

A first year mathematics degree

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Chapter 1

Analysis

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Introduction

The following are core reading:

- M W Liebeck, A concise intro to pure mathematics, 2016
- F M Hart, Guide to analysis, 2001
- D A Brannan, A first course in mathematical analysis, 2006
- K G Binmore, Mathematical analysis: a straightforward approach, 1982

The following are supplementary reading:

- S R Lay, Analysis with an introduction to proof, 1994
- S Abbott, Understanding Analysis, 2015

The following are advanced reading:

- M Spivak, Calculus, 2010
- T Tao, Analysis 1, 2022

Lecture 1
Thursday
10/01/19

1 Number systems

1.1 Naturals, integers and rationals

Definition 1.1.1 (Natural numbers). As in IUM, we define the **natural numbers**, \mathbb{N} , from the Peano axioms:

P1 0 is a natural number,

P6 if n is a natural number then $S(n)$ is a natural number where $S(n)$ is the successor of n ,

P9 the principle of mathematical induction.

Clearly, there are many Peano axioms not included, these are however not particularly relevant to this course. Addition and multiplication is defined as expected and will descend to our other number systems

Definition 1.1.2 (Integers). The **integers** are defined as $\mathbb{Z} := \mathbb{N} \times \mathbb{N} / \sim$ where \sim is the equivalence relation given by $(a, b) \sim (c, d)$ iff $a + d = b + c$. Subtraction is defined as expected and will also descend to our other number systems.

Definition 1.1.3 (Rationals). The **rationals** are defined as $\mathbb{Q} := \mathbb{Z} \times \mathbb{N}^{>0} / \sim$ where \sim is the equivalence relation given by $(a, b) \sim (c, d)$ iff $ad = bc$. The equivalence class (p, q) will be written as $\frac{p}{q}$. There is an element of each equivalence class $\frac{p'}{q'}$ with $\gcd(p', q') = 1$, we say that $\frac{p'}{q'}$ is in **lowest terms**.

Theorem 1.1.4 (Axioms of the rationals). With the usual operations descended from \mathbb{N} and \mathbb{Z} , \mathbb{Q} satisfies the following axioms with $a, b, c \in \mathbb{Q}$ throughout:

Q1 $a + (b + c) = (a + b) + c$ ($+$ is associative),

Q2 $\exists 0 \in \mathbb{Q}$ such that $a + 0 = a$ (0 is the additive identity of \mathbb{Q}),

Q3 $\forall a \in \mathbb{Q}, \exists (-a) \in \mathbb{Q}$ such that $a + (-a) = 0$ (\mathbb{Q} is closed under additive inverses),

Q4 $a + b = b + a$ ($+$ is commutative),

Q5 $a \times (b \times c) = (a \times b) \times c$ (\times is associative),

Q6 $\exists 1 \in \mathbb{Q}$ such that $a \times 1 = a$ (1 is the multiplicative identity of \mathbb{Q}),

Q7 $a \times (b + c) = (a \times b) + (a \times c)$ (\times is left distributive over $+$),

Q8 $(a + b) \times c = (a \times c) + (b \times c)$ (\times is right distributive over $+$),

Q9 $a \times b = b \times a$ (\times is commutative),

Q10 $\forall a \in \mathbb{Q}, \exists a^{-1} \in \mathbb{Q}$ such that $a \times a^{-1} = 1$ (\mathbb{Q} is closed under multiplicative inverses),

Q11 for all $a \in \mathbb{Q}$ either $x < 0$, $x = 0$ or $x > 0$ (Trichotomy),

Q12 for all $x, y \in \mathbb{Q}$ we have $x > 0, y > 0 \implies x + y > 0$,

Q13 for all $x \in \mathbb{Q}$ there exists a $n \in \mathbb{N}$ such that $x < n$ (Archimedean axiom).

1-4 says $(\mathbb{Q}, +)$ is an abelian group, 1-9 says $(\mathbb{Q}, +, \times)$ is a commutative ring, 1-10 says $(\mathbb{Q}, +, \times)$ is a field.

1.2 Decimal expansions

Definition 1.2.1. For $a_0 \in \mathbb{N}$ and $a_i \in [1, 9]$ for $i > 0 \in \mathbb{N}$, define the **periodic decimal**

$$a_0.a_1a_2 \dots \overline{a_ia_{i+1} \dots a_j},$$

to be equal to the rational number

$$a_0 + \frac{a_1}{10} + \frac{a_2}{100} + \dots + \frac{a_i}{10^i} + \left(\frac{a_{i+1}a_{i+2} \dots a_j}{10^j} \right) \left(\frac{1}{1 - 10^{i-j}} \right).$$

Theorem 1.2.2. If $x \in \mathbb{Q}$ has 2 decimal expansions, then they will be of the form

$$x = a_0.a_1a_2 \dots a_n\overline{9} = a_0.a_1a_2 \dots (a_n + 1), a_n \in [0, 8].$$

Definition 1.2.3 (Real numbers). The **real numbers**, \mathbb{R} , can be defined as:

$$\mathbb{R} := \{a_0.a_1a_2 \dots : a_0 \in \mathbb{Z}, a_i \in [0, 9], \exists N \in \mathbb{N} \text{ such that } a_i = 9 \forall i \geq N\}.$$

1.3 Countability

Definition 1.3.1 (Countability). A set S is **countably infinite** iff there exists a bijection $f : \mathbb{N} \rightarrow S$, a set is **countable** if it is finite or countable infinite.

Theorem 1.3.2. All $S \subseteq \mathbb{N}$ are countable, \mathbb{Z} and \mathbb{Q} are both countable, \mathbb{R} is uncountable.

2 Bounded sets

2.1 Supremums and infimums

Definition 2.1.1 (Maximum and minimum). $s \in \mathbb{R}$ is the **maximum** of a set $S \subset \mathbb{R}$ iff $\forall s' \in S, s \geq s'$. **Minimums** are defined similarly. Maximums and minimums are unique.

Definition 2.1.2 (Bounded). A non-empty set $S \subset \mathbb{R}$ is **bounded above** if there exists some $M \in \mathbb{R}$ such that $\forall s \in S, s \leq M$ with **bounded below** defined similarly. S is **bounded** if it is both bounded above and bounded below.

Theorem 2.1.3. If S is bounded then $\exists R > 0$ such that $|s| < R$ for all $s \in S$.

Definition 2.1.4 (Supremum and infimum). If $S \subset \mathbb{R}$ is bounded above, we say $x \in \mathbb{R}$ is the **least upper bound** or **supremum** iff x is an upper bound for S and for all $y \in \mathbb{R}$ such that y is an upper bound of S , $x \leq y$. The **infimum** is defined similarly.

2.2 Completeness

Theorem 2.2.1 (Completeness axiom). For all non-empty $S \subset \mathbb{R}$, if S is bounded above then S has a supremum, and similarly for S bounded below.

2.3 Dedekind cuts

Definition 2.3.1 (Dedekind cut). A non-empty set $S \subset \mathbb{Q}$ is a **Dedekind cut** if it satisfies:

1. $s \in S$ and $s > t \in \mathbb{Q} \implies t \in S$ (S is a semi-infinite interval to the left),
2. S is bounded above with no maximum.

Dedekind cuts are in the form $S_r := (-\infty, r) \cap \mathbb{Q}$.

Theorem 2.3.2 (Real numbers). We can redefine the reals as the set of Dedekind cuts, $\mathbb{R} := \{S_r \subset \mathbb{Q}\}$. All operations and orderings as well as the completeness axiom are held by this new Dedekind cut definition.

Theorem 2.3.3 (Triangle inequality). For all $a, b \in \mathbb{R}$ we have $|a + b| \leq |a| + |b|$.

3 Sequences

Definition 3.0.1 (Real sequence). A **real sequence** is a function $a : \mathbb{N} \rightarrow \mathbb{R}$ written (a_n) . Sequences of other number systems are defined similarly.

3.1 Convergence

Definition 3.1.1 (Convergence of sequences). A real sequence (a_n) **converges** to some $a \in \mathbb{R}$ as $n \rightarrow \infty$ iff

$$\forall \epsilon > 0, \exists N_\epsilon \text{ such that } \forall n \geq N_\epsilon, |a_n - a| < \epsilon.$$

For complex series the definition is the same just with $|\cdot|$ referring to the modulus instead of the absolute value. This is written $a_n \rightarrow a$ (as $n \rightarrow \infty$).

3.2 Divergence

Definition 3.2.1 (Divergence). A sequence (a_n) **diverges** iff it doesn't converge.

Definition 3.2.2 (Divergence to infinity). A sequence (a_n) **diverges to ∞** iff $\forall R > 0, \exists N \in \mathbb{N}$, such that $\forall n \geq N, a_n > R$. And similarly for a sequence diverging to $-\infty$.

3.3 Limits

Theorem 3.3.1 (Uniqueness of limits). Given a sequence (a_n) if $a_n \rightarrow a$ and $a_n \rightarrow b$, $a = b$.

Theorem 3.3.2. If a sequence (a_n) is convergent then (a_n) is bounded.

Theorem 3.3.3 (Algebra of limits). Given two sequences $a_n \rightarrow a$ and $b_n \rightarrow b$ the following hold:

- $a_n + b_n \rightarrow a + b$,
- $a_n b_n \rightarrow ab$ (a special case of this is $ca_n \rightarrow ca$ for a constant c),
- $\frac{a_n}{b_n} \rightarrow \frac{a}{b}$ given $b \neq 0$.

Theorem 3.3.4. If (a_n) is a positive sequence then $a_n \rightarrow 0 \iff \frac{1}{a_n} \rightarrow +\infty$, and similarly for negative sequences.

Theorem 3.3.5 (Ratio test). If a sequence (a_n) satisfies $\left| \frac{a_{n+1}}{a_n} \right| \rightarrow L < 1$ then $a_n \rightarrow 0$.

3.4 Monotone sequences

Definition 3.4.1 (Monotonically increasing sequence). A sequence, (a_n) , is **monotonically increasing** iff $\forall m, n \in \mathbb{N}$ with $n > m$ we have $a_n \geq a_m$, and similarly for monotonically decreasing and their strict equivalents.

Theorem 3.4.2 (Monotone convergence). If a sequence (a_n) is monotone increasing and bounded above then $a_n \rightarrow a := \sup\{a_i : i \in \mathbb{N}\}$ written $a_n \uparrow a$. This holds similarly for monotone decreasing sequences.

3.5 Cauchy sequences

Definition 3.5.1 (Cauchy sequence). A sequence (a_n) is a **Cauchy sequence** iff $\forall \epsilon > 0 \in \mathbb{R}, \exists N \in \mathbb{N}$ such that $\forall n, m < N, |a_n - a_m| < \epsilon$.

Theorem 3.5.2 (Cauchy convergence criterion). A sequence (a_n) converges iff it is a Cauchy sequence.

3.6 Subsequences

Definition 3.6.1 (Subsequence). Given a strictly monotonically increasing function $n : \mathbb{N} \rightarrow \mathbb{N}$ and a sequence (a_n) , the sequence (b_n) defined by $b_i := a_{n(i)}$ is a **subsequence** of (a_n) .

Theorem 3.6.2. Given a subsequence of (a_n) , $(a_{n(i)})$, if $a_n \rightarrow a$ then $a_{n(i)} \rightarrow a$ as $i \rightarrow \infty$.

Theorem 3.6.3 (Bolzano-Weierstrass). If a sequence (a_n) is bounded then it has a convergent subsequence.

Note 3.6.4 (Sketch of the Bolzano-Weierstrass theorem proof). The proof of the Bolzano-Weierstrass theorem is an equally valuable point as the statement of the theorem itself. The idea of the proof considers the “peak points” of the sequence: if there are infinitely many peak points, then the peak points themselves form a monotonically decreasing subsequence; if there are finitely many, then the points after the final peak must have a monotonically increasing subsequence bounded above by the final peak. By the monotone convergence theorem both of these subsequences must converge.

4 Series

Definition 4.0.1 (Infinite series). An **(infinite) series** is an expression of the form $\sum_{i=1}^{\infty} a_i$ of $a_1 + a_2 + \dots$ for some sequence (a_n) . The sequence **partial sums** of the series (S_n) is given by

$$S_n := \sum_{i=1}^n a_i = a_1 + a_2 + \dots + a_n.$$

4.1 Convergence

Definition 4.1.1 (Convergence of series). The series $\sum_{i=1}^{\infty} a_i$ of (a_n) **converges** iff $S_n \rightarrow A \in \mathbb{R}$, written $\sum_{n=1}^{\infty} a_n = A$.

Theorem 4.1.2. For a sequence (a_n) , $\sum_{n=1}^{\infty} a_n$ converges if $a_n \rightarrow 0$ (the converse is not true).

Theorem 4.1.3. Given a sequence non-negative sequence (a_n) , the convergence of the infinite series and the boundedness of (S_n) are equivalent.

Theorem 4.1.4 (Algebra of limits for series). A similar algebra of limits for series can be established from the algebra of limits for sequences acting on the partial sums of the series.

Theorem 4.1.5 (Comparison I test). Given sequences $(a_n), (b_n)$ if $0 \leq a_n \leq b_n$ then:

- If $\sum_{n=1}^{\infty} a_n$ and $\sum_{n=1}^{\infty} b_n$ converge, $0 \leq \sum_{n=1}^{\infty} a_n \leq \sum_{n=1}^{\infty} b_n$,
- If $\sum_{n=1}^{\infty} a_n$ diverges, $\sum_{n=1}^{\infty} b_n$ also diverges.

Theorem 4.1.6 (Comparison II test (Sandwich theorem)). Given sequences $(a_n), (b_n), (c_n)$ with $a_n \leq b_n \leq c_n$, if $\sum_{n=1}^{\infty} a_n$ and $\sum_{n=1}^{\infty} c_n$ both converge, $\sum_{n=1}^{\infty} b_n$ converges.

Theorem 4.1.7. If $\alpha > 1 \in \mathbb{R}$, $\sum_{n=1}^{\infty} \frac{1}{n^\alpha}$ converges.

Definition 4.1.8 (Alternating sequence). A sequence (a_n) is **alternating** iff $a_{2n} \geq 0$ and $a_{2n-1} \leq 0$ of vice versa for all $n \in \mathbb{N}^{>0}$.

Theorem 4.1.9. If (a_n) is alternating with $|a_n| \downarrow 0$, a_n converges and $\sum_{n=1}^{\infty} a_n$ converges.

4.2 Absolute convergence

Definition 4.2.1 (Absolute convergence). Given a sequence (a_n) the series $\sum_{n=1}^{\infty} a_n$ is **absolutely convergent** iff $\sum_{n=1}^{\infty} |a_n|$ converges.

Theorem 4.2.2. Absolute convergence \implies convergence.

Theorem 4.2.3 (Comparison III test). Given sequences $(a_n), (b_n)$ with $\frac{a_n}{b_n} \rightarrow L \in \mathbb{R}$ if $\sum_{n=1}^{\infty} b_n$ is absolutely convergent then $\sum_{n=1}^{\infty} a_n$ is also absolutely convergent.

Theorem 4.2.4 (Ratio test). If the sequence (a_n) is such that $\left| \frac{a_{n+1}}{a_n} \right| \rightarrow r < 1$ then $\sum_{n=1}^{\infty} a_n$ is absolutely convergent or divergent if $r > 1$.

Theorem 4.2.5 (Root test). If the sequence (a_n) is such that $|a_n|^{\frac{1}{n}} \rightarrow r < 1$ then $\sum_{n=1}^{\infty} a_n$ is absolutely convergent or divergent is $r > 1$.

Remark 4.2.6. Both the ratio test and the root test are inconclusive if $r = 1$.

4.3 Rearrangement of series

Sometimes, series are easier to deal with and have cancellations when their terms are rearranged. However, the rearrangement of terms will only preserve limits under certain conditions.

Definition 4.3.1 (Reordering). Given a bijection $n : \mathbb{N} \rightarrow \mathbb{N}$ and a sequence (a_n) , the sequence (b_n) with $b_i := a_{n(i)}$ is a **rearrangement** or **reordering** of (a_n) .

Theorem 4.3.2. If (a_n) is a sequence satisfying $a_n \rightarrow 0$, $\sum_{n: a_n \geq 0} a_n = \infty$ and $\sum_{n: a_n \leq 0} a_n = -\infty$ then $\sum_{n=1}^{\infty} a_n$ can be rearranged to converge to any $r \in \mathbb{R}$.

Theorem 4.3.3. If (a_n) is a sequence with absolutely convergent series, $\sum_{n: a_n \geq 0} a_n = A$ and $\sum_{n: a_n \leq 0} a_n = B$ with all arrangements of (a_n) converging to $A + B$.

4.4 Power series

Throughout this subsection $[0, \infty] := [0, \infty) \cup \{+\infty\}$.

Definition 4.4.1 (Power series). For $z \in \mathbb{C}$ and a complex sequence (a_n) , a **power series** is an expression in the form $\sum_{n=1}^{\infty} a_n z^n$.

Definition 4.4.2 (Radius of convergence). Given the power series $\sum_{n=1}^{\infty} a_n z^n$, there exists some $R \in [0, \infty]$ such that:

- $|z| < R \implies \sum_{n=1}^{\infty} a_n z^n$ converges,
- $|z| > R \implies \sum_{n=1}^{\infty} a_n z^n$ diverges.

We cannot tell what happens when $|z| = R$ so this has to be checked separately. R is the **radius of convergence** of the power series.

Corollary 4.4.3. Given the same power series $\sum_{n=1}^{\infty} a_n z^n$, have $S := \{|z| \in \mathbb{R}^{\geq 0} : a_n z^n \rightarrow 0\}$ then

$$R := \begin{cases} \sup(S) & \text{if } S \text{ is bounded} \\ \infty & \text{otherwise} \end{cases}.$$

is the radius of convergence for the power series.

Theorem 4.4.4 (Evaluating radius of convergence from tests). For the power series $\sum_{n=1}^{\infty} a_n z^n$:

- if $\left| \frac{a_{n+1}}{a_n} \right| \rightarrow a \in [0, \infty]$ then $R = \frac{1}{a}$ is the radius of convergence for the power series,
- if $|a_n|^{\frac{1}{n}} \rightarrow a \in [0, \infty]$ then $R = \frac{1}{a}$ is the radius of convergence for the power series,

Definition 4.4.5 (Cauchy product). Given two series $\sum_{n=1}^{\infty} a_n$, $\sum_{n=1}^{\infty} b_n$; their **Cauchy product** is the series

$$\sum_{n=0}^{\infty} \sum_{i=0}^n a_i b_{n-i}.$$

Remark 4.4.6. If (a_n) , (b_n) are the coefficients for a power series, then the Cauchy product of their series will be the coefficients of the product of the power series.

Theorem 4.4.7. If $\sum_{n=1}^{\infty} a_n$, $\sum_{n=1}^{\infty} b_n$ are absolutely convergent their Cauchy product converges absolutely to $\left(\sum_{n=1}^{\infty} a_n\right) \left(\sum_{n=1}^{\infty} b_n\right)$.

Theorem 4.4.8. If the power series $\sum_{n=1}^{\infty} a_n z^n$, $\sum_{n=1}^{\infty} b_n z^n$ have radii of convergence R_a, R_b respectively then their Cauchy product has radius of convergence $R_c \geq \min(R_a, R_b)$.

4.5 Exponential series

Definition 4.5.1. For $z \in \mathbb{C}$, its **exponential series** is

$$E(z) := \sum_{n=0}^{\infty} \frac{z^n}{n!} = 1 + z + \frac{z^2}{2!} + \frac{z^3}{3!} + \dots,$$

with $E(z)$ converging absolutely for all $z \in \mathbb{C}$.

Theorem 4.5.2 (Properties of exponential series). For all $z, w \in \mathbb{C}$:

1. $E(z)E(w) = E(z+w)$, 2. $\frac{1}{E(z)} = E(-z)$, 3. $E(z) \neq 0$.

