Part II — Probability and Measure

Based on lectures by Dr Sarkar

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§0 Holes in classical theory

Analysis

1. What is the "volume" of a subset of $\mathbb{R}^d$.
2. Integration (Riemann Integration has holes)
   - $\{f_n\}$ a sequence of continuous functions on $[0,1]$ s.t.
     - $0 \leq f_n(x) \leq 1 \; \forall \; x \in [0,1]$.
     - $f_n(x)$ is monotonically decreasing on $n \to \infty$, i.e. $f_n(x) \geq f_{n+1}(x) \; \forall \; x/$

   So, $\lim_{n\to\infty} f_n(x)$ exists. But f is not Riemann integrable . We want a theory of integration s.t. f is integrable and $\lim_{n\to\infty} \int_0^1 f_n(x)\,\mathrm{d}x = \int_0^1 f(x)\,\mathrm{d}x$.
3. $L^1 = ()$ If $f \in L^1$ is f Riemann integrable? Will have to change the definition of integral. L^2 a hilbert space

Probability

1. Discrete probability has its limitations,
   - Toss a unbiased coin 5 times. What is the probability if getting 3 heads?
   - Take an infinite sequence of coin tosses ($E = \{0,1\}^{\mathbb{N}}$ which is uncountable) and an event A that depends on that infinite sequence. How do you define $\mathbb{P}(A)$? E.g. $X_i \sim \text{Ber}\left(\frac{1}{2}\right)$ and $A = \frac{\sum_{i=1}^{n} X_i}{n}$, the average number of heads. By strong law of large numbers $\mathbb{P}\left(\frac{\sum_{i=1}^{n} X_i}{n} \to \frac{1}{2}\right) = 1$.
   - How to draw a point uniformly at random from $[0,1]$? $U \sim U[0,1]$. Probability needs axioms to be made rigorous.
2. Define Expectation for a r.v.. Also would want the following if $0 \leq X_n \leq 1$ and $X_n \downarrow X$ then $\mathbb{E}X_n \to \mathbb{E}X$.

§1 Introduction

Notation. $A_n \uparrow A$ means that the sequence A_n is increasing ($A_1 \subseteq A_2 \subseteq \dots$) and $\bigcup_n A_n = A$.

§1.1 Definitions

Definition 1.1 (σ-algebra)
Let E be a (nonempty) set. A collection $\mathcal{E}$ of subsets of E is called a **σ-algebra** if the following properties hold:

- $\varnothing \in \mathcal{E}$;
- $A \in \mathcal{E} \implies A^c = E \setminus A \in \mathcal{E}$;
- if $(A_n)_{n\in\mathbb{N}}$ is a countable collection of sets in $\mathcal{E}$, $\bigcup_{n\in\mathbb{N}} A_n \in \mathcal{E}$.

Example 1.1
Let $\mathcal{E} = \{\varnothing, E\}$. This is a σ-algebra. Also, $\mathcal{P}(E) = \{A \subseteq E\}$ is a σ-algebra.

Remark 1. Since $\bigcap_n A_n = (\bigcup_n A_n^c)^c$, any σ-algebra $\mathcal{E}$ is closed under countable intersections as well as under countable unions. Note that $B \setminus A = B \cap A^c \in \mathcal{E}$, so σ-algebras are closed under set difference.

Definition 1.2 (Measurable Space and Set)
A set E with a σ-algebra $\mathcal{E}$ is called a **measurable space**. The elements of $\mathcal{E}$ are called **measurable sets**.

Definition 1.3 (Measure)
A **measure** μ is a set function $\mu : \mathcal{E} \to [0, \infty]$, such that $\mu(\varnothing) = 0$, and for a sequence $(A_n)_{n\in\mathbb{N}}$ such that the A_n are disjoint, we have

$$\mu\left(\bigcup_{n\in\mathbb{N}} A_n\right) = \sum_{n\in\mathbb{N}} \mu(A_n)$$

This is the **countable additivity** property of the measure.

Remark 2. $(E, \mathcal{E}, \mu)$ is a measure space.

Remark 3. If E is countable, then for any $A \in \mathcal{P}(E)$ and measure μ, we have

$$\mu(A) = \mu\left(\bigcup_{x \in A} \{x\}\right) = \sum_{x \in A} \mu(\{x\})$$

Hence, measures are uniquely defined by the measure of each singleton.

Define $m : E \to [0, \infty]$ s.t. $m(x) = \mu(\{x\})$, such an m is called a "mass function", and measures μ are in 1-1 correspondence with the mass function m. This corresponds to the notion of a probability mass function.

Here $\mathcal{E} = \mathcal{P}(E)$ and this is the theory in elementary discrete prob. (when $\mu(\{x\}) = 1 \; \forall \; x \in E$, μ is called the counting measure. Here $\mu(A) = |A| \; \forall \; A \subset E$).

For uncountable E however, the story is not so simple and $\mathcal{E} = \mathcal{P}(E)$ is generally not feasible. Indeed measures are defined on σ-algebra "generated" by a smaller class $\mathcal{A}$ of simple subsets of E.

Definition 1.4 (Generated $\sigma-$algebra)

For a collection $\mathcal{A}$ of subsets of E, we define the σ-algebra $\sigma(A)$ **generated by** $\mathcal{A}$ by

$$\sigma(\mathcal{A}) = \{A \subseteq E \colon A \in \mathcal{E} \text{ for all } \sigma\text{-algebras } \mathcal{E} \supseteq \mathcal{A}\}$$

So it is the smallest σ-algebra containing $\mathcal{A}$. Equivalently,

$$\sigma(\mathcal{A}) = \bigcap_{\mathcal{E} \supseteq \mathcal{A},\, \mathcal{E} \text{ a } \sigma\text{-algebra}} \mathcal{E}$$

Question

Why is $\sigma(A)$ a σ-algebra? See Sheet 1, Q1.

§1.2 Rings and algebras

The class $\mathcal{A}$ will usually satisfy some properties too, let E be a set and $\mathcal{A}$ a collection of subsets of E. To construct good generators, we define the following.

Definition 1.5 (Ring)

$\mathcal{A} \subseteq \mathcal{P}(E)$ is called a **ring** over E if $\varnothing \in \mathcal{A}$ and $A, B \in \mathcal{A}$ implies $B \setminus A \in \mathcal{A}$ and $A \cup B \in \mathcal{A}$.

Rings are easier to manage than σ-algebras because there are only finitary operators.

Definition 1.6 (Algebra)

$\mathcal{A}$ is called an **algebra** over E if $\varnothing \in \mathcal{A}$ and $A, B \in \mathcal{A}$ implies $A^c \in \mathcal{A}$ and $A \cup B \in \mathcal{A}$.

Remark 4. Rings are closed under symmetric difference $A \triangle B = (B \setminus A) \cup (A \setminus B)$, and are closed under intersections $A \cap B = A \cup B \setminus A \triangle B$. Algebras are rings, because $B \setminus A = B \cap A^c = (B^c \cup A)^c$. Not all rings are algebras, because rings do not need to include the entire space.

