and “not in \( T \)”. 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.

How about words without repeating entries? 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 2.11.** 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

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

□

Some further things to think about:

- 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?
- Using that $(n)_k = n \cdot (n - 1)_{k-1}$, can you find a proof for Theorem 2.11 that uses induction?

2.5. Choice problems. We finish up with some related counting problems. Recall we showed that an $n$-element set has exactly $2^n$ subsets. We can refine this problem by asking about subsets of a given size.

Theorem 2.12. The number of $k$-element subsets of $[n]$ is

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

There are many ways to prove this, but we’ll just do one for now:

Proof. In the last section on words, we identified subsets of $[n]$ with words of length $n$ on $\{0,1\}$, with a 1 in position $i$ if and only if $i$ belongs to the subset. So the number of subsets of size $k$ are exactly the number of words with exactly $k$ instances of 1. This is the same as arranging $n-k$ 0’s and $k$ 1’s from the section on permutations. In that case, we saw the answer is $\frac{n!}{(n-k)!k!}$. □

Corollary 2.13. $\sum_{k=0}^{n} \binom{n}{k} = 2^n$.

Proof. The left hand side counts the number of subsets of $[n]$ of some size $k$ where $k$ ranges from 0 to $n$. But all subsets of $[n]$ are accounted for and we’ve seen that $2^n$ is the number of all subsets of $[n]$. □

Here’s an important identity for binomial coefficients (we interpret $\binom{n}{-1} = 0$):

Proposition 2.14 (Pascal’s identity). For any $k \geq 0$, we have

$$\binom{n}{k-1} + \binom{n}{k} = \binom{n+1}{k}.$$ 

Proof. The right hand side is the number of subsets of $[n+1]$ of size $k$. There are 2 types of such subsets: those that contain $n+1$ and those that do not. Note that the subsets that do contain $n+1$ are naturally in bijection with the subsets of $[n]$ of size $k-1$: to get such a subset, delete $n+1$. Those that do not contain $n+1$ are naturally already in bijection with the subsets of $[n]$ of size $k$. The two sets don’t overlap and their sizes are $\binom{n}{k-1}$ and $\binom{n}{k}$, respectively. □
An important variation of subset is the notion of a multiset. Given a set $S$, a **multiset** of $S$ is like a subset, but we allow elements to be repeated. Said another way, a subset of $S$ can be thought of as a way of assigning either a 0 or 1 to an element, based on whether it gets included. A multiset is then a way to assign some non-negative integer to each element, where numbers bigger than 1 mean we have picked them multiple times.

**Example 2.15.** There are 10 multisets of $[3]$ of size 3:

\[
\{1,1,1\}, \{1,1,2\}, \{1,1,3\}, \{1,2,2\}, \{1,2,3\}, \\
\{1,3,3\}, \{2,2,2\}, \{2,2,3\}, \{2,3,3\}, \{3,3,3\}.
\]

Aside from exhaustively checking, how do we know that’s all of them? Here’s a trick: given a multiset, add 1 to the second smallest values (including ties) and add 2 to the largest value. What happens to the above:

\[
\{1,2,3\}, \{1,2,4\}, \{1,2,5\}, \{1,3,4\}, \{1,3,5\}, \\
\{1,4,5\}, \{2,3,4\}, \{2,3,5\}, \{2,4,5\}, \{3,4,5\}.
\]

We get all of the 3-element subsets of $[5]$. The process is reversible using subtraction, so there is a more general fact here.

**Theorem 2.16.** The number of $k$-element multisets of $[n]$ is

\[
\binom{n+k-1}{k}.
\]

**Proof.** We adapt the example above to find a bijection between $k$-element multisets of $[n]$ and $k$-element subsets of $[n+k-1]$. Given a multiset $S$, sort the elements as $s_1 \leq s_2 \leq \cdots \leq s_k$. From this, we get a subset $\{s_1, s_2+1, s_3+2, \ldots, s_k+(k-1)\}$ of $[n+k-1]$. On the other hand, given a subset $T$ of $[n+k-1]$, sort the elements as $t_1 < t_2 < \cdots < t_k$. From this, we get a multiset $\{t_1, t_2-1, t_3-2, \ldots, t_k-(k-1)\}$ of $[n]$. I will omit the details that these are well-defined and inverse to one another. (But you should make sure that you could do this if asked.)

Some additional things:

- From the formula, we see that $\binom{n}{k} = \binom{n}{n-k}$. This would also be implied if we could construct a bijection between the $k$-element subsets and the $(n-k)$-element subsets of $[n]$. Can you find one?
- Given variables $x, y, z$, we can form polynomials. A monomial is a product of the form $x^ay^bz^c$, and its degree is $a + b + c$. How many monomials in $x, y, z$ are there of degree $d$? What if we have $n$ variables $x_1, x_2, \ldots, x_n$?

3. Binomial theorem and generalizations

3.1. Binomial theorem. The binomial theorem is about expanding powers of $x + y$ where we think of $x, y$ as variables. For example:

\[
(x + y)^2 = x^2 + 2xy + y^2, \\
(x + y)^3 = x^3 + 3x^2y + 3xy^2 + y^3.
\]
**Theorem 3.1** (Binomial theorem). For any \( n \geq 0 \), we have

\[
(x + y)^n = \sum_{i=0}^{n} \binom{n}{i} x^i y^{n-i}.
\]

Here’s the proof given in the book.

**Proof.** Consider how to expand the product \((x + y)^n = (x + y)(x + y) \cdots (x + y)\). To get a term, from each expression \((x + y)\), we have to either pick \(x\) or \(y\). The final term we get is \(x^i y^{n-i}\) if the number of times we chose \(x\) is \(i\) (and hence the number of times we’ve chosen \(y\) is \(n-i\)). The number of times this term appears is therefore the number of different ways we could have chosen \(x\) exactly \(i\) times. For each way of doing this, we can associate to it a subset of \([n]\) of size \(i\): the number \(j\) is in the subset if and only if we chose \(x\) in the \(j\)th copy of \((x + y)\). We have already seen that the number of subsets of \([n]\) of size \(i\) is \(\binom{n}{i}\). \(\square\)

Here’s a proof using induction.

**Proof.** For \(n = 0\), the formula becomes \((x + y)^0 = 1\) which is valid.

Now suppose the formula is valid for \(n\). Then we have

\[
(x + y)^{n+1} = (x + y)(x + y)^n = (x + y) \sum_{i=0}^{n} \binom{n}{i} x^i y^{n-i}.
\]

For a given \(k\), there are at most 2 ways to get \(x^k y^{n+1-k}\) on the right side: either we get it from \(x \cdot \binom{n}{k} x^{k-1} y^{n-k+1}\) or from \(y \cdot \binom{n}{k} x^k y^{n-k}\). If we add these up, then we get \(\binom{n+1}{k}\) by Pascal’s identity. \(\square\)

We can manipulate the binomial theorem in a lot of different ways (taking derivatives with respect to \(x\) or \(y\), or doing substitutions). This will give us a lot of new identities. Here are a few of particular interest (some are old):

**Corollary 3.2.** \(2^n = \sum_{i=0}^{n} \binom{n}{i}\).

**Proof.** Substitute \(x = y = 1\) into the binomial theorem. \(\square\)

This says that the total number of subsets of \([n]\) is \(2^n\) which is a familiar fact from before.

**Corollary 3.3.** For \(n > 0\), we have \(0 = \sum_{i=0}^{n} (-1)^i \binom{n}{i}\).

**Proof.** Substitute \(x = -1\) and \(y = 1\) into the binomial theorem. \(\square\)

If we rewrite this, it says that the number of subsets of even size is the same as the number of subsets of odd size. It is worth finding a more direct proof of this fact which does not rely on the binomial theorem.

**Corollary 3.4.** \(n 2^{n-1} = \sum_{i=0}^{n} i \binom{n}{i}\).

**Proof.** Take the derivative of both sides of the binomial theorem with respect to \(x\) to get

\[
n(x + y)^{n-1} = \sum_{i=0}^{n} i \binom{n}{i} x^{i-1} y^{n-i}.
\]

Now substitute \(x = y = 1\). \(\square\)
It is possible to interpret this formula as the size of some set so that both sides are different ways to count the number of elements in that set. Can you figure out how to do that? How about if we took the derivative twice with respect to \( x \)? Or if we took it with respect to \( x \) and then with respect to \( y \)?

### 3.2. Multinomial theorem.

**Theorem 3.5** (Multinomial theorem). For \( n, k \geq 0 \), we have

\[
(x_1 + x_2 + \cdots + x_k)^n = \sum_{\substack{(a_1, a_2, \ldots, a_k) \in \mathbb{N}_0^k \colon \sum a_i = n}} \binom{n}{a_1, a_2, \ldots, a_k} x_1^{a_1} x_2^{a_2} \cdots x_k^{a_k}.
\]

**Proof.** The proof is similar to the binomial theorem. Consider expanding the product \( (x_1 + \cdots + x_k)^n \). To do this, we first have to pick one of the \( x_i \) from the first factor, pick another one from the second factor, etc. To get the term \( x_1^{a_1} x_2^{a_2} \cdots x_k^{a_k} \), we need to have picked \( x_1 \) exactly \( a_1 \) times, picked \( x_2 \) exactly \( a_2 \) times, etc. We can think of this as arranging \( n \) objects, where \( a_i \) of them have “type \( i \)”. In that case, we’ve already discussed that this is counted by the multinomial coefficient \( \binom{n}{a_1, a_2, \ldots, a_k} \).

By performing substitutions, we can get a bunch of identities that generalize the one from the previous section. I’ll omit the proofs, try to fill them in.

\[
k^n = \sum_{\substack{(a_1, a_2, \ldots, a_k) \in \mathbb{N}_0^k \colon \sum a_i = n}} \binom{n}{a_1, a_2, \ldots, a_k},
\]

\[
0 = \sum_{\substack{(a_1, a_2, \ldots, a_k) \in \mathbb{N}_0^k \colon \sum a_i = n}} (1 - k)^{a_1} \binom{n}{a_1, a_2, \ldots, a_k},
\]

\[
nk^{n-1} = \sum_{\substack{(a_1, a_2, \ldots, a_k) \in \mathbb{N}_0^k \colon \sum a_i = n}} a_1 \binom{n}{a_1, a_2, \ldots, a_k}.
\]

### 4. Inclusion-exclusion

**Example 4.1.** Suppose we have a room of students, and 14 of them play basketball, 10 of them play football. How many students play at least one of these? We can’t answer the question because there might be students who play both. But we can say that the total number is 24 minus the amount in the overlap.

![Venn Diagram](attachment:venn.png)

Alternatively, let \( B \) be the set who play basketball and let \( F \) be the set who play football. Then what we’ve said is:

\[
|B \cup F| = |B| + |F| - |B \cap F|.
\]
New situation: there are additionally 8 students who play hockey. Let $H$ be the set of students who play hockey. What information do we need to know how many total students there are?

Here the overlap region is more complicated: it has 4 regions, which suggest that we need 4 more pieces of information. The following formula works:

$$|B \cup F \cup H| = |B| + |F| + |H| - |B \cap F| - |B \cap H| - |F \cap H| + |B \cap F \cap H|.$$

To see this, the total diagram has 7 regions and we need to make sure that students in each region get counted exactly once in the right side expression. For example, consider students who play basketball and football, but don’t play hockey. They get counted in $B$, $F$, $B \cap F$ with signs $+1$, $+1$, $-1$, which sums up to 1. How about students who play all 3? They get counted in all terms with 4 $+1$ signs and 3 $-1$ signs, again adding up to 1. You can check the other 5 to make sure the count is right. □

The examples above have a generalization to $n$ sets, though the diagram is harder to draw beyond 3.

What’s the pattern so far? We have to add up all of the sizes of the sets involved, then we subtract off the sizes of all ways of intersecting two of them, and then we add back the sizes of all ways of intersecting three of them. How does this continue? In general, the signs continue to alternate (add, subtract, add, subtract, ...) and at the $j$th step, we have to consider all sizes of intersecting $j$ different sets.

**Theorem 4.2** (Inclusion-Exclusion). Let $A_1, \ldots, A_n$ be finite sets. Then

$$|A_1 \cup \cdots \cup A_n| = \sum_{j=1}^{n} (-1)^{j-1} \sum_{\{i_1, i_2, \ldots, i_j\}} |A_{i_1} \cap A_{i_2} \cap \cdots \cap A_{i_j}|,$$

where the second sum is over all $j$-element subsets of $[n]$, i.e., we add up the sizes of all possible ways of intersecting $j$ of the sets $A_1, \ldots, A_n$.

**Proof.** We just need to make sure that every element $x \in A_1 \cup \cdots \cup A_n$ is counted exactly once on the right hand side. Let $S = \{s_1, \ldots, s_k\}$ be all of the indices such that $x \in A_{s_r}$. Then $x$ belongs to $A_{i_1} \cap \cdots \cap A_{i_j}$ if and only if $\{i_1, \ldots, i_j\} \subseteq S$. So the relevant contributions for $x$ is a sum over all of the nonempty subsets of $S$:

$$\sum_{T \subseteq S} (-1)^{|T|-1} = -\sum_{n=1}^{|S|} \binom{|S|}{n} (-1)^n.$$

However, since $|S| > 0$, we have shown before that $\sum_{n=0}^{|S|} \binom{|S|}{n} (-1)^n = 0$, so the sum above is $\binom{|S|}{0} = 1$. □
We can also prove this by induction on \( n \). Can you see how?