Theorem 4.5.3. For all $x \in \mathbb{Q}$, $E(x) = e^x$, with $e := E(1)$.

5 Continuity

5.1 Continuous functions

Definition 5.1.1 (Limit of real functions). For a function $f : \mathbb{R} \rightarrow \mathbb{R}$ and some $a, b \in \mathbb{R}$ we have $f(x) \rightarrow b$ as $x \rightarrow a$ of $\lim_{x \rightarrow a} f(x) = b$ iff:

$$\forall \epsilon > 0, \exists \delta > 0 \text{ such that } |x - a| < \delta \iff |f(x) - b| < \epsilon.$$

Definition 5.1.2 (Continuity of real functions). Given the function $f : \mathbb{R} \rightarrow \mathbb{R}$

1. f is **continuous at a point** $a \in \mathbb{R}$ iff $\lim_{x \rightarrow a} f(x) = f(a)$,
2. f is **continuous (on \mathbb{R})** iff f is continuous at all $a \in \mathbb{R}$.

Definition 5.1.3 (Discontinuity of real functions). The function $f : \mathbb{R} \rightarrow \mathbb{R}$ is **discontinuous** at a point if it is not continuous at that point.

Definition 5.1.4 (Sequential continuity). The function $f : \mathbb{R} \rightarrow \mathbb{R}$ is continuous at $a \in \mathbb{R} \iff f(a_n) \rightarrow f(a)$ as $n \rightarrow \infty$ for all sequences (a_n) converging to a .

Remark 5.1.5. The definition for limits and continuity of complex functions is similar with $|\cdot|$ being the modulus instead of the absolute values. The same definition also applies for functions that are continuous on certain subsets of \mathbb{R} or \mathbb{C} .

Theorem 5.1.6. $E : \mathbb{C} \rightarrow \mathbb{C}$ given by $E(z) := \sum_{n=0}^{\infty} \frac{z^n}{n!}$ is continuous on \mathbb{C} .

Theorem 5.1.7 (Properties of the real exponential function). Given the exponential function $E : \mathbb{R} \rightarrow (0, \infty)$:

1. for all $x \in \mathbb{R}$, $E(x) > 0$,
2. $x > 0 \implies E(x) > 1$,
3. $E(x)$ is a strictly increasing function,

4. For $|x| < 1$, $|E(x) - 1| \leq \frac{|x|}{1 - |x|}$,
5. E is a continuous bijection.

Theorem 5.1.8. The inverse of $E(x) = e^x$ is the **natural logarithm** function $\ln : (0, \infty) \rightarrow \mathbb{R}$ satisfying $y = \ln x \iff x = e^y$ for all $x, y \in \mathbb{R}$.

Definition 5.1.9 (Exponentiation of positive bases). For $a \in (0, \infty)$, for all $x \in \mathbb{R}$ define $a^x := E(x \ln a)$.

Definition 5.1.10 (Trigonometric functions). The **sine** and **cosine** functions are defined as:

$$\sin(\theta) := \Im[E(i\theta)], \quad \cos(\theta) := \Re[E(i\theta)].$$

and are both continuous functions from $\mathbb{R} \rightarrow [-1, 1]$.

Theorem 5.1.11 (Continuity of piecewise functions). For $a, c \in \mathbb{R}$ with functions $f_1 : (-\infty, a) \rightarrow \mathbb{R}$ and $f_2 : (a, \infty) \rightarrow \mathbb{R}$, the **piecewise function** $f : \mathbb{R} \rightarrow \mathbb{R}$, defined as,

$$f(x) := \begin{cases} f_1(x) & \text{if } x < a \\ c & \text{if } x = a \\ f_2(x) & \text{if } x > a \end{cases}$$

is continuous on \mathbb{R} iff both f_1 and f_2 are continuous on their respective domains and

$$\lim_{x \uparrow a} f_1(x) = \lim_{x \downarrow a} f_2(x) = c.$$

5.2 Properties of continuity

Theorem 5.2.1. For $f, g : \mathbb{R} \rightarrow \mathbb{R}$ continuous at $a \in \mathbb{R}$ the following functions are also continuous at a :

1. αf for all $\alpha \in \mathbb{R}$;
2. $f + g, f \cdot g$;
3. $\frac{f}{g}$, given $g(a) \neq 0$.

Theorem 5.2.2. The following functions (all by their well known definitions) are continuous:

1. $f(x) = x^n$, for $n \in \mathbb{N}_0$ (**monomials**);
2. $p(x) = \sum_{i=1}^n a_i x^i$, given (a_n) is a real sequence (**polynomials**);
3. $\frac{p(x)}{q(x)}$ at $a \in \mathbb{R}$ given p, q are polynomials with $q(a) \neq 0$ (**rational functions**);
4. $\sin(x)$, $\cos(x)$ on \mathbb{R} and $\tan(x)$ whenever $\cos(x) \neq 0$, plus their reciprocals under similar conditions;
5. $f \circ g$ at $a \in \mathbb{R}$ when g is continuous at a and f is continuous at $g(a)$.

Theorem 5.2.3 (Intermediate value theorem). Given $a, b \in \mathbb{R}$ with $a \leq b$, if $f : [a, b] \rightarrow \mathbb{R}$ is continuous, then for all c between $f(a)$ and $f(b)$ there exists some $x \in [a, b]$ such that $f(x) = c$.

Definition 5.2.4 (Boundedness of real functions). Given some $S \subseteq \mathbb{R}$ a function $f : S \rightarrow \mathbb{R}$ is **bounded above** iff $\exists M \in \mathbb{R}$ such that $f(x) \leq M$ for all $x \in S$. The definitions for **bounded below** and **bounded** extend naturally from this.

Theorem 5.2.5 (Extreme value theorem). Given $a, b \in \mathbb{R}$ with $a \leq b$, if $f : [a, b] \rightarrow \mathbb{R}$ is continuous then f is bounded.

6 Properties of subsets

6.1 Open sets

Definition 6.1.1 (Open sets). A set $S \subset \mathbb{R}$ is **open** iff $\forall x \in S, \exists \delta$ such that $(x - \delta, x + \delta) \subset S$.

Theorem 6.1.2 (Union of open sets). For a collection of open sets in \mathbb{R} , $\{S_i\}$, given the indexing set \mathcal{I} (could be countable or uncountable), $\bigcup_{i \in \mathcal{I}} S_i$ is open in \mathbb{R} .

Theorem 6.1.3 (Finite intersections of open sets). The intersection of finitely many open sets in \mathbb{R} is open in \mathbb{R} .

6.2 Closed and compact sets

Definition 6.2.1 (Closed sets). A set $S \subset \mathbb{R}$ is **closed** in \mathbb{R} if all convergent subsequences of S have a limit in S .

Definition 6.2.2 (Compact sets). A set $S \subset \mathbb{R}$ is **compact** in \mathbb{R} if it is closed and bounded in \mathbb{R} .

Theorem 6.2.3. The complement of an open set is closed.

Remark 6.2.4. Not every set in \mathbb{R} is either open or closed. Half-open intervals are neither open nor closed while \mathbb{R} and \emptyset are both open and closed.

Theorem 6.2.5. The finite union or any intersection of closed sets in \mathbb{R} is closed.

Theorem 6.2.6. A set $S \subset \mathbb{R}$ is compact iff every subsequence of S has as convergent subsequence $x_{n(i)} \rightarrow x \in S$.

Theorem 6.2.7 (Extreme value theorem for compact sets). If $S \subset \mathbb{R}$ is compact with $f : S \rightarrow \mathbb{R}$ continuous, there exists some $c, d \in S$ with $f(c) = \inf_{x \in S} f(x)$ and $f(d) = \sup_{x \in S} f(x)$.

7 Uniform continuity and convergence

7.1 Uniform continuity

Definition 7.1.1 (Uniform continuity). A function $f : S \rightarrow \mathbb{R}$ is **uniformly continuous** iff

$$\forall \epsilon > 0, \exists \delta > 0 \text{ such that } \forall x, y \in S, |x - y| < \delta \implies |f(x) - f(y)| < \epsilon.$$

Uniform continuity is a more powerful notion than continuity with f is uniformly continuous $\implies f$ is continuous.

Theorem 7.1.2. If $S \subset \mathbb{R}$ is compact and $f : S \rightarrow \mathbb{R}$ continuous then f is uniformly continuous.

7.2 Convergence of sequences of functions

Definition 7.2.1 (Pointwise convergence). For some $S \subset \mathbb{R}$ with the sequence $f_1, f_2, \dots : S \rightarrow \mathbb{R}$, f_n **converges pointwise** to some $f : S \rightarrow \mathbb{R}$ if

$$\forall x \in S, \forall \epsilon > 0, \exists N \in \mathbb{N} \text{ such that } \forall n \geq N, |f(x) - f_n(x)| < \epsilon.$$

Written $\forall x \in S, \lim_{n \rightarrow \infty} f_n(x) = f(x)$.

Definition 7.2.2 (Uniform convergence). For some $S \subset \mathbb{R}$, the sequence $f_1, f_2, \dots : S \rightarrow \mathbb{R}$ **uniformly converges** to some $f : S \rightarrow \mathbb{R}$ if

$$\forall \epsilon > 0, \exists N \in \mathbb{N} \text{ such that } \forall x \in S, \text{ and } \forall n > N, |f(x) - f_n(x)| < \epsilon.$$

Theorem 7.2.3. If a sequence of (uniformly) continuous functions converges uniformly to a function f then f is (uniformly) continuous.

Theorem 7.2.4. If, given $S \subset \mathbb{R}$, $(f_n) : S \rightarrow \mathbb{R}$ is a uniformly convergent sequence of continuous functions with $a \in S$ open in S , $\lim_{n \rightarrow \infty} \lim_{x \rightarrow a} f_n(x) = \lim_{x \rightarrow a} \lim_{n \rightarrow \infty} f_n(x)$.

7.3 Convergence of series of functions

Definition 7.3.1 (Convergence of series of functions). Given $(f_n) : S \rightarrow \mathbb{R}$ defined on $S \subset \mathbb{R}$, the series $\sum_{n=1}^{\infty} f_n(x)$ **converges (uniformly)** iff the sequence of partial sums $S_n(x) = \sum_{n=1}^n f_n(x)$ converges (uniformly).

Theorem 7.3.2 (Weierstrass M-test). Given continuous $(f_n) : S \rightarrow \mathbb{R}$ defined on $S \subset \mathbb{R}$,

$$\begin{aligned} &\forall x \in S \text{ and } \forall i \in \mathbb{N}, \exists M_1, M_2, \dots \in \mathbb{R} \text{ such that } |f_i(x)| \leq M_i \text{ and } \sum_{i=1}^{\infty} M_i \text{ converges} \\ &\implies \sum_{n=1}^{\infty} f_n(x) \text{ converges uniformly to some continuous } g : S \rightarrow \mathbb{R}. \end{aligned}$$

Theorem 7.3.3. If a power series $f(x) = \sum_{n=1}^{\infty} f_n(x)$ has radius of convergence $R > 0$ then f is continuous on $(-R, R)$.

8 Differentiation

8.1 Differentiability

Definition 8.1.1 (Differentiability). A function $f : \mathbb{R} \rightarrow \mathbb{R}$ is **differentiable** at $a \in \mathbb{R}$, with **derivative**

$$f'(a) = \left. \frac{d}{dx} f(x) \right|_a \text{ iff}$$

$$\lim_{x \rightarrow a} \frac{f(x) - f(a)}{x - a} \text{ exists, which we set to } f'(a).$$

f is differentiable on $S \subseteq \mathbb{R}$, with derivative $\frac{d}{dx} f = \frac{df}{dx} = f' : \mathbb{R} \rightarrow \mathbb{R}$, if it is differentiable at every $x \in S$.

Examples 8.1.2. The following functions are all differentiable,

- $f(x) = x^n$, for $n \in \mathbb{N}$ on \mathbb{R} with $f'(x) = nx^{n-1}$,
- $f(x) = e^x$ on \mathbb{R} with $f'(x) = e^x$,
- $f(x) = \ln x$ on $\mathbb{R}^{>0}$ with $f'(x) = \frac{1}{x}$.

Theorem 8.1.3. f is differentiable $\implies f$ is continuous.

Theorems 8.1.4 (Operations on derivatives). If $f, g : \mathbb{R} \rightarrow \mathbb{R}$ are both differentiable at $x = a \in \mathbb{R}$ then,

1. for all $c, d \in \mathbb{R}$, $h(x) := c \cdot f(x) + d \cdot g(x)$ is differentiable at $x = a$ with $h'(a) = c \cdot f'(a) + d \cdot g'(a)$,
2. $p(x) := f(x) \cdot g(x)$ is differentiable at $x = a$ with $p'(a) = f(a) \cdot g'(a) + f'(a) \cdot g(a)$,
3. if $f(a) \neq 0$, $q(x) := \frac{1}{f(x)}$ is differentiable at $x = a$ with $q'(a) = -\frac{f'(a)}{[f(a)]^2}$,
4. if $g(a) \neq 0$ $r(x) := \frac{f(x)}{g(x)}$ is differentiable at $x = a$ with $r'(a) = \frac{f'(a) \cdot g(a) - f(a) \cdot g'(a)}{[g(a)]^2}$.

Theorem 8.1.5 (Chain rule). If $g, f : \mathbb{R} \rightarrow \mathbb{R}$ are differentiable at $x = a \in \mathbb{R}$ and $x = g(a)$ respectively then $s(x) := f \circ g(x)$ is differentiable at $x = a$ with $s'(a) = g'(a) \cdot f'(g(a))$.

8.2 Local extrema and mean values

Definition 8.2.1 (Local extrema). For a function $f : S \rightarrow \mathbb{R}$, f has a **local minimum** at $a \in \mathbb{R}$ iff $\exists \delta > 0$ such that $\forall y \in S$ with $|y - a| < \delta$, $f(y) \geq f(a)$, and similarly for a **local maximum**.

Theorem 8.2.2. If $f : [a, b] \rightarrow \mathbb{R}$ is differentiable on (a, b) and has a local maximum or minimum at $c \in (a, b)$, $f'(c) = 0$.

Theorem 8.2.3 (Rolle's theorem). If $f : [a, b] \rightarrow \mathbb{R}$ is differentiable on (a, b) with $f(a) = f(b)$, $\exists c \in (a, b)$ such that $f'(c) = 0$.

Theorem 8.2.4 (Mean value theorem). If $f : [a, b] \rightarrow \mathbb{R}$ is differentiable on (a, b) , $\exists c \in (a, b)$ such that $f'(c) = \frac{f(b) - f(a)}{b - a}$.

Theorem 8.2.5. If $f : [a, b] \rightarrow \mathbb{R}$ is differentiable on (a, b) with $f'(x) \geq 0$ for all $x \in (a, b)$ then f is monotone increasing. Similar holds for monotone/strictly increasing/decreasing or constant.

Theorem 8.2.6 (Cauchy's MVT). A similar but slightly more general statement than the MVT: if $f, g : [a, b] \rightarrow \mathbb{R}$ are differentiable on (a, b) , $\exists c \in (a, b)$ with $(f(b) - f(a))g'(c) = (g(b) - g(a))f'(c)$.

8.3 L'Hôpital's rule

Theorem 8.3.1 (L'Hôpital's rule). Given $f, g : [c, d] \rightarrow \mathbb{R}$ are differentiable on (c, d) except possibly at some $a \in (c, d)$ with $g'(x) \neq 0$ on $(c, d) \setminus \{a\}$:

$$\text{if } \lim_{x \rightarrow a} f(x) = \lim_{x \rightarrow a} g(x) = 0 \text{ or } \infty \text{ and } \lim_{x \rightarrow a} \frac{f'(x)}{g'(x)} = L \text{ then } \lim_{x \rightarrow a} \frac{f(x)}{g(x)} = L.$$

This also applies when taking $\lim_{x \rightarrow \infty}$.

Definition 8.3.2 (Higher derivatives). **Higher derivatives** of $f : \mathbb{R} \rightarrow \mathbb{R}$ are defined inductively as

$$f^{(n)}(x) := \begin{cases} f(x) & \text{if } x = 0 \\ f^{(n-1)'}(x) & \text{otherwise} \end{cases}.$$

The existence of the n th derivative of f requires all lower order derivatives of f also exist and be differentiable.

Theorem 8.3.3 (Second derivative test). For a second differentiable function $f : \mathbb{R} \rightarrow \mathbb{R}$ with $f'(a) = 0$ for some $a \in \mathbb{R}$,

- $f''(a) > 0 \implies f$ has a local minimum at $x = a$,
- $f''(a) < 0 \implies f$ has a local maximum at $x = a$,
- the test is inconclusive if $f''(a) = 0$.

8.4 Taylor's theorem

Definition 8.4.1 (Taylor polynomial of a function). Given $f : [c, d] \rightarrow \mathbb{R}$ has an order $n \in \mathbb{N}_0$ derivative at $x = a \in (c, d)$, the **Taylor polynomial** of order n at $x = a$ is

$$P_n(x) := \sum_{i=0}^n \frac{f^{(i)}(a)}{i!} (x-a)^i = f(a) + \frac{f'(a)}{1!} (x-a) + \frac{f''(a)}{2!} (x-a)^2 + \dots + \frac{f^{(n)}(a)}{n!} (x-a)^n.$$

Theorem 8.4.2 (Taylor's theorem). Given $f : [c, d] \rightarrow \mathbb{R}$ has an order $n+1$, for some $n \in \mathbb{N}_0$, derivative for all $x \in (c, d)$. For $a, b \in [c, d]$ with $a \neq b$ there exists some t between a and b such that,

$$f(b) = P_n(b) + \frac{f^{(n+1)}(t)}{(n+1)!} (b-a)^{n+1}.$$

This is a further, massive generalisation of the MVT (the case when $n = 0$).

Definition 8.4.3 (Taylor series of a function). The **Taylor series**, $P(x)$, for a function $f : \mathbb{R} \rightarrow \mathbb{R}$ at $x = a$ exists if $f^{(n)}(a)$ exists for all $n \in \mathbb{N}$ and is given by

$$P(x) := \sum_{n=0}^{\infty} \frac{f^{(n)}(a)}{n!} (x-a)^n = f(a) + \frac{f'(a)}{1!} (x-a) + \frac{f''(a)}{2!} (x-a)^2 + \dots$$

Definition 8.4.4 (Analytic function). A function $f : \mathbb{R} \rightarrow \mathbb{R}$ is **analytic** if it equals its Taylor series.

8.5 Convexity

Definition 8.5.1 (Convexity of functions). A function $f : [a, b] \rightarrow \mathbb{R}$ is **convex** iff

$$\forall c, t, d \in [a, b] \text{ with } c < t < d, f(c) + \frac{f(d) - f(c)}{d - c} (t - c) \geq f(t).$$

Theorem 8.5.2. Given the function $f : [a, b] \rightarrow \mathbb{R}$ with $f''(x)$ existing on (a, b) , f is convex $\iff f''(x)$ non-negative on (a, b) .

8.6 Exchange of limits and derivatives

Theorem 8.6.1 (Criteria for exchange of limits and derivatives). Given (f_n) is a sequence of functions with $f_n : [a, b] \rightarrow \mathbb{R}$ differentiable, if $\lim_{n \rightarrow \infty} f_n(c)$ exists for some $c \in [a, b]$ and $(f'_n(x))$ converges uniformly on $[a, b]$: (f_n) converges uniformly to some differentiable f satisfying $f'(x) = \lim_{n \rightarrow \infty} f'_n(x)$.

Theorem 8.6.2 (Derivatives of power series). Given a power series $f(x) = \sum_{n=0}^{\infty} a_n x^n$ with radius of convergence $R > 0$, f has a continuous derivative on $(-R, R)$ with $f'(x) = \sum_{n=0}^{\infty} n a_n x^{n-1}$.

Corollary 8.6.3. Given a power series $f(x) = \sum_{n=0}^{\infty} a_n x^n$ with radius of convergence $R > 0$, the Taylor series of f centered at $x = 0$ is $\sum_{n=0}^{\infty} a_n x^n$.