The idea:

- Define a set function on a suitable collection $\mathcal{A}$.
- Extend the set function to a measure on $\sigma(\mathcal{A})$. (Carathéodory's Extension theorem)
- Such an extension is unique. (Dynkin's Lemma)

Goal: Start with a "measure" on $\mathcal{A}$ that has some nice properties and then extend it to $\sigma(A)$.

Definition 1.7 (Set Function)

A **set function** on a collection $\mathcal{A}$ of subsets of E, where $\varnothing \in \mathcal{A}$, is a map $\mu \colon \mathcal{A} \to [0, \infty]$ such that $\mu(\varnothing) = 0$.

- We say μ is **increasing** if $\mu(A) \leq \mu(B)$ for all $A \subseteq B$ in $\mathcal{A}$.
- We say μ is **additive** if $\mu(A \cup B) = \mu(A) + \mu(B)$ for disjoint $A, B \in \mathcal{A}$ and $A \cup B \in \mathcal{A}$.
- We say μ is **countably additive** if $\mu(\bigcup_n A_n) = \sum_n \mu(A_n)$ for disjoint sequences A_n where $\bigcup_n A_n$ and each A_n lie in $\mathcal{A}$.
- We say μ is **countably subadditive** if $\mu(\bigcup_n A_n) \leq \sum_n \mu(A_n)$ for arbitrary sequences A_n under the above conditions.

Remark 5. If μ is countably additive set function on $\mathcal{A}$ and $\mathcal{A}$ is a ring then μ satisfies all the previous listed properties.

Proposition 1.1 (Disjointification of countable unions)

Consider $\bigcup_n A_n$ for $A_n \in \mathcal{E}$, where $\mathcal{E}$ is a σ-algebra (or a ring, if the union is finite). Then there exist $B_n \in \mathcal{E}$ that are disjoint such that $\bigcup_n A_n = \bigcup_n B_n$.

Proof. Define $\widetilde{A}_n = \bigcup_{j \leq n} A_j$, then $B_n = \widetilde{A}_n \setminus \widetilde{A}_{n-1}$. □

Remark 6. A measure satisfies all four of the above conditions. Countable additivity implies the other conditions. Proof on Sheet 1.

Theorem 1.1 (Carathéodory's theorem)
Let μ be a countably additive set function on a ring $\mathcal{A}$ of subsets of E. Then there exists a measure $\mu^\star$ on $\sigma(\mathcal{A})$ such that $\mu^\star|_{\mathcal{A}} = \mu$.

We will later prove that this extended measure is unique.

Proof. For $B \subseteq E$, we define the **outer measure** $\mu^\star$ as

$$\mu^\star(B) = \inf\left\{\sum_{n\in\mathbb{N}} \mu(A_n) : A_n \in \mathcal{A}, B \subseteq \bigcup_{n\in\mathbb{N}} A_n\right\}$$

If there is no sequence A_n such that $B \subseteq \bigcup_{n\in\mathbb{N}} A_n$, we declare the outer measure $\mu^\star(B)$ to be ∞. Clearly, $\mu^\star(\varnothing)$ and $\mu^\star$ is increasing, so $\mu^\star$ is an increasing set fcn on $\mathcal{P}(E)$.

Definition 1.8 ($\mu^\star$ measurable)
A set $A \subseteq E$ $\mu^\star$ **measurable** if $\forall\, B \subseteq E$ $\mu^\star(B) = \mu^\star(B \cap A) + \mu^\star(B \cap A^c)$.

We define the class

$$\mathcal{M} = \{A \subseteq E : A \text{ is } \mu^\star \text{ measurable}\}$$

We shall show that M is a σ-algebra that contains $\mathcal{A}$, $\mu^\star|_M$ is a measure on M that extends μ (i.e. $\mu^\star|_{\mathcal{A}} = \mu$).

Step 1. $\mu^\star$ is countably sub-additive on $\mathcal{P}(E)$: It suffices to prove that for $B \subseteq E$ and $B_n \subseteq E$ such that $B \subseteq \bigcup_n B_n$ we have

$$\mu^\star(B) \leq \sum_n \mu^\star(B_n) \tag{†}$$

We can assume without loss of generality that $\mu^\star(B_n) < \infty$ for all n, otherwise there is nothing to prove. For all $\varepsilon > 0$ there exists a collection $A_{n,m} \in \mathcal{A}$ such that $B_n \subseteq \bigcup_m A_{n,m}$ and

$$\mu^\star(B_n) + \frac{\varepsilon}{2^n} \geq \sum_m \mu(A_{n,m})$$

as we took an infimum. Now, since $\mu^\star$ is increasing, and $B \subseteq \bigcup_n B_n \subseteq \bigcup_n \bigcup_m A_{n,m}$, we have

$$\mu^\star(B) \leq \mu^\star\left(\bigcup_{n,m} A_{n,m}\right) \leq \sum_{n,m} \mu(A_{n,m}) \leq \sum_n \mu^\star(B_n) + \sum_n \frac{\varepsilon}{2^n} = \sum_n \mu^\star(B_n) + \varepsilon$$

Since ε was arbitrary in the construction, (†) follows by construction.

Step 2. $\mu^\star$ extends μ: Let $A \in \mathcal{A}$, and we want to show $\mu^\star(A) = \mu(A)$.

We can write $A = A \cup \varnothing \cup \ldots$, hence $\mu^\star(A) \leq \mu(A) + 0 + \cdots = \mu(A)$ by definition of $\mu^\star$.

If $\mu^\star$ is infinite, there is nothing to prove.

We need to prove the converse, that $\mu(A) \leq \mu^\star(A)$. For the finite case, suppose there is a sequence A_n where $\mu(A_n) < \infty$ and $A \subseteq \bigcup_n A_n$. Then, $A = \bigcup_n (A \cap A_n)$, which is a union of elements of the ring $\mathcal{A}$. As μ is countably additive on $\mathcal{A}$ and $\mathcal{A}$ is a ring, μ is countably subadditive on $\mathcal{A}$ and increasing by remark 6. Hence $\mu(A) \leq \sum_n \mu(A \cap A_n) \leq \sum_n \mu(A_n)$. Since the A_n were arbitrary taking the infimum over A_n, we have $\mu(A) \leq \mu^\star(A)$ as required.

Step 3. $\mathcal{M} \supseteq \mathcal{A}$: Let $A \in \mathcal{A}$. We must show that for all $B \subseteq E$, $\mu^\star(B) = \mu^\star(B \cap A) + \mu^\star(B \cap A^c)$.

We have $B \subseteq (B \cap A) \cup (B \cap A^c) \cup \varnothing \cup \ldots$, hence by countable subadditivity (†), $\mu^\star(B) \leq \mu^\star(B \cap A) + \mu^\star(B \cap A^c)$.