We use this to solve the derangements problem. Here is a version of that problem: suppose we have \( n \) people and they all put their hat into a box. The hats are redistributed to the people at random. What is the chance that nobody gets their own hat back? (We won’t solve this exactly, but see how to get a close approximation to the answer.)

First, we can think of a permutation of \([n]\) as the same thing as a bijection \( f : [n] \to [n] \) (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 4.3.** The number of derangements of size \( n \) is

\[
\sum_{i=0}^{n} (-1)^{i} \frac{n!}{i!}.
\]

**Proof.** It turns out to be easier to count the number of permutations which are *not* derangements and then subtract that from the total number of permutations. For \( i = 1, \ldots, n \), let \( A_i \) be the set of bijections \( f \) such that \( f(i) = i \). Then the set of non-derangements is \( A_1 \cup \cdots \cup A_n \). To apply inclusion-exclusion, we need to count the size of \( A_{i_1} \cap \cdots \cap A_{i_j} \) for some choice of indices \( i_1, \ldots, i_j \). This is the set of bijections \( f : [n] \to [n] \) such that \( f(i_1) = i_1, \ldots, f(i_j) = i_j \). The remaining information to specify \( f \) are its values outside of \( i_1, \ldots, i_j \), which we can interpret as a bijection of \([n] \setminus \{i_1, \ldots, i_j\} \) to itself. So there are \((n - j)! \) of them. So we get

\[
|A_1 \cup \cdots \cup A_n| = \sum_{j=1}^{n} (-1)^{j-1} \sum_{\{i_1, \ldots, i_j\}} |A_{i_1} \cap \cdots \cap A_{i_j}|
\]

\[
= \sum_{j=1}^{n} (-1)^{j-1} \sum_{\{i_1, \ldots, i_j\}} (n - j)!
\]

\[
= \sum_{j=1}^{n} (-1)^{j-1} \binom{n}{j} (n - j)!
\]

\[
= \sum_{j=1}^{n} (-1)^{j-1} \frac{n!}{j!}.
\]

Remember that we have to subtract this from \( n! \). So the final answer simplifies as so:

\[
n! - \sum_{j=1}^{n} (-1)^{j-1} \frac{n!}{j!} = \sum_{j=0}^{n} (-1)^{j} \frac{n!}{j!}.
\]

\[\square\]

The problem with formulas coming from inclusion-exclusion is the alternating sign. It can generally be hard to estimate the behavior of the quantity as \( n \) grows. For example, binomial coefficients \( \binom{n}{i} \) (for fixed \( i \)) limit to infinity as \( n \) goes to infinity. However, the alternating sum

\[
\sum_{i=0}^{n} (-1)^{i} \binom{n}{i}
\]
is 0. For derangements, we can use the following observation. We have a formula for the exponential function

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

If we plug in $x = -1$ and 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. From calculus, taking the first $n$ terms of a Taylor expansion is supposed to be a good approximation for a function, so for $n \to \infty$, the proportion of permutations that are derangements is $e^{-1} \approx 0.368$, or roughly 36.8%.

## 5. Graph theory, introduction

The origin of graph theory is said to be Euler’s solution of the bridges of Königsberg problem.

(Taken from [https://commons.wikimedia.org/wiki/File:Konigsberg_bridges.png](https://commons.wikimedia.org/wiki/File:Konigsberg_bridges.png))

The problem was to find a path which starts and ends at the same point, and crosses each bridge exactly once.

There are a few things to note here: the lengths of the bridges are not important, nor are the sizes of the landmasses. All that really matters is: between any two landmasses, how many bridges connect them? So we can simplify the picture dramatically:

![Konigsberg bridges](https://commons.wikimedia.org/wiki/File:Konigsberg_bridges.png)

This notion can be abstracted with the definition of a graph. Given a set $V$, let $\binom{V}{2}$ denote the set of 2-element subsets of $V$.

**Definition 5.1.** 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 use it 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 of lengths in this definition. So while we can draw a graph as we have done above, such a pictorial representation is not unique.

Example 5.2. The picture above represents the graph with $V = \{A, B, C, D\}$ and $E = \{\{A, B\}, \{B, D\}, \{B, C\}, \{A, C\}, \{C, D\}, \{A, C\}\}$.  

5.1. Eulerian trails.

Definition 5.3. 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 $v_k$. A walk is closed if $v_0 = v_k$. The length of the walk is $k$.

A trail is a walk that does not repeat edges (we treat multiple edges between two vertices as distinct). An Eulerian trail is a trail that uses every edge exactly once.

A path is a trail that does not repeat vertices (in particular, it has no loops).

Now we can phrase the Königsberg problem in our new language: does the Königsberg graph have a closed Eulerian trail?

We need a few more definitions before answering this question.

Definition 5.4. Given a vertex $v$ in a graph $G$, its degree $\deg(v)$ is the number of edges connected to $v$, except that loops at $v$ must be counted twice. In other words, it is 2 times the number of loops at $v$ plus the number of non-loop edges at $v$.

Heuristically, we can think of an edge as being comprised of two “halves”, each one connected to one of the endpoints. This perspective is partially why we want to count loops twice: we really want to know how many half-edges are connected to $v$. Alternatively, imagine zooming in really close to a vertex. In that case, the two pieces of a loop would look like two separate edges.

Definition 5.5. A graph is connected if, for any two vertices $v$ and $w$, there exists a walk that begins at $v$ and ends at $w$.

In general, we can put an equivalence relation on the vertices of $G$, declaring that $v \sim w$ if there exists a walk beginning at $v$ and ending at $w$. The equivalence classes of this relation are the connected components of $G$. Each connected component is a connected graph.

The following will be convenient to use repeatedly. I’ll leave it as homework to prove.

Proposition 5.6. If there is a walk from $x$ to $y$, then there is also a path from $x$ to $y$.

Theorem 5.7 (Euler). Let $G$ be a connected graph. Then $G$ has a closed Eulerian trail if and only if every vertex has even degree.

Proof. If $G$ has a closed Eulerian trail, then along this path, each vertex is entered the same number of times that it is exited. Since none of these edges are being repeated, we see that the number of “half-edges” being used must be even, and so $\deg(v)$ is even for all $v$.

Now we have to prove the opposite implication. So let $G$ be a connected graph such that $\deg(v)$ is even for all vertices $v$. We need to show that $G$ has a closed Eulerian trail. We will prove this statement by induction on the number of edges of $G$. If the number is 0, then $G$ is a single vertex, and there is nothing to prove so our base case is done. Now suppose we have proven it for all connected graphs such that $\deg(v)$ is even for all $v$ that have $\leq n$ edges and let $G$ be a graph such that $\deg(v)$ is even for all $v$ and has $n + 1$ edges.
Step 1. Find any closed trail. To do this, begin with any edge. If it is a loop, we are done. Otherwise, it connects \( v \) and \( w \) with \( v \neq w \). Since \( \deg(w) \) is even, there must be another edge that we can follow. Continue doing this until we hit a vertex we’ve seen before. Then our path will contain a closed trail. Let \( H \) be the graph consisting of the vertices and edges used on this trail.

Step 2. Delete the edges in \( H \) from \( G \). From the reasoning in the first part of the proof, this will change the degree of each vertex by an even amount, so the resulting graph \( G' \) still has all vertices with even degree. However, it may not be connected. But by induction, each connected component has a closed Eulerian trail.

Step 3. Each connected component of \( G' \) must share a vertex with one in \( H \). If not, pick a connected component that violates this condition and pick a vertex \( v \) in it and pick a vertex \( w \) in \( H \). In the original graph \( G \), there is a path between \( v \) and \( w \) since \( G \) is connected. If we follow this path from \( v \), there is a first time that it hits a vertex in \( H \), and the edges up to this point aren’t in \( H \) (because all edges of \( H \) go between two vertices in \( H \)), which means that some vertex in \( H \) is in the same connected component as \( v \).

So we can attach the closed Eulerian trails from each component to \( H \). The resulting path is a closed Eulerian trail for \( G \). \( \square \)

If a graph \( G \) has a closed Eulerian trail, we can “rotate” it so that it begins and ends at any particular vertex that we like. Formally, if our closed Eulerian trail is \( v_0, e_1, v_1, \ldots, v_k \), then we get another closed Eulerian trail \( v_1, e_2, \ldots, v_k, e_1, v_1 \). The important thing is that \( v_k = v_0 \) since the trail is closed. We can repeat this rotation as many times as needed to get the desired starting vertex and starting edge.

Corollary 5.8. Let \( G \) be a connected graph and let \( v, w \in V \) be different vertices. Then there is an Eulerian trail starting at \( v \) and ending at \( w \) if and only if \( \deg(v), \deg(w) \) are both odd, and \( \deg(x) \) is even for all other vertices \( x \in V \).

Proof. Add a new edge \( f = \{v, w\} \) to \( G \) to get a graph \( G' \). Then the degree of every vertex of \( G' \) has even degree, so by the previous theorem, it has a closed Eulerian trail. By the above comments, we can rotate this trail until \( f \) is the starting edge. Delete \( f \) from this closed trail, and the result is an Eulerian trail of \( G \) which starts at one of \( v, w \) and ends at the other one. If it’s backwards, we can always reverse the trail. \( \square \)

Finally, a related statement about vertices of odd degree.

Theorem 5.9. Let \( G \) be a graph. There are an even number of vertices with odd degree.

Proof. Let \( v_1, \ldots, v_n \) be the vertices of \( G \) and let \( d_i = \deg(v_i) \). Note that \( d_1 + \cdots + d_n \) is twice the number of edges, since each edge contributes 1 to the degree of each of its endpoints (still true for loops). In particular, \( d_1 + \cdots + d_n \) is even. If there are an odd number of vertices with odd degree, then the sum would be odd, and hence we know that it is not the case. \( \square \)

5.2. Directed graphs. Our definition of graph models when things are related by putting an edge between them. For example, our nodes could represent places and the edges could represent whether or not they are connected by a road. This suggests an equal relation between the two, but we might want to be able to talk about roads that only go in one direction, for example. For that, we can use the notion of a directed graph. Intuitively, this is the same as a graph except each edge now has an orientation, i.e., it has a direction placed on it. Formally, instead of thinking of (non-loop) edges as a 2-element subset of the
vertex set, they are now elements of $V \times V$ (where the entries are unequal; loops then become the case when the entries are equal).

The definition of a walk in a directed graph is changed in the following way: it is now a sequence $v_0, e_1, \ldots, v_k$ where $e_i = (v_{i-1}, v_i)$, i.e., it has to be compatible with the orientations of the edges. All of the other definitions can then be adapted once we’ve made this change. We define a directed graph $G$ to be **strongly connected** if for all $v, w \in V$, there is a walk from $v$ to $w$.

Given a vertex $v$, its **in-degree** $\text{indeg}(v)$ is the number of edges of the form $(x, v)$ for some $x \in V$ and its **out-degree** $\text{outdeg}(v)$ is the number of edges of the form $(v, x)$ for some $x \in V$. Note that each loop at $v$ contributes 1 to its in-degree and 1 to its out-degree.

The directed version of Euler’s theorem is the following.

**Theorem 5.10.** Let $G$ be a strongly connected directed graph. Then $G$ has a closed Eulerian trail if and only if $\text{indeg}(v) = \text{outdeg}(v)$ for all $v \in V$.

The proof is fairly similar to the proof of the undirected version, so we won’t repeat it here.

### 5.3. Hamiltonian cycles

The last section was about walks that use all of the edges exactly once. Now we consider the dual idea where we use all of the vertices exactly once (but not necessarily all edges).

**Definition 5.11.** Let $G$ be a graph. A **cycle** is a closed trail such that each vertex is used at most once. In other words, it’s a closed walk that doesn’t repeat vertices or edges. A **Hamiltonian cycle** is a closed trail such that each vertex is used exactly once. A **Hamiltonian path** is a path that uses every vertex exactly once (the difference is that it is not required to be closed).

**Remark 5.12.** We would like to have a simple criterion which determines if a graph has a Hamiltonian cycle or path like Euler’s theorem. However, it is known that determining if a graph has a Hamiltonian cycle or path is an “NP-complete problem”. This means that there is unlikely to even be an efficient algorithm to determine this, so we probably don’t expect any simple condition.

While we don’t have a good condition for the existence of Hamiltonian cycles, we can give stronger conditions which imply their existence.

**Theorem 5.13.** Let $n \geq 3$. Let $G$ be a simple graph with $n$ vertices and assume that $\text{deg}(v) \geq n/2$ for all vertices $v$. Then $G$ has a Hamiltonian cycle.