8.7 Trigonometric properties

Definition 8.7.1 (π). Let $S = \{y > 0 : \sin(y) = 0\}$, $\pi := \inf S$.

Definition 8.7.2 (Periodic function). A function $f : \mathbb{R} \rightarrow \mathbb{R}$ is **$2L$ -periodic** iff $f(x + 2L) = f(x)$ for all $x \in \mathbb{R}$.

Theorem 8.7.3. \sin and \cos satisfy the following important properties: 1. $\sin(x)$ is odd, 2. $\cos(x)$ is even, 3. $\cos^2(x) + \sin^2(x) = 1$ for all $x \in \mathbb{R}$, 4. \sin and \cos are 2π -periodic functions.

9 Integration

9.1 Partitions

Definition 9.1.1 (Partition). A **partition**, P , of the interval $[a, b] \subset \mathbb{R}$ is a finite collection of points $x_0, x_1, \dots, x_n \in [a, b]$ such that $a = x_0 < x_1 < \dots < x_n = b$. A partition naturally splits the domain $[a, b]$ into finitely many closed intervals.

Definition 9.1.2 (Refinement). Given partitions Q, P , Q is a **refinement** of P , written $Q \prec P$, iff every point of P is also in Q .

Definition 9.1.3 (Common refinement). Given partitions P, Q the **common refinement** of P and Q is the partition R containing all points in P or Q . $R \prec P$ and $R \prec Q$.

9.2 Darboux sums

Definition 9.2.1 (Darboux sums). Given the bounded function $f : [a, b] \rightarrow \mathbb{R}$ and the partition $P = \{x_0, x_1, \dots, x_n\}$ of $[a, b]$, we will assign to each subintervals generated by P :

- a length, $\Delta x_i := x_{i+1} - x_i$,
- an infimum, $m_i := \inf_{x_i \leq t \leq x_{i+1}} f(t)$,
- a supremum, $M_i := \sup_{x_i \leq t \leq x_{i+1}} f(t)$.

Now define the **lower Darboux sum** and **upper Darboux sum** of f w.r.t. P as:

$$L(f, P) := \sum_{i=0}^{n-1} m_i \Delta x_i, \quad U(f, P) := \sum_{i=0}^{n-1} M_i \Delta x_i \quad \text{respectively.}$$

If $f : [a, b] \rightarrow \mathbb{R}$ is continuous then $L(f, P)$ and $U(f, P)$ exist. $L(f, P)$ is always less than or equal to $U(f, P)$.

Theorem 9.2.2 (Boundedness of refined Darboux sums). If $f : [a, b] \rightarrow \mathbb{R}$ is bounded with $Q \prec P$ partitions of $[a, b]$, $L(f, P) \leq L(f, Q) \leq U(f, Q) \leq U(f, P)$.

Theorem 9.2.3. Given some bounded $f : [a, b] \rightarrow \mathbb{R}$, the set $\{L(f, P) : P \text{ is a partition of } [a, b]\}$ is bounded above by any upper Darboux sum on $[a, b]$ w.r.t. f .

9.3 Darboux integral

Definition 9.3.1 (Darboux integrals). Given a bounded function $f : [a, b] \rightarrow \mathbb{R}$, the **lower Darboux integral** and **upper Darboux integral** are:

$$\int_a^b f(x) dx := \sup_P L(f, P), \quad \overline{\int_a^b} f(x) dx := \inf_P U(f, P) \quad \text{respectively.}$$

Definition 9.3.2 (Darboux integrability). If the upper and lower Darboux integral of a bounded function $f : [a, b] \rightarrow \mathbb{R}$ are equal, f is **Darboux integrable** on $[a, b]$ with

$$\int_a^b f(x) dx := \int_a^b f(x) dx = \overline{\int_a^b} f(x) dx.$$

We will now refer to Darboux integrable functions simply as **integrable**.

Theorem 9.3.3. A bounded function $f : [a, b] \rightarrow \mathbb{R}$ is integrable iff $\forall \epsilon > 0$ there exists a partition P with $U(f, P) - L(f, P) < \epsilon$. Furthermore, given a sequence of partitions (P_n) if $\lim_{n \rightarrow \infty} (U(f, P_n) - L(f, P_n)) = 0$ then

$$\int_a^b f(x) dx = \lim_{n \rightarrow \infty} (L(f, P_n)) = \lim_{n \rightarrow \infty} (U(f, P_n)).$$

Remark 9.3.4. For a bounded function $f : [a, b] \rightarrow \mathbb{R}$, f is integrable if it is, continuous, differentiable, monotone, or discontinuous at finitely many points.

9.4 Properties of integration

Theorem 9.4.1 (Monotonicity). If $f, g : [a, b] \rightarrow \mathbb{R}$ are integrable with $f(x) \leq g(x)$ for all $x \in \mathbb{R}$,

$$(1) \quad \int_a^b f(x) dx \leq \int_a^b g(x) dx. \quad (2) \quad m \cdot (b - a) \leq \int_a^b f(x) dx \leq M \cdot (b - a).$$

Theorem 9.4.2 (Boundedness). If $f : [a, b] \rightarrow \mathbb{R}$ is integrable with $m \leq f(x) \leq M$ for all $x \in \mathbb{R}$,

Theorem 9.4.3 (Linearity). If $f, g : [a, b] \rightarrow \mathbb{R}$ are integrable, for all $c, d \in \mathbb{R}$,

$$(3) \quad \int_a^b (cf(x) + dg(x)) dx = c \int_a^b f(x) dx + d \int_a^b g(x) dx. \quad (4) \quad \int_a^b f(x) dx = \int_a^c f(x) dx + \int_c^b f(x) dx.$$

Theorem 9.4.4 (Integrability on subdomains). $f : [a, b] \rightarrow \mathbb{R}$ is integrable iff $\forall c \in [a, b]$, f is integrable on $[a, c]$ and $[c, b]$ with,

Theorem 9.4.5 (Composition). If $f : [a, b] \rightarrow [m, M] \subset \mathbb{R}$, $g : [m, M] \rightarrow \mathbb{R}$ are integrable and continuous respectively, $h(x) := g \circ f(x)$ is integrable on $[a, b]$.

Theorem 9.4.6 (Triangle inequality). If $f : [a, b] \rightarrow \mathbb{R}$ is integrable then $|f|$ is integrable on $[a, b]$ with,

$$(6) \quad \left| \int_a^b f(x) dx \right| \leq \int_a^b |f(x)| dx. \quad (7) \quad \int_a^b f(x) dx = \int_a^b g(x) dx.$$

Theorem 9.4.7 (Finite point differences). If $f, g : [a, b] \rightarrow \mathbb{R}$ are integrable with $f(x) = g(x)$ except at finitely many points,

Theorem 9.4.8 (Products). If $f, g : [a, b] \rightarrow \mathbb{R}$ are integrable then $f \cdot g : [a, b] \rightarrow \mathbb{R}$ is integrable.

Theorem 9.4.9 (Maxima and minima). If $f, g : [a, b] \rightarrow \mathbb{R}$ are integrable then $\max(f, g), \min(f, g) : [a, b] \rightarrow \mathbb{R}$ are integrable.

9.5 Fundamental theorems of calculus

Theorem 9.5.1 (Fundamental theorem of calculus 1). Given continuous $f : [a, b] \rightarrow \mathbb{R}$, have $F : [a, b] \rightarrow \mathbb{R}$ with $F(x) := \int_a^x f(t) dt$. F is continuous on $[a, b]$ and differentiable on (a, b) . $F'(x) = f(x)$ for all $x \in [a, b]$.

Theorem 9.5.2 (Fundamental theorem of calculus 2). Given continuous $f : [a, b] \rightarrow \mathbb{R}$ with continuous derivative on (a, b) , $\int_a^b f'(x) dx = f(b) - f(a)$.

9.6 Methods of integration

Theorem 9.6.1 (MVT). If $f : [a, b] \rightarrow \mathbb{R}$ is continuous, $\exists c \in [a, b]$ such that $\int_a^b f(x) dx = f(c)(b - a)$.

Theorem 9.6.2 (Integration by parts). If $f, g : [a, b] \rightarrow \mathbb{R}$ have continuous first derivatives,

$$(2) \quad \int_a^b f(x)g'(x) dx = \left[f(x)g(x) \right]_a^b - \int_a^b f'(x)g(x) dx. \quad (3) \quad \int_{u(c)}^{u(d)} f(x) dx = \int_c^d f(u(x))u'(x) dx.$$

Theorem 9.6.3 (Integration by substitution). Given continuous $f : [a, b] \rightarrow \mathbb{R}$ if $u : [a, b] \rightarrow [c, d]$ has a continuous derivative on (c, d) ,

9.7 Limits and integrals

Theorem 9.7.1 (Exchanging limits and integrals). If $f_n : [a, b] \rightarrow \mathbb{R}$ is a sequence of integrable functions converging uniformly to $f : [a, b] \rightarrow \mathbb{R}$, then f is integrable with,

$$\int_a^b f(x) \, dx = \int_a^b \lim_{n \rightarrow \infty} f_n(x) \, dx = \lim_{n \rightarrow \infty} \int_a^b f_n(x) \, dx.$$

Theorem 9.7.2 (Power series integration). If the power series $f(x) = \sum_{n=0}^{\infty} a_n x^n$ has radius of convergence $R > 0$, f is integrable on all closed subintervals of $(-R, R)$ with

$$\int_0^x f(t) \, dt = \sum_{n=0}^{\infty} \frac{a_n}{n+1} x^{n+1} \text{ for all } x \in (-R, R).$$

9.8 Improper integrals

Definition 9.8.1 (Improper integral). Given $f : (a, b] \rightarrow \mathbb{R}$ integrable on all $[c, b] \subset (a, b]$, the **improper integral**,

$$\int_a^b f(x) \, dx = \lim_{c \downarrow a} \int_c^b f(x) \, dx,$$

if the limit exists, otherwise the integral **diverges**; and similarly for other non-closed intervals or those with $\pm\infty$ as bounds.

Remark 9.8.2. When integrating over intervals with multiple undefined points, the integral is split into sums of multiple integrals each with single undefined points on their boundaries.

Chapter 2

Linear Algebra

Lectured by Dr Charlotte Kestner
Typed by Yu Coughlin
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Introduction

The following are supplementary reading:

- S Lang, Linear algebra, 1987
- G Strang, Introduction to linear algebra, 2023
- S Axler, Linear Algebra Done Right, 2015

Lecture 1
Thursday
10/01/19

1 Linear Systems and matrices

1.1 Linear systems

Definition 1.1.1 (Linear system). A **linear system** is a set of linear equations in the same variables.

Notation 1.1.2. The follow are all equivalent notation for the same linear system:

$$\begin{aligned}
 & \begin{array}{ccccccccc}
 a_{11}x_1 & + & a_{12}x_2 & + & \dots & + & a_{1n}x_n & = & b_1 \\
 a_{21}x_1 & + & a_{22}x_2 & + & \dots & + & a_{2n}x_n & = & b_2 \\
 \vdots & & \vdots & & \ddots & & \vdots & & \vdots \\
 a_{m1}x_1 & + & a_{m2}x_2 & + & \dots & + & a_{mn}x_n & = & b_m
 \end{array} \\
 & \iff \begin{pmatrix} a_{11} & a_{12} & \dots & a_{1n} \\ a_{21} & a_{22} & \dots & a_{2n} \\ \vdots & \vdots & \ddots & \vdots \\ a_{m1} & a_{m2} & \dots & a_{mn} \end{pmatrix} \begin{pmatrix} x_1 \\ x_2 \\ \vdots \\ x_n \end{pmatrix} = \begin{pmatrix} b_1 \\ b_2 \\ \vdots \\ b_m \end{pmatrix} \\
 & \iff \left(\begin{array}{cccc|c}
 a_{11} & a_{12} & \dots & a_{1n} & b_1 \\
 a_{21} & a_{22} & \dots & a_{2n} & b_2 \\
 \vdots & \vdots & \ddots & \vdots & \vdots \\
 a_{m1} & a_{m2} & \dots & a_{mn} & b_m
 \end{array} \right).
 \end{aligned}$$

1.2 Matrix algebra

Definition 1.2.1 (Matrix by elements). An $m \times n$ matrix written as $A = [a_{ij}]_{m \times n}$ has the element a_{ij} in the i th row and j th column.

Definition 1.2.2 (Matrix addition). If $A = [a_{ij}]_{m \times n}$ and $B = [b_{ij}]_{m \times n}$ then $A + B := [a_{ij} + b_{ij}]_{m \times n}$.

Definition 1.2.3 (Scalar multiplication). If $A = [a_{ij}]_{m \times n}$ then $\lambda A := [\lambda a_{ij}]_{m \times n}$.

Definition 1.2.4 (Matrix multiplication). If $A = [a_{ij}]_{p \times q}$ and $B = [b_{ij}]_{q \times r}$ then $AB := C = [c_{ij}]_{p \times r}$ where $c_{ij} = \sum_{k=1}^q a_{ik}b_{kj}$.

Theorem 1.2.5. Matrix multiplication is associative.

Remark 1.2.6. Matrix multiplication is not commutative.

1.3 EROs

Definition 1.3.1 (Elementary row operations). The three **elementary row operations (EROs)** that can be performed on augmented matrixes are as follows:

1. Multiply a row by a non-zero scalar.
2. Swap two rows.
3. Add a scalar multiple of a row to another row.

Remark 1.3.2. EROs preserve the set of solutions of a linear system. Each ERO has an inverse.

Definition 1.3.3 (Equivalence of linear systems). Two systems of linear equations are equivalent iff either:

1. They are both inconsistent.
2. (wlog) The augmented matrix of the first system can be transformed to the augmented matrix of the second system with just EROs.

Definition 1.3.4 (Row echelon form / Echelon form / REF). A matrix is in **row echelon form** if it satisfies the following:

1. All of the zero rows are at the bottom of the matrix,
2. The first non-zero entry in any row is **1**,
3. The first non-zero entry in row i is strictly to the left of the first non-zero entry in row $i + 1$.

Definition 1.3.5 (Reduced row echelon form / Row reduced echelon form / rREF). A matrix is in **reduced row echelon form** if it is in REF and the first non-zero entry in a row is the only non-zero entry in its column.

1.4 Matrices of note

Definition 1.4.1 (Square matrix). A matrix is **square** iff it has the same number of rows and columns.

Definition 1.4.2. A square matrix $(A = [a_{ij}]_{n \times n})$ is: 1. **Upper triangular** iff $i > j \implies a_{ij} = 0$. 2. **Lower triangular** iff $i < j \implies a_{ij} = 0$. 3. **Diagonal** iff $i \neq j \implies a_{ij} = 0$.

Definition 1.4.3 (Identity matrix). The **identity matrix** of size n written I_n , is the square diagonal matrix of size n with all diagonal entries equal 1.

Definition 1.4.4 (Elementary matrix). An **elementary matrix** is a matrix that can be achieved by applying a single ERO to the identity matrix.

Definition 1.4.5 (Inverse). For a square matrix B if there exists a matrix B^{-1} such that $BB^{-1} = I = B^{-1}B$ then B^{-1} is the **inverse** of B and vice versa.

Definition 1.4.6 (Singular). A matrix without an inverse is **singular**.

Theorem 1.4.7. The inverse of a matrix is unique.

Definition 1.4.8. A **transpose** of the matrix $A = [a_{ij}]_{m \times n}$ is $A^T := [a_{ji}]_{n \times m}$.

Theorem 1.4.9. If the matrix A has an inverse then its transpose has an inverse with $(A^T)^{-1} = (A^{-1})^T$.

Theorem 1.4.10. If a matrix $A \in M_{m \times n}$ can be reduced to I_n by a sequence of EROs then A is invertible with A^{-1} given by applying the same sequence of EROs to I_n .

Definition 1.4.11. A matrix A is **orthogonal** if it has an inverse with $A^{-1} = A^T$.

Theorem 1.4.12. An orthogonal matrix A satisfies the condition $(Ax) \cdot (Ay) = x \cdot y$, where \cdot is the dot product.

2 Vector Spaces

The notion of a vector space is a structure designed to generalise that of real vectors, so before developing them we must first produce a generalisation of the real numbers.

2.1 Fields

Definition 2.1.1 (Field). A **field** is a set \mathbb{F} equipped with the binary operations **addition** $+: \mathbb{F} \times \mathbb{F} \rightarrow \mathbb{F}$ and **multiplication** $\cdot: \mathbb{F} \times \mathbb{F} \rightarrow \mathbb{F}$ satisfying the follow axioms:

- A1 $\forall x, y \in \mathbb{F} : x + y = y + x$ (commutativity of addition),
- A2 $\forall x, y, z \in \mathbb{F} : x + (y + z) = (x + y) + z$ (associativity of addition),
- A3 $\exists 0_{\mathbb{F}} \in \mathbb{F}$ such that $\forall x \in \mathbb{F} : x + 0_{\mathbb{F}} = x$, (additive identity element),
- A4 $\forall x \in \mathbb{F}, \exists (-x) \in \mathbb{F}$ such that $x + (-x) = 0_{\mathbb{F}}$, (additive inverse);
- M1 $\forall x, y \in \mathbb{F} : x \cdot y = y \cdot x$ (commutativity of multiplication),
- M2 $\forall x, y, z \in \mathbb{F} : x \cdot (y \cdot z) = (x \cdot y) \cdot z$ (associativity of multiplication),
- M3 $\exists 1_{\mathbb{F}} \in \mathbb{F}$ such that $\forall x \in \mathbb{F} : x \cdot 1_{\mathbb{F}} = x$, (multiplicative identity element),
- M4 $\forall x \in \mathbb{F} \setminus \{0_{\mathbb{F}}\}, \exists x^{-1} \in \mathbb{F}$ such that $x \cdot x^{-1} = 1_{\mathbb{F}}$, (multiplicative inverse);
- D $\forall x, y, z \in \mathbb{F} : x \cdot (y + z) = x \cdot y + x \cdot z$ (distributivity of multiplication over addition).

The field $(\mathbb{F}, +, \cdot)$ is often referred to as just \mathbb{F} .

Example 2.1.2. The familiar sets $\mathbb{R}, \mathbb{Q}, \mathbb{C}$ are all fields.

Theorem 2.1.3. If $p \in \mathbb{N}$ is prime with $\mathbb{F}_p = \{0, 1, \dots, p-1\}$ then $(\mathbb{F}_p, +_p, \cdot_p)$ is a field.

2.2 Vector spaces

Definition 2.2.1 (Vector space). A **vector space** over a field \mathbb{F} is a set V equipped with the binary operations **vector addition** $\oplus : V \times V \rightarrow V$ and **scalar multiplication** $\odot : \mathbb{F} \times V \rightarrow V$ satisfying the follow axioms:

- A1 $\forall u, v, w \in V : u \oplus (v \oplus w) = (u \oplus v) \oplus w$ (associativity of addition),
- A2 $\forall u, v \in V : u \oplus v = v \oplus u$ (commutativity of vector addition),
- A3 $\exists 0_V \in V$ such that $\forall v \in V : v \oplus 0_V = v$, (vector additive identity element),
- A4 $\forall v \in V, \exists (-v) \in V$ such that $v \oplus (-v) = 0_V$, (vector additive inverse),
- A5 $\forall x \in \mathbb{F}, \forall u, v \in V : x \odot (u \oplus v) = (x \odot u) \oplus (x \odot v)$ (vector distributivity 1),
- A6 $\forall x, y \in \mathbb{F}, \forall v \in V : x \cdot (y \odot v) = (x \cdot y) \odot v$ (vector distributivity 2),
- A7 $\forall x, y \in \mathbb{F}, \forall v \in V : (x \cdot y) \odot v = x \odot (y \odot v)$ (scalar multiplication associativity),
- A8 $\forall v \in V : 1_{\mathbb{F}} \odot v = v$, (scalar multiplication identity element).

If V is a vector space over \mathbb{F} we say V is an \mathbb{F} -vector space with $v \in V$ a **vector** and $x \in \mathbb{F}$ a **scalar**.