It now suffices to prove the converse, that $\mu^\star(B) \geq \mu^\star(B \cap A) + \mu^\star(B \cap A^c)$.
We can assume $\mu^\star(B)$ is finite, and so $\forall\ \varepsilon > 0\ \exists\ A_n \in \mathcal{A}$ s.t. $B \subseteq \bigcup_n A_n$ and $\mu^\star(B) + \varepsilon \geq \sum_n \mu(A_n)$. Now, $B \cap A \subseteq \bigcup_n (A_n \cap A)$, and $B \cap A^c \subseteq \bigcup_n (A_n \cap A^c)$. All of the members of these two unions are elements of $\mathcal{A}$, since $A_n \cap A^c = A_n \setminus A$. Therefore,

$$\begin{aligned}
\mu^\star(B \cap A) + \mu^\star(B \cap A^c) &\leq \sum_n \mu(A_n \cap A) + \sum_n \mu(A_n \cap A^c) \\
&\leq \sum_n [\mu(A_n \cap A) + \mu(A_n \cap A^c)] \\
&\leq \sum_n \mu(A_n) \leq \mu^\star(B) + \varepsilon
\end{aligned}$$

Since ε was arbitrary, $\mu^\star(B) = \mu^\star(B \cap A) + \mu^\star(B \cap A^c)$ as required.

Step 4. $\mathcal{M}$ is an algebra: Clearly $\varnothing$ lies in $\mathcal{M}$, and by the symmetry in the definition of $\mathcal{M}$, complements lie in $\mathcal{M}$. We need to check $\mathcal{M}$ is stable under finite intersections. Let $A_1, A_2 \in \mathcal{M}$ and let $B \subseteq E$. We have

$$\mu^\star(B) = \mu^\star(B \cap A_1) + \mu^\star(B \cap A_1^c) \text{ as } A_1 \in M$$

$$= \mu^\star(B \cap A_1 \cap A_2) + \mu^\star(B \cap A_1 \cap A_2^c) + \mu^\star(B \cap A_1^c) \text{ taking } \tilde{B} = B \cap A_1$$

We can write $A_1 \cap A_2^c = (A_1 \cap A_2^c)^c \cap A_1$, and $A_1^c = (A_1 \cap A_2)^c \cap A_1^c$. Hence

$$\begin{aligned}
\mu^\star(B) &= \mu^\star(B \cap A_1 \cap A_2) + \underbrace{\mu^\star(B \cap (A_1 \cap A_2)^c \cap A_1) + \mu^\star(B \cap (A_1 \cap A_2)^c \cap A_1^c)}_{\mu^\star(B \cap (A_1 \cap A_2)^c) \text{ as } A_1 \in M} \\
&= \mu^\star(B \cap A_1 \cap A_2) + \mu^\star(B \cap (A_1 \cap A_2)^c)
\end{aligned}$$

which is the requirement for $A_1 \cap A_2$ to lie in $\mathcal{M}$.

Step 5. $\mathcal{M}$ is a σ-algebra and $\mu^\star$ is a measure on $\mathcal{M}$:
It suffices now to show that $\mathcal{M}$ has countable unions and the measure respects these countable unions. Let $A = \bigcup_n A_n$ for $A_n \in \mathcal{M}$. Without loss of generality, let the A_n be disjoint. We want to show $A \in \mathcal{M}$, and that $\mu^\star(A) = \sum_n \mu^\star(A_n)$.

By (†), we have for any $B \subseteq E$ $\mu^\star(B) \leq \mu^\star(B \cap A) + \mu^\star(B \cap A^c) + 0 + \dots$ so we need to check only the converse of this inequality. Also, $\mu^\star(A) \leq \sum_n \mu^\star(A_n)$, so we need only check the converse of this inequality as well. Similarly to before,

$$\begin{aligned}
\mu^\star(B) &= \mu^\star(B \cap A_1) + \mu^\star(B \cap A_1^c) \\
&= \mu^\star(B \cap A_1) + \mu^\star(B \cap \underbrace{A_1^c \cap A_2}_{A_2 \text{ as } A_1, A_2 \text{ disjoint}}) + \mu^\star(B \cap A_1^c \cap A_2^c) \\
&= \mu^\star(B \cap A_1) + \mu^\star(B \cap A_2) + \mu^\star(B \cap A_1^c \cap A_2^c) \\
&= \mu^\star(B \cap A_1) + \mu^\star(B \cap A_2) + \mu^\star(B \cap A_1^c \cap A_2^c \cap A_3) + \mu^\star(B \cap A_1^c \cap A_2^c \cap A_3^c) \\
&= \mu^\star(B \cap A_1) + \mu^\star(B \cap A_2) + \mu^\star(B \cap A_3) + \mu^\star(B \cap A_1^c \cap A_2^c \cap A_3^c) \\
&= \cdots \\
&= \sum_{n \leq N} \mu^\star(B \cap A_n) + \mu^\star(B \cap A_1^c \cap \cdots \cap A_N^c)
\end{aligned}$$

Since $\bigcup_{n \leq N} A_n \subseteq A$, we have $\bigcap_{n \leq N} A_n^c \supseteq A^c$. $\mu^\star$ is increasing, hence, taking limits,

$$\mu^\star(B) \geq \sum_{n=1}^{\infty} \mu^\star(B \cap A_n) + \mu^\star(B \cap A^c)$$

By (†),

$$\mu^\star(B) \geq \mu^\star(B \cap A) + \mu^\star(B \cap A^c)$$

as required. Hence $\mathcal{M}$ is a σ-algebra. For the other inequality, we take the above result for $B = A$.

$$\mu^\star(A) \geq \sum_{n=1}^{\infty} \mu^\star(A \cap A_n) + \mu^\star(A \cap A^c) = \sum_{n=1}^{\infty} \mu^\star(A_n)$$

So $\mu^\star$ is countably additive on $\mathcal{M}$ and is hence a measure on $\mathcal{M}$. $\square$

§1.3 Uniqueness of extension

To address uniqueness of extension, we introduce further subclasses of $\mathcal{P}(E)$. Let $\mathcal{A}$ be a collection of subsets of E.

Definition 1.9 (π-system)
A collection $\mathcal{A}$ of subsets of E is called a **π-system** if $\varnothing \in \mathcal{A}$ and $A, B \in \mathcal{A} \implies A \cap B \in \mathcal{A}$.

Definition 1.10 (d-system)
A collection $\mathcal{A}$ of subsets of E is called a **d-system** if

- $E \in \mathcal{A}$;
- $A, B \in \mathcal{A}$ and $A \subseteq B$ then $B \setminus A \in \mathcal{A}$;
- $A_n \in \mathcal{A}$ is an increasing sequence of sets then $\bigcup_n A_n \in \mathcal{A}$.

Remark 7. Equivalently, $\mathcal{A}$ is a d-system if

- $\varnothing \in \mathcal{A}$;
- $A \in \mathcal{A} \implies A^c \in \mathcal{A}$
- $A_n \in \mathcal{A}$ is a sequence of disjoint sets then $\bigcup_n A_n \in \mathcal{A}$.

The difference between this and a σ-algebra is the requirement for disjoint sets.

Proof on Sheet 1.

Proposition 1.2
A d-system which is also a π-system is a σ-algebra.

Proof. Sheet 1. □

Lemma 1.1 (Dynkin's Lemma/π-λ/π-d theorem)
Let $\mathcal{A}$ be a π-system. Then any d-system that contains $\mathcal{A}$ also contains $\sigma(\mathcal{A})$.