**Proof.** Assume that the theorem is false. Then there exists a simple graph $G$ with $n \geq 3$ vertices such that $\text{deg}(v) \geq n/2$ for all vertices and $G$ does not have a Hamiltonian cycle. We want a maximal counterexample in the following sense: if there is an edge we can add to $G$ such that the result does not have a Hamiltonian cycle, then do it and keep doing it until it’s not possible. We end up with a counterexample, call it $G'$ with $n$ vertices which still has $\text{deg}(v) \geq n/2$ and no Hamiltonian cycle, with the additional property that adding any new edge will introduce a Hamiltonian cycle.

First, $G'$ must be missing at least one edge: if every edge is present, then there is clearly a Hamiltonian cycle. Let $x, y$ be a pair of vertices such that $\{x, y\}$ is not an edge. Let $G''$ be the result of adding $\{x, y\}$. By construction, $G''$ has a Hamiltonian cycle, and it must use $\{x, y\}$ (otherwise $G'$ has one too). In particular, if we consider the rest of the edges, we have
a Hamiltonian path starting at \( x \) and ending at \( y \) using only edges from \( G' \). Let \( z_1, \ldots, z_n \) be the order in which the vertices are visited on this path (so \( z_1 = x \) and \( z_n = y \)).

Now define

\[
S = \{ i \mid 2 \leq i \leq n - 1, \ \{ z_{i+1}, x \} \text{ is an edge in } G' \} \\
T = \{ i \mid 2 \leq i \leq n - 1, \ \{ z_i, y \} \text{ is an edge in } G' \}.
\]

Then \( |S| \geq \deg(x) - 1 \geq n/2 - 1 \) (since we aren’t counting \( z_2 \) though it could be connected by an edge to \( x \)) and \( |T| = \deg(y) \geq n/2 \). Note that \( S \cup T \subseteq \{2, \ldots, n - 1\} \) and so \( |S \cup T| \leq n - 2 \), so we have (by inclusion-exclusion)

\[
|S \cap T| = |S| + |T| - |S \cup T| \geq \frac{n}{2} - 1 + \frac{n}{2} - (n - 2) = 1
\]

and so \( S \cap T \neq \emptyset \). Pick \( i \in S \cap T \). Then \( \{x, z_{i+1}\} \) and \( \{y, z_i\} \) are both edges of \( G' \). That means that \( G' \) does in fact have a Hamiltonian cycle: the order is

\[
x = z_1, z_2, \ldots, z_i, y = z_n, z_{n-1}, \ldots, z_{i+2}, z_{i+1},
\]

which is a contradiction. \( \Box \)

The condition \( \deg(v) \geq n/2 \) is way too strong though. Consider the graph with vertices \([n]\) and edges \( \{\{1, 2\}, \{2, 3\}, \{3, 4\}, \ldots, \{n, n+1\}, \{1, n\}\} \), so if you draw it, it is a cycle of length \( n \). Then this has a Hamiltonian cycle but the degree of each vertex is only 2. You’ll see another example in homework where the condition is too strong to be applicable.

5.4. **Graph isomorphisms.** Graph isomorphisms formalize what it means for two graphs to “look the same” without literally being equal to each other. For example, consider the following 3 graphs

\[
\begin{array}{ccc}
1 & \rightarrow & 2 \\
\mid & \mid & \mid \\
3 & \rightarrow & 4 \\
\end{array}
\]

Formally, both graphs have vertex set \([4]\), the first graph has edges \( \{\{1, 2\}, \{1, 3\}, \{2, 4\}, \{3, 4\}\} \) and the second graph has edges \( \{\{1, 2\}, \{1, 4\}, \{3, 4\}, \{2, 3\}\} \), so they are not the same. However, consider the bijection \( f: [4] \rightarrow [4] \) defined by \( f(1) = 1, f(2) = 4, f(3) = 2, f(4) = 3 \). Applying this function to the first graph gives the second graph.

**Definition 5.14.** Let \( G = (V, E) \) and \( H = (V', E') \) be graphs. A function \( f: V \rightarrow V' \) is a **graph isomorphism** if it is a bijection, and for all \( x, y \in V \), the number of edges between \( x \) and \( y \) equals the number of edges of \( f(x), f(y) \in V' \).

If a graph isomorphism exists between \( G \) and \( H \), then \( G \) and \( H \) are **isomorphic**.

If \( G = H \), then a graph isomorphism is called an **automorphism**. \( \Box \)

If \( f \) defines a graph isomorphism between \( G \) and \( H \), then \( f^{-1}: V' \rightarrow V \) defines a graph isomorphism between \( H \) and \( G \). Furthermore, if we have an isomorphism from \( H \) to a third graph \( I \), then we can compose these isomorphisms to get one between \( G \) and \( I \).

Many properties of graphs are preserved by isomorphism, meaning that if \( G \) has that property then so does any graph which is isomorphic to \( G \). Heuristically, any property or quantity which does not depend on the specific way that the vertices are named will be preserved by isomorphism. Such properties are called **isomorphism invariants**. Some examples:

- Number of vertices
• Number of edges
• The multiset of degrees of vertices
• Whether or not a graph has a Hamiltonian cycle.

Remark 5.15. Given two graphs, there is a naive algorithm for determining whether or not they are isomorphic: first, check if they have the same number of vertices, and if so, test all bijections between their vertex sets. But this is a pretty bad algorithm. One can ask whether or not there is an efficient (polynomial-time) algorithm and this is an open problem. □

Automorphisms of $G$ capture its symmetries and can be used to simplify proofs and provide justification for the phrase “without loss of generality...”. Here is an example which illustrates what I mean.

Example 5.16. Consider the graph $G$ with vertices $[n]$ ($n \geq 3$) and edges $\{i, i+1\}$ for $i = 1, \ldots, n-1$ and $\{1, n\}$. Consider the following statement: every path $v_0, e_1, v_1, e_2, v_2$ of length 2 can be extended uniquely to a Hamiltonian cycle. When $v_0 = 1$, $v_1 = 2$, $v_2 = 3$, we can prove this directly since the unique way to extend to a Hamiltonian cycle is to continue traveling around the graph, i.e., take $v_i = i+1$ for $i = 0, \ldots, n-1$ and $v_{n+1} = 1$.

But there are other paths of length 2. However, up to applying an automorphism, they can be turned into the one we just analyzed. Namely, “rotation” gives us an automorphism, i.e., $f(i) = i - 1$ for $i = 2, \ldots, n$ and $f(1) = n$ as does “reflection”, i.e., $f(1) = 1$ and $f(i) = n + 2 - i$ for $i \neq 1$. If we repeatedly apply rotation, we can turn $v_0$ into 1. Then either $v_1 = 2$ (in which case $v_2 = 3$ and we’re done) or $v_1 = n$, in which case one application of reflection will turn this into what we want. □

How many simple graphs with $n$ vertices are there up to isomorphism? The numbers for $n = 1, \ldots, 8$ are 1, 2, 4, 11, 34, 156, 1044, 12346. There probably isn’t a simple formula in general.

6. Trees

6.1. Definition and basic properties.

Definition 6.1. A **tree** is a connected simple graph that does not contain a cycle. A **forest** is a simple graph that does not contain a cycle. □

The difference is that forests are not required to be connected, but they can be. So every tree is a forest, but not vice versa. The connected components of a forest are all trees, so we can think of forests as being made up of trees (hence the name).

Example 6.2. Here are two examples of trees on 5 vertices:

![Tree Diagram]

We’ll give a few different ways to characterize trees.

Theorem 6.3. Let $G$ be a connected simple graph. Then $G$ is a tree if and only if for all vertices $x, y \in V$, there is exactly one path going from $x$ to $y$. □
Proof. Suppose that $G$ is a tree. Since $G$ is connected, for any vertices $x, y$, there is at least one walk going from $x$ to $y$, and hence there is at least one path going from $x$ to $y$ (by a homework problem). We claim that there is always exactly one path. If not, then let $x, y$ be a pair that has at least 2 paths between them. Furthermore, among all such pairs, pick $x, y$ so that the length of a path between them is as small as possible. Let $e_1, \ldots, e_k$ be the shortest possible path between them (in case of ties, pick any) and let $f_1, \ldots, f_\ell$ be the next shortest path between them. I claim that $e_1, \ldots, e_k, f_\ell, f_{\ell-1}, \ldots, f_1$ is a cycle, i.e., it doesn’t visit the same vertex more than once (except $x$). This is left as a homework problem.

Now suppose that $G$ is not a tree. Since $G$ is connected, it contains a cycle. We can split this into two pieces to get two different paths that go between the same pair of vertices. □

Definition 6.4. Let $G$ be a simple connected graph. Then $G$ is minimally connected if deleting any edge causes it to become disconnected.

Theorem 6.5. Let $G$ be a connected simple graph. Then $G$ is a tree if and only if $G$ is minimally connected.

Proof. First suppose that $G$ is a tree. If $G$ has 1 vertex, there’s nothing to say. If it has more, then it has at least one edge, pick any of them $\{x, y\}$. This edge gives a path from $x$ to $y$, so by uniqueness (previous theorem), there isn’t another one. So if we delete this edge, there’s no way to get from $x$ to $y$, which means the result is disconnected. Since this is true for any edge, $G$ is minimally connected.

Now suppose that $G$ is not a tree. Since $G$ is connected, that means there is a cycle. If we delete any edge $\{x, y\}$ from that cycle, the resulting graph is still connected: for any two vertices $a, b$, if we had to use that edge to get from $a$ to $b$ before, then we can instead replace that part with the rest of the cycle. So $G$ is not minimally connected. □

Proposition 6.6. Let $G$ be a connected simple graph with $n$ vertices. Then $G$ has at least $n - 1$ edges.

Proof. Consider deleting all of the edges of $G$ and then re-adding them one at a time. At first, we have $n$ connected components since all of the vertices are isolated. Every time we add an edge, we are either connecting two vertices in the same component (in which case the number of components stays the same), or we are connecting two vertices in different components (in which case the two components get merged into one). We conclude that to go from $n$ connected components to 1 connected component, we will need to add at least $n - 1$ edges. □

Corollary 6.7. Let $F$ be a forest with $n$ vertices and $m$ edges. Then it has $n - m$ connected components.

Proof. Go through the proof above. Each time we add an edge, we are decreasing the number of connected components by 1 (otherwise we’d introduce a cycle). □

Theorem 6.8. Let $G$ be a connected simple graph with $n$ vertices. Then $G$ is a tree if and only if $G$ has exactly $n - 1$ edges.

Proof. Suppose that $G$ is a tree. In particular, it is a forest so by the previous corollary, it has $n - m = 1$ connected components where $m$ is the number of edges. Solving for $m$ gives $m = n - 1$. 

□
Now suppose that $G$ has exactly $n - 1$ edges. Pick one and remove it. By the previous proposition, the result cannot be connected since it has too few edges. This means that $G$ is minimally connected, and hence a tree.

We have given several different properties of connected simple graph to be a tree. Here is the summary:

**Theorem 6.9.** Let $G$ be a connected simple graph with $n$ vertices. The following conditions are all equivalent:

1. $G$ is a tree.
2. $G$ has no cycles.
3. For all vertices $x, y \in V$, there is exactly one path starting at $x$ and ending at $y$.
4. $G$ is minimally connected.
5. $G$ has $n - 1$ edges.

In other words, if any of these properties hold, then all of them hold. If any of them fail, then they all fail.

6.2. **Adjacency matrix.** 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 $a_{ij}$ 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 $a_{ij}$. Note that $A_G$ is a symmetric matrix since $a_{ij} = a_{ji}$.

Going back to the Königsberg graph, it has the following adjacency matrix (using the ordering $A, B, C, D$):

$$
\begin{pmatrix}
0 & 1 & 2 & 0 \\
1 & 0 & 1 & 1 \\
2 & 1 & 0 & 2 \\
0 & 1 & 1 & 0
\end{pmatrix}
$$

If $G$ is a directed graph with vertices $v_1, \ldots, v_n$, then set $a_{ij}$ 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 $a_{ij}$. In general, $A_G$ need not be symmetric.

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

$$
v_1 \rightarrow v_2 \\
v_3 \\
\begin{pmatrix}
0 & 1 & 1 \\
0 & 0 & 1 \\
0 & 0 & 0
\end{pmatrix}
$$

**Theorem 6.10.** 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 $k > 0$, the number of walks of length $k$ starting at $v_i$ and ending at $v_j$ is $(A^k)_{i,j}$.

**Proof.** We prove this by induction on $k$. 
When $k = 1$, the $(i,j)$ entry of $A$ is defined so that it is the number of edges going from $v_i$ to $v_j$ (in both the directed and undirected cases).

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

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

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 $k + 1$. \hfill \Box

**Remark 6.11.** 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 $A_G = BDB^{-1}$ where $D$ is a diagonal matrix whose entries are the eigenvalues of $A_G$. In particular, $A_G^k = BD^k B^{-1}$. 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. \hfill \Box

**6.3. Deletion-contraction.** Let $G$ be a graph and $e$ an edge of $G$ which is not a loop. There are two important operations (deletion and contraction) that we can perform on $G$ using $e$ and which are useful for certain kinds of induction proofs.

The **deletion** of $e$ is denoted $G \setminus e$ and is a graph with the same vertices as $G$, and the same edges, except we don’t use $e$.