2.3 Subspaces

Definition 2.3.1 (Subspace). A subset $W \subseteq V$ is a **subspace** of V , denoted $W \leq V$ iff:

- S1 $W \neq \emptyset$,
- S2 $\forall w_1, w_2 \in W : w_1 \oplus w_2 \in W$,
- S3 $\forall x \in \mathbb{F}, \forall w \in W : x \odot w \in W$.

If $W = \{0_V\}$ then W is the **trivial subspace**.

Theorem 2.3.2. Every subspace of V contains 0_V .

Theorem 2.3.3. If U and W are subspaces of V , $U \cap W$ is a subspace of V .

3 Spanning and Linear Independence

Throughout this section, assume V is an \mathbb{F} -vector space.

3.1 Spanning

Definition 3.1.1 (Span). Given some set $\{v_1, v_2, \dots, v_n\} \subseteq V$ define the **span** by,

$$\text{Span}(\{v_1, v_2, \dots, v_n\}) := \{u \in V : u = \sum_{i=1}^n \alpha_i v_i \text{ with } \alpha_i \in \mathbb{F}\}.$$

Note that the span of a subset of V is always a subspace of V .

Remark 3.1.2. If $S \subseteq V$ is infinite, $\text{Span}(S)$ is the set of all **finite** linear combinations of elements of S .

Definition 3.1.3 (Spanning sets). If $S \subseteq V$ and $\text{Span}(S) = V$, we say S **spans** V or S is a **spanning set** for V .

3.2 Linear independence

Definition 3.2.1. The set $\{v_1, v_2, \dots, v_n\} \subseteq V$ is **linearly independent** in V iff:

$$\sum_{i=1}^n \alpha_i v_i = 0_V \iff \alpha_i = 0_{\mathbb{F}} \text{ for all } i \in [1, n].$$

Theorem 3.2.2. If $S = \{v_1, v_2, \dots, v_n\} \subseteq V$ is linearly independent in V with $v_{n+1} \in V \setminus \text{Span}(S)$ then $S \cup \{v_{n+1}\}$ is also linearly independent in V .

4 Bases

4.1 Definition

Again, assume V is an \mathbb{F} -vector space throughout this section.

Definition 4.1.1 (Bases). A **basis** for V is linearly independent, spanning set of V . If V has a finite bases then V is said to be a **finite dimensional** vector space.

Theorem 4.1.2. Any $S \subseteq V$ is a basis for V iff every vector in V can be uniquely expressed as a linear combination of the elements of S .

Theorem 4.1.3. If V is non-trivial and S is a finite spanning set of V then S contains a basis for V .

Lemma 4.1.4 (Steinitz Exchange Lemma). Given some $X \subseteq V$ with $u \in \text{Span}(X)$ but $u \notin \text{Span}(X \setminus \{v\})$ for some $v \in X$, let $Y = (X \setminus \{v\}) \cup \{u\}$ then $\text{Span}(X) = \text{Span}(Y)$.

Theorem 4.1.5. Given a LI $S \subseteq V$ and spanning set $T \subseteq V$, $|S| \leq |T|$.

Corollary 4.1.6. If S and T are both bases for V , $|S| = |T|$.

4.2 Dimension

Definition 4.2.1 (Dimension of a vector space). If V is finite dimensional then the **dimension** of V , $\dim V$, is the size of any basis of V .

Definition 4.2.2 (Notable subspaces). Let U and W both be subspaces of V , the **intersection** of U and W :

$$U \cap W := \{v \in V : v \in U \text{ and } v \in W\}$$

is a subspace of V , and the **sum** of U and W :

$$U + W := \{u + w : u \in U, w \in W\}$$

is also a subspace of V .

Theorem 4.2.3. Let U and W both be subspaces of V , we have:

$$\dim(U + W) = \dim U + \dim W - \dim(U \cap W).$$

5 Matrix rank

Definition 5.0.1. Given a field \mathbb{F} and a matrix $A \in M_{m \times n}(\mathbb{F})$ we have:

- the **row space** of A , $\text{RSp}(A)$, as the span of the rows of A , this is a subspace of \mathbb{F}^n ,
- the **row rank** of A , is $\dim(\text{RSp}(A))$,
- the **column space** of A , $\text{CSp}(A)$, as the span of the columns of A , this is a subspace of \mathbb{F}^m ,
- the **column rank** of A , is $\dim(\text{CSp}(A))$.

Theorem 5.0.2. For any matrix A , the row rank of A is equal to the column rank of A .

Definition 5.0.3 (Rank of a matrix). The **rank** of a matrix A , $\text{rank}(A)$, is equal to the row/column rank of A .

Theorem 5.0.4. Given a field \mathbb{F} and a matrix $A \in M_{n \times n}(\mathbb{F})$ with $\text{rank}(A) = n$:

- the rows of A for a basis for \mathbb{F}^n ,
- the columns of A for a basis for \mathbb{F}^n ,
- A is invertible.

6 Linear transformations

6.1 Definition

Definition 6.1.1 (Linear transformation). Given \mathbb{F} -vector spaces V and W , let $T : V \rightarrow W$ be a function, T is a **linear transformation** iff the following two properties hold:

1. T **preserves vector addition**: $\forall v_1, v_2 \in V$ we have $T(v_1 +_V v_2) = T(v_1) +_W T(v_2)$,
2. T **preserves scalar multiplication**: $\forall v \in V$ and $\forall \lambda \in \mathbb{F}$ we have $\lambda T(v) = T(\lambda v)$.

Definition 6.1.2 (Identity transformation). The **identity transformation** of the vector space V is the linear transformation $\text{Id}_V : V \rightarrow V$ with $\text{Id}_V(v) := v$ for all $v \in V$.

Definition 6.1.3 (Linear transformation of a matrix). If $A \in M_{m \times n}(\mathbb{F})$ then we can define $T : \mathbb{F}^n \rightarrow \mathbb{F}^m$ by $T(v) := Av$, T is a linear transformation.

Theorem 6.1.4. If V and W are \mathbb{F} -vector spaces, $T(0_V) = 0_W$ and

$$v = \lambda_1 v_1 + \dots + \lambda_n v_n \iff T(v) = \lambda_1 T(v_1) + \dots + \lambda_n T(v_n).$$

Theorem 6.1.5. If V and W are \mathbb{F} -vector spaces, $\text{Hom}(V, W)$ is the set of linear transformations from V to W , with pointwise addition and scalar multiplication $\text{Hom}(V, W)$ is a \mathbb{F} -vector space.

6.2 Image and kernel

Throughout, assume $T : V \rightarrow W$ is a linear transformation and V, W are \mathbb{F} -vector spaces

Definition 6.2.1 (Image). We define the **image** of T , denoted $\text{Im } T$, as

$$\text{Im } T := \{w \in W : \exists v \in V, T(v) = w\},$$

with $\text{Im } T$ being a subspace of W .

Definition 6.2.2 (Kernel). We define the **kernel** of T , denoted $\ker T$, as

$$\ker T := \{v \in V : T(v) = 0_W\},$$

with $\ker T$ being a subspace of V .

Theorem 6.2.3. If $v_1, v_2 \in V$ then $T(v_1) = T(v_2) \iff v_1 - v_2 \in \ker T$.

Theorem 6.2.4. If $\{v_1, v_2, \dots, v_n\}$ is a basis for V , then $\text{Im } T = \text{Span}(\{T(v_1), T(v_2), \dots, T(v_n)\})$.

Remark 6.2.5. If T is the linear transformation for some matrix $A \in M_{m \times n}(\mathbb{F})$ then, $\ker T$ is the set of solutions for $Av = 0$, $\text{Im } T$ is the column space of A , and $\dim(\text{Im } T) = \text{rank } A$

6.3 Rank nullity

Theorem 6.3.1 (Rank Nullity Theorem). If V and W are finite dimensional \mathbb{F} -vector spaces and $T : V \rightarrow W$ is a linear transformation, we have:

$$\dim(\text{Im } T) + \dim(\ker T) = \dim V.$$

7 Representations

7.1 Matrices of transformations

Throughout this subsection let V be an n -dimensional \mathbb{F} -vector space and $B = \{e_1, e_2, \dots, e_n\}$ be a basis for V .

Definition 7.1.1 (Representation of a vector). Given some $v \in V$ with $v = \lambda_1 v_1 + \lambda_2 v_2 + \cdots + \lambda_n v_n$ for $\lambda_i \in \mathbb{F}$, we define the v **with respect to** (w.r.t.) B as

$$[v]_B := \begin{pmatrix} \lambda_1 \\ \lambda_2 \\ \vdots \\ \lambda_n \end{pmatrix} \in \mathbb{F}^n.$$

Remark 7.1.2. This must be well defined as all vectors have a unique representation in terms of every basis.

Definition 7.1.3 (Linear isomorphism). A **linear isomorphism** is a bijective linear transformation.

Theorem 7.1.4. The linear transformation $T : V \rightarrow \mathbb{F}^n$ given by $T(v) := [v]_B$ is a linear isomorphism.

7.2 Matrices of transformations

Definition 7.2.1 (Representation of a linear transformation). Given finite dimensional \mathbb{F} -vector spaces V and W with bases $B = \{v_1, v_2, \dots, v_n\}$, $C = \{w_1, w_2, \dots, w_n\}$ respectively, the **matrix of T w.r.t. B and C** denoted ${}_C[T]_B$ is $m \times n$ matrix with the i th column given by $[T(v_i)]_C$. ${}_B[T]_B$ is often shortened to $[T]_B$.

Remark 7.2.2. ${}_C[T]_B[v]_B = [T(v)]_C$, for all $v \in V$.

Theorem 7.2.3. Given a finite dimensional \mathbb{F} -vector space V with bases $B = \{v_1, v_2, \dots, v_n\}$ and $C = \{w_1, w_2, \dots, w_n\}$, if $v_i = \lambda_{1i} w_1 + \lambda_{2i} w_2 + \cdots + \lambda_{ni} w_n$ and P is the matrix given by $P = [\lambda_{ij}]_{n \times n}$, we have:

- $P = [X]_C$ where X is the unique linear transformation given by $X(w_i) = v_i$ for all $i \in [1, n]$,
- $P([v]_B) = [v]_C$ for all $v \in V$,
- $P = {}_C[\text{Id}_V]_B$.

P is often also called the **change of basis matrix** from B to C .

Corollary 7.2.4. P is invertible with $(P)^{-1} = ({}_C[\text{Id}_V]_B)^{-1} = {}_B[\text{Id}_V]_C$.

Theorem 7.2.5. If $T : V \rightarrow V$ is a linear transformation $[T]_C = ({}_C[\text{Id}_V]_B)[T]_B({}_B[\text{Id}_V]_C)$.

8 Determinants

8.1 Definition

Definition 8.1.1 (Minor of matrix). Given a matrix $A \in M(\mathbb{F})_n$ the **ij th-minor** of the matrix A , $A_{ij} \in M(\mathbb{F})_n$, is A with row i and column j removed.

Definition 8.1.2 (Determinant). The **determinant** of the matrix A is defined recursively by

$$\det(A) := \begin{cases} a_{11} & \text{if } A \text{ is a matrix with a single row and column} \\ \sum_{j=1}^n (-1)^{j+1} a_{1j} \det(A_{1j}) & \text{otherwise.} \end{cases}.$$

The determinant is only a function on square matrices.

Theorem 8.1.3. If a matrix A is singular, $\det(A) = 0$.

Theorem 8.1.4. If A is invertible then the columns of A are LI.

8.2 Properties

Definition 8.2.1 (EROs). The three ERO's on the matrix A to form A' have the following effects on the determinant:

- multiplying a row by $\lambda \neq 0$, $\det(A') = \lambda \det(A)$;
- swapping two rows, $\det(A') = -\det(A)$;
- adding a scalar multiple of one row to another, $\det(A') = \det(A)$.

Definition 8.2.2 (Other miscellaneous properties). For obvious types:

- If A , B and C all only differ in the i th row with the i th row of C being the sum of the i th row of A and B , $\det(C) = \det(A) + \det(B)$,
- if a matrix A has two identical rows, $\det(A) = 0$,
- $\det(AB) = \det(A)\det(B)$,
- $\det(A^T) = \det(A)$,
- $\det(I_n) = 1$.

Definition 8.2.3 (Cofactor). The ij th cofactor of a matrix A is defined as,

$$c_{ij} := (-1)^{i+j} \det(A_{ij}).$$

The **matrix of cofactors** of A is defined as $C = [c_{ij}]_{n \times n}$ where c_{ij} is the ij th cofactor of A .

Theorem 8.2.4. For a matrix A with matrix of cofactors C , $C^T A = \det(A) I_n$.

Theorem 8.2.5 (Cramer's Rule). ugh

Definition 8.2.6. The **determinant** of a linear transformation $T : V \rightarrow V$ where B is a basis for T , $\det(T) = \det([T]_B)$. This definition says, rather importantly, that the determinant of the matrix of linear transformation is independent of the basis that linear transformation is represented in.

9 Eigen-things

The prefix “eigen” comes from the German word “eigen” which can be roughly translated to mean “proper” or “characteristic”.

9.1 Eigenvectors and eigenvalues

Definition 9.1.1 (Eigenvectors and eigenvalues). Given the finite dimensional \mathbb{F} -vector space, V , and the linear transformation $T : V \rightarrow V$, we say $v \in V \setminus \{0_V\}$ is an **eigenvector** of T if it satisfies the equation $T(v) = \lambda v$ for some $\lambda \in \mathbb{F}$, we call λ the corresponding **eigenvalue**.

Definition 9.1.2 (Eigenspace). The **eigenspace** of an eigenvalue of a given linear transformation $T : V \rightarrow V$ is the set of eigenvectors that correspond to said eigenvalue. The eigenspace of any eigenvalue λ of T is a subspace of V .

Remark 9.1.3. The eigenvectors, eigenvalues and eigenspaces of a matrix are defined obviously and do not depend on which basis the linear transformation is represented in.

9.2 Characteristic polynomial

Theorem 9.2.1. If D is a square diagonal matrix, D^k is D with its entries raised to the power of k .

Definition 9.2.2 (Characteristic polynomial). Given the finite dimensional \mathbb{F} -vector space V with basis B and the linear transformation $T : V \rightarrow V$, we define the **characteristic polynomial** of T , $\chi_T : \mathbb{F} \rightarrow \mathbb{F}$ by $\chi_T(\lambda) := \det(\lambda I_n - T_B)$.

Theorem 9.2.3. The characteristic polynomial of a linear transformation is independent of the basis it is represented in.

Remark 9.2.4. Therefore, the characteristic polynomial of a matrix can be defined as the characteristic polynomial of the linear transformation it represents.

9.3 Diagonalisation

Definition 9.3.1 (Diagonalisability). Given a finite dimensional \mathbb{F} -vector space V , a linear transformation $T : V \rightarrow V$ is **diagonalisable** if there exists a basis for V consisting of eigenvectors of T . Similarly the matrix $A \in M_n(\mathbb{F})$ is **diagonalisable** if there exists a basis to \mathbb{F}^n as eigenvectors of A .

Theorem 9.3.2. If V is a n -dimensional vector space and $T : V \rightarrow V$ has n distinct eigenvalues, T is diagonalisable.

Theorem 9.3.3. If a matrix $A \in M_n(\mathbb{F})$ is diagonalisable, let P be the matrix with columns as eigenvectors of A and D be the diagonal matrix with i th entry as the corresponding eigenvalue for the i th column of P , then $A = PDP^{-1}$.

Theorem 9.3.4 (Precursor of Cayley-Hamilton). If a matrix $A \in M_n(\mathbb{F})$ is diagonalisable, $\chi_A(A) = 0$ where 0 is the zero-matrix in $M_n(\mathbb{F})$.

Note 9.3.5. The full **Cayley-Hamilton theorem**, which applies to all $A \in M_n(\mathbb{F})$, requires more advanced reasoning, this simpler version only requires the expansion of the polynomial and diagonalisation.

10 Orthogonality

10.1 Inner product spaces*

Definition 10.1.1 (Inner product). Let V be an n -dimensional \mathbb{F} -vector space, a **inner product** on V is a bilinear map $\langle \cdot, \cdot \rangle : V \times V \rightarrow \mathbb{F}$ satisfying the following:

- $\langle u, v \rangle = \langle v, u \rangle$ for all $u, v \in V$ (Symmetry),
- $\langle u + v, w \rangle = \langle u, w \rangle + \langle v, w \rangle$ and $\langle \lambda u, v \rangle = \lambda \langle u, v \rangle$ for all $u, v, w \in V$ and $\lambda \in \mathbb{F}$ (Bilinearity),
- $\langle v, v \rangle \geq 0$ for all $v \in V$ with equality when $v = 0_V$ (Positive-definite).

Bilinearity must hold in both arguments however it can be derived from a single argument with the symmetry property.

Definition 10.1.2 (Norm). Given a real-vector space V with an inner product $\langle \cdot, \cdot \rangle$ the **norm** induced by the inner product of $v \in V$ is:

$$\|v\| := \sqrt{\langle v, v \rangle}$$

Definition 10.1.3 (Orthogonality). Two vectors u, v in an real or complex vector space V with inner product $\langle \cdot, \cdot \rangle$ are **orthogonal** iff: $\langle u, v \rangle = 0$.

10.2 Orthonormal sets

Throughout the remainder of this section we will assume all vector spaces are over the real or complex numbers and will use the dot product as our inner product with its induced norm.

Definition 10.2.1 (Orthogonal sets). A set of vectors $\{v_1, v_2, \dots, v_n\}$ in a vector space is **orthogonal** if it is pairwise orthogonal.

Definition 10.2.2 (Orthonormal sets). A set of vectors $\{v_1, v_2, \dots, v_n\}$ in a vector space is **orthonormal** if it is orthogonal and satisfies $\|u_i\| = 1$ for all $i \in [1, n]$.

Theorem 10.2.3. The columns of an orthogonal matrix $P \in M_n(\mathbb{R})$ form an orthonormal set.

10.3 Gram-Schmidt process

The Gram-Schmidt process is a method of producing orthonormal bases.

Algorithm 10.3.1 (Gramm-Shmidt process). Given a LI set $\{v_1, v_2, \dots, v_r\} \in \mathbb{R}^n$ the **Gramm-Shmidt process** will produce the set of vectors $\{w_1, w_2, \dots, w_r\} \in \mathbb{R}^n$ by the following:

$$\begin{aligned} w_1 &= v_1, \\ w_2 &= v_2 - \frac{w_1 \cdot v_2}{\|w_1\|^2} w_1, \\ w_3 &= v_3 - \left(\frac{w_1 \cdot v_3}{\|w_1\|^2} w_1 + \frac{w_2 \cdot v_3}{\|w_2\|^2} w_2 \right), \\ &\vdots \\ w_r &= v_r - \sum_{j=1}^{r-1} \frac{w_j \cdot v_r}{\|w_j\|^2} w_j. \end{aligned}$$

Note that each vector is the original vector v_i with its projection along all of the previous w_j s subtracted and therefore $\{w_1, w_2, \dots, w_r\}$ is orthogonal. Finally, $\{u_1, u_2, \dots, u_r\}$, where $u_i = \frac{w_i}{\|w_i\|}$ for all $i \in [1, r]$, is an orthonormal set with $\text{Span}(\{u_1, u_2, \dots, u_r\}) = \text{Span}(\{v_1, v_2, \dots, v_r\})$.

Corollary 10.3.2. Given some vector $u \in \mathbb{R}^n$ there exists an orthogonal matrix in $M_n(\mathbb{R})$ with u as its first column.

11 Real symmetric matrices

Throughout this section, unsurprisingly, all matrices will be assumed to be real.

11.1 Introduction

Definition 11.1.1 (Self-adjoint matrices). If a matrix $A \in M_n(\mathbb{R})$ is symmetric and satisfies $A(u \cdot v) = (Au) \cdot v$ for all vectors $u, v \in \mathbb{R}^n$, we say A is **self-adjoint** w.r.t. the usual scalar product.