Proof. We define

$$\mathcal{D} = \bigcap_{\mathcal{D}' \text{ is a } d\text{-system};\ \mathcal{D}' \supseteq \mathcal{A}} \mathcal{D}'$$

We can show this is a d-system (proof same as in $\sigma(\mathcal{A})$ on Sheet 1). It suffices to

prove that $\mathcal{D}$ is a π-system, because then it is a σ-algebra[a].

We now define

$$\mathcal{D}' = \{B \in \mathcal{D} : \forall A \in \mathcal{A}, B \cap A \in \mathcal{D}\}$$

We can see that $\mathcal{A} \subseteq \mathcal{D}'$, as $\mathcal{A}$ is a π-system.

We now show that $\mathcal{D}'$ is a d-system, fix $A \in \mathcal{A}$.

- Clearly $E \cap A = A \in \mathcal{A} \subseteq \mathcal{D}'$ hence $E \in \mathcal{D}'$.
- Let $B_1, B_2 \in \mathcal{D}'$ such that $B_1 \subseteq B_2$. Then $(B_2 \setminus B_1) \cap A = (B_2 \cap A) \setminus (B_1 \cap A)$, and since $B_i \cap A \in \mathcal{D}$ this difference also lies in $\mathcal{D}$, so $B_2 \setminus B_1 \in \mathcal{D}'$.
- Now, suppose B_n is an increasing sequence converging to B, and $B_n \in \mathcal{D}'$. Then $B_n \cap A \in \mathcal{D}$, and $\mathcal{D}$ is a d-system, we have $B \cap A \in \mathcal{D}$, so $B \in \mathcal{D}'$.

Hence $\mathcal{D}'$ is a d-system. Also, $\mathcal{D}' \subseteq \mathcal{D}$ by construction of $\mathcal{D}'$. But also $\mathcal{A} \subseteq \mathcal{D}'$ and $\mathcal{D}'$ is a d-system so $\mathcal{D} \subset \mathcal{D}'$ as $\mathcal{D}$ is the smallest d-system containing $\mathcal{A}$. Thus $\mathcal{D} = \mathcal{D}'$, i.e $\forall\ B \in \mathcal{D}$ and $A \in \mathcal{A}, B \cap A \in \mathcal{D}$ $(*)$.

We then define

$$\mathcal{D}'' = \{B \in \mathcal{D} : \forall A \in \mathcal{D}, B \cap A \in \mathcal{D}\}$$

Note that $\mathcal{A} \subseteq \mathcal{D}''$ by $(*)$. Running the same argument as before, we can show that $\mathcal{D}''$ is a d-system. So $\mathcal{D}'' = \mathbb{D}$. But then (by the definition of $\mathcal{D}''$), $\forall\ B \in \mathcal{D}, A \in \mathcal{D} \implies B \cap A \in \mathcal{D}$, i.e. $\mathcal{D}$ is a π-system (check that $\varnothing \in \mathcal{D}$).

So $\mathcal{D}$ is a σ-algebra containing $\mathcal{A}$, hence $\mathcal{D} \supseteq \sigma(\mathcal{A})$. $\square$

[a] As $\mathcal{D} \supseteq \mathcal{A}$ and $\sigma(\mathcal{A})$ the intersection of all σ-algebras containing $\mathcal{A}$, $\mathcal{D} \supseteq \sigma(\mathcal{A})$.

Theorem 1.2 (Uniqueness of Extension)

Let μ_1, μ_2 be measures on a measurable space $(E, \mathcal{E})$, such that $\mu_1(E) = \mu_2(E) < \infty$. Suppose that μ_1 and μ_2 coincide on a π-system $\mathcal{A}$, such that $\mathcal{E} \subseteq \sigma(\mathcal{A})$. Then $\mu_1 = \mu_2$ on $\sigma(\mathcal{A})$, and hence on $\mathcal{E}$.

Proof. We define

$$\mathcal{D} = \{A \in \mathcal{E} : \mu_1(A) = \mu_2(A)\}$$

This collection contains $\mathcal{A}$ by assumption. By Dynkin's lemma, it suffices to prove $\mathcal{D}$ is a d-system, because then $\mathcal{D} \supseteq \sigma(\mathcal{A}) \supseteq \mathcal{E}$ giving $\mathcal{D} = \mathcal{E}$ as $\mathcal{D} \subseteq \mathcal{E}$.

- $\varnothing \in \mathcal{D}$, since $\mu_1(\varnothing) = \mu_2(\varnothing) = 0$;
- $A \in \mathcal{D} \implies \mu_1(A) = \mu_2(A)$, thus $\mu_1(A^c) = \mu_1(E) - \mu_1(A) = \mu_2(E) - \mu_2(A) = \mu_2(A^c)$, so $A^c \in \mathcal{D}$ (μ_1, μ_2 finite so this works);

- Let $A_n \in \mathcal{D}$ be a disjoint sequence then, $\mu_1(\bigcup_n A_n) = \sum \mu_1(A_n) = \sum \mu_2(A_n) = \mu_2(\bigcup_n A_n)$ by countable additivity. So $\bigcup_n A_n \in \mathcal{D}$.

So $\mathcal{D}$ is a d-system. □

Remark 8. If $A_n \in \mathcal{A}$ an increasing sequence, then $\mu(\mathcal{A}) = \lim_{n\to\infty} \mu(A_n)$. Use this to show that $\mathcal{D}$ is a d-system satisfying conditions in d-system.

The above theorem applies to finite measures (μ such that $\mu(E) < \infty$) only. However, the theorem can be extended to measures that are σ-finite, for which $E = \bigcup_{n\in\mathbb{N}} E_n$ where $\mu(E_n) < \infty$.

Question

How to show all sets of a σ-algebra $\mathcal{E}$ generated by $\mathcal{A}$ has a certain property $\mathcal{P}$?

Answer

Consider set $\mathcal{G} = \{A \subseteq E : A \text{ has the property } \mathcal{P}\}$ and have that all elements of $\mathcal{A}$ have the property $\mathcal{P}$.

Method 1: Show that $\mathcal{G}$ is a σ-algebra, as it then must contain $\sigma(\mathcal{A}) = \mathcal{E}$.

Method 2: Show that $\mathcal{G}$ is a d-system and pick $\mathcal{A}$ s.t. it is a π-system and use Dynkin's Lemma/π-λ/π-d theorem.

Method 3: Monotone Convergence Theorem, we will see it shortly.

§1.4 Borel measures

Definition 1.11 (Borel Sets)

Let (E, τ) be a Hausdorff topological space. The σ-algebra generated by the open sets of E, i.e. $\sigma(\mathcal{A})$ where $\mathcal{A} = \{A \subseteq E : A \text{ open}\}$, is called the **Borel σ-algebra** on E, denoted $\mathcal{B}(E)$.
A measure μ on $(E, \mathcal{B}(E))$ is called a **Borel measure on E**.

Members of $\mathcal{B}(E)$ are called **Borel sets**.

Notation. We write $\mathcal{B} = \mathcal{B}(\mathbb{R})$.

Definition 1.12 (Radon Measure)

A **Radon measure** is a Borel measure μ on E such that $\mu(K) < \infty$ for all $K \subseteq E$ compact.

Note that in a Hausdorff space, compact sets are closed and hence measurable.