The **contraction** of $e$ is denoted $G/e$. Let $e = \{x,y\}$. To define it, take the vertices of $G$, replace the two vertices $x,y$ with a single vertex that we will call $z$. For each edge in $G$ that does not use $x$ or $y$, add it into $G/e$. For each vertex $a$ different from $x$ and $y$, the number of edges between $a$ and $z$ in $G/e$ is the number of edges between $a$ and $x$ plus the number of edges between $a$ and $y$. The number of loops at $z$ is the number of loops at $x$ plus the number of loops at $y$ plus the number of edges between $x$ and $y$ different from $e$.

To visualize this, pretend we are shrinking $e$ until $x$ and $y$ become the same point (hence the use of the word contraction).

Here’s a small example to illustrate. Say our graph is as follows (I put numbers on the edges to denote multiple edges):

$$
\begin{array}{c}
v_3 & v_4 \\
4 & 3 & 2 \\
v_1 & 3 & \cdots & v_2
\end{array}
$$
Let $e$ be one of the edges between the bottom two vertices. Then

$$G \setminus e = \begin{array}{c|c|c|c}
 v_1 & v_2 & v_3 & v_4 \\
 4 & 3 & 2 & \end{array}$$

Visually, $G/e$ is the result of shrinking the bottom edge of $G$ towards its midpoint. The two other bottom edges become loops.

6.4. **Matrix-tree theorem.**

**Definition 6.12.** Let $G = (V, E)$ be a graph. A subgraph of $G$ is a graph $G' = (V', E')$ such that $V' \subseteq V$ and $E' \subseteq E$.

A subgraph $G'$ is called a spanning tree of $G$ if $V' = V$ and $G'$ is a tree. □

Let $\tau(G)$ (that letter is TAU) be the number of spanning trees of $G$. Our goal in this section is to find a formula for $\tau(G)$.

Let $v_1, \ldots, v_n$ be the vertices of $G$. Define $D_G$ be the $n \times n$ diagonal matrix whose $i$th diagonal is the number of edges connected to $v_i$ (loops only get counted once). The Laplacian matrix of $G$ is defined to be $L_G = D_G - A_G$. Next, given any $n \times n$ matrix $B$ and $1 \leq i \leq n$, let $B[i]$ be the matrix obtained from $B$ by deleting the $i$th column and $i$th row.

**Theorem 6.13** (Kirchhoff’s matrix tree theorem). For any $1 \leq i \leq n$, we have $\tau(G) = \det(L_G[i])$. (By convention, the determinant of the $0 \times 0$ matrix is 1.)

**Example 6.14.** Taking $G$ from the previous section, $A_G$ and $L_G$ are as follows:

$$G = \begin{array}{c|c|c|c}
 v_1 & v_2 & v_3 & v_4 \\
 4 & 3 & 2 & \end{array}$$

$$A_G = \begin{bmatrix}
 0 & 3 & 4 & 3 \\
 3 & 0 & 0 & 2 \\
 4 & 0 & 0 & 0 \\
 3 & 2 & 0 & 0 \\
\end{bmatrix}, \quad L_G = \begin{bmatrix}
 10 & -3 & -4 & -3 \\
 -3 & 5 & 0 & -2 \\
 -4 & 0 & 4 & 0 \\
 -3 & -2 & 0 & 5 \\
\end{bmatrix}$$

Taking $i = 1$ in the theorem, the determinant of the bottom right $3 \times 3$ matrix is 84. We can also check directly that $G$ has 84 spanning trees. □

**Remark 6.15.** Loops are annoying to deal with in this definition, but they are forced on us in some situations when we contract an edge. A spanning tree cannot use any loops, so for the purposes of computing spanning trees, we can delete any loops a graph has and forget about them.

Although it sounds more general to let $i$ be anything between 1 and $n$, if we remember that we’re allowed to order the vertices any way we please, then it’s enough to prove the
theorem when \( i = 1 \) since the taking any other value of \( i \) would correspond to \( i = 1 \) for a different way to label the vertices. We’ll usually assume \( i = 1 \) since it makes the notation a little simpler to explain (but remember it’s not less general).

To prove the matrix-tree theorem, we will prove that both sides satisfy the same recurrence, namely a deletion-contraction type recurrence.

**Proposition 6.16.** Let \( e \) be an edge \( \{v_1, v_r\} \) for some \( 2 \leq r \leq n \). Then

\[
\det(L_G[1]) = \det(L_{G\setminus e}[1]) + \det(L_{G/e}[1]).
\]

When we contract the edge \( e \), the vertices \( v_1, v_r \) become the same vertex (which we will name \( z \)) and we use the ordering for \( G/e \) where \( z \) comes before all of the other vertices.

**Example 6.17.** Continue taking \( G \) as in the previous example and \( r = 2 \), let \( e \) be one of the edges connecting \( v_1 \) and \( v_2 \) as in the previous section. Here are the matrices mentioned in the proposition:

\[
L_G[1] = \begin{bmatrix} 5 & 0 & -2 \\ 0 & 4 & 0 \\ -2 & 0 & 5 \end{bmatrix}, \quad L_{G\setminus e}[1] = \begin{bmatrix} 4 & 0 & -2 \\ 0 & 4 & 0 \\ -2 & 0 & 5 \end{bmatrix}, \quad L_{G/e}[1] = \begin{bmatrix} 4 & 0 \\ 0 & 5 \end{bmatrix}
\]

We can check by computing that their determinants are 84, 64, 20, so the equation in the proposition checks for this example. However, here is a better way to check the equation, which can be used to prove the identity in general. Introduce the following matrix

\[
M = \begin{bmatrix} 1 & 0 & 0 \\ 0 & 4 & 0 \\ -2 & 0 & 5 \end{bmatrix}.
\]

The general pattern for \( M \): we replace the first row of \( L_G[1] \) by \((1, 0, 0, \ldots, 0)\). Then \( \det M = \det(L_{G\setminus e}[1]) \) which we can see by expanding along the first row. Furthermore, \( L_G[1], L_{G\setminus e}[1], L_{G/e}[1] \) agree in all rows except the first, and the first row of \( L_G[1] \) is the sum of the first rows of the other two. By the multilinearity property of determinants, this tells us that \( \det L_G[1] = \det L_{G\setminus e}[1] + \det M \).

**Proof.** We first assume \( r = 2 \). Define an \((n-1) \times (n-1)\) matrix \( M \) as follows. The bottom right \((n-2) \times (n-2)\) matrix will be a copy of \( L_{G\setminus e}[1] \). The first row will be \((1, 0, 0, \ldots, 0)\). The remaining entries are \( M_{i,1} \) where \( 2 \leq i \leq n - 1 \) and we set it to be \( L_{G\setminus e}[1]_{i,1} \). Then \( \det(M) = \det(L_{G\setminus e}[1]) \) (see this by expanding \( \det(M) \) along the first row).

Next, we claim that \( M, L_{G\setminus e}[1], L_G[1] \) are all identical in rows 2, \ldots, \( n - 1 \). This is clear outside of the first column because they all record information about vertices \( v_3, \ldots, v_n \) which are exactly the same in \( G, G \setminus e, G/e \). For the first column, these involve information about edges between \( v_2 \) and \( v_3, \ldots, v_n \). Again, this is the same between \( G \) and \( G \setminus e \) and we have defined \( M \) to have the same entries as \( L_{G\setminus e}[1] \).

Finally, the first row of \( L_G[1] \) is the sum of the first row of \( M \) and the first row of \( L_{G\setminus e}[1] \) (I’ll leave this to you to fill in the details why). By multilinearity of determinants, this implies

\[
\det(L_G[1]) = \det(L_{G\setminus e}[1]) + \det(M).
\]

But \( \det M = \det(L_{G\setminus e}[1]) \), so we get the identity we want.

For the case of general \( r \), the same argument works (though the notation is a little more complicated), or we note that swapping the rows 2 and \( r \) and also columns 2 and \( r \) does not affect the determinant. \( \square \)
Proposition 6.18. Let $e$ be any edge of $G$. Then
\[ \tau(G) = \tau(G \setminus e) + \tau(G/e). \]

Proof. Write $e = \{x, y\}$. $\tau(G \setminus e)$ counts the number of spanning trees in $G$ that do not use the edge $e$ while $\tau(G/e)$ counts the number of spanning trees in $G$ that do use the edge $e$.

The second requires some more explanation. First, let $n$ be the number of vertices of $G$ and let $z$ be the new vertex formed by merging $x$ and $y$. If we have a spanning tree $T$ of $G$ that uses $e$ and we contract $e$, the remaining edges of $T$ become a spanning tree of $G/e$: $T$ has $n - 1$ edges and is connected, $T/e$ has $n - 2$ edges and $n - 1$ vertices and is still connected. We can reverse this: by the way we defined it, there is a bijection between the edges of $G/e$ that aren’t loops at $z$ and the edges of $G$ whose endpoints aren’t $\{x, y\}$. So if we have a spanning tree $T'$ of $G/e$, take the edges in $G$ that correspond to the edges of $T'$ and add $e$ to get a spanning tree of $G$. We just defined a bijection between spanning trees of $G$ using $e$ and spanning trees of $G/e$. One thing to note: spanning trees never use more than one edge with the same endpoints, and never use loops; edges with the same endpoints as $e$ correspond to loops in $G/e$.

Every spanning tree of $G$ either uses $e$ or doesn’t, so we get the desired identity. $\Box$

Now we can finish the proof.

Proof of Theorem 6.13. We prove the theorem by induction on the number of non-loop edges of $G$. For the base case, let $G$ be a graph with no non-loop edges. If the number of vertices is 1, then $\tau(G) = 1$, and if it has more than 1 vertex, then $\tau(G) = 0$. In either case, $L_G$ is the $n \times n$ matrix consisting of 0’s. If $n = 1$, then $L_G[1]$ is the 0 $\times$ 0 matrix and so has determinant 1, which is what we want. Otherwise, $L_G[i]$ is a 0 matrix of positive size so has determinant 0 for each $i$, again which is what we want.

Now assume that we have proven the statement for all graphs with $\leq n$ edges, and let $G$ be a graph with $n + 1$ edges. If $v_1$ has no non-loop edges touching it, then $G$ is disconnected and hence has no spanning trees. On the other hand, $L_G[1]$ is the Laplacian for the graph obtained from $G$ by deleting $v_1$. In particular, the sum of its rows is 0, so $\det(L_G[1]) = 0$. Otherwise, there is some $r$ so that $\{v_1, v_r\} = e$ is an edge in $G$.

Then both $G/e$ and $G \setminus e$ have less non-loop edges than $G$, so know by induction that
\[ \tau(G \setminus e) = \det(L_{G\setminus e}[1]), \quad \tau(G/e) = \det(L_{G/e}[1]). \]

However, we have previously proven the identities
\[ \tau(G) = \tau(G \setminus e) + \tau(G/e), \quad \det(L_G[1]) = \det(L_{G\setminus e}[1]) + \det(L_{G/e}[1]). \]

Combining this with the previous identities gives us
\[ \tau(G) = \det(L_G[1]). \]$

There is also a version of the matrix-tree theorem that uses eigenvalues instead of determinants. Below we will use the fact that if $G$ is connected, then exactly one of the eigenvalues of $L_G$ is 0.

Theorem 6.19 (Matrix tree theorem, eigenvalue version). Let $G$ be a connected graph with $n$ vertices and let $\lambda_1, \ldots, \lambda_{n-1}$ be the nonzero eigenvalues of $L_G$. Then $\tau(G) = \lambda_1 \lambda_2 \cdots \lambda_{n-1}/n$.

Remark 6.20. This is a purely linear algebra proof once we have the determinant version of the matrix-tree theorem. Since this isn’t a linear algebra class, I’ll just tell you what goes into it. First, given a square matrix $A$ with eigenvalues $\mu_1, \ldots, \mu_n$, a principal minor of size
$k$ is the determinant of a $k \times k$ square submatrix where you choose the same indices for the rows and columns. A general fact is that the sum of all principal minors of size $k$ is the sum of $\mu_{i_1} \mu_{i_2} \cdots \mu_{i_k}$ over all $k$-element subsets $\{i_1, \ldots, i_k\}$ of $[n]$ (this generalizes the statement that the trace of $A$ is the sum of the eigenvalues when $k = 1$ and that the determinant of $A$ is the product of all of them when $k = n$). Now apply this to $A = L_G$ and $k = n - 1$. All of the principal minors are $\det L_G[i]$ for some $i$ and each one is $\tau(G)$ by the matrix-tree theorem. The sum of the product of eigenvalues (over all choices of $n - 1$ of them) is $\lambda_1 \cdots \lambda_{n-1}$ (all other subsets use the 0 eigenvalue) so we get $n \tau(G) = \lambda_1 \cdots \lambda_{n-1}$.

We finish with two general examples.

**Example 6.21.** Let $G = K_n$ be the complete graph on $n$ vertices $v_1, \ldots, v_n$. This means that there is exactly 1 edge between any pair of distinct vertices. The number of spanning trees counts the number of labeled trees: these are trees on $n$ vertices whose vertices are numbered $1, \ldots, n$. We will use both versions of the matrix-tree theorem to compute $\tau(K_n)$.