Theorem 11.1.2. If $A \in M_n(\mathbb{R})$ is symmetric with $\lambda \in \mathbb{C}$ a root of $\chi_A(x) = 0$, $\lambda \in \mathbb{R}$.

Corollary 11.1.3. Real symmetric matrices have at least 1 real eigenvalue.

Theorem 11.1.4. If $A \in M_n(\mathbb{R})$ is symmetric with discrete eigenvalues $\lambda, \mu \in \mathbb{R}$, their corresponding eigenvectors $u, v \in \mathbb{R}^n$ satisfy $u \cdot v = 0$.

11.2 Spectral theorem

Theorem 11.2.1 (Spectral theorem). If $A \in M_n(\mathbb{R})$ is symmetric, then there exists an orthonormal matrix P such that $P^{-1}AP$ is diagonal.

Corollary 11.2.2. Appropriately scaled eigenvectors of a symmetric matrix $A \in M_n(\mathbb{R})$ form an orthonormal basis for \mathbb{R}^n .

Chapter 3

Groups

Lectured by Dr Michele Zordan

Typed by Yu Coughlin

Spring 2024

Introduction

The following are supplementary reading:

- J B Fraleigh, A first course in abstract algebra, 2014
- R B J T Allenby, Rings, field and groups: an introduction to abstract algebra, 1991
- A W Knap, Basic Algebra, 2006

Lecture 1

Thursday

10/01/19

1 Binary operations and groups

Definition 1.0.1 (Binary operation). Given a set G a **binary operation** on G is a mapping $\cdot : G \times G \rightarrow G$ written $\cdot(g, h) = g \cdot h$ (and sometimes gh) for all $g, h \in G$.

Definition 1.0.2 (Group). A **group** is a pair $G = (G, \cdot)$, for some set G and a binary operation \cdot , satisfying the following properties:

- G1 $(a \cdot b) \cdot c = a \cdot (b \cdot c)$ for all $a, b, c \in G$ - the binary operation is **associative**,
- G2 $\exists e \in G$ such that $\forall g \in G, g \cdot e = e \cdot g = g$ - there is an **identity** element,
- G3 $\forall g \in G, \exists g^{-1} \in G$ such that $g \cdot g^{-1} = g^{-1} \cdot g = e$ - every element has an **inverse**.

In some literature, the condition of **closure** is also required however this is given in the fact that \cdot is a binary operation on G .

Theorem 1.0.3 (Uniqueness). The identity element for some group G is unique. The inverse, g^{-1} , of any element $g \in G$ is also unique.

Lemma 1.0.4 (Inverse of product). Given a group G and the elements $g_1, g_2, \dots, g_n \in G$ we have,

$$(g_1 g_2 \dots g_n)^{-1} = g_n^{-1} g_{n-1}^{-1} \dots g_1^{-1}.$$

Definition 1.0.5 (Abelian Group). If a group G also satisfies the condition $g \cdot h = h \cdot g$ for all $g, h \in G$ - **commutativity**, then G is said to be an **abelian group**.

Definition 1.0.6 (Powers of elements). Given a group G and some $g \in G$ the n th **power** of g in G is defined recursively as,

$$g^n := \begin{cases} e & \text{if } n = 0 \\ g^{n-1}g & \text{if } n > 0 \\ (g^n)^{-1} & \text{if } n < 0 \end{cases}.$$

Definition 1.0.7 (Order of group). The **order** of a group G , written $|G|$, is the cardinality of the underlying set of G .

Example 1.0.8 (Symmetric group). The **symmetric group of size n** , denoted S_n , is the set of bijections on the interval $[1, n]$, for $n \in \mathbb{N}$, under function composition.

2 Subgroups

2.1 Subgroups

Definition 2.1.1 (Subgroup). Given a group (G, \cdot) and a subset $H \subseteq G$ we say (H, \cdot) is a **subgroup** of G , written $H \leq G$, if (H, \cdot) forms a group and

$$\forall h_1, h_2 \in H : h_1 \cdot h_2 \in H.$$

A subgroup, H , is a **proper subgroup** if $H \neq G$. $\{e\}$ is the trivial subgroup.

Theorem 2.1.2 (Subgroup test). Given a group (G, \cdot) , (H, \cdot) is a subgroup iff:

- S1 H is non-empty - **existence**,
- S2 for all $h_1, h_2 \in H$ we have $h_1 \cdot h_2 \in H$ - **closure under group operation**,
- S3 for all $h \in H$ we have $h^{-1} \in H$ - **closure under inverses**.

2.2 Cyclic groups and orders

Definition 2.2.1 (Cyclic group). We say a group G is **cyclic** if there is an element $g \in G$ such that

$$G = \langle g \rangle := \{g^n : n \in \mathbb{N}\}.$$

We say that G is **generated** by g or g is a **generator** of G .

Definition 2.2.2 (Order of elements). Given a group G and some $g \in G$, the **order** of g in G , written $\text{ord } g$, is the smallest positive integer n such that $g^n = e$ or ∞ if no such n exists.

Theorem 2.2.3. Suppose G is a cyclic group generated by g with $|G| = n$, $\text{ord } g = |\{e, g, g^2, \dots, g^{n-1}\}| = |G| = n$.

Theorem 2.2.4. Suppose G is a cyclic group with $G = \langle g \rangle$, the three statements:

1. $H \leq G \implies H$ is cyclic,
2. suppose $|G| = n$ and $m \in \mathbb{Z}$ with $f = \gcd(m, n)$,

$$\langle g^m \rangle = \langle g^d \rangle \text{ and } |\langle g^m \rangle| = \frac{n}{d}.$$

In particular, $\langle g^m \rangle = G$ iff $\gcd(m, n) = 1$,

3. if $|G| = n$ and $k \leq n$, then G has a subgroup of order k iff $k|n$, this subgroup is $\langle g^{n/k} \rangle$.

Definition 2.2.5 (Euler totient). The **Euler totient** function ϕ is defined as $\phi(n) := |\{k \in \mathbb{N} : k \leq n \text{ and } \gcd(k, n) = 1\}|$.

Corollary 2.2.6. For $n \in \mathbb{N}$:

$$\sum_{d|n} \phi(d) = n.$$

2.3 Cosets

Definition 2.3.1 (Coset). Given a group G with $H \leq G$ and $g \in G$ then

$$gH := \{gh : h \in H\},$$

is a **left coset** of H in G (the definition of a **right coset** follows clearly).

Note 2.3.2. For the rest of this section, unless specified otherwise, a coset is assumed to be a left-coset.

Theorem 2.3.3. Given a group G with $H \leq G$, all cosets of H in G have the same size.

Theorem 2.3.4. If G is a finite group with $H \leq G$, the left cosets of H form a partition of G .

2.4 Lagrange's theorem

Theorem 2.4.1 (Lagrange's theorem). If G is a finite group and $H \leq G$, $|H|$ divides $|G|$.

Corollary 2.4.2. Given a group G with $H \leq G$, the relation \sim on G given by: $g \sim k$ iff $g^{-1}k \in H$, is an equivalence relation with equivalence classes given by cosets of H .

Corollary 2.4.3. Given a group G of order n , for all $g \in G$, $\text{ord } g|n$ and $g^n = e$.

Corollary 2.4.4 (Fermat's little theorem). Let p be prime. If $x \in \mathbb{Z}$ and $p \nmid x$, then $x^{p-1} \equiv 1 \pmod{p}$.

2.5 Generating groups

Definition 2.5.1. Given a group G with $S \subseteq G$, $S^{-1} := \{g^{-1} \in G : g \in S\}$.

Definition 2.5.2 (Subgroup generated by a set). Let G be a group with non-empty $S \subseteq G$. The **subgroup generated by S** is defined as

$$\langle S \rangle := \{g_1 g_2 \dots g_k \in G : k \in \mathbb{N} \text{ and } g_i \in S \cup S^{-1} \text{ for all } i \in [1, k]\}.$$

Lemma 2.5.3. Given a group G with non-empty $S \subseteq G$, $\langle S \rangle \leq G$ and, $H \leq G$, $S \subseteq H \implies \langle S \rangle \leq H$. This is equivalent to saying “ $\langle S \rangle$ is the smallest subgroup of G containing S ”.

3 Group homomorphisms

Definition 3.0.1 (Group homomorphism). If (G, \cdot) and $(H, *)$ are groups, $\phi : G \rightarrow H$ is a **group homomorphism** iff $\phi(g_1) * \phi(g_2) = \phi(g_1 \cdot g_2)$ for all $g_1, g_2 \in G$. If ϕ is bijective then it is called a **group isomorphism** with G and H being **isomorphic**, written $G \cong H$.

Example 3.0.2. The **determinant** is a group homomorphism, suppose \mathbb{F} is a field:

$$\det : \text{GL}(n, \mathbb{F}) \rightarrow (\mathbb{F}^*, \times).$$

Lemma 3.0.3. If G, H are groups with $\phi : G \rightarrow H$,

1. $\phi(e_G) = e_H$,
2. $\phi(g^{-1})(\phi(g))^{-1}$ for all $g \in G$.

Definition 3.0.4 (Image and kernel of group homomorphism). If G, H are groups with $\phi : G \rightarrow H$, the **image** of ϕ is:

$$\text{im } \phi := \{h \in H : \exists g \in G, h = \phi(g)\}.$$

and the **kernel** of ϕ is

$$\ker \phi := \{g \in G : \phi(g) = e_H\}.$$

These are each subgroups of H and G respectively.

Lemma 3.0.5. A group homomorphism, $\phi : G \rightarrow H$, is injective iff $\ker \phi = \{e_H\}$.

Theorem 3.0.6. The composition of two compatible group homomorphisms is also a group homomorphism.

Theorem 3.0.7. All cyclic groups of the same order are isomorphic.

4 Symmetric groups

4.1 Disjoint cycle decomposition

Definition 4.1.1. If $f, g \in S_n$ and $x \in [1, n]$ then f **fixes** x if $f(x) = x$ and f **moves** x otherwise.

Definition 4.1.2. The **support** of $f \in S_n$ is the set of points f moves, $\text{supp}(f) := \{x \in [1, n] : f(x) \neq x\}$.

Definition 4.1.3. If $f, g \in S_n$ satisfy $\text{supp}(f) \cap \text{supp}(g) = \emptyset$, f and g are **disjoint**.

Lemma 4.1.4. If $f, g \in S_n$ are disjoint, $fg = gf$.

Definition 4.1.5 (Cycles). If $f \in S_n$ with $i_1, i_2, \dots, i_r \in [1, n]$ for some $r \leq n$ such that,

$$f(i_s) = i_{s+1 \bmod r} \text{ for all } s \in [1, r],$$

with f fixing all other elements of $[1, n]$, then f is a **cycle of length r** or an **r -cycle** and we write $f = (i_1 i_2 \dots i_r)$.

Theorem 4.1.6 (Disjoint cycle form). if $f \in S_n$ then there exists $f_1, f_2, \dots, f_k \in S_n$ all with disjoint supports such that $f = f_1 f_2 \dots f_k$. If we further have, for all $i \in [1, k]$, both f_i is not a 1-cycle when $f \neq \text{id}$ and $\text{supp}(f_i) \subseteq \text{supp}(f)$. We say f is in **disjoint cycle form** or **d.c.f.**

Theorem 4.1.7 (Uniqueness of disjoint cycles). The disjoint cycle form of some $f \in S_n$ is unique up to rearrangement.

Theorem 4.1.8. If $f \in S_n$ is written in d.c.f as $f = f_1 f_2 \dots f_k$ where f_i is an r_i -cycle for $i \in [1, k]$ then,

1. $f^m = \text{id}$ iff $f_i^m = \text{id}$ for all $i \in [1, k]$,
2. $\text{ord}(f) = \text{lcm}(r_1, r_2, \dots, r_k)$.

4.2 Alternating groups

Theorem 4.2.1. Every permutation in S_n can be written as the product of 2-cycles.

Definition 4.2.2 (Sign of a permutation). We define the **sign** of a permutation with the group homomorphism, $\text{sgn} : S_n \rightarrow \{-1, 1\}$ with $\text{sgn}(i\ j) := -1$ for all $i, j \in [1, n]$ with $i \neq j$. This is defined over all permutations by the decomposition into 2-cycles, the sign of a permutation is unique. We say $f \in S_n$ is **even** if $f \in \ker(\text{sgn})$ and **odd** otherwise.

Definition 4.2.3 (Alternating group). The **alternating group** of size n is $A_n := \ker(\text{sgn})$ with $A_n \leq S_n$.

4.3 Dihedral groups

Definition 4.3.1 (Dihedral group). The **dihedral group** of order $2n$, denoted D_{2n} , is the group of symmetries of a regular n -gon in \mathbb{R}^3 centered at the origin, it is often written at

$$D_{2n} = \{e, r, r^2, \dots, r^{n-1}, s, sr, sr^2, \dots, sr^{n-1}\},$$

where r is a rotation by $\frac{2\pi}{n}$ and s is the reflection along the centre of the polygon and the first vertex.

Theorem 4.3.2. The elements of D_{2n} can be written as elements of S_n giving $D_{2n} \leq S_n$. Specifically, $r = (1\ 2\ \dots\ n)$ and $s = (1)(2\ n)(3\ n-1)\dots$ or $(1\ n)(2\ n-1)\dots$ if n is odd or even respectively.

Chapter 4

Calculus

Lectured by Professor Demetrios Papageorgiou
Typed by Yu Coughlin
Autumn 2023

Introduction

The following are suggested textbooks:

- G F Simmons, Calculus with Analytic Geometry, 1995
- J Stewart, Calculus, 2011
- S Lang, A First Course in Calculus, 1986
- S Lang, Undergraduate Analysis, 1997
- J Marsden and A Weinstein, Calculus I and Calculus II, 1985

Note. The actual majority of MATH40004A Calculus was a less formal and more example / application based derivation of the entirety of MATH40002 Analysis. As all of this content can be found in the corresponding document for Analysis, it isn't included in here.

Lecture 1
Thursday
10/01/19

1 Lengths, volumes and surfaces

1.1 Lengths

Theorem 1.1.1 (Arc length). The **arc length** of the curve $y = f(x)$ along $[a, b]$ is given by

$$\int_a^b \sqrt{1 + (f'(x))^2} \, dx$$

Theorem 1.1.2 (Distance and velocity of parameterised curves). If a curve is parameterised by $(x(t), y(t), z(t))$, the **distance travelled** from time t_0 to t is given by:

$$L(t) = \int_{t_0}^t \sqrt{\left(\frac{dx}{dt}\right)^2 + \left(\frac{dy}{dt}\right)^2 + \left(\frac{dz}{dt}\right)^2} \, dt$$

which naturally leads to the velocity at t :

$$v(t) = \sqrt{\left(\frac{dx}{dt}\right)^2 + \left(\frac{dy}{dt}\right)^2 + \left(\frac{dz}{dt}\right)^2}$$

1.2 Volumes and volumes of revolution

Theorem 1.2.1 (Volume). If the cross sectional area of a shape when cut by a plane at $x = x_0$ is given by $A(x_0)$ for all $x_0 \in [a, b]$, the volume of the shape is given by

$$V = \int_a^b A(x) \, dx$$

Theorem 1.2.2 (Disk method). The **volume of revolution** of $y = f(x)$ about the x -axis from $x = a$ to $x = b$ is given by,

$$V_x = \int_a^b \pi (f(x)^2) \, dx$$

Theorem 1.2.3 (Shell method). The **volume of revolution** of $y = f(x)$ about the y -axis from $y = a$ to $y = b$ is given by,

$$V_y = \int_a^b \pi (f^{-1}(x)^2) \, dx = \int_a^b 2\pi x f(x) \, dx$$

1.3 Surfaces

Theorem 1.3.1. The **surface area of revolution** of $y = f(x)$ about the x -axis from $x = a$ to $x = b$ is given by,

$$S_x = \int_a^b 2\pi f(x) \sqrt{1 + (f'(x))^2} \, dx$$

Theorem 1.3.2. The **surface area of revolution** of $y = f(x)$ about the y -axis from $y = a$ to $y = b$ is given by,

$$S_y = \int_a^b 2\pi x \sqrt{1 + (f'(x))^2} \, dx$$

1.4 Centres of mass

Theorem 1.4.1 (1D discrete case). If we have a system of n particles each with mass m_k and position x_k we can define the **centre of mass** at \bar{x} by

$$\bar{x} = \frac{\sum_{k=1}^n m_k x_k}{\sum_{k=1}^n m_k}$$

Theorem 1.4.2 (2D continuous case). If we have a region limited by $f(x)$ and $g(x)$, give $g(x) \leq f(x)$ for all $x \in [a, b]$, with uniform mass, the coordinates of the **centre of mass**, (\bar{x}, \bar{y}) is

$$\bar{x} = \frac{\int_a^b x(f(x) - g(x)) \, dx}{\int_a^b f(x) - g(x) \, dx} \quad \bar{y} = \frac{\int_a^b f(x)^2 - g(x)^2 \, dx}{2 \int_a^b f(x) - g(x) \, dx}$$

Theorem 1.4.3 (Pappus's theorem). If R is a region with area A lying on one side of the line l , $V = Ad$ is the volume obtained by rotation R about l , where d is the distance travelled by the **com** when R is rotated about l .

1.5 Moments of inertia

Theorem 1.5.1. Given a curve $y = f(x)$ in the interval $[a, b]$, this is representing a wire in a given shape, and have the density per unit length of the wire at a given x be $\rho(x)$, the **moment of inertia** of the curve about the x and y axis respectively is given by

$$I_x = \int_a^b \rho(x) f(x)^2 \sqrt{1 + f'(x)^2} \, dx \quad I_y = \int_a^b \rho(x) x^2 \sqrt{1 + f'(x)^2} \, dx$$

1.6 Polar coordinates

Definition 1.6.1 (Polar coordinates). A parameterisation of x, y is r, θ with $x = r \cos(\theta)$ and $y = r \sin(\theta)$.

Theorem 1.6.2 (Polar arc length). The arc length of a curve, $r = f(\theta)$ in polar coordinates between angles α, β is given by

$$L = \int_{\alpha}^{\beta} \sqrt{\left(\frac{dr}{d\theta}\right)^2 + r^2} \, d\theta$$

Theorem 1.6.3 (Polar area). The area of a polar curve, $r = f(\theta)$ between angles α, β is given by

$$A = \frac{1}{2} \int_{\alpha}^{\beta} f(\theta)^2 \, d\theta$$

2 Fourier series

2.1 Orthogonal and orthonormal function spaces

Definition 2.1.1 (Inner product of functions). If $f, g : [a, b] \rightarrow \mathbb{R}$ are integrable on $[a, b]$, the their **inner product** is defined as

$$\langle f, g \rangle := \int_a^b f(x)g(x) \, dx$$

Definition 2.1.2 (Orthogonal and orthonormal system). If $\mathcal{S} = \{\phi_0, \phi_1, \dots\}$ is a collection of integrable real functions on $[a, b]$, iff $\langle \phi_n, \phi_m \rangle = 0$ for all $n \neq m$ then \mathcal{S} is an **orthogonal system** on $[a, b]$. Furthermore, \mathcal{S} is a **orthonormal system** on $[a, b]$ iff $\|\phi_n\| := \langle \phi_n, \phi_n \rangle = 1$ for all n .

Theorem 2.1.3. The system

$$\mathcal{S} = \left\{ \frac{1}{\sqrt{2\pi}}, \frac{\cos(x)}{\sqrt{2\pi}}, \frac{\sin(x)}{\sqrt{2\pi}}, \frac{\cos(2x)}{\sqrt{2\pi}}, \frac{\sin(2x)}{\sqrt{2\pi}}, \dots \right\}$$

is orthonormal on all closed intervals of length 2π .

2.2 Periodic functions

Definition 2.2.1 (Periodic function). A function $f : \mathbb{R} \rightarrow \mathbb{R}$ is **periodic** with period T iff $f(x+T) = f(x)$ for all $x \in \mathbb{R}$.