Definition 1.13 (Probability Measure)
If $\mu(E) = 1$, μ is called a **probability measure** on E, and $(E, \mathcal{E}, \mu)$ is called a probability space, typically denoted instead by $(\Omega, \mathcal{F}, \mathcal{P})$.

Definition 1.14 (Finite Measure)
If $\mu(E) < \infty$, μ is a **finite measure** on E.

Definition 1.15 (σ-finite Measure)
If $\exists$ sequence $E_n \in \mathcal{E}$ s.t. $\mu(E_n) < \infty \ \forall\ n$ and $E = \bigcup_n E_n$, then μ is called a **σ-finite measure**.

Remark 9. Arguments that hold for finite measures can usually be extended to σ-finite measures.

§1.5 Lebesgue measure

One of the main goals for this course is to define a notion of volume for arbitrary sets, we can do this by constructing a Borel measure μ on $\mathcal{B}(\mathbb{R}^d)$ s.t $\mu\left(\prod_{i=1}^d (a_i, b_i)\right) = \prod_{i=1}^d (b_i - a_i)$ where $a_i < b_i$ corresponding to the usual notion of volume of rectangles.

Initially, we will perform this construction for $d = 1$, and later we will consider product measures to extend this to higher dimensions.

Theorem 1.3 (Construction of the Lebesgue measure)
There exists a unique Borel measure μ on $\mathbb{R}$ such that

$$a < b \implies \mu\left((a, b]\right) = b - a. \tag{†}$$

μ is called the Lebesgue measure on $\mathbb{R}$.

Proof. First we shall prove the existence of the measure and then uniqueness.

Consider the ring $\mathcal{A}$ of finite unions of disjoint intervals[a] of the form

$$\mathcal{A} = (a_1, b_1] \cup \cdots \cup (a_n, b_n]$$

where $a_1 \leq b_1 \leq a_2 \leq \cdots \leq a_n \leq b_n$. Note that $\sigma(\mathcal{A}) = \mathcal{B}$ (see Example Sheets[b]).

Define for each $A \in \mathcal{A}$

$$\mu(A) = \sum_{i=1}^{n} (b_i - a_i).$$

This agrees with (†) for $(a, b]$. This is additive and well-defined (check).

So, the existence of μ on $\sigma(\mathcal{A}) = \mathcal{B}$ follows from Carathéodory's theorem if we can show that μ is *countable additive* on $\mathcal{A}$.

Remark 10. Suppose μ a finitely additive set function on a ring $\mathcal{A}$. Then μ is countable additive iff

- $A_n \uparrow$ [c] A; $A_n, A \in \mathcal{A} \implies \mu(A_n) \uparrow \mu(A)$.
- In addition, if μ is finite and $A_n \downarrow A$ s.t. $A_n, A \in \mathcal{A}$ then $\mu(A_n) \downarrow \mu(A)$[d].

See Example Sheet for proof.

So showing μ is countably additive on $\mathcal{A}$ is equivalent to showing the following
If $A_n \in \mathcal{A}, A_n \downarrow \varnothing$ then $\mu(A_n) \downarrow 0$. We require that μ is finite, as A_n decreasing we require A_1 to have finite measure. ?????

We shall prove this by contradiction.

Suppose this is not the case, so there exist $\varepsilon > 0$ and $B_n \in \mathcal{A}$ such that $B_n \downarrow \varnothing$ but $\mu(B_n) \geq 2\varepsilon$ for infinitely many n (and so wlog for all n).
We can approximate B_n from within by a sequence $\overline{C}_n$[e] $\in \mathcal{A}$ s.t. $C_n \subseteq B_n$ and $\mu(B_n \setminus C_n) \leq \varepsilon/2^n$. Suppose $B_n = \bigcup_{i=1}^{N_n} (a_{ni}, b_{ni}]$, then define $C_n = \bigcup_{i=1}^{N_n} (a_{ni} + \frac{2^{-n}\varepsilon}{N_n}, b_{ni}]$. Note that the C_n lie in $\mathcal{A}$, and $\mu(B_n \setminus C_n) \leq 2^{-n}\varepsilon$. Since B_n is decreasing, we have $B_N = \bigcap_{n \leq N} B_n$, and

$$B_N \setminus (C_1 \cap \cdots \cap C_N) = B_n \cap \left(\bigcup_{n \leq N} C_n^c \right) = \bigcup_{n \leq N} B_N \setminus C_n \subseteq \bigcup_{n \leq N} B_n \setminus C_n$$

Since μ is increasing and finitely additive and thus subadditive on $\mathcal{A}$,

$$\mu(B_N \setminus (C_1 \cap \cdots \cap C_N)) \leq \mu \left(\bigcup_{n \leq N} B_n \setminus C_n \right) \leq \sum_{n \leq N} \mu(B_n \setminus C_n) \leq \sum_{n \leq N} 2^{-N}\varepsilon \leq \varepsilon$$

Since $\mu(B_N) \geq 2\varepsilon$, additivity implies that $\mu(C_1 \cap \cdots \cap C_N) \geq \varepsilon$. This means that $C_1 \cap \cdots \cap C_N$ cannot be empty. We can add the left endpoints of the intervals, giving $K_N = \overline{C}_1 \cap \cdots \cap \overline{C}_N \neq \varnothing$. By Analysis I, K_N is a nested sequence of bounded nonempty closed intervals and therefore there is a point $x \in \mathbb{R}$ such that $x \in K_N$ for all N[f]. But $K_N \subseteq \overline{C}_N \subseteq B_N$, so $x \in \bigcap_N B_n$, which is a contradiction since $\bigcap_N B_N$ is empty. Therefore, a measure μ on $\mathcal{B}$ exists.

Now we prove uniqueness. Suppose μ, λ are measures such that the measure of an interval $(a, b]$ is $b - a$. We define truncated measures for $A \in \mathcal{B}$

$$\mu_n(A) = \mu\left(A \cap (n, n+1)\right)$$
$$\lambda_n(A) = \lambda\left(A \cap (n, n+1]\right)$$

Then μ_n, λ_n are *probability measures* on $\mathcal{B}$ and $\mu_n = \lambda_n$ on the π-system of intervals of the form $(a, b]$ with $a < b$[g]. This π-system generates $\mathcal{B}$, so by the uniqueness theorem for finite measures (theorem 1.2) $\mu_n = \lambda_n$ on $\mathcal{B}$. Hence $\forall\ A \in \mathcal{B}$

$$\begin{aligned}
\mu(A) &= \mu\left(\bigcup_n A \cap (n, n+1]\right) \\
&= \sum_{n \in \mathbb{Z}} \mu(A \cap (n, n+1]) \\
&= \sum_{n \in \mathbb{Z}} \mu_n(A) \\
&= \sum_{n \in \mathbb{Z}} \lambda_n(A) = \cdots = \lambda(A)
\end{aligned}$$

□

[a] We take semi intervals as for $\mathcal{A}$ to be a ring, we require the set difference to be in $\mathcal{A}$.
[b] as all open intervals are in $\sigma(\mathcal{A})$ and open intervals generate open sets
[c] increasing sequence tending to A
[d] E.g. let $A_n = [n, \infty)$ with the Lebesgue measure then $A_n \downarrow \varnothing$. But $\mu(A_n) = \infty$ whilst $\mu(\varnothing) = 0$
[e] $\overline{C}_n$ means the closure of C_n, i.e. make it a closed set by including the left endpoint
[f] As completeness of $\mathbb{R}$ implies $\bigcap_n K_n$ is closed and non empty.
[g] As $(a, b] \cap (c, d] = \varnothing$ or $(e, f]$.