First, $L_G$ has $n - 1$ in the diagonals and $-1$ in all off-diagonal entries. Now delete the first row and column. The result is an $(n - 1) \times (n - 1)$ matrix which looks like:

$$L_G[1] = \begin{bmatrix}
1 & 1 & 1 & \cdots & 1 \\
1 & n-1 & 1 & \cdots & 1 \\
1 & 1 & n-1 & \cdots & 1 \\
\vdots & & & & \\
-1 & -1 & -1 & \cdots & n-1
\end{bmatrix}$$

Remember that adding a row to another row does not change the determinant of a matrix. Our goal is to do these row operations to make it so the first column only has one nonzero entry. Add each of the rows $2, \ldots, n-1$ to the first row to get the matrix:

$$\begin{bmatrix}
1 & 1 & 1 & \cdots & 1 \\
1 & n-1 & 1 & \cdots & 1 \\
1 & 1 & n-1 & \cdots & 1 \\
\vdots & & & & \\
-1 & -1 & -1 & \cdots & n-1
\end{bmatrix}$$

Finally, add the first row to all of the other rows to get

$$\begin{bmatrix}
1 & 1 & 1 & \cdots & 1 \\
0 & n & 0 & \cdots & 0 \\
0 & 0 & n & \cdots & 0 \\
\vdots & & & & \\
0 & 0 & 0 & \cdots & n
\end{bmatrix}$$

Now expand the determinant along the first column to see that $\tau(G) = n^{n-2}$. This is also known as Cayley’s theorem.

Now let’s compute $\tau(G)$ using the eigenvalues of $L_G$. Instead of computing them from the characteristic polynomial, let’s instead find a basis of eigenvectors. Let $e_i$ be the column vector with a 1 in position $i$ and 0’s elsewhere. Then the following are all eigenvectors:

$$e_1 + e_2 + \cdots + e_n, e_1 - e_2, e_2 - e_3, e_3 - e_4, \ldots, e_{n-1} - e_n.$$ 

The first one has eigenvalue 0 and the others have eigenvalue $n$. We can also check they are linearly independent (but I’ll leave the details to you). So this tells us that the eigenvalues
are 0 and n (multiplicity \( n - 1 \)). The eigenvalue version then tells us that \( \tau(G) = \frac{n^{n-1}}{n} = n^{n-2} \).

\[ \square \]

**Example 6.22.** For our second example, pick two integers \( n, m \geq 1 \) and let \( G = K_{m,n} \) be the complete bipartite graph. This has vertices \( v_1, \ldots, v_m \) and \( w_1, \ldots, w_n \) and there is one edge between \( v_i \) and \( w_j \) for all choices of \( i \) and \( j \), and no other edges. (This is an example of a bipartite graph, which we will introduce later.) To write down the adjacency matrix, we will order the vertices as \( v_1, \ldots, v_m, w_1, \ldots, w_n \). Then \( L_G \) has a natural block structure:

\[
L_G = \begin{bmatrix}
  n & 0 & \cdots & 0 & -1 & -1 & \cdots & -1 \\
  0 & n & \cdots & 0 & -1 & -1 & \cdots & -1 \\
  \vdots \\
  0 & 0 & \cdots & n & -1 & -1 & \cdots & -1 \\
  -1 & -1 & \cdots & -1 & m & 0 & \cdots & 0 \\
  -1 & -1 & \cdots & -1 & 0 & m & \cdots & 0 \\
  \vdots \\
  -1 & -1 & \cdots & -1 & 0 & 0 & \cdots & m \\
\end{bmatrix}
\]

This is a matrix of size \((m + n) \times (m + n)\). It will be convenient to delete the last row and column instead of the first. It looks the same, except now there are only \( n - 1 \) \( m \)'s in the bottom right corner. We proceed as before. Add rows 2, \ldots, \( m + n - 2 \) to the first row:

\[
\begin{bmatrix}
  1 & 1 & \cdots & 1 & 0 & 0 & \cdots & 0 \\
  0 & n & \cdots & 0 & -1 & -1 & \cdots & -1 \\
  \vdots \\
  0 & 0 & \cdots & n & -1 & -1 & \cdots & -1 \\
  -1 & -1 & \cdots & -1 & m & 0 & \cdots & 0 \\
  -1 & -1 & \cdots & -1 & 0 & m & \cdots & 0 \\
  \vdots \\
  -1 & -1 & \cdots & -1 & 0 & 0 & \cdots & m \\
\end{bmatrix}
\]

Expand along the first column to get \( \tau(G) = \frac{n^{m-1}m^{n-1}}{n} \). We can also use eigenvalues. Again, let \( e_i \) be the column vector with 1 in position \( i \) and 0's elsewhere. The following are an eigenbasis for \( L_G \):

\[
\sum_{i=1}^{n+m} e_i, \sum_{i=1}^{n} e_i - \sum_{j=n+1}^{n+m} e_j, \{e_i - e_{i+1} \mid 1 \leq i \leq m - 1\}, \{e_j - e_{j+1} \mid m + 1 \leq j \leq m + n\}.
\]
I’ll leave the details of computing the eigenvalues of these vectors and showing they are
linearly independent to you.

6.5. Minimum-weight spanning trees.

**Proposition 6.23.** If $G$ is connected, then it has at least one spanning tree.

*Proof.* If $G$ is a tree, then we are done. Otherwise, $G$ is not minimally connected, so we can
remove some edge in such a way that the resulting graph is still connected. Continue doing
this until we get down to a tree.

Suppose we are given a graph $G = (V, E)$ together with a real-valued function on the edge
set $w: E \to \mathbb{R}_{\geq 0}$ which takes non-negative values. Given a spanning tree $T \subseteq G$, define the
**weight** $w(T)$ to be the sum of $w(e)$ over all edges $e \in T$. Our goal is to find a tree that
has the smallest possible weight (there could be multiple solutions). Such a tree is called a
**minimum-weight spanning tree.**

**Example 6.24.** The vertices might represent cities and the edges represent possible railroad
tracks that could be built between them. The function $w$ could represent the cost for building
that track. If we wanted to connect all of the cities by rail by spending as little money as
possible, then we want a minimum-weight spanning tree.

[⋆ Steven: put in picture ⋆]

Here is Kruskal’s greedy algorithm. Let $G = (V, E)$ be a connected graph with $n$ vertices.

- Set $T = (V, \emptyset)$.
- If $T$ is a spanning tree, then terminate. Otherwise, consider all possible edges that
we can add to $T$ so that the result does not have a cycle and add one that has the
smallest possible weight to $T$. Repeat this step until $T$ is a spanning tree.

We have to be careful that the second step is well-defined. At each step, we have some
subgraph $T$ that has no cycles and isn’t a spanning tree. Why is it always possible that we
can find an edge to add to it so that we don’t create a cycle? (Left for homework.)

**Theorem 6.25.** The result of Kruskal’s greedy algorithm is a minimum-weight spanning
tree.

In order to prove this, we need to know the “exchange property”.

**Lemma 6.26** (Exchange property). Let $F$ and $F'$ be forests with the same set of vertices.
Assume that $F'$ has strictly more edges than $F$ does. Then there exists an edge $e$ of $F'$ not
belonging to $F$ such that adding $e$ to $F$ results in a forest.

*Proof.* Suppose this is false. Then for all edges $e$ of $F'$, $F \cup \{e\}$ has a cycle. In particular, if $x$
and $y$ are connected by an edge in $F'$, then they are in the same connected component of $F$.
Being in the same connected component is a transitive property, so this also implies that if $x$
and $y$ are in the same connected component of $F'$, then they are also in the same connected
component of $F$. Let $c, c'$ be the number of connected components of $F, F'$, respectively.
This means that $c' \geq c$. However, for a forest, the number of connected components is
the number of vertices minus the number of edges by Corollary 6.7. If $n$ is the number of
vertices, then we see that $n - c' \leq n - c$ but this contradicts that $F'$ has strictly more edges
than $F$.

Now we can prove the theorem about Kruskal’s greedy algorithm.
Proof of Theorem 6.25. Let $T$ be the output of Kruskal’s greedy algorithm. Suppose it is not a minimum-weight spanning tree. Then there is another spanning tree $H$ such that $w(H) < w(T)$. Let $h_1, \ldots, h_{n-1}$ be the edges of $H$ ordered so that $w(h_1) \leq \cdots \leq w(h_{n-1})$ and similarly let $t_1, \ldots, t_{n-1}$ be the edges of $T$ ordered so that $w(t_1) \leq \cdots \leq w(t_{n-1})$.

Let $S$ be the set of $1 \leq k \leq n - 1$ such that $\sum_{i=1}^{k} w(h_i) < \sum_{i=1}^{k} w(t_i)$. Note that $S$ is non-empty because $n-1 \in S$. Let $m$ be the minimal element of $S$. Note that $m > 1$ because otherwise it means that $w(h_1) < w(t_1)$ which contradicts how we picked the first edge in the algorithm. By definition of $m$, we have

$$\sum_{i=1}^{m-1} w(h_i) \geq \sum_{i=1}^{m-1} w(t_i) \quad \text{and} \quad \sum_{i=1}^{m} w(h_i) < \sum_{i=1}^{m} w(t_i).$$

This is only possible if $w(h_m) < w(t_m)$.

Let $T'$ be the forest with all vertices of $G$ and the edges $t_1, \ldots, t_{m-1}$ and let $H'$ be the forest with all vertices of $G$ and the edges $h_1, \ldots, h_m$. By the exchange property, there is some $h_i$ such that $h_i$ is different from the $t_1, \ldots, t_{m-1}$ and the edges $t_1, \ldots, t_{m-1}, h_i$ form a forest, i.e., has no cycle. However, $w(h_i) \leq w(h_m) < w(t_m)$, which contradicts our choice of edge to add in the $m$th step of the greedy algorithm. In particular, $H$ cannot exist, so $T$ is a minimum-weight spanning tree. □

Remark 6.27. The context of the greedy algorithm can be expanded quite a bit. The keypoint is the exchange property. Here is another familiar situation where it holds: let $S, T$ be linearly independent sets of vectors in a vector space and assume that $|T| > |S|$. Then there is always a choice of vector $v \in T$ such that $S \cup \{v\}$ is linearly independent (try to prove it).

In particular, imagine we had a collection of vectors that spans our vector space and each one has a weight and our goal is to find a basis from this collection that minimizes the total weight. The greedy algorithm, adapted suitably, tells us that we can successively pick vectors of smallest possible weight that don’t create a linear dependence with the vectors we’ve already chosen and the result will have the desired property.

Structures that satisfy the exchange property are called matroids. □

7. Coloring and matching

7.1. Colorings and chromatic polynomials.

Definition 7.1. Let $G$ be a graph and $k$ a positive integer. A proper $k$-coloring is a function $f : V \rightarrow [k]$ such that for each pair $\{v, v'\}$, $f(v) \neq f(v')$.

The chromatic number $\chi(G)$ is the smallest $k$ such that $G$ has a proper $k$-coloring.

We let $\chi_G(k)$ be the number of proper $k$-colorings of $G$. □

Heuristically, we think of the numbers $1, \ldots, k$ as colors, like red, blue, green, etc. A proper $k$-coloring is a way to assign colors so that adjacent vertices never have the same color. We are free to use colors multiple times and we don’t have to use all of them.

Note that if $G$ has loops, then it has no proper colorings, so this section will only be interesting for loopless graphs.

As an example, the vertices of our graph could be wireless routers and an edge is present if the routers are close enough to interfere if they broadcast at the same frequency. So the colors would be different frequencies to use and a proper coloring corresponds to assignments
so that no interference occurs, and the chromatic number is the least number of different frequencies that we need to use.

**Example 7.2.** Here is a proper 3-coloring of a length 5 cycle:

![Cycle Graph](image)

Note that a proper 2-coloring isn’t possible. So this graph has chromatic number 3. □

Here are some easy properties of chromatic numbers. They all follow from the definitions, so I won’t write out the proofs.

**Proposition 7.3.** (1) If $G$ has a proper $k$-coloring, then $\chi(G) \leq k$.

(2) If $G$ has a clique of size $r$, meaning there are $r$ vertices of $G$ all of which are connected to each other by an edge, then $\chi(G) \geq r$.

(3) If $H$ is a subgraph of $G$, then $\chi(H) \leq \chi(G)$.

(4) If $G$ has $n$ vertices (and no loops), then $\chi(G) \leq n$.