Definition 2.2.2 (Discontinuity). When periodically extending a function, if $\lim_{x \rightarrow \xi+} f(x) \neq \lim_{x \rightarrow \xi-} f(x)$, we

$$\text{set } f(\xi) := \frac{1}{2} \left[\lim_{x \rightarrow \xi+} f(x) + \lim_{x \rightarrow \xi-} f(x) \right]$$

Theorem 2.2.3 (Integral over period). If $f(x)$ is a T periodic function, for all $a, b \in \mathbb{R}$ we have

$$\int_{a+T}^{b+T} f(x) dx = \int_a^b f(x) dx$$

2.3 Trigonometric polynomials

Definition 2.3.1 (Trigonometric polynomial). A **trigonometric polynomial** is a function in the form

$$S_n(x) = \frac{1}{2}a_0 + \sum_{k=1}^n (a_k \cos(kx) + b_k \sin(kx))$$

Theorem 2.3.2. Using euler's identity we can rewrite a trigonometric polynomial

$$S_n(x) = \frac{1}{2}a_0 + \sum_{k=1}^n (a_k \cos(kx) + b_k \sin(kx)) \text{ as } \sum_{k=-n}^n (\gamma_k e^{ikx}) \text{ where } \gamma_k = \begin{cases} \frac{1}{2}a_0 & \text{if } k = 0 \\ \frac{1}{2}(a_k - ib_k) & \text{if } k \in [1, n] \\ \gamma_k^* & \text{otherwise} \end{cases}$$

2.4 Fourier series

Definition 2.4.1 (Fourier series). If $f(x)$ is $2L$ periodic then its **Fourier series** is given by

$$f(x) := \frac{1}{2}a_0 + \sum_{n=1}^{\infty} \left[a_n \cos\left(\frac{n\pi x}{L}\right) + b_n \sin\left(\frac{n\pi x}{L}\right) \right] \text{ where}$$

$$a_n := \frac{1}{\pi} \int_{-L}^L f(x) \cos\left(\frac{n\pi x}{L}\right) dx, \quad b_n := \frac{1}{\pi} \int_{-L}^L f(x) \sin\left(\frac{n\pi x}{L}\right) dx$$

Lemma 2.4.2 (Riemann-Lebesgue). If the function $f(x)$ is integrable on $[a, b]$ then

$$I_\lambda := \int_a^b g(x) \sin(\lambda x) dx \rightarrow 0 \text{ as } \lambda \rightarrow \infty$$

Theorem 2.4.3 (Parseval's). If $f(x)$ is periodic on 2π and is represented by its Fourier series,

$$\frac{1}{\pi} \int_{-\pi}^{\pi} f^2(x) dx = \frac{1}{2}a_0^2 + \sum_{n=0}^{\infty} (a_n^2 + b_n^2)$$

3 Laplace transform

3.1 Definition

Definition 3.1.1 (Laplace transform). The **Laplace transform** is a linear operator that when applied to a function $f(x)$ gives

$$F(p) := \mathcal{L}[f(x)] := \int_0^{\infty} e^{-px} f(x) \, dx$$

Theorem 3.1.2 (Common Laplace transformations). These are some common functions with their Laplace transforms and the conditions for which they converge:

$f(x) = 1$	$F(p) = \frac{1}{p}$	Converges for $p > 0$
$f(x) = x$	$F(p) = \frac{1}{p^2}$	Converges for $p > 0$
$f(x) = x^n$	$F(p) = \frac{n!}{p^{n+1}}$	Converges for $p > 0$
$f(x) = e^{ax}$	$F(p) = \frac{1}{p-a}$	Converges for $p > a$
$f(x) = \sin(ax)$	$F(p) = \frac{a}{p^2 + a^2}$	Converges for $p > 0$
$f(x) = \cos(ax)$	$F(p) = \frac{p}{p^2 + a^2}$	Converges for $p > 0$
$f(x) = \sinh(ax)$	$F(p) = \frac{a}{p^2 - a^2}$	Converges for $p > a$
$f(x) = \cosh(ax)$	$F(p) = \frac{p}{p^2 - a^2}$	Converges for $p > a$

Theorem 3.1.3 (Existence of Laplace transform). The Laplace transform for a function $f(x)$ exists iff there exists constants $M, c \in \mathbb{R}$ with $|f(x)| \leq Me^{cx}$ ($f(x)$ is of **exponential order**).

3.2 Differentiating

Theorem 3.2.1 (Derivatives of Laplace transforms). By performing DUTIS $n \in \mathbb{N}$ times we have

$$F^{(n)}(p) = \mathcal{L}[(-1)^n x^n f(x)]$$

3.3 Convolution theorem

Theorem 3.3.1 (Convolution theorem for Laplace transforms). For integrable functions $f, g : \mathbb{R} \rightarrow \mathbb{R}$:

$$\mathcal{L} \left[\int_0^x f(x-t)g(t) \, dt \right] = F(p)G(p)$$

Chapter 5

Differential Equations

Lectured by Professor Demetrios Papageorgiou

Typed by Yu Coughlin

Spring 2024

Introduction

The following are suggested textbooks:

- G F Simmons, Calculus with Analytic Geometry, 1995
- J Stewart, Calculus, 2011
- S Lang, A First Course in Calculus, 1986
- S Lang, Undergraduate Analysis, 1997
- J Marsden and A Weinstein, Calculus I and Calculus II, 1985

Lecture 1
Thursday
10/01/19

- 1 Fourier transform
- 2 Ordinary differential equations
- 3 Qualitative analysis
- 4 Bifurcations
- 5 Multivariate calculus
- 6 Partial differential equations

Chapter 6

Probability

Lectured by Professor Almut Veraart
Typed by Yu Coughlin
Autumn 2023

Introduction

The following are complementary reading for the course.

- G. Grimmett and D. J. A. Welsh, Probability: An Introduction, 1986
- J. K. Blitzstein and J. Hwang, Introduction to Probability, 2019
- D. F. Anderson et al, Introduction to Probability, 2018
- S. M. Ross, Introduction to Probability Models, 2014
- G. Grimmett and D. Stirzaker, Probability and Random Processes, 2001
- G. Grimmett and D. Stirzaker, One Thousand Exercises in Probability, 2009

1 Introduction

1.1 Sample spaces and set theory

Definition 1.1.1. The **sample space** Ω is the set of all possible outcomes of an experiment. An element of the sample space $\omega \in \Omega$ is a **sample point**.

Examples 1.1.2. When flipping a coin $\Omega = \{H, T\}$. When rolling a standard die $\Omega = \{1, 2, 3, 4, 5, 6\}$.

Definition 1.1.3. Subsets of Ω are collections of sample points and called **events**.

Suppose events $A, B \subseteq \Omega$:

- $A \cup B$ is the event that A or B or both occur,
- $A \cap B$ is the event that A and B both occur,
- $A^c = \bar{A}$ is the event that occurs iff A does not occur.

Let \mathcal{I} be a general index set with $A_i \subseteq \Omega, \forall i \in \mathcal{I}$ and $B \subseteq \Omega$. The following identities hold.

$$\left(\bigcup_{i \in \mathcal{I}} A_i\right)^c = \bigcap_{i \in \mathcal{I}} A_i^c, \quad \left(\bigcap_{i \in \mathcal{I}} A_i\right)^c = \bigcup_{i \in \mathcal{I}} A_i^c, \quad B \cap \left(\bigcup_{i \in \mathcal{I}} A_i\right) = \bigcap_{i \in \mathcal{I}} (A_i \cup B), \quad B \cup \left(\bigcap_{i \in \mathcal{I}} A_i\right) = \bigcap_{i \in \mathcal{I}} (A_i \cap B).$$

These are **De Morgan's Laws** and **Distributivity** respectively.

1.2 Interpretation of probability

Definition 1.2.1. The **Cardinality** of a set, denoted $\text{card}(A)$ or $|A|$ is the number of elements in the set A .

Definition 1.2.2. Two sets have the same cardinality iff there exists a bijection between the them.

Definition 1.2.3. A is **finite** if it has as finite numbers of elements, A is **countably infinite** if there exists a bijection $f: A \rightarrow \mathbb{N}$, A is **countable** if it is finite or countable infinite, A is **uncountable** or **uncountable infinite** if it isn't countable.

Samples spaces can be countable or uncountable.

Definition 1.2.4 (Naive probability). Suppose $|A| < \infty$ and we want to assign a probability to $A \subseteq \Omega$.

$$P_{\text{Naive}}(A) := \frac{|A|}{|\Omega|} \implies P(A^c) = 1 - P(A).$$

This Naive example does not consider when $|A|$ is infinite but of finite area.

Example 1.2.5. Let $\Omega = \{(x, y) \in \mathbb{R}^2, x^2 + y^2 = 1\}$ and $A \subseteq \Omega$. Define:

$$P(A) := \frac{\text{area of } A}{\text{area of } \Omega}$$

In the case where $A = \{(x, y) \in \mathbb{R}^2, x^2 + y^2 = 0.5^2\}$ we have $P(A) = 0.25$

Remark 1.2.6. For classical / naive probability we require $|\Omega| < \infty$ or the "area" of Ω be finite.

Definition 1.2.7 (Limiting frequency). Consider n_{total} repetitions of an experiment and n_A the number of time A occurs.

$$P(A) := \lim_{n_{\text{total}} \rightarrow \infty} \frac{n_A}{n_{\text{total}}}$$

Unfortunately, $n_{\text{total}} \rightarrow \infty$ is often hard to conceive with finite representations not necessarily being representative.

Definition 1.2.8 (Subjective probability). For an event A assign the probability $P(A)$ based on our own personal beliefs. The subjective probability need not be the same for different individuals, and despite its appearance it remains a valid interpretation of probability.

Remark 1.2.9. All three interpretations of probability depend of assumptions about the experiment.

2 Counting

2.1 Multiplication principle

Computing naive probabilities often requires some combinatorics.

Definition 2.1.1 (Multiplication principle). If we perform an experiment A that has a possible outcomes and an experiment B with b possible outcomes (in any order) then the number of outcomes of the **compound experiment** will be ab .

Remark 2.1.2. When dealing with repetitions of the same experiment (with sample space Ω , the sample space is given by the Cartesian product of the individual samples spaces.

$$\Omega_1 \times \Omega_2 \times \cdots \times \Omega_n := \{(\omega_1, \omega_2, \dots, \omega_n) : \omega_i \in \Omega_i\}.$$

The cardinality of this samples space follows from the multiplication principle.

2.2 Power sets

Definition 2.2.1 (Power Set). Given a set A its **power set** is defined as:

$$\mathcal{P}(A) := \{X : X \subseteq A\}.$$

Theorem 2.2.2. If A is a finite set, $|\mathcal{P}(A)| = 2^{|A|}$.

2.3 Combinatorial coefficients

Definition 2.3.1 (Factorial). Let $n \in \mathbb{N}$ the **factorial** of n is defined as:

$$n! := \prod_{i=1}^n i.$$

Definition 2.3.2 (Descending factorial). Let $k, n \in \mathbb{N}$ with $k \leq n$ the **descending factorial** denoted $(n)_k$ is defined as:

$$(n)_k := n(n-1) \cdots (n-k+1) = \prod_{i=0}^{k-1} (n-i) = \prod_{j=n-k+1}^n j = \frac{n!}{(n-k)!}.$$

Definition 2.3.3 (Binomial coefficient). Let $k, n \in \mathbb{N}_0$ the **binomial coefficient** is the number of subsets of size k of a set n :

$$\binom{n}{k} := \begin{cases} \frac{n(n-1) \cdots (n-(k-1))}{k!} = \frac{(n)_k}{k!} = \frac{n!}{(n-k)!k!} & \text{if } k \leq n \\ 0 & \text{otherwise.} \end{cases}$$

2.4 Sampling with and without replacement

“Definitions” given in the context of drawing balls from an urn, $S = \{1, 2, \dots, n\}$.

Definition 2.4.1 (Ordered sampling with replacement). Take out a ball from S , note its number, put it back; repeat this k times. The sample space for this experiment is $\Omega = S^k$.

Definition 2.4.2 (Ordered sampling without replacement). Take out a ball from S , note its number but **do not** put it back; repeat $k < n$ times. There are $|\Omega| = (n)_k$ possible outcomes.

Definition 2.4.3 (Unordered sampling without replacement). We take k balls out of the urn, there are $\binom{n}{k}$ possibilities.

Definition 2.4.4 (Unordered sampling with replacement). We take k balls out of the urn, with the stars and bars argument: there must be k stars divided by $n-1$ bars giving us:

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

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06/11/2023

Lecture 5
Tuesday
07/11/2023

3 Axiomatic probability

3.1 Event space

We do not always want to consider all subsets of Ω so denote $\mathcal{F} \subseteq \mathcal{P}(\Omega)$ the **event space**, which contains the events we are allowed to consider. \mathcal{F} must always be a σ -algebra.

Definition 3.1.1 (Algebra). \mathcal{F} is an **algebra** (or a field) on Ω iff:

1. $\emptyset \in \mathcal{F}$,
2. $A \in \mathcal{F} \implies A^c \in \mathcal{F}$,
3. $A, B \in \mathcal{F} \implies A \cup B \in \mathcal{F}$.

Definition 3.1.2 (σ -algebra). \mathcal{F} is a **σ -algebra** (or a σ -field) on Ω iff:

1. $\emptyset \in \mathcal{F}$;
2. $A \in \mathcal{F} \implies A^c \in \mathcal{F}$,
3. For all i in some countable indexing set \mathcal{I} , $A_i \in \mathcal{F} \implies \bigcup_{i \in \mathcal{I}} A_i \in \mathcal{F}$.

Remark 3.1.3. 1. Any algebra is closed under finite unions and finite intersections,

2. Any σ -algebra is closed under countable intersections,

3. Any (σ -)algebra on Ω contains Ω .

Definition 3.1.4 (Trivial sigma algebra). The **trivial sigma algebra** on Ω is defined as $\mathcal{F}_{trivial} := \{\emptyset, \Omega\}$.

Example 3.1.5 (Smallest σ -algebra of an element). For some $A \subseteq \Omega$, the sigma algebra $\mathcal{F}_A := \{\emptyset, A, A^c, \Omega\}$ is the smallest σ -algebra on Ω (smallest cardinality) that contains A .

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3.2 Probability measure

Definition 3.2.1 (Probability measure). A mapping $P : \mathcal{F} \rightarrow \mathbb{R}$ is a **probability measure** on (Ω, \mathcal{F}) iff:

1. $P(A) \geq 0$ for all $A \in \mathcal{F}$;
2. $P(\Omega) = 1$;
3. for a countable, disjoint sequence of events $(A_i)_{i \in \mathcal{I}}$ on an indexing set \mathcal{I} :

$$P\left(\bigcup_{i \in \mathcal{I}} A_i\right) = \sum_{i \in \mathcal{I}} P(A_i).$$

3.3 Probability space

Definition 3.3.1 (Probability space). A **probability space** is a triple (Ω, \mathcal{F}, P) , with Ω a sample space, \mathcal{F} a σ -algebra on Ω , and P a probability measure on (Ω, \mathcal{F}) .

Corollary 3.3.2. For $A, B \in \mathcal{F}$:

1. $P(A^c) = 1 - P(A)$,
2. $A \subseteq B \implies P(A) \leq P(B)$,
3. $P(A \cup B) = P(A) + P(B) - P(A \cap B)$.

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4 Conditional probability

Definition 4.0.1 (Conditional probability measure). Consider the probability space (Ω, \mathcal{F}, P) and some event $B \in \mathcal{F}$ with $P(B) > 0$, we construct the probability measure Q on (Ω, \mathcal{F}) by

$$Q(A) := \frac{P(A \cap B)}{P(B)}.$$

Denote the **conditional probability** of A given B by $P(A|B) = Q(A)$.

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4.1 Bayes' rule and total probability

Theorem 4.1.1 (Bayes' rule). For $A, B \in \mathcal{F}$ with $P(A) > 0, P(B) > 0$ we have,

$$P(A|B) = \frac{P(B|A)P(A)}{P(B)}.$$

Definition 4.1.2 (Partition of a set). A partition of some set Ω is a collection $\{B_i, i \in \mathcal{I}\}$ for some countable index set \mathcal{I} with $B_i \cap B_j = \emptyset$ for all $i, j \in \mathcal{I}$ with $i \neq j$ and $\bigcup_{i \in \mathcal{I}} B_i = \Omega$.

Theorem 4.1.3 (Total probability). Given some partition $\{B_i, i \in \mathcal{I}\}$ of Ω with $P(B_i) > 0$ for all $i \in \mathcal{I}$ and some event $A \in \mathcal{F}$,

$$P(A) = \sum_{i \in \mathcal{I}} P(A \cap B_i) = \sum_{i \in \mathcal{I}} P(A|B_i)P(B_i).$$

These two theorems can then be combined to form the following.

Theorem 4.1.4 (Bayes' rule with extra conditioning). For events $A, B, E \in \mathcal{F}$ with $P(A \cap E) > 0, P(B \cap E) > 0$ we have

$$P(A|B \cap E) = \frac{P(B|A \cap E)P(A|E)}{P(B|E)}.$$

Theorem 4.1.5 (Total probability with extra conditioning). Given events $A, E \in \mathcal{I}$ with $P(E) > 0$ and some partition $\{B_i, i \in \mathcal{I}\}$ of Ω with $P(B_i \cap E) > 0$ for all $i \in \mathcal{I}$,

$$P(A|E) = \sum_{i \in \mathcal{I}} \frac{P(A \cap B_i \cap E)}{P(E)} = \sum_{i \in \mathcal{I}} P(A|B_i \cap E)P(B_i|E).$$

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5 Independence

5.1 Event independence

Two events $A, B \in \mathcal{F}$ will be independent iff the occurrence of one does not effect the probability the other occurs, i.e $P(A|B) = P(A)$ and vice versa.

Definition 5.1.1 (Independent events). Two events $A, B \in \mathcal{F}$ are said to be **independent** iff

$$P(A \cap B) = P(A)P(B),$$

and **dependent** otherwise.

Corollary 5.1.2. If A and B are independent then so are all pairs of their complements.

Definition 5.1.3 (General independence). A finite collection of events $\{A_1, A_2, \dots, A_n\}$ is independent iff

$$P(A_1 \cap A_2 \cap \dots \cap A_n) = P(A_1)P(A_2) \dots P(A_n),$$

and similarly a countably or uncountably infinite collection of events is independent iff each finite subcollection is independent.

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5.2 Conditional independence

Definition 5.2.1 (Conditional independence). Given the events $A, B, C \in \mathcal{F}$ with $P(C) > 0$ we say A and B are **conditional independent** given C iff,

$$P(A \cap B|C) = P(A|C)P(B|C).$$

5.3 Product rule for general independence

The upcoming subsection may seem disparate, they are however necessary parts to the omitted proof of the product rule for general independence and therefore deemed relevant.

Definition 5.3.1 (Set difference). Given two set $A, B \in \Omega$ the **set difference** of A and B is defined as, $A \setminus B := A \cap B^c$.

Lemma 5.3.2. Any countable union of sets can be written as a countable union of disjoint sets.

Definition 5.3.3 (Increasing and decreasing sets). A sequence of sets $(A_i)_{i=1}^{\infty}$ is said to **increase** to A (written $A_i \uparrow A$) iff $A_1 \subseteq A_2 \subseteq \dots$ and $\bigcup_{i=1}^{\infty} A_i = A$. The definition for a sequence of sets $(B_i)_{i=1}^{\infty}$ to **decrease** to a set B ($B_1 \downarrow B$) is defined similarly.

Theorem 5.3.4 (Continuity property of probability measures). If $A_1, A_2, \dots \in \mathcal{F}$ with $A_i \uparrow A$ or $A_i \downarrow A$ for some $A \in \mathcal{F}$,

$$\lim_{i \rightarrow \infty} P(A_i) = P(\lim_{i \rightarrow \infty} A_i) = P(A).$$

Theorem 5.3.5 (Product rule for general independence). Given a countably infinite set of independent events $A_1, A_2, \dots \in \mathcal{F}$,

$$P\left(\bigcap_{i=1}^{\infty} A_i\right) = \prod_{i=1}^{\infty} P(A_i).$$

6 Discrete random variables

6.1 Images and their properties

throughout this subsection we will be considering the function $f : \mathcal{X} \rightarrow \mathcal{Y}$.