Definition 1.16 (Lebesgue null set)

A Borel set $B \in \mathcal{B}$ is called a **Lebesgue null set** if $\lambda(B) = 0$ where λ is the Lebesgue measure.

Remark 11. A singleton $\{x\}$ can be written as $\bigcap_n \left(x - \frac{1}{n}, x\right]$, hence $\lambda(x) = \lim_n \frac{1}{n} = 0$. Hence singletons are null sets. In particular, $\lambda((a, b)) = \lambda((a, b]) = \lambda([a, b)) = \lambda([a, b])$. Any countable set $Q = \bigcup_q \{q\}$ is a null set. Not all null sets are countable; the Cantor set is an example.

The Lebesgue measure is *translation-invariant*. Let $x \in \mathbb{R}$, then the set $B + x = \{b + x : b \in B\}$ lies in $\mathcal{B}$ iff $B \in \mathcal{B}$, and in this case, it satisfies $\lambda(B + x) = \lambda(B)$. We can define the translated Lebesgue measure $\lambda_x(B) = \lambda(B + x)$ for all $B \in \mathcal{B}$, then $\lambda_x((a, b]) = \lambda((a, b] + x) = \lambda((a + x, b + x]) = b - a = \lambda((a, b])$. So $\lambda_x = \lambda$ on the π-system of intervals and so $\lambda_x = \lambda$ on the sigma algebra $\mathcal{B}$ (i.e. $\forall\, B \in \mathcal{B}, \lambda(B + x) = \lambda(B)$).

Question

Is the Lebesgue measure the only such translation invariant measure on $\mathcal{B}$?

Carathéodory's theorem extends λ from $\mathcal{A}$ to not just $\sigma(\mathcal{A}) = \mathcal{B}$, but actually to $\mathcal{M}$, the set of outer-measurable sets $M \supseteq \mathcal{B}$, but how large is $\mathcal{M}$?

The class of outer measurable sets $\mathcal{M}$ used in Carathéodory's extension theorem is here called the class of Lebesgue measurable sets. This class, the Lebesgue σ-algebra, can be shown to be

$$\mathcal{M} = \{A \cup N : A \in \mathcal{B}, N \subseteq B, B \in \mathcal{B}, \lambda(B) = 0\} \supsetneq \mathcal{B}$$

§1.6 Existence of non-measurable sets

We now show that $\mathcal{B} \subsetneq \mathcal{P}(\mathbb{R})$ (in fact $\mathcal{M}_{leb} \subsetneq \mathcal{P}(\mathbb{R})$).

Consider $E = [0, 1)$ with addition defined modulo one. By the same argument as before, the Lebesgue measure is translation-invariant modulo one. Consider the subgroup $Q = E \cap \mathbb{Q}$ of $(E, +)$. We define $x \sim y$ for $x, y \in E$ if $x - y \in Q$. Assuming the axiom of choice (uncountable version), we can select a representative from each equivalence class, and denote by S the set of such representatives. We shall show that $S \notin \mathcal{B}$.

We can partition E into the union of its cosets, so $E = \bigcup_{q \in Q}(S + q)$ is a disjoint[1] union.

Suppose S is a Borel set. Then $S + q$ is also a Borel set[2]. Therefore by translation invariance of λ and by countably additivity,

$$\lambda([0, 1)) = 1 = \lambda\left(\bigcup_{q \in Q}(S + q)\right) = \sum_{q \in Q} \lambda(S + q) = \sum_{q \in Q} \lambda(S)$$

But no value for $\lambda(S) \in [0, \infty]$ can be assigned to make this equation hold. Therefore S is not a Borel set.

Remark 12. We can extend this proof to show that $S \notin \mathcal{M}_{leb}$.

One can further show that λ cannot be extended to all subsets $\mathcal{P}(E)$.

[1] Suppose $s_1 + q_1 = s_2 + q_2$ then $s_1 - s_2 = q_1 - q_2 \in \mathbb{Q}$ but then $s_1, s_2 \in S$ by definition $\lightning$.
[2] Consider $\mathcal{G} = \{B \in \mathcal{B} : B + x \in \mathcal{B}\}$ we can show this is a σ-algebra, see page 12.

Theorem 1.4 (Banach - Kuratowski)
Assuming the continuum hypothesis, there exists no measure μ on the set $\mathcal{P}([0,1))$ such that $\mu([0,1)) = 1$ and $\mu(\{x\}) = 0$ for $x \in [0,1)$.

Henceforth, whenever we are on a metric space E, we will work with $\mathcal{B}(E)$, which will be perfectly satisfactory.

§1.7 Probability spaces

Definition 1.17
If a measure space $(E, \mathcal{E}, \mu)$ has $\mu(E) = 1$, we call it a **probability space**, and instead write $(\Omega, \mathcal{F}, \mathbb{P})$. We call Ω the outcome space or sample space, $\mathcal{F}$ the set of events, and $\mathbb{P}$ the probability measure.

The axioms of probability theory (Kolmogorov, 1933), are

1. $\mathbb{P}(\Omega) = 1, \mathbb{P}(\varnothing) = 0$;
2. $0 \leq \mathbb{P}(E) \leq 1$ for all $E \in \mathcal{F}$;
3. if A_n are a disjoint sequence of events in $\mathcal{F}$, then $\mathbb{P}\left(\bigcup_n A_n\right) = \sum_n \mathbb{P}(A_n)$.

This is exactly what is required by our definition: $\mathbb{P}$ is a measure on a σ-algebra.

Remark 13.

- $\mathbb{P}\left(\bigcup_n A_n\right) \leq \sum_n \mathbb{P}(A_n)$ for all sequences $A_n \in \mathcal{F}$;
- $A_n \uparrow A \implies \mathbb{P}(A_n) \uparrow \mathbb{P}(A)$;
- $A_n \downarrow A \implies \mathbb{P}(A_n) \downarrow \mathbb{P}(A)$ as $\mathbb{P}$ a finite measure.

This definition is what separates probability from analysis.

Definition 1.18 (Independent)
Events $(A_i, i \in I)$, $A_i \in \mathcal{F}$ are **independent** if for all finite $J \subseteq I$, we have

$$\mathbb{P}\left(\bigcap_{j \in J} A_j\right) = \prod_{j \in J} \mathbb{P}(A_j).$$

σ-algebras $(\mathcal{A}_i, i \in I)$, $A_i \subseteq \mathcal{F}$ are **independent** if for any $A_j \in \mathcal{A}_j$, where $J \subseteq I$ is finite, the A_j are independent.

Kolmogorov showed that these definitions are sufficient to derive the law of large numbers.