**Lemma 7.4.** Let $x, y$ be two vertices of $G$ with exactly one edge $e$ between them. Then

$$\chi_G(k) = \chi_{G \setminus e}(k) - \chi_{G/e}(k).$$

**Proof.** By the definitions, a proper $k$-coloring of $G$ is the same thing as a proper $k$-coloring of $G \setminus e$ where $x$ and $y$ get different labels. On the other hand, proper $k$-colorings of $G \setminus e$ where $x$ and $y$ receive the same color are naturally in bijection with proper $k$-colorings of $G/e$: if $z$ is the result of contracting $x$ and $y$, make its color the common color of $x$ and $y$. The identity $\chi_G(k) = \chi_{G \setminus e}(k) - \chi_{G/e}(k)$ is a translation of what we just said: proper $k$-colorings of $G$ are the same thing as proper $k$-colorings of $G \setminus e$ once we subtract off all of those that give $x$ and $y$ the same color. □

The assumption about $x$ and $y$ having exactly one edge between them is a little bit annoying, but it’s easy to get around. Let $G$ be a graph. Construct a simple graph $\overline{G}$ whose vertices are the same as $G$ and where $x$ and $y$ have an edge in $\overline{G}$ if they have at least one edge in $G$. In other words, multiple edges in $G$ get replaced by a single edge in $\overline{G}$.

**Lemma 7.5.** $\chi_G(k) = \chi_{\overline{G}}(k)$.

**Proof.** The definition of proper $k$-coloring only involves labeling vertices and the conditions on them only depend on whether or not two vertices have the same color if they’re connected by an edge (but we don’t care how many edges). □

Now we’re ready to prove the main result:

**Theorem 7.6.** If $G$ is a graph with $n$ vertices, then $\chi_G(k)$ is a polynomial in $k$ of degree $n$ (more precisely, there is a unique polynomial of degree $n$ whose values agree with $\chi_G(k)$ at all non-negative integer inputs $k$).
Proof. By Lemma 7.5, it is enough to prove this for simple graphs \( G \). We proceed by induction on the number of edges. If there are no edges in \( G \), then any labeling of the vertices is a proper \( k \)-coloring, so \( \chi_G(k) = k^n \) which is certainly a polynomial of degree \( n \).

Now assume we've proved this for graphs with \( < m \) edges and let \( G \) be a graph with \( m \) edges. Let \( e \) be an edge of \( G \). Then \( G \setminus e \) and \( G/e \) both have \( < m \) edges. So \( \chi_{G \setminus e}(k) \) is a polynomial in \( k \) of degree \( n \) and \( G/e \) is a polynomial in \( k \) of degree \( n - 1 \). By Lemma 7.4, \( \chi_G(k) = \chi_{G \setminus e}(k) - \chi_{G/e}(k) \), so \( \chi_G(k) \) is a polynomial in \( k \) of degree \( n \).

For notation, we will write \( \chi_G(z) \) for this polynomial (\( z \) is now a variable) and we will use \( k \) to denote non-negative integers. This is the chromatic polynomial of \( G \). Then the chromatic number is the smallest positive integer \( k \) such that \( \chi_G(k) \neq 0 \).

Here are some easy properties:

**Proposition 7.7.** (1) \( G \) has at least one vertex if and only if \( \chi_G(0) = 0 \).

(2) \( G \) has at least one edge if and only if \( \chi_G(1) = 0 \). (The converse is clearly true.)

(3) If \( G \) has an odd length cycle, then \( \chi_G(2) = 0 \). (The converse is also true, as we will see when we discuss bipartite graphs.)

How about a property that determines if \( \chi_G(3) = 0 \)? This is an NP-complete problem, so there likely isn’t a simple criterion to determine this for a general graph.

**Example 7.8.** Let’s compute \( \chi_G(z) \) for the square:

\[
\begin{array}{c}
2 & 3 \\
\hline
1 & 4
\end{array}
\]

It will follow that \( \chi(G) = 2 \) (or you can figure that out by staring at the square). Here are some different approaches:

(1) For the first way, we just use the definition. If we want to properly \( k \)-color \( G \), then 1 can be colored anything, so there are \( k \) choices for it. Now the color on 2 and 4 have to be different from the color assigned to 1, so there are \( k - 1 \) choices for each. There are two cases to consider: if the colors of 2 and 4 are the same, then the color for 3 has \( k - 1 \) choices. If they’re not the same, then the color for 3 has \( k - 2 \) choices. So the total number of colorings is: \( k(k-1)^2 + k(k-1)(k-2)^2 \). (The first term counts the number of colorings where 2 and 4 have the same color and the second counts the number of colorings where 2 and 4 have different colors.) We can simplify it to get

\[
\chi_G(k) = k(k-1)(k^2 - 3k + 3).
\]

(2) For the second way, we’ll use deletion-contraction. Let \( e = \{1, 4\} \). Then

\[
\begin{array}{c}
2 & 3 \\
\hline
1 & 4
\end{array}
\]
Its chromatic polynomial is simple to compute: for a proper \( k \)-coloring, 1 has \( k \) choices, 2 has \( k-1 \) choices (any color different from the one given to 1), similarly 3 has \( k-1 \) choices, and similarly, 4 has \( k-1 \) choices. So
\[
\chi_{G \setminus e}(k) = k(k-1)^3.
\]
The contraction by \( e \) is

\[
\begin{array}{c}
\begin{tikzpicture}
  \node (1) at (0,0) {2};
  \node (3) at (1,0) {3};
  \node (4) at (0,-1) {5};
  \draw (1) -- (3);
  \draw (1) -- (4);
  \draw (3) -- (4);
\end{tikzpicture}
\end{array}
\]

I called the new vertex 5. This is also easy to compute: for a proper \( k \)-coloring, 5 has \( k \) choices, 2 has \( k-1 \) choices, and 3 has \( k-2 \) choices (any color different from the one given to 2 and 5 which are different from each other). So
\[
\chi_{G/e}(k) = k(k-1)(k-2).
\]

So using Lemma 7.4, we get
\[
\chi_G(k) = \chi_{G \setminus e}(k) - \chi_{G/e}(k)
= k(k-1)^3 - k(k-1)(k-2)
= k(k-1)(k^2 - 3k + 3).
\]

\[\square\]

**Example 7.9.** Here are some families of graphs where we can give explicit formulas for \( \chi_G(z) \) and \( \chi(G) \). I won’t explain how to get the derivation, you should see if you can figure out how to do it.

1. The **complete graph** on \( n \) vertices is denoted \( K_n \) and is defined so that every pair of vertices has an edge between them. Then
\[
\chi_{K_n}(z) = z(z-1)(z-2) \cdots (z-n+1),
\]
\[
\chi(K_n) = n.
\]

2. The cycle \( C_n \) of length \( n \) has vertices \( v_1, \ldots, v_n \) and edges \( \{i, i+1\} \) for \( i = 1, \ldots, n-1 \) and \( \{1, n\} \). Then
\[
\chi_{C_n}(z) = \begin{cases} 
(z-1)^n + (z-1) & \text{if } n \text{ is even} \\
(z-1)^n - (z-1) & \text{if } n \text{ is odd}
\end{cases},
\]
\[
\chi(C_n) = \begin{cases} 
2 & \text{if } n \text{ is even} \\
1 & \text{if } n = 1 \\
3 & \text{if } n \text{ is odd and } n \geq 3
\end{cases}.
\]

3. If \( G \) is a tree with \( n \) vertices, then
\[
\chi_G(z) = z(z-1)^{n-1}
\]
\[
\chi(G) = \begin{cases} 
1 & \text{if } n = 1 \\
2 & \text{if } n > 1
\end{cases}.
\]
\[\square\]
7.2. Bipartite graphs.

**Definition 7.10.** A graph $G$ is **bipartite** if it has a proper 2-coloring. Equivalently, $G$ is bipartite if there is a subset $X \subseteq V$ such that no 2 vertices in $X$ are joined by an edge, and no 2 vertices in $Y = V \setminus X$ are joined by an edge. We will call $(X,Y)$ a **bipartition** of $G$ if this holds.

Note that a bipartite graph could have many bipartitions. For example, if $G$ has no edges at all, then any choice of $X$ will work.

**Example 7.11.** Every forest is bipartite. A cycle of even length is bipartite. However, a cycle of odd length is not bipartite.

Given two vertices $v, w$, let $d(v, w)$ be the smallest length of a path from $v$ to $w$. This is the **distance** between $v$ and $w$. We can give a nice criterion for a graph to be bipartite. But first we need a preliminary lemma.

**Lemma 7.12.** Suppose a graph $G$ has a closed walk of odd length. Then $G$ has a cycle of odd length.

**Proof.** Pick a closed walk of odd length of shortest possible size $v_0, e_1, v_1, \ldots, v_k$. Suppose that this is not a cycle, so that there are $0 < i < j < k$ such that $v_i = v_j$. Then $v_i, e_{i+1}, v_{i+1}, \ldots, e_j, v_j$ is a closed walk and so is the result of deleting $v_i, e_{i+1}, v_{i+1}, \ldots, e_j$ from the first walk. Since both of their lengths add up to $k$, which is odd, one of these closed walks must have odd length. But that contradicts the fact that our original choice is the shortest possible size.

Note that if you make the above claim with “odd” replaced with “even”, the proof above does not work (can you see where it goes wrong?). In fact, the modified statement is false: let $G$ be a triangle. Then it has a closed walk of length 6 (go around the triangle twice) but no cycle of even length.

**Theorem 7.13.** A graph $G$ is bipartite if and only if it does not contain any cycles of odd length.

**Proof.** If $G$ has a cycle of odd length, let $H$ be that cycle. We know that $\chi(H) = 3$, and $\chi(G) \geq \chi(H) = 3$, so $G$ cannot be bipartite.

Now assume that $G$ has no cycles of odd length. We will show that $G$ has a proper 2-coloring using “red” and “blue”. We first assume that $G$ is connected. Start by picking a vertex $v$. We color the vertices as follows: for each vertex $w$, if $d(v, w)$ is even, color $w$ red, and if $d(v, w)$ is odd, color $w$ blue. We claim that this is a proper 2-coloring. If not, then there is some edge $\{w, w'\}$ such that $w$ and $w'$ have the same color. Then we have a closed walk of odd length: take a path of shortest length from $v$ to $w$, then use the edge $\{w, w'\}$, then take a path of shortest length from $w'$ to $v$. These two paths have the same parity, so the total number of edges used is odd. From the previous lemma, this implies that $G$ has an odd cycle, so $\{w, w'\}$ can’t exist, which means we have a proper 2-coloring.

Finally, if $G$ is not connected, then none of its connected components has a cycle of odd length, so each of them has a proper 2-coloring by the above argument. Coloring each component then gives a proper 2-coloring of the whole graph.
7.3. Matchings.

**Definition 7.14.** Let \( G = (V, E) \) be a graph. A subset \( S \subseteq E \) is a **matching** if every vertex is the endpoint of at most one edge of \( S \). In other words, no two edges in \( S \) share a common vertex. A matching is a **perfect matching** if every vertex is the endpoint of exactly one edge of \( S \).

For example, let’s say the vertices are students in a class and edges represent pairs of students who are willing to work together on a project. A matching then consists of pairing off some of the students in teams of two (who are willing to work together), and a perfect matching means that we’ve paired off all of the students. We’ll be mostly interested in matchings in bipartite graphs. We have a modified definition in this case.

**Definition 7.15.** Let \( G \) be a bipartite graph with bipartition \((X, Y)\). A **perfect matching** of \( X \) into \( Y \) is a matching \( S \) of \( G \) such that every vertex of \( X \) is the endpoint of some edge in \( S \).

This situation comes up a lot. For example, \( X \) might represent doctors and \( Y \) might represent positions in hospitals. An edge between a doctor and a position exists if the hospital is willing to hire that person for that job. A perfect matching of \( X \) into \( Y \) in this situation means finding a job for each doctor. On the other hand, a perfect matching of \( Y \) into \( X \) would mean filling each job with a qualified doctor. Our goal is to give a theorem which guarantees the existence of a perfect matching of \( X \) into \( Y \).

**Definition 7.16.** Given a subset \( T \) of the vertices of a graph, define \( N_G(T) \), the neighbor set of \( T \), to be the set of vertices which are connected by an edge to some vertex in \( T \).

**Theorem 7.17** (Hall’s marriage theorem). Let \( G \) be a bipartite graph with bipartition \((X, Y)\). Then there is a perfect matching of \( X \) into \( Y \) if and only if for all subsets \( T \subseteq X \), we have \(|T| \leq |N_G(T)| \).

In our doctor example above, this says there is a perfect matching if for every group of \( k \) doctors, the number of jobs for which at least one of them is qualified is at least \( k \). This is clearly a necessary condition, so the real content is that this is also enough.

**Proof.** If there is a perfect matching of \( X \) into \( Y \), then clearly \(|T| \leq |N_G(T)| \) for all subsets \( T \subseteq X \): for \( t \in T \), let \( t' \in Y \) be the vertex matched with it. Then \( t' \in N_G(T) \) and \( t' \neq s' \) for \( t \neq s \) by definition of matching.

Now suppose that \(|T| \leq |N_G(T)| \) for all subsets \( T \subseteq X \). We need to show that \( G \) has a perfect matching of \( X \) into \( Y \). We will prove this statement by induction on \(|X| \). The base case is \(|X| = 1 \). In this case, \( X = \{x\} \) and our assumption says that \( x \) has at least one neighbor in \( Y \), so picking any of them gives the perfect matching. For the induction step, assume that the statement is true whenever \(|X| < d \) and consider the setup with \(|X| = d \). We consider two cases.