Definition 6.1.1 (Image). For some subset $A \subseteq \mathcal{X}$ we define the **image** of A under f by,

$$f(A) := \{y \in \mathcal{Y} : \exists x \in A, y = f(x)\} = \{f(x) : x \in A\}.$$

When $A = \mathcal{X}$, $f(\mathcal{X}) = \text{im } f$.

Definition 6.1.2 (Pre-image). For some subset $B \subseteq \mathcal{Y}$ we now define the **pre-image** of B under f by,

$$f^{-1}(B) := \{x \in \mathcal{X} : f(x) \in B\}.$$

Despite the similar notation to the inverse function of f they are not the same thing. Notably, the pre-image under f always exists while the inverse function need not exist.

Lemma 6.1.3. For a collection of subsets $B_i \in \mathcal{F}$ for i in some indexing set \mathcal{I} we have,

$$f^{-1}\left(\bigcup_{i \in \mathcal{I}} B_i\right) = \bigcup_{i \in \mathcal{I}} f^{-1}(B_i).$$

6.2 DRVs and their distributions

Definition 6.2.1 (Discrete random variable). A **discrete random variable** (DRV) on the probability space (Ω, \mathcal{F}, P) is a function $X : \Omega \rightarrow \mathbb{R}$ that satisfies the following properties:

- $\text{im } X = \{X(\omega) : \omega \in \Omega\}$ must be a countable subset of \mathbb{R} ,
- $X^{-1}(x) \in \mathcal{F}$ for all $x \in \mathbb{R}$.

Remark 6.2.2. The nomenclature of X being discrete stems from the fact that its image is a countable subset of \mathbb{R} and so can be mapped to \mathbb{N} which we see as being discrete.

Definition 6.2.3 (Probability mass function). The **probability mass function** (pmf) of a DRV X is defined as a function $p_X : \mathbb{R} \rightarrow [0, 1]$ such that,

$$p_X(x) := P(X^{-1}(x)).$$

This is commonly denoted by $p_X(x) = P(X = x)$.

Remark 6.2.4. Some useful properties of the pmf extending from the definition are:

- $x \notin \text{im } X \implies p_X(x) = 0$,
- For $x_1, x_2 \in \text{im } X$ with $x_1 \neq x_2$, $X^{-1}(x_1) \cap X^{-1}(x_2) = \emptyset$,
- $\sum_{x \in \text{im } X} p_X(x) = \sum_{x \in \mathbb{R}} p_X(x) = 1$.

Theorem 6.2.5. Suppose \mathcal{I} is some indexing set and $S = \{s_i \in \mathbb{R} : i \in \mathcal{I}\}$ is countable and $\{\pi_i : i \in \mathcal{I}\}$ is a collection such that $\pi_i \geq 0$ for all $i \in \mathcal{I}$ and $\sum_{i \in \mathcal{I}} \pi_i = 1$. Then there exists some probability space (Ω, \mathcal{F}, P) and a DRV X on said probability space such that $p_X(s_i) = \pi_i$ for all $i \in \mathcal{I}$ and $p_X(s) = 0$ otherwise.

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7 Common DRV's

All DRV's within this section will be considered over the probability space (Ω, \mathcal{F}, P) .

7.1 Bernoulli distribution

Definition 7.1.1 (Bernoulli distribution). A DRV X is said to have **Bernoulli distribution** with parameter $p \in (0, 1)$ if $\text{im } X = \{0, 1\}$ with $p_X(1) = p$. This is denoted by $X \sim \text{Bern}(p)$.

Definition 7.1.2 (Indicator variable). Given some event $A \in \mathcal{F}$ the **indicator variable** of the event A is given by,

$$\mathbb{I}_A(\omega) := \begin{cases} 1 & \text{if } \omega \in A \\ 0 & \text{if } \omega \notin A \end{cases}.$$

Remark 7.1.3. $\mathbb{I}_A \sim \text{Bern}(P(A))$.

7.2 Binomial distribution

Definition 7.2.1 (Binomial distribution). Consider a sequence of $n \in \mathbb{N}$ iid Bernoulli trials with parameter p , count the number of successes and denote this by the random variable X then $\text{im } X = [0, n]$ and,

$$p_X(x) = \binom{n}{x} p^x (1-p)^{n-x} \quad \text{for } x \in [0, n].$$

We say X follows a **binomial distribution** and this is denoted by $X \sim \text{Bin}(n, p)$.

7.3 Hypergeometric distribution

As we have done previously, consider of urn of $N \in \mathbb{N}$ balls with $K \in \mathbb{N}$ of these being white and the remainder being black from which we will draw $n \in \mathbb{N}$ balls and want to consider the DRV (X) for the number of white balls drawn. When drawing with replacement we have $X \sim \text{Bin}(n, K/N)$. However, when drawing without replacement X follows the hypergeometric distribution.

Definition 7.3.1 (Hypergeometric distribution). A DRV X follows the **hypergeometric distribution** with three parameters $N \in \mathbb{N}_0, K \in \mathbb{N}, n \in [0, N]$ if $\text{im } X = [0, \min(n, K)]$ and,

$$p_X(x) = \frac{\binom{K}{x} \binom{N-K}{n-x}}{\binom{N}{n}} \quad \text{for } x \in [0, K].$$

Lemma 7.3.2 (Vandemonde's identity). **Vandemonde's identity** is an important tool in the derivation of the pmf for the hypergeometric distribution and so is included here. The identity is as follows, for $k, m, n \in \mathbb{N}$ with $k \leq m + n$, we have:

$$\binom{m+n}{k} = \sum_{i=0}^k \binom{m}{i} \binom{n}{k-i}.$$

7.4 Discrete uniform distribution

Definition 7.4.1 (Discrete uniform distribution). A DRV X follows the **discrete uniform distribution** over a nonempty set of numbers C , denoted $X \sim \text{DUnif}(C)$, if $\text{im } X = C$ and,

$$p_X(x) = \begin{cases} \frac{1}{\text{card}(C)} & \text{for } x \in C \\ 0 & \text{otherwise} \end{cases}.$$

7.5 Poisson distribution

The poisson distribution is commonly used for modelling the number of events occurring in a certain time period. Its pdf is derived by taking the $\lim_{n \rightarrow \infty} p_X(x)$ where $X \sim \text{Bin}(n, \frac{\lambda}{n})$ for some $\lambda \in \mathbb{R}$.

Definition 7.5.1 (Poisson distribution). A DRV X follows the **poisson distribution** with parameter $\lambda \in \mathbb{R}^{>0}$, denoted $X \sim \text{Poi}(\lambda)$, if $\text{im } X = \mathbb{N}_0$ and,

$$p_X(x) = \frac{\lambda^x}{x!} e^{-\lambda} \quad \text{for } x \in \mathbb{N}_0.$$

7.6 Geometric distribution

Definition 7.6.1 (Geometric distribution). A DRV X follows the **geometric distribution** with parameter $p \in (0, 1)$, denoted $X \sim \text{Geom}(p)$, if $\text{im } X = \mathbb{N}$ and,

$$p_X(x) = (1 - p)^x p \quad \text{for } x \in \mathbb{N}.$$

This can be seen as counting the number of Bernoulli trials with parameter p that occur before a failure.

7.7 Negative binomial distribution

Definition 7.7.1 (Generalised binomial coefficient). Let $\alpha \in \mathbb{C}$ and $k \in \mathbb{N}$ and define the **generalised binomial coefficient** by,

$$\binom{\alpha}{k} := \frac{\alpha(\alpha - 1) \dots (\alpha - k + 1)}{k!}.$$

Lemma 7.7.2. For $x \in \mathbb{N}_0$ and $r \in \mathbb{N}$ the following identity holds,

$$\binom{x + r - 1}{r - 1} = (-1)^x \binom{-r}{x}.$$

The generalised binomial coefficient as well as this lemma are necessary to have a well defined and valid pdf for the negative binomial distribution.

Definition 7.7.3 (Negative binomial distribution). A DRV X follows the **negative binomial distribution** with parameters $r \in \mathbb{N}$ and $p \in (0, 1)$, denoted $X \sim \text{NBin}(r, p)$, if $\text{im } X = \mathbb{N}_0$ and,

$$p_X(x) = \binom{x + r - 1}{r - 1} p^r (1 - p)^x \quad \text{for } x \in \mathbb{N}_0.$$

This is the distribution of the number of failed ii Bernoulli trials with parameter p before r successes have occurred.

8 Continuous random variables

8.1 General random variables and their distributions

Definition 8.1.1 (Random variable). A **random variable (RV)** on the probability space (Ω, \mathcal{F}, P) is a mapping $X : \Omega \rightarrow \mathbb{R}$ such that $X^{-1}((-\infty, x]) = \{\omega \in \Omega : X(\omega) \leq x\} \in \mathcal{F}$ for all $x \in \mathbb{R}$. By taking the countable union of pre-images of all $\omega \leq x$ in \mathcal{F} , it can be seen that a DRV satisfies this condition.

Definition 8.1.2 (Cumulative distribution function). For some RV X on the probability space (Ω, \mathcal{F}, P) , the **cumulative distribution function (CDF)** of X is defined as the mapping $F_X : \mathbb{R} \rightarrow [0, 1]$ given by,

$$F_X(x) = P(X^{-1}((-\infty, x])),$$

often denoted $F_X(x) = P(X \leq x)$.

Theorem 8.1.3 (cdf properties). For some RV X on the probability space (Ω, \mathcal{F}, P) the following properties hold:

1. F_X is monotonically non-decreasing,
2. F_X is right-continuous ($(x_n) \downarrow x \implies F_X(x_n) \rightarrow F_X(x)$ as $n \rightarrow \infty$),
3. $\lim_{x \rightarrow -\infty} F_X(x) = 0$ and $\lim_{x \rightarrow \infty} F_X(x) = 1$.

Theorem 8.1.4. For $a, b \in \mathbb{R}$ if $a < b$, then $P(a < X \leq b) = F_X(b) - F_X(a)$.

Remark 8.1.5. In general the cdf of an RV is not left continuous.

8.2 CRVs and pdfs

Definition 8.2.1 (Continuous random variable). A random variable X on the probability space (Ω, \mathcal{F}, P) is called a **continuous random variable (CRV)** iff its cdf can be written as:

$$F_X(x) = \int_{-\infty}^x f_X(u) du \quad \text{for all } x \in \mathbb{R},$$

where $f_X : \mathbb{R} \rightarrow \mathbb{R}$ satisfies: $f_X(u) \geq 0$ for all $u \in \mathbb{R}$ and $\int_{-\infty}^{\infty} f_X(u) du = 1$. We call f_X the **probability density function (pdf)** of X .

Theorem 8.2.2. If X is a CRV on the probability space (Ω, \mathcal{F}, P) with pdf f_X , $P(X = x) = 0$ for all $x \in \mathbb{R}$.

Theorem 8.2.3. With the same conditions, $P(a \leq X \leq b) = \int_a^b f_X(u) du$ for all $a, b \in \mathbb{R}$ with $a \leq b$.

Remark 8.2.4. Combining the results from this section leads to the conclusion that the cdf of a CRV is continuous.

9 Common CRVs

All CRVs X within this section will be considered over the probability space (Ω, \mathcal{F}, P) with the natural notation for their pdf and cdf. These distributions will be uniquely identified by their pdfs.

9.1 Uniform distribution

Definition 9.1.1 (Uniform distribution). A CRV X follows the **uniform distribution** on the interval (a, b) for $a, b \in \mathbb{R}$ with $a < b$, denoted $X \sim U(a, b)$ if it satisfies:

$$f_X(x) = \begin{cases} \frac{1}{b-a} & \text{if } a < x < b \\ 0 & \text{otherwise} \end{cases}, \quad F_X(x) = \begin{cases} 0 & \text{if } x \leq a \\ \frac{x-a}{b-a} & \text{if } a < x < b \\ 1 & \text{if } x \geq b \end{cases}.$$

9.2 Exponential distribution

Definition 9.2.1 (Exponential distribution). A CRV X follows the **exponential distribution** with parameter $\lambda \in \mathbb{R}^{>0}$, denoted $X \sim \text{Exp}(\lambda)$ if it satisfies:

$$f_X(x) = \begin{cases} \lambda e^{-\lambda x} & \text{if } x > 0 \\ 0 & \text{otherwise} \end{cases}, \quad F_X(x) = \begin{cases} 0 & \text{if } x \leq 0 \\ 1 - e^{-\lambda x} & \text{if } x > 0 \end{cases}.$$

9.3 Gamma distribution

Definition 9.3.1 (Gamma function). For $t \in \mathbb{R}$ with $t > 0$ we define the **gamma function** by,

$$\Gamma(t) := \int_0^{\infty} x^{t-1} e^{-x} dx.$$

One of the gamma function's many interesting properties is that $\Gamma(t) = (t-1)\Gamma(t-1)$ for $t > 1$.

Definition 9.3.2 (Gamma distribution). A CRV X follows the **gamma distribution** with shape and rate parameter $\alpha, \beta \in \mathbb{R}^{>0}$ respectively, denoted $X \sim \text{Gamma}(\alpha, \beta)$ if it satisfies:

$$f_X(x) = \begin{cases} \frac{\beta^\alpha}{\Gamma(\alpha)} x^{\alpha-1} e^{-\beta x} & \text{if } x > 0 \\ 0 & \text{otherwise} \end{cases}.$$

Its cdf cannot be written in a closed form so must be left as an integral of the pdf or approximated.

9.4 Chi-squared distribution

Definition 9.4.1 (Chi-squared distribution). A CRV X follows the **chi-squared distribution** with $n \in \mathbb{N}$ degrees of freedom, denoted $X \sim \chi^2(n)$ if it satisfies:

$$f_X(x) = \begin{cases} \frac{1}{2\Gamma(n/2)} \left(\frac{x}{2}\right)^{n/2-1} e^{-x/2} & \text{if } x > 0 \\ 0 & \text{otherwise} \end{cases}.$$

Its cdf can also not be written in a closed form. The $\chi^2(n)$ distribution is the same as the **Gamma** $(\frac{n}{2}, \frac{1}{2})$ distribution.

9.5 F-distribution

These pdfs are getting tough.

Definition 9.5.1 (F-distribution). A CRV X follows the **f-distribution** with $d_1, d_2 \in \mathbb{R}^{>0}$ degrees of freedom, denoted $X \sim F(d_1, d_2)$ if it satisfies:

$$f_X(x) = \begin{cases} \frac{\Gamma\left(\frac{d_1+d_2}{2}\right) \left(\frac{d_1}{d_2}\right)^{d_1/2} x^{d_1/2-1}}{\Gamma\left(\frac{d_1}{2}\right) \Gamma\left(\frac{d_2}{2}\right) \left(1 + \frac{d_1}{d_2}x\right)^{(d_1+d_2)/2}} & \text{if } x > 0 \\ 0 & \text{otherwise} \end{cases}.$$

Its cdf can also not be written in a closed form. It is important to note that d_1, d_2 are not restricted to integer values, and that $X = \frac{X_1/d_1}{X_2/d_2}$ where $X_1 \sim \chi^2(d_1)$ and $X_2 \sim \chi^2(d_2)$.

9.6 Beta distribution

Definition 9.6.1 (Beta function). For $\alpha, \beta \in \mathbb{R}^{>0}$ we define the **beta function** by,

$$B(\alpha, \beta) := \int_0^1 x^{\alpha-1} (1-x)^{\beta-1} dx = \frac{\Gamma(\alpha)\Gamma(\beta)}{\Gamma(\alpha+\beta)}.$$

Definition 9.6.2 (Beta distribution). A CRV X follows the **beta distribution** with parameters $\alpha, \beta \in \mathbb{R}^{>0}$, denoted $X \sim \text{Beta}(\alpha, \beta)$ if it satisfies:

$$f_X(x) = \begin{cases} \frac{1}{B(\alpha, \beta)} x^{\alpha-1} (1-x)^{\beta-1} & \text{if } 0 \leq x \leq 1 \\ 0 & \text{otherwise} \end{cases}.$$

Its cdf can also not be written in a closed form.

9.7 Normal distribution

Definition 9.7.1 (Standard normal distribution). A CRV X follows the **standard normal distribution** or **Gaussian distribution**, denoted $X \sim N(0, 1)$ if it satisfies,

$$f_X(x) = \phi(x) := \frac{1}{\sqrt{2\pi}} e^{-x^2/2} \quad \text{for } x \in \mathbb{R}, \quad F_X(x) = \frac{1}{\sqrt{2\pi}} \int_{-\infty}^x e^{-t^2/2} dt \quad \text{for } x \in \mathbb{R}.$$

Where, once again, there is no explicit formula for the cdf.

Definition 9.7.2 (Normal distribution). A CRV X follows the **normal distribution** with mean $\mu \in \mathbb{R}$ and variance σ^2 for $\sigma \in \mathbb{R}^{>0}$ denoted $X \sim N(\mu, \sigma^2)$ if it satisfies,

$$f_X(x) = \phi(x) := \frac{1}{\sqrt{2\pi\sigma^2}} e^{-\frac{(x-\mu)^2}{2\sigma^2}} \quad \text{for } x \in \mathbb{R}.$$

9.8 Cauchy distribution

Definition 9.8.1 (Cauchy distribution). A CRV X follows the **Cauchy distribution** if it satisfies,

$$f_X(x) = \frac{1}{\pi(1+x^2)} \quad \text{for } x \in \mathbb{R}, \quad F_X(x) = \frac{1}{\pi} \arctan(x) + \frac{1}{2} \quad \text{for } x \in \mathbb{R}.$$

If $X, Y \sim N(0, 1)$, then $Z = X/Y$ follows the Cauchy distribution.

9.9 Student t-distribution

Definition 9.9.1 ((Student's) t-distribution). A CRV X follows the **Student t-distribution** with $\nu \in \mathbb{R}^{>0}$ degrees of freedom if it satisfies,

$$f_X(x) = \frac{\Gamma\left(\frac{\nu+1}{2}\right)}{\sqrt{\nu\pi}\Gamma\left(\frac{\nu}{2}\right)} \left(1 + \frac{x^2}{\nu}\right)^{-\frac{\nu+1}{2}} \quad \text{for } x \in \mathbb{R}.$$

Its cdf cannot be written in a closed form.

Remark 9.9.2. Not all RVs are either discrete or continuous.

10 Transformations of random variables

10.1 DRVs

Theorem 10.1.1. Let X be a DRV on (Ω, \mathcal{F}, P) and let $g : \mathbb{R} \rightarrow \mathbb{R}$ be a deterministic function, then $Y = g(X)$ is a DRV with pmf:

$$p_Y(y) = \sum_{\{x \in \text{im } X : g(x)=y\}} p_X(x) \quad \text{if } y \in \text{im } Y \text{ and } 0 \text{ otherwise.}$$

10.2 CRVs

Theorem 10.2.1. Let X be a CRV on (Ω, \mathcal{F}, P) and let $g : \mathbb{R} \rightarrow \mathbb{R}$ be a strictly monotonic and differentiable function with inverse $g^{-1} : \mathbb{R} \rightarrow \mathbb{R}$, then $Y = g(X)$ is a CRV with pdf:

$$f_Y(y) = f_X(g^{-1}(y)) \left| \frac{d}{dy} [g^{-1}(y)] \right| \quad \text{for all } y \in \mathbb{R}.$$

Remark 10.2.2. The term $\left| \frac{d}{dy} [g^{-1}(y)] \right|$ is often called the **Jacobian** of the transformation.

11 Expectation of random variables

Throughout this section, unless otherwise specified, all infinite sums will be assumed to converge absolutely and all integrals will be assumed to be $< \infty$.