Proposition 1.3

Let $\mathcal{A}_1, \mathcal{A}_2$ be π-systems of sets in $\mathcal{F}$. Suppose $\mathbb{P}(A_1 \cap A_2) = \mathbb{P}(A_1)\mathbb{P}(A_2)$ for all $A_1 \in \mathcal{A}_1, A_2 \in \mathcal{A}_2$. Then the σ-algebras $\sigma(\mathcal{A}_1), \sigma(\mathcal{A}_2)$ are independent.

Proof. Fix $A_1 \in \mathcal{A}_1$, and define for all $A \in \sigma(\mathcal{A}_2)$.

$$\mu(A) = \mathbb{P}(A_1 \cap A), \nu(A) = \mathbb{P}(A_1)(A).$$

Then μ, ν are finite measures and they agree on the π-system $\mathcal{A}_2$. Hence by Uniqueness of Extension, $\mu(A) = \nu(A)\ \forall\ A \in \sigma(\mathcal{A}_2)$, i.e. $\mathbb{P}(A_1 \cap A) = \mathbb{P}(A_1)\mathbb{P}(A)\ \forall\ A_1 \in \mathcal{A}_1, A_2 \in \sigma(\mathcal{A}_2)$.

Now repeat same argument, but now by fixing $A_2 \in \sigma(\mathcal{A}_2)$ define for all $A \in \sigma(\mathcal{A}_1)$

$$\mu'(A) = \mathbb{P}(A \cap A_2), \nu'(A) = \mathbb{P}(A)(A_2).$$

Then μ', ν' are finite measures and they agree on the π-system $\mathcal{A}_1$. Hence by Uniqueness of Extension, $\mu'(A) = \nu'(A)\ \forall\ A \in \sigma(\mathcal{A}_1)$, i.e. $\mathbb{P}(A_1 \cap A) = \mathbb{P}(A_1)\mathbb{P}(A)\ \forall\ A_1 \in \sigma(\mathcal{A}_1), A_2 \in \sigma(\mathcal{A}_2)$.

□

This follows by uniqueness.

§1.8 Borel–Cantelli lemmas

Definition 1.19

Let $A_n \in \mathcal{F}$ be a sequence of events. Then the **limit superior** of A_n is

$$\limsup_n A_n = \bigcap_n \bigcup_{m \geq n} A_m = \{A_n \text{ infinitely often}\}^a$$

The **limit inferior** of A_n is

$$\liminf_n A_n = \bigcup_n \bigcap_{m \geq n} A_m = \{A_n \text{ eventually}\}^b$$

[a] Consider ω, if $\omega \in \limsup_n A_n$ then $\forall\ n, \omega \in \bigcup_{m \geq n} A_m$ thus ω must be in an infinite number of A_ns.

[b] ω is in all but finitely many A_n.

Lemma 1.2 (First Borel–Cantelli lemma)
Let $A_n \in \mathcal{F}$ be a sequence of events such that $\sum_n \mathbb{P}(A_n) < \infty$. Then $\mathbb{P}(A_n \text{ infinitely often}) = 0$.

Proof. For all n, we have

$$\mathbb{P}\left(\limsup_n A_n\right) = \mathbb{P}\left(\bigcap_n \bigcup_{m \geq n} A_m\right) \leq \mathbb{P}\left(\bigcup_{m \geq n} A_m\right) \leq^{a} \sum_{m \geq n} \mathbb{P}(A_m) \to 0$$

□

[a] By countable subadditivity

This proof did not require that $\mathbb{P}$ be a probability measure, just that it is a measure. Therefore, we can use this for arbitrary measures.

Lemma 1.3 (Second Borel–Cantelli lemma)
Let $A_n \in \mathcal{F}$ be a sequence of independent events with $\sum_n \mathbb{P}(A_n) = \infty$. Then $\mathbb{P}(A_n \text{ infinitely often}) = 1$.

Proof. By independence, for all $N \geq n \in \mathbb{N}$ and using $1 - a \leq e^{-a}$, we find

$$\mathbb{P}\left(\bigcap_{m=n}^{N} A_m^c\right) = \prod_{m=n}^{N} (1 - \mathbb{P}(A_m)) \leq \prod_{m=n}^{N} e^{-\mathbb{P}(A_m)} = e^{-\sum_{m=n}^{N} \mathbb{P}(A_m)}$$

As $N \to \infty$, this approaches zero.
Since $\bigcap_{m=n}^{N} A_m^c$ decreases to $\bigcap_{m=n}^{\infty} A_m^c$, $\mathbb{P}(\bigcap_{m=n}^{\infty} A_m^c) = 0$ as $\mathbb{P}(\bigcap_{m=n}^{\infty} A_m^c) \leq \mathbb{P}\left(\bigcap_{m=n}^{N} A_m^c\right) \leq e^{-\sum_{m=n}^{N} \mathbb{P}(A_m)} \to 0$. So by taking complements $\mathbb{P}(\bigcup_{m=n}^{\infty} A_n) = 1 \; \forall \, n (\dagger)$.

Let $B_n = \bigcup_{m=n}^{\infty} A_m$, B_n decreasing and so $B_n \downarrow \bigcap_n B_n = \bigcap_n \bigcup_{m \geq n} A_m = \{A_n \text{ i.o}\}$[a]. As $\mathbb{P}(B_n) = 1$ by $(\dagger)$, $\mathbb{P}\{A_n \text{ i.o}\} = \lim_{n \to \infty} \mathbb{P}(B_n) = 1$ as probabilities are a finite measure. □

[a] A_n occurs infinitely often

Remark 14. If A_n independent, then $\{A_n \text{ i.o}\}$ has either probability 0 or 1 and is called a "tail event". Kolmogorov 0-1 law shows this is true for all "tail events".

§2 Measurable Functions

§2.1 Definition

Definition 2.1 (Measurable)

Let $(E, \mathcal{E}), (G, \mathcal{G})$ be measurable spaces. A function $f\colon E \to G$ is called **measurable** if $f^{-1}(A) \in \mathcal{E}\ \forall\ A \in \mathcal{G}$, where $f^{-1}(A)$ is the preimage of A under f i.e. $f^{-1}(A) = \{x \in E : f(x) \in A\}$.

If $G = \mathbb{R}$ and $\mathcal{G} = \mathcal{B}$, we can just say that $f\colon (E, \mathcal{E}) \to G$ is measurable. Moreover, if E is a topological space and $\mathcal{E} = \mathcal{B}(E)$, we say f is Borel measurable.

Note that preimages f^{-1} commute with many set operations such as intersection, union, and complement. This implies that $\{f^{-1}(A) : A \in \mathcal{G}\}$ is a σ-algebra over E, and likewise, $\{A : f^{-1}(A) \in \mathcal{E}\}$ is a σ-algebra over G. Hence, if $\mathcal{A}$ is a collection of subsets s.t. $G \supset \sigma(\mathcal{A})$ then if $f^{-1}(A) \in \mathcal{E}$ for all $A \in \mathcal{A}$, the class $\{A : f^{-1} \in \mathcal{E}\}$ is a σ-algebra that contains $\mathcal{A}$ and so $\sigma(\mathcal{A})$. So f is measurable.