- **Case 1:** Suppose that for all subsets \( T \subset X \) with \( T \neq X \), we have \(|T| < |N_G(T)| \).

  In this case, pick any \( x \in X \) and then pick \( y \in Y \) such that \( \{x, y\} \) is an edge (this exists because \(|N_G(\{x\})| \geq 1 \). Let \( H \) be the subgraph of \( G \) with \( x, y \) deleted and \( \{x, y\} \) deleted. Then \( H \) is a bipartite graph with bipartition \((X', Y')\) where \( X' = X \setminus \{x\} \) and \( Y' = Y \setminus \{y\} \). For any subset \( T \subseteq X' \), we have \(|N_H(T)| \geq |N_G(T)| - 1 \) since, at worst, \( y \) might have been in \( N_G(T) \). However, \(|N_G(T) - 1 \geq |T| \) by the assumption in this case (note \( T \neq X \) because \( T \) does not contain \( x \)). So \(|N_H(T)| \geq |T| \), and by
induction, $H$ has a perfect matching of $X'$ into $Y'$. Now add the edge $\{x, y\}$ back to get a perfect matching of $X$ into $Y$.

- **Case 2:** There is some nonempty subset $B \subset X$ with $B \neq X$ and $|B| = |N_G(B)|$.

  In this case, form two graphs $H_1$ and $H_2$: $H_1$ is the subgraph of $G$ with vertices $B \cup N_G(B)$ and whatever edges go between them, and $H_2$ is the result of deleting $B \cup N_G(B)$ from $G$ along with any edges that touch those vertices.

  First, $H_1$ is bipartite with bipartition $(B, N_G(B))$. Since $B \neq X$, we have $|B| < d$. Furthermore, if $T \subseteq B$, then $N_{H_1}(T) = N_G(T)$, so by assumption we have $|T| \leq |N_{H_1}(T)|$ for all $T$. By induction, $H_1$ has a perfect matching of $B$ into $N_G(B)$.

  Next, $H_2$ is bipartite with bipartition $(X', Y')$ where $X' = X \setminus B$ and $Y' = Y \setminus N_G(B)$. Let $T \subseteq X'$ be a subset. Then $N_G(B \cup T) = N_G(B) \cup N_{H_2}(T)$ and $N_G(B) \cap N_{H_2}(T) = \emptyset$ by definition of $H_2$. Furthermore, $B \cap T = \emptyset$, so we have

  \[ |B| + |T| = |B \cup T| \leq |N_G(B \cup T)| = |N_G(B) + N_{H_2}(T)| = |N_G(B)| + |N_{H_2}(T)|. \]

  In particular, $|B| + |T| \leq |N_G(B)| + |N_{H_2}(T)|$. Since $|B| = |N_G(B)|$, we conclude that $|T| \leq |N_{H_2}(T)|$.

  So we have verified the hypothesis of the theorem for $H_2$. Since $B$ is nonempty, $|X'| < |X|$, so we may apply the induction hypothesis to conclude that $H_2$ has a perfect matching of $X'$ into $Y'$.

  Combining these two perfect matchings gives one for $G$. \qed

If know that a perfect matching of $X$ into $Y$ exists, how do we find one? Or more generally, how do we find a matching with as many edges as possible? We could check all possibilities, but let’s discuss one way which isn’t just brute force.

**Definition 7.18.** Let $G$ be a bipartite graph and let $M$ be a matching.

A path in $G$ is called **$M$-alternating** if its edges alternate between edges in $M$ and edges not in $M$ (we don’t specify which one it starts or ends with).

An $M$-alternating path is called **$M$-augmenting** if its first and last edges are not in $M$.

Finally, a matching is called a **maximum matching** if there is no other matching with more edges than it. \qed

**Theorem 7.19.** Let $G$ be a bipartite graph and let $M$ be a matching. Then $M$ is a maximum matching if and only if there are no $M$-augmenting paths.

**Proof.** First, given an $M$-augmenting path, we can construct a bigger matching: just swap whether or not the edges in this path are in $M$. So if $M$ is maximum, then there are no $M$-augmenting paths.

Now assume that there are no $M$-augmenting paths in $G$. We need to prove that $M$ is a maximum matching. Suppose that $M'$ is a maximum matching of $G$. We need to show that $M$ and $M'$ have the same number of edges. Let $H$ be the subgraph of $G$ with the same vertices but only the edges which appear in exactly one of $M$ or $M'$. Each connected component of $H$ has to be a path or a cycle (an isolated vertex is a length 0 cycle): Start with any edge (say in $M$), if there is another edge touching it, it comes from $M'$. Any new edge touching that one must be in $M$, etc. Eventually we can’t continue finding more edges, and then go back to the first edge and go in “reverse” if you can. You either end up with a path or something which has revisited a vertex. The second case has to be a cycle, otherwise you’d have 3 edges touching at one vertex which contradicts that both $M$ and $M'$ are matchings.
Since $G$ is bipartite, all of these cycles have even length. Each path is $M$-alternating and also $M'$-alternating as well. If any path had odd length, then it would be either $M$-augmenting (but we assumed they didn’t exist) or $M'$-augmenting (but they also don’t exist because $M'$ is a maximum matching and we can use the argument in the first paragraph), so they don’t appear. This means that $H$ consists of even length cycles and even length paths, which means that $M$ and $M'$ have the same number of edges: the number of edges of $M$ is the number of edges that $M$ and $M'$ have in common plus half the number of edges of $H$ (and same for $M'$). Hence $M$ is also maximum. 

This gives a better way to build maximum matchings than just brute force: start with any matching $M$, then look for $M$-augmenting paths. If they exist, swap the edges in the path like in the proof and continue. If they don’t exist, we know we’re done.

7.4. Stable matchings. We end this section with a discussion of stable matchings, which take into account preferences. Continuing with the doctor example, we don’t want to just put people into jobs just because they are qualified to do them. Instead, we should try to take into account individual rankings. First, let’s assume that $|X| = |Y| = n$. Furthermore, attached to each vertex $x \in X$ is a ranking of the vertices of $Y$, and similarly, each vertex $y \in Y$ has a ranking of the vertices of $X$. For the doctors, a ranking is their ranked list of the available jobs in order of preference, and similarly, for the jobs.

With the above setup, a stable matching is a perfect matching of $X$ into $Y$ such that for every $x \in X$ and $y \in Y$ that are not matched, either $x$ prefers its current match over $y$, or $y$ prefers its current match over $x$. In other words, a matching fails to be stable if there is some pair $(x, y)$ which prefer each other over their current matches. In our example, a perfect matching has placed each doctor into a job. Some doctors might prefer jobs they didn’t get over the one they get. Stable here means that whenever doctor $x$ preferred job $y$ over their current job, the person filling job $y$ was better for the job than $x$ is, and vice versa: if job $y$ has doctor $x$ working in it but doctor $x'$ has a better fit, then doctor $x'$ has a better job than $y$ (in their own personal ranking).

The Gale-Shapley algorithm finds a stable matching. We’ll use doctors and hospitals to explain it. It proceeds in rounds:

- In the first round, each doctor sends an application to their favorite hospital. Each hospital ranks the applications and says “maybe” to the best candidate and “no” to the rest.
- In all future rounds, any doctor with a “maybe” does nothing. All others send an application to their favorite hospital that have not already rejected them. Each hospital ranks new applicants against their current “maybe” and then says “maybe” to top candidate and “no” to the rest (possibly including one that used to have a “maybe”).
- This keeps going until each doctor has a “maybe”. In that case, each “maybe” becomes a “yes”.

I’ll leave you to think about why this algorithm terminates and correctly produces a stable matching.

Note the asymmetry in the problem: we could have instead had hospitals soliciting their favorite doctors to work for them. The result can be different, and in fact, this setup gives an advantage to the hospitals. Furthermore, this algorithm actually has been used before to assign jobs to graduating medical students. One more thing: there are definitely flaws
with this algorithm in the sense that it is possible for hospitals to lie about their preferences in order to get a better candidate in some cases at the expense of making things worse for other hospitals.

8. Planarity

There's an old puzzle related to the topic of planarity: there are 3 houses in a row and 3 utilities (gas, water, electric) in a row below them. Can we connect each house to each utility without crossing any lines? We'll formulate this in terms of graph theory and use it to answer the question.

8.1. Definitions. A planar graph is, roughly speaking, a graph which can be drawn in the plane (for example, a piece of paper) in such a way that edges do not overlap each other (except at their endpoints). The basic instructive example is the complete graph $K_4$. Sometimes it is drawn as follows:

and the two diagonal edges overlap. However, here are two different ways to draw it so that none of the edges overlap:

In the left drawing, one of edges is “curved”: this is allowed. In the right drawing, all of the edges are straight lines. It is a theorem that planar graphs can always be drawn so that the edges are all straight lines, but we won't use this and it takes a bit of effort to prove this, so we won't say anything more about it.

It's easy enough to see that if $\overline{G}$ is the simple graph associated to $G$ (i.e., multiple edges are replaced with a single edge and loops are deleted), then $G$ is planar if and only if $\overline{G}$ is planar. So we will restrict our attention to simple graphs since it's no less general but might make notation easier to deal with.

When you draw a graph in the plane, you have separated the plane into different regions. Another way to say this: if you delete the graph from the plane, the regions are the different connected pieces. These are called faces. Note that the boundary of every face is a cycle of the graph. In the right drawing of $K_4$ above, the faces are the 3 inside triangles and then there is 1 “outside” face, so 4 in total. You can also see there are 4 faces in the left drawing. We'll see shortly that the number of faces is an isomorphism-invariant of the graph.

Once you draw a graph $G$ in the plane, there is something called the dual graph $G^*$: the vertices are the faces of $G$, and two faces are connected by an edge if they share an edge. This notion is useful for translating the problem of coloring maps into problems of coloring graphs (the countries are faces), but otherwise we won’t do much with this definition.
8.2. Some equations and inequalities.

**Theorem 8.1.** (Euler) Let $G$ be a connected planar graph with $n$ vertices, $m$ edges, and $f$ faces. Then $n - m + f = 2$.

**Proof.** We prove this by induction on $m$, the number of edges. If $m = 0$, then $G$ is a single vertex since it is connected. In that case, $n = f = 1$, and the equation is satisfied.

Now assume the equation holds for all graphs with $< m$ edges and let $G$ be a graph with $m$ edges. We consider 2 separate cases:

- If $G$ is a tree, then $m = n - 1$ and $f = 1$ (there’s just the outside face) and we see again the identity holds.
- If $G$ is not a tree, then it fails to be minimally connected, so there is some edge $e = \{x, y\}$ such that $G' = G \setminus \{e\}$ is still connected. In particular, there is a cycle of $G$ that uses $e$. Relative to our planar drawing for $G$, pick such a cycle that has no edges going inside of it. Then this is a face of $G$. There is also another face of $G$ which is on the other side of $G$. When we remove $e$, these two faces merge into one.

So, $G'$ has $n$ vertices, $m - 1$ edges, and $f - 1$ faces. By induction, $n - (m - 1) + (f - 1) = 2$, but the left hand side simplifies to $n - m + f$, so the identity is also valid for $G$.

**Corollary 8.2.** Let $G$ be a planar graph with $c$ connected components, $n$ vertices, $m$ edges, and $f$ faces. Then $n - m + f = c + 1$.

**Proof.** We’ll do induction on $c$, the case $c = 1$ being Theorem 8.1. Let $G$ be a graph with $c$ connected components $G_1, \ldots, G_c$. Say $G_c$ has $n'$ vertices, $m'$ edges, and $f'$ faces. Let $H$ be the result of removing $G_c$. Then $H$ is still planar and has $n - n'$ vertices, $m - m'$ edges, but $f - (f' - 1)$ faces: we subtract $f'$ - 1 because $G_c$ and $H$ have the same “outside” face. By induction: $n - n' - (m - m') + f - (f' - 1) = c$. By Theorem 8.1, we also know that $n' - m' + f' = 2$. Adding these two equations gives $n - m + f = c + 1$.

A simple planar graph with $n$ vertices can’t have arbitrarily many edges. We can give a bound in terms of $n$, but can do much better if we use the girth. For $n \geq 3$, let $C_n$ be the cycle graph with $n$ vertices (so it has $n$ vertices $v_1, \ldots, v_n$ and $n$ edges $\{v_1, v_2\}, \{v_2, v_3\}, \ldots, \{v_n, v_1\}$).

**Definition 8.3.** A simple graph $G$ has girth at least $g$ if it does not contain a subgraph isomorphic to $C_3, C_4, \ldots, C_{g-1}$. (If $g = 3$, this condition is vacuous, so every simple graph has girth $\geq 3$.) We say that $G$ has girth equal to $g$ if it has girth at least $g$ and also contains a subgraph isomorphic to $C_g$. If $G$ has no cycles, then it has infinite girth.

If $H$ is a subgraph of $G$, then the girth of $H$ is at least as big as the girth of $G$.