11.1 Definition

Definition 11.1.1 (Expectation of a DRV). Let X be a DRV on (Ω, \mathcal{F}, P) then the **expectation** of X is defined by,

$$E(X) := \sum_{x \in \text{im } X} x p_X(x).$$

Definition 11.1.2 (Expectation of a CRV). Let X be a CRV on (Ω, \mathcal{F}, P) then the **expectation** of X is defined by,

$$E(X) := \int_{-\infty}^{\infty} x f_X(x) dx.$$

11.2 LOTUS

Theorem 11.2.1 (Discrete LOTUS). Let X be a DRV on (Ω, \mathcal{F}, P) and $g : \mathbb{R} \rightarrow \mathbb{R}$, we have,

$$E(g(X)) = \sum_{x \in \text{im } X} g(x) p_X(x).$$

Theorem 11.2.2 (Continuous LOTUS). Let X be a CRV on (Ω, \mathcal{F}, P) and $g : \mathbb{R} \rightarrow \mathbb{R}$, we have,

$$E(g(X)) = \int_{-\infty}^{\infty} g(x) f_X(x) dx.$$

Note that this is one of the few theorems throughout the course given without proof.

Theorem 11.2.3. If X is non-negative then $E(X) \geq 0$.

11.3 Variance

Definition 11.3.1 (Variance). Let X be a discrete/continuous random variable, then the **variance** of X is defined by,

$$\text{Var}(X) := E[X - E(X)]^2.$$

Theorem 11.3.2. For a discrete/continuous random variable with finite variance,

$$\text{Var}(X) = E(X^2) - [E(X)]^2.$$

12 Multivariate random variables

12.1 Multivariate distributions

Definition 12.1.1 (Joint distribution function). Consider the sequence of random variables X_1, X_2, \dots, X_n all on (Ω, \mathcal{F}, P) and write $\mathbf{X} = (X_1, X_2, \dots, X_n)$ and $\mathbf{x} = (x_1, x_2, \dots, x_n) \in \mathbb{R}^n$. Then the **joint distribution function** of \mathbf{X} is $F_{\mathbf{X}} : \mathbb{R}^n \rightarrow [0, 1]$ defined by:

$$F_{\mathbf{X}}(\mathbf{x}) := P(X_1 \leq x_1, X_2 \leq x_2, \dots, X_n \leq x_n) \quad \text{for all } \mathbf{x} \in \mathbb{R}^n.$$

12.2 Independence

Definition 12.2.1 (Pairwise independence for n random variables). We call the sequence of RVs, X_1, X_2, \dots, X_n , **pairwise independent** iff,

$$F_{X_i, X_j}(x_i, x_j) = F_{X_i}(x_i)F_{X_j}(x_j) \quad \text{for all } x_i, x_j \in \mathbb{R} \text{ with } i \neq j.$$

Definition 12.2.2 (Independence of a family of random variables). Given some indexing set $\mathcal{I} \subset \mathbb{R}$, a family of random variables $\{X_i : i \in \mathcal{I}\}$ is **independent** iff for all finite $\mathcal{J} \subseteq \mathcal{I}$:

$$P\left(\bigcap_{j \in \mathcal{J}} \{X_j \leq x_j\}\right) = \prod_{j \in \mathcal{J}} P(\{X_j \leq x_j\}) \quad \text{for all } (x_j)_{j \in \mathcal{J}} \text{ with } x_j \in \mathbb{R}.$$

(All finite subfamilies of the family of random variables is independent by the natural definition)

12.3 Multivariate DRVs

Definition 12.3.1 (Joint probability mass functions). Let X_1, X_2, \dots, X_n all be DRVs on (Ω, \mathcal{F}, P) that form the random vector \mathbf{X} , their **joint probability mass function**, $p_{\mathbf{X}} : \mathbb{R}^n \rightarrow [0, 1]$ is defined as,

$$p_{\mathbf{X}}(x_1, x_2, \dots, x_n) := P(\{\omega \in \Omega : X_1(\omega) = x_1, X_2(\omega) = x_2, \dots, X_n(\omega) = x_n\}) = P(X_1 = x_1, X_2 = x_2, \dots, X_n = x_n).$$

The **marginal probability mass function** of $X_i \in \mathbf{X}$ is given by,

$$p_{X_i}(k) = \sum_{(x_1, x_2, \dots, x_n) \in \mathbb{R}^n} \dots \sum p_{\mathbf{X}}(x_1, x_2, \dots, x_{i-1}, k, x_{i+1}, \dots, x_n).$$

It can be obtained that for any sufficiently "nice" set $A \in \mathbb{R}^n$,

$$P(\mathbf{X} \in A) = \sum_{(x_1, x_2, \dots, x_n) \in A} \dots \sum P(X_1 = x_1, X_2 = x_2, \dots, X_n = x_n).$$

Definition 12.3.2 (Independence of DRVs). Given some indexing set $\mathcal{I} \in \mathbb{R}$ a family of DRVs, $\{X_i : i \in \mathcal{I}\}$ with joint pmf $p_{\mathbf{X}}$, is **independent** iff for all finite $\mathcal{J} \subseteq \mathcal{I}$:

$$p_{\mathbf{X}}(\mathbf{x}) = \prod_{i=1}^n p_{X_i}(x_i) \quad \text{for all } \mathbf{x} = (x_1, x_2, \dots, x_n) \in \mathbb{R}^n.$$

12.4 Multivariate CRVs

Definition 12.4.1 (Continuous random vector). The random vector $\mathbf{X} = (X_1, X_2, \dots, X_n)$ is a **continuous random vector** iff,

$$F_{\mathbf{X}}(\mathbf{x}) = \int \cdots \int_B f_{\mathbf{X}}(\mathbf{x}) dx_1 dx_2 \cdots dx_n \quad \text{with } B = (\infty, x_1] \times (\infty, x_2] \times \cdots \times (\infty, x_n], \quad \text{for all } \mathbf{x} \in \mathbb{R}^n;$$

for some $f_{\mathbf{X}} : \mathbb{R}^n \rightarrow \mathbb{R}$ satisfying: $f_{\mathbf{X}}(\mathbf{x}) \geq 0$ for all $\mathbf{x} \in \mathbb{R}^n$ and $\int \cdots \int_{\mathbb{R}^n} f_{\mathbf{X}}(\mathbf{x}) dx_1 dx_2 \cdots dx_n = 1$.

Note that $f_{\mathbf{X}}(\mathbf{x}) = \frac{\partial^n}{\partial x_1 \partial x_2 \cdots \partial x_n} F_{\mathbf{X}}(\mathbf{x})$ and $P(\mathbf{X} \in A) = \int \cdots \int_A f_{\mathbf{X}}(\mathbf{x}) d^n \mathbf{x}$.

Definition 12.4.2 (Independence of CRVs). Given some indexing set $\mathcal{I} \in \mathbb{R}$ a family of CRVs, $\{X_i : i \in \mathcal{I}\}$ with joint pdf $f_{\mathbf{X}}$, is **independent** iff for all finite $\mathcal{J} \in \mathcal{I}$:

$$f_{\mathbf{X}}(\mathbf{x}) = \prod_{i=1}^n f_{X_i}(x_i) \quad \text{for all } \mathbf{x} = (x_1, x_2, \dots, x_n) \in \mathbb{R}^n.$$

12.5 Transformations of random vector

Definition 12.5.1 (Transformation). We are going to **transform** the random vector \mathbf{X} with joint pdf $f_{\mathbf{X}}$ to $\mathbf{U} = (u_1(\mathbf{X}), u_2(\mathbf{X}), \dots, u_n(\mathbf{X}))$ with $u_i : \mathbb{R}^n \rightarrow \mathbb{R}^n$ for all $i \in [1, n]$. Firstly, define $T : \mathbb{R}^n \rightarrow \mathbb{R}^n$ by $T(\mathbf{x}) = (u_1(\mathbf{x}), u_2(\mathbf{x}), \dots, u_n(\mathbf{x}))$ and assume T is bijective on $D = \{\mathbf{x} \in \mathbb{R}^n : f_{\mathbf{X}}(\mathbf{x}) > 0\}$ with range $S \subseteq \mathbb{R}^n$. Secondly, have the Jacobian determinant of $T^{-1} : S \rightarrow D$, $J(\mathbf{u}) = \det([a_{ij}]_{m \times n})$ with $a_{ij} = \frac{\partial x_i}{\partial u_j}$. Finally, define:

$$f_{\mathbf{U}}(\mathbf{u}) := \begin{cases} f_{\mathbf{X}}(T^{-1}(\mathbf{u})) |J(\mathbf{u})| & \text{if } \mathbf{u} \in S \\ 0 & \text{otherwise} \end{cases}.$$

12.6 Multivariate LOTUS

Theorem 12.6.1 (Discrete multivariate LOTUS). If X_1, X_2, \dots, X_n are DRVs on (Ω, \mathcal{F}, P) and form the random vector \mathbf{X} with $g : \mathbb{R}^n \rightarrow \mathbb{R}$, then $Y = g(\mathbf{X})$ is a DRV on (Ω, \mathcal{F}, P) with expectation,

$$E(g(\mathbf{X})) = \sum_{x_i \in \text{im } X_i} \cdots \sum g(\mathbf{X}) P(\mathbf{X} = \mathbf{x}) \quad \text{for all } \mathbf{x} = (x_1, x_2, \dots, x_n) \in \mathbb{R}^n.$$

Theorem 12.6.2 (Continuous multivariate LOTUS). If X_1, X_2, \dots, X_n are DRVs on (Ω, \mathcal{F}, P) and form the random vector \mathbf{X} with $h : \mathbb{R}^n \rightarrow \mathbb{R}$ we have,

$$E(h(\mathbf{X})) = \int \cdots \int_{\mathbb{R}^n} g(\mathbf{x}) f_{\mathbf{X}}(\mathbf{x}) dx_1 dx_2 \cdots dx_n \quad \text{for all } \mathbf{x} = (x_1, x_2, \dots, x_n) \in \mathbb{R}^n.$$

12.7 Covariance

Definition 12.7.1 (Covariance). Given two random variable X and Y on the same probability space with expectations μ_X and μ_Y respectively. The **covariance** of X and Y is defined as,

$$\text{Cov}(X, Y) := E[(X - \mu_X)(Y - \mu_Y)] \quad \text{assuming both expectation take finite values.}$$

Definition 12.7.2 (Correlation). Given the same X and Y the **correlation** of X and Y is defined as,

$$\text{Cor}(X, Y) := \frac{\text{Cov}(X, Y)}{\sqrt{\text{Var}(X)\text{Var}(Y)}}.$$

Theorem 12.7.3. For jointly discrete/continuous RVs with finite expectations the following hold:

1. when $X = Y$, $\text{Cov}(X, Y) = E[(X - \mu_X)^2] = \text{Var}(X)$,
2. $\text{Cov}(X, Y) = E(XY) - E(X)E(Y)$,
3. when X and Y are independent, $E(XY) = E(X)E(Y)$,
4. when variances are also finite, $\text{Var}(X + Y) = \text{Var}(X) + \text{Var}(Y) + 2\text{Cov}(X, Y)$.

13 Generating functions

13.1 Probability generating functions

Definition 13.1.1 (Probability generating functions). Given a DRV X with $\text{im}(X) \subseteq \mathbb{N}_0$, denote,

$$\mathcal{S}_X := \left\{ s \in \mathbb{R} : \sum_{x=0}^{\infty} |s|^x \mathbb{P}(X = x) < \infty \right\},$$

and define the **probability generating function (pgf)** of X as the function $G_X : \mathcal{S}_X \rightarrow \mathbb{R}$ given by,

$$G_X(s) := \mathbb{E}(s^X) = \sum_{x=0}^{\infty} s^x \mathbb{P}(X = x),$$

noting that the pgf is well defined for $|s| < 1$ and $G_X(0) = \mathbb{P}(X = 0)$ and $G_X(1) = 1$.

Theorem 13.1.2 (Uniqueness of pgfs). Given two DRVs X and Y with pgfs G_X and G_Y respectively,

$$G_X(s) = G_Y(s) \quad \text{for all } s \in \mathcal{S}_X \cap \mathcal{S}_Y \iff p_X(x) = p_Y(x) \quad \text{for all } x \in \mathbb{N}_0.$$

Theorem 13.1.3. Let X, Y be independent DRVs with $\text{im } X, \text{im } Y \in \mathbb{N}_0$, then

$$G_{X+Y}(s) = G_X(s)G_Y(s) \quad \text{for all } s \in \mathcal{S}_X \cap \mathcal{S}_Y.$$

Theorem 13.1.4 (Pgfs of sum of independent DRVs). Given a collection of n independent DRVs X_1, X_2, \dots, X_n ,

$$G_{\sum_{i=1}^n X_i}(s) = \prod_{i=1}^n G_{X_i}(s) \quad \text{for all } s \in \bigcap_{i=1}^n \mathcal{S}_{X_i}.$$

Theorem 13.1.5 (Moments). Given a DRV X with $\text{im } X \subseteq \mathbb{N}_0$, the k th derivative of G_X , for $k \in \mathbb{N}$ is given by,

$$\left. \frac{d^k}{ds^k} G_X(s) \right|_{s=1} = G_X^{(k)}(1) = \mathbb{E}[X(X-1)\dots(X-k+1)].$$

13.2 Common pgfs

Example 13.2.1 (Bernoulli distribution). Let $X \sim \text{Bern}(p)$, then $G_X(s) = 1 - p + sp$ for all $s \in \mathbb{R}$.

Example 13.2.2 (Binomial distribution). Let $X \sim \text{Bin}(n, p)$, then $G_X(s) = (1 - p + sp)^n$ for all $s \in \mathbb{R}$.

Example 13.2.3 (Poisson distribution). Let $X \sim \text{Poi}(\lambda)$, then $G_X(s) = \exp(\lambda(s-1))$ for all $s \in \mathbb{R}$.

13.3 Moment generating functions

Definition 13.3.1 (Moment generating function). Let X be a RV, its **moment generating function (mgf)** is defined as,

$$M_X(t) = \mathbb{E}(e^{tX}),$$

assuming the expectation exists in some neighbourhood of zero.

Remark 13.3.2. If X is a RV with a mgf, $M_X(t) = \mathbb{E}(e^{tX}) = G_X(e^t)$

Theorem 13.3.3. If X is a RV with a mgf, the k th moment of X is $\mathbb{E}(X^k) = M_X^{(k)}(0)$.

Theorem 13.3.4. If X_1, X_2, \dots, X_n are a family of independent RVs on the same probability space with mgfs $M_{X_1}, M_{X_2}, \dots, M_{X_n}$ respectively, we have,

$$M_{\sum_{i=1}^n X_i}(t) = \prod_{i=1}^n M_{X_i}(t).$$

Theorem 13.3.5 (Characterisation by mgf). If the RVs X, Y have existent mgfs M_X, M_Y respectively such that $M_X(t) = M_Y(t)$ for all t in some neighbourhood of 0, we have,

$$F_X(u) = F_Y(u) \quad \text{for all } u.$$

14 Conditional distribution and expectation

14.1 Discrete: Conditional expectation and total expectation

Definition 14.1.1 (Condition distribution and expectation of a DRV). Given a DRV Y on (Ω, \mathcal{F}, P) and some event $B \in \mathcal{F}$ with $P(B) > 0$, the **conditional distribution** of Y given B is defined as,

$$P(Y = y|B) := \frac{P(\{Y = y\} \cap B)}{P(B)} \quad \text{for } y \in \mathbb{R};$$

with the **conditional expectation** of Y given B defined as,

$$E(Y|B) := \sum_{i \in \text{im } Y} eP(Y = y|B).$$

Theorem 14.1.2 (Discrete law of total expectation). Given a DRV Y on (Ω, \mathcal{F}, P) and some partition $\{B_i : i \in \mathcal{I}\}$ of Ω with $P(B_i) > 0$ for all $i \in \mathcal{I}$ we have,

$$E(Y) = \sum_{i \in \mathcal{I}} E(Y|B_i)P(B_i).$$

14.2 Conditioning on a DRV

Definition 14.2.1 (Conditional probability mass function). Given two joint DRVs (X, Y) , the **conditional probability mass function** of Y given $X = x$ is given by,

$$p_{Y|X}(y|x) := \frac{p_{X,Y}(x, y)}{p_X(x)} \quad \text{for } y \in \mathbb{R}.$$

Theorem 14.2.2 (Conditional expectation). Given two joint DRVs (X, Y) , the **conditional expectation** of Y given $X = x$ is given by,

$$E(Y|X = x) = \sum_{y \in \text{im } Y} yp_{Y|X}(y|x).$$

Independence, LOTUS and a Bayes' rule formulation all follow naturally from this as they do for the non-distribution settings.

14.3 Continuous: Conditional density, distribution and expectation

Definition 14.3.1 (Conditional distribution and conditional density). For two joint CRVs (X, Y) the **conditional density** of Y given $X = x$ is define as,

$$f_{Y|X}(y|x) := \frac{f_{X,Y}(x, y)}{f_X(x)} \quad \text{for all } y, x \in \mathbb{R} \text{ with } f_X(x) > 0;$$

with the corresponding **conditional distribution** of Y given $X = x$ defined as,

$$F_{Y|X=x}(y|x) := \frac{1}{f_X(x)} \int_{-\infty}^y f_{X,Y}(x, t) dt \quad \text{for all } y, x \in \mathbb{R} \text{ with } f_X(x) > 0.$$

Where, once again, familiar formulations for independence and Bayes' rule can be easily derived.

Definition 14.3.2 (Conditional expectation). Given two joint CRVs (X, Y) , the **conditional expectation** of Y given $X = x$ is defined as,

$$E(Y|X = x) := \int_{-\infty}^{\infty} yf_{Y|X}(y|x)dy \quad \text{provided } f_X(x) > 0.$$

Theorem 14.3.3 (Continuous law of total expectation). Given two joint CRVs (X, Y) with $E(|Y|) < \infty$, we have,

$$E(Y) = \int_{\{x: f_X(x) > 0\}} E(Y|X = x)f_X(x)dx.$$

Chapter 7

Statistics

1 Introduction

The following are references.

- E Artin, Galois theory, 1994
- A Grothendieck and M Raynaud, Revêtements étales et groupe fondamental, 2002
- I N Herstein, Topics in algebra, 1975
- M Reid, Galois theory, 2014

Notation. If K is a field, or a ring, I denote the **ring of polynomials** with coefficients in K .

Lecture 1
Thursday
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2 Central tendency and dispersion

- 2.1 Mean, variance and moments
- 2.2 Parameter estimation
- 2.3 Other measures of central tendency
- 2.4 Sampling from normal RVs

3 Hypothesis testing

- 3.1 Introduction
- 3.2 Single sample hypothesis testing
- 3.3 Distribution of p-values
- 3.4 Errors
- 3.5 Two sample hypothesis testing
- 3.6 Multiple hypothesis testing

4 Covariance and Correlations

- 4.1 Covariance
- 4.2 Correlation

5 Statistical models

- 5.1 Definitions
- 5.2 Likelihood
- 5.3 Linear regression

6 Bayesian inference

- 6.1 Definitions
- 6.2 Conjugate pair distributions
- 6.3 Intractable posteriors
- 6.4 Choosing a prior

7 Bootstrap

- 7.1 Empirical distribution
- 7.2 Bootstrap procedure

Chapter 8

Computation

1 Introduction

The following are references.

- E Artin, Galois theory, 1994
- A Grothendieck and M Raynaud, Revêtements étales et groupe fondamental, 2002
- I N Herstein, Topics in algebra, 1975
- M Reid, Galois theory, 2014

Notation. If K is a field, or a ring, I denote the **ring of polynomials** with coefficients in K .

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Chapter 9

Applied Mathematics

1 Introduction

The following are references.

- E Artin, Galois theory, 1994
- A Grothendieck and M Raynaud, Revêtements étales et groupe fondamental, 2002
- I N Herstein, Topics in algebra, 1975
- M Reid, Galois theory, 2014

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