**Theorem 8.4.** Let $G$ be a simple planar graph with $n$ vertices, $m$ edges, and finite girth $\geq g$. Then

$$m \leq \frac{g}{g-2}(n-2).$$

**Proof.** Let $c$ be the number of connected components of $G$. Let $N$ be the number of pairs $(e, F)$ where $e$ is an edge on the boundary of the face $F$. Since $G$ has finite girth $g$, the boundary of every face of $G$ has $\geq g$ edges (the boundary of the outside face is the boundary of $G$), so $N \geq gf$. Next, each edge is on the boundary of at most 2 faces, so $N \leq 2m$. We conclude that $2m \geq gf$. By Corollary 8.2,

$$f = c + 1 - n + m \geq 2 - n + m,$$
so \(2m \geq g(2 - n + m)\). Rearranging terms we get \(g(n - 2) \geq (g - 2)m\). Now divide by \(g - 2\) (note that \(g - 2 \geq 1\) because \(g \geq 3\) by the way we defined girth).

If \(G\) has infinite girth, then it’s a forest, and we already know that \(m \leq n - 1\), so the restriction of finite girth isn’t a big deal.

**Corollary 8.5.** Let \(G\) be a simple planar graph with \(n\) vertices and \(m\) edges. If \(n \geq 3\), then \(m \leq 3n - 6\).

**Proof.** If \(G\) is a forest, then \(m \leq n - 1 \leq 3n - 6\) (the second inequality holds because \(n \geq 3\)). Otherwise, \(G\) has finite girth \(\geq 3\) so \(m \leq 3(n - 2)\) by Theorem 8.4.

The following result will be useful when we study colorings of planar graphs:

**Corollary 8.6.** If \(G\) is a simple planar graph, then it has a vertex with degree \(\leq 5\).

**Proof.** Suppose not. Then every vertex of \(G\) has degree \(\geq 6\) (in particular, \(G\) has at least 7 vertices). Let \(m\) be the number of edges and \(n\) be the number of vertices. By the handshake lemma (a problem from HW4), \(2m = \sum_v \deg(v) \geq 6n\) where the sum is over all vertices \(v\). In particular, \(m \geq 3n\), which contradicts Corollary 8.5 since \(n \geq 7\).

You can’t improve the above statement to 4 because here is a planar graph where every vertex has degree 5 (this is the projection of an icosahedron):

8.3. **Obstructions to planarity.** Recall that \(K_n\) is the complete graph on \(n\) vertices: every pair of vertices has an edge between them. Also recall the complete bipartite graph \(K_{n,m}\). It has vertices \(\{x_1, \ldots, x_n, y_1, \ldots, y_m\}\) and the edges \(\{x_i, y_j\}\) for all \(i = 1, \ldots, n\) and \(j = 1, \ldots, m\).

Here are two important examples of non-planar graphs:

**Example 8.7.** \(K_5\) has 5 vertices and 10 edges, so is not planar: if it were, then Corollary 8.5 implies \(10 \leq 3(5 - 2)\) which is false.

**Example 8.8.** \(K_{3,3}\) has 6 vertices and 9 edges. We claim it is not planar. However, if you test the inequality in Corollary 8.5, you get \(9 \leq 3(6 - 2)\) which is valid, so doesn’t help us. We instead use the stronger Theorem 8.4. Since \(K_{3,3}\) is bipartite, it has no odd cycles, and in particular, its girth is \(\geq 4\). Then the sharper inequality would be \(9 \leq 2(6 - 2)\) which is false, and allows us to conclude \(K_{3,3}\) is not planar.

**Example 8.9.** We’ve seen already that \(K_4\) is planar. In particular, so is \(K_n\) for \(n \leq 4\) (subgraphs of planar graphs are clearly planar). \(K_{3,2}\) is also planar (can you find a planar drawing?). In fact, \(K_{n,2}\) is planar for any \(n\).
Given a graph $G$ and an edge $e = \{x, y\}$, the subdivision of $G$ along $e$ is a new graph obtained as follows: add a new vertex $z$, and edges $\{x, z\}$ and $\{y, z\}$ and remove $e$ (pictorially, we've replaced $e$ with two edges. Here's an example where we're subdividing the left edge:

Starting with any graph $G$, we can subdivide edges all we like (including the new ones), the resulting set of graphs are called subdivisions of $G$.

It's clear that if $G$ isn't planar, then neither is any subdivision of it (all we really did is add vertices along edges). Furthermore, if $G$ contains a subgraph which isn't planar, then $G$ also can't be planar.

Combining what we know, if $G$ contains a subgraph which is isomorphic to a subdivision of $K_5$ or $K_{3,3}$, then it isn't planar. The converse is also true, but we won't prove it:

**Theorem 8.10** (Kuratowski, 1930). A simple graph $G$ is planar if and only if it does not have a subgraph which is isomorphic to a subdivision of $K_5$ or $K_{3,3}$.

There is a variant which is also convenient for testing planarity. Given a graph $G$ and an edge $e$, we can delete it to get $G \setminus e$ or contract it to get $G/e$. We will call both of these **graph minors** of $G$, and more generally, a graph minor of $G$ is any graph which can be obtained from $G$ by repeatedly deleting edges and vertices, and also contracting edges.

If $G$ is planar, then so is any graph minor (consider how these operations affect a planar drawing). Equivalently, if $G$ has a non-planar graph minor, then $G$ is non-planar. So we conclude that if $G$ has a graph minor isomorphic to $K_5$ or $K_{3,3}$, then $G$ is not planar. Again, the converse is also true, but we won't prove it:

**Theorem 8.11** (Wagner, 1937). A simple graph $G$ is planar if and only if it does not have a graph minor which is isomorphic to $K_5$ or $K_{3,3}$.

In this way, we see that $K_5$ and $K_{3,3}$ are the basic obstructions to a graph being planar.

### 8.4. Chromatic numbers of planar graphs.

The following result is a famous result in graph theory which was first proven by Appel and Haken in 1976:

**Theorem 8.12** (4-color theorem). Let $G$ be a simple planar graph. Then $\chi(G) \leq 4$, i.e., $G$ has a proper 4-coloring.

The proof is very difficult so we won’t say much about it. The method of proof roughly consisted of reducing the check to a finite (but very large) number of different cases and have the result checked by a computer program. This is historically one of the early uses of computer verification of proofs in math.

However, it turns out to be not that hard to prove weaker versions. Our goal will be to prove the 5-color theorem. But first, let’s warm-up with the 6-color theorem.

**Theorem 8.13.** Let $G$ be a simple planar graph. Then $\chi(G) \leq 6$.

**Proof.** We do induction on the number of vertices. If $G$ has 1 vertex, there’s nothing to say. Otherwise, assume that all simple planar graphs with $< n$ vertices have a proper 6-coloring, and let $G$ be a simple planar graph with $n$ vertices. By Corollary 8.6, there is a vertex $v$ of $G$ with $\deg(v) \leq 5$. Let $H$ be the graph obtained by deleting $v$ and all of the edges...
that touch it. By induction, $H$ has a proper 6-coloring. The neighbors of $v$ use at most 5 different colors, so we are free to use the 6th remaining one for $v$, and hence extend the proper coloring to all of $G$. \hfill \Box

Now let’s do the 5-color theorem. We follow the same outline, but need to consider more situations in the induction step.

**Theorem 8.14.** Let $G$ be a simple planar graph. Then $\chi(G) \leq 5$.

**Proof.** Same as before: we do induction on the number of vertices of $G$. Let $v$ be a vertex of $G$ with $\deg(v) \leq 5$, $H = G \setminus \{v\}$ has a proper 5-coloring by induction. If the neighbors don’t use all 5 colors, then we can use the same argument from before to finish. So we may assume $\deg(v) = 5$ and that all 5 neighbors of $v$ must use different colors.

Let $x_1, \ldots, x_5$ be the neighbors of $v$ in clockwise order and let $c_i$ be the name of the color of $x_i$. For each $i < j$, let $H(i, j)$ be the subgraph of $H$ on all vertices that use colors $c_i$ or $c_j$ (with all edges between them). If $x_1$ and $x_3$ are in different connected components of $H(1, 3)$, then swap the colors $c_1$ and $c_3$ of all vertices in the component containing $x_1$. Then this gives a new proper 5-coloring of $H$, and now the neighbors of $v$ only use 4 different colors, so we can assign $c_1$ to $v$ and finish. Otherwise, $x_1$ and $x_3$ are in the same connected component of $H(1, 3)$, which means there is a path from $x_1$ to $x_3$ consisting of vertices alternating between $c_1$ and $c_3$:

![Diagram](attachment:diagram.png)

But this means that $x_2$ and $x_4$ have to be in different connected components of $H(2, 4)$ (the red-blue path above separates them). So we can do the swapping we mentioned above to reduce the number of colors used to 4. \hfill \Box

It’s tempting to keep pushing this argument to bring down the number of colors needed to 4. Actually, something like this was done by Kempe in 1879. It turned out there was a subtle mistake in Kempe’s argument which wasn’t found until 11 years later. I’ve posted an article to the course webpage which explains the details of Kempe’s argument and why it was wrong (this is purely optional further reading).

9. **Pigeon-hole principle**

9.1. **Basic version.** The following is really obvious, but is a very important tool. The proof illustrates how to make “obvious” things rigorous. It is important to always keep this in mind especially in this course when many things you might want to use sound obvious. There are many interesting ways to use this theorem which are not obvious.

**Theorem 9.1** (Pigeon-hole principle (PHP)). Let $n, k$ be positive integers with $n > k$. If $n$ objects are placed into $k$ boxes, then there is a box that has at least 2 objects in it.
Proof. We will do proof by contradiction. So suppose that the statement is false. Then each box has either 0 or 1 object in it. Let \( m \) be the number of boxes that have 1 object in it. Then there are \( m \) objects total and hence \( n = m \). However \( m \leq k \) since there are \( k \) boxes, but this contradicts our assumption that \( n > k \). \( \square \)

Note that the objects can be anything and the boxes don’t literally have to be boxes.

Example 9.2.  
- Simple example: If we have 4 flagpoles and we put up 5 flags, then there is some flagpole that has at least 2 flags on it.
- Draw 10 points in a square with unit side length. Then there is some pair of them that are less than .48 distance apart. There’s some content here since the corners on opposite ends have distance \( \sqrt{2} \approx 1.4 \). Also, if we only have 9 points, we could arrange them like so:

\[
\begin{array}{cccc}
\cdot & \cdot & \cdot & \\
\cdot & & & \\
\cdot & & & \\
\end{array}
\]

The pairs of points that are closest are .5 away from each other, so it is important that we have at least 10 points.

To see why the statement holds, divide the square into 9 equal parts:

Then some little square has to contain at least 2 points in it (is it ok if the points are on the boundary segments?). Each square has side length \( 1/3 \), and so the maximum distance between 2 points in the same square is given by the length of its diagonal (why?) which is \( \sqrt{(1/3)^2 + (1/3)^2} = \sqrt{2}/3 \approx 0.4714 \). \( \square \)

Here are some more to think about:

- At least 2 of the students in this class were born in the same month.
- If you have 10 white socks and 10 grey socks, and you grabbed 3 of them without looking, you automatically have a matching pair.
- Pick 5 different integers between 1 and 8. Then there must be a pair of them that add up to 9.
- Given 5 points on a sphere, there is a hemisphere that contains at least 4 of the points.
- There is a party with 1000 people. Some pairs of people have a conversation at this party. There must be at least 2 people who talked to the same number of people.
- Given an algorithm for compressing data, if there exist files whose length strictly decreases, then there exist files whose length strictly increases!

In mathematical terms: let’s represent a file by a sequence of 0’s and 1’s. Then an algorithm for compressing data can be thought of as a function that takes each
sequence to some other sequence in such a way that different inputs must result in different outputs.

9.2. **General version.** Here’s a more general version of the PHP:

**Theorem 9.3** (General pigeon-hole principle). Let \( n, m, r \) be positive integers and suppose that \( n > rm \). If \( n \) objects are placed into \( m \) boxes, then there is a box that contains at least \( r + 1 \) objects in it.

If you set \( r = 1 \), then this is exactly the first version of the PHP.

**Proof.** We can again do this via proof by contradiction. Suppose the statement is false and label the boxes 1 up to \( m \). Let \( b_i \) be the number of objects in box number \( i \). Then \( b_i \leq r \) since the conclusion is false. Furthermore, we have \( n = b_1 + b_2 + \cdots + b_m \leq r + r + \cdots + r = rm \).

But this contradicts the assumption that \( n > rm \).

**Example 9.4.**
- Simple example: If we have 4 flagpoles and 9 flags distributed to them, then some flagpole must have at least 3 flags on it.
- Continuing from our geometry example from before: draw 9 points in a square of unit side length. Then there must be a triple of them that are contained in a single semicircle of radius 0.5. (Is this true if we only have 8 points?)

For the solution, we divide up the square into 4 triangles as follows:

Then some triangle must contain at least 3 points. Furthermore, each triangle fits into a semicircle of radius 0.5.

\[ \square \]

10. **Ramsey theory**

10.1. **Ramsey’s theorem for graphs.**

10.2. **Ramsey’s theorem for hypergraphs.**

10.3. **Turán’s theorem.**

10.4. **Lower bounds on Ramsey numbers.**