If $f\colon (E, \mathcal{E}) \to \mathbb{R}$, the collection $\mathcal{A} = \{(-\infty, y]\colon y \in \mathbb{R}\}$ generates $\mathcal{B}$ (Sheet 1). Hence f is Borel measurable iff $f^{-1}((-\infty, y]) = \{x \in E : f(x) \leq y\} \in \mathcal{E}$ for all $y \in \mathbb{R}$.

If E is a topological space and $\mathcal{E} = \mathcal{B}(E)$, then if $f\colon E \to \mathbb{R}$ is continuous, the preimages of open sets B are open, and hence Borel sets. The open sets in $\mathbb{R}$ generate the σ-algebra $\mathcal{B}$. Hence, continuous functions to the real line are measurable.

Example 2.1

Consider the indicator function 1_A of a set $A \subset E$. $1_A^{-1}(1) = A$ and $1_A^{-1}(0) = A^c$ hence measurable iff $A \in \mathcal{E}$.

Example 2.2

The composition of measurable functions is measurable. Note that given a collection of maps $\{f_i\colon E \to (G, \mathcal{G}) : i \in I\}$, we can make them all measurable by taking $\mathcal{E}$ to be a large enough σ-algebra, for instance $\sigma\left(\left\{f_i^{-1}(A) : A \in \mathcal{G}, i \in I\right\}\right)$ called the σ-algebra generated by $\{f_i\}_{i \in I}$.

Proposition 2.1

If $f_1, f_2, \ldots$ are measurable $\mathbb{R}$-valued. Then $f_1 + f_2$, $f_1 f_2$, $\inf_n f_n$, $\sup_n f_n$, $\liminf f_n$, $\limsup f_n$ are all measurable.

Proof. See Sheet 1. □

§2.2 Monotone Class Theorem

Theorem 2.1 (Monotone Class Theorem)
Let $(E, \mathcal{E})$ be a measurable space and $\mathcal{A}$ be a π-system that generates the σ-algebra $\mathcal{E}$. Let $\mathcal{V}$ be a vector space of bounded maps from E to $\mathbb{R}$ s.t.

1. $1_E \in \mathcal{V}$;
2. $1_A \in \mathcal{V}$ for all $A \in \mathcal{A}$;
3. if f is bounded and $f_n \in \mathcal{V}$ are nonnegative functions that form an increasing sequence that converge pointwise to f on E, then $f \in \mathcal{V}$.

Then $\mathcal{V}$ contains all bounded measurable functions $f\colon E \to \mathbb{R}$.

Proof. Define $\mathcal{D} = \{A \in \mathcal{E} : 1_A \in \mathcal{V}\}$. Then $\mathcal{D}$ is a d-system as $1_E \in \mathcal{V}$ and for $A \subseteq B$, $1_{B \setminus A} = 1_B - 1_A \in \mathcal{V}$ as $\mathcal{V}$ a vector space so $B \setminus A \in \mathcal{D}$.
If $A_n \in \mathcal{D}$ increases to A, we have 1_{A_n} increases pointwise to 1_A, which lies in $\mathcal{V}$ by the (3.) so $A \in \mathcal{D}$.

$\mathcal{D}$ contains $\mathcal{A}$ by (2.), as well as E itself. So by Dynkin's lemma $\mathcal{D}$ contains $\sigma(\mathcal{A}) = \mathcal{E}$ so $\mathcal{E} = \mathcal{D}$ i.e. $1_A \in V \; \forall \; A \in \mathcal{E}$.

Since V a vector space it contains all finite linear combinations of indicators of measurable sets. Let $f\colon E \to \mathbb{R}$ be a bounded measurable function, which we will assume at first is nonnegative. We define

$$\begin{aligned} f_n(x) &= 2^{-n} \lfloor 2^n f(x) \rfloor \\ &= 2^{-n} \sum_{j=0}^{\infty} 1_{A_{n,j}}(x) \\ A_{n,j} &= \{2^n f(x) \in [j, j+1)\} \\ &= f^{-1}\left(\left[\frac{j}{2^n}, \frac{j+1}{2^n}\right)\right) \in \mathcal{E}. \end{aligned}$$

As f is bounded we do not need an infinite sum but only a finite one. Then $f_n \leq f \leq f_n + 2^{-n}$. Hence $|f_n - f| \leq 2^{-n} \to 0$ and $f_n \uparrow f$.

So $0 \leq f_n \uparrow f, f_n \in \mathcal{V}$ and f is bounded non-negative so $f \in \mathcal{V}$ by (3.).

Finally, for any f bounded and measurable, $f = f^{+}$[a] $- f^{-}$[b]. f^+, f^- are bounded, nonnegative and measurable, so in $\mathcal{V}$ and $\mathcal{V}$ a vector space thus $f \in \mathcal{V}$. □

[a] $\max(f, 0)$
[b] $\max(-f, 0)$

§2.3 Image measures

Definition 2.2 (Image Measure)
Let $f\colon (E,\mathcal{E}) \to (G,\mathcal{G})$ be a measurable function and μ a measure on $(E,\mathcal{E})$. Then the **image measure** $\nu = \mu \circ f^{-1}$ is obtained from assigning $\nu(A) = \mu(f^{-1}(A))$ for all $A \in \mathcal{G}$.

Remark 15. This is well defined as $f^{-1}(A) \in \mathcal{E}$ as f measurable. ν is countably additive because the preimage satisfies set operations and μ countably additive (See Sheet 1).

Starting from the Lebesgue measure, we can get all probability measures (in fact we can get all Radon measures) in this way.

Definition 2.3 (Right-Continuous)
A function f is **right-continuous** if $x_n \downarrow x \implies f(x_n) \to f(x)$.

Lemma 2.1
Let $g\colon \mathbb{R} \to \mathbb{R}$ be a non-constant, increasing, right-continuous function, and set $g(\pm\infty) = \lim_{z\to\pm\infty} g(z)$. On $I = (g(-\infty), g(+\infty))$ we define the **generalised inverse** $f : I \to \mathbb{R}$ by

$$f(x) = \inf\{y \in \mathbb{R} : g(y) \geq x\}.$$

Then f is increasing, left-continuous, and $f(x) \leq y$ iff $x \leq g(y)$ for all $x \in I, y \in \mathbb{R}$.

Remark 16. f and g form a Galois connection.

Proof. Fix $x \in I$.
Let $J_x = \{y \in \mathbb{R} : g(y) \geq x\}$. Since $x > g(-\infty)$, J_x is nonempty and bounded below. Hence $f(x)$ is a well-defined real number.
If $y \in J_x$, then $y' \geq y$ implies $y' \in J_x$ since g is increasing. Since g is right-continuous, if $y_n \downarrow y$, and all $y_n \in J_x$, then $g(y) = \lim_n g(y_n) \geq x$ so $y \in J_x$.
So $J_x = [f(x), \infty)$. Hence $f(x) \leq y \iff x \leq g(y)$ as required.

If $x \leq x'$, we have $J_x \supseteq J_{x'}$ (as $y \in J_x \impliedby y \in J'_x$), i.e. $[f(x),\infty) \supseteq [f(x'),\infty)$ so $f(x) \leq f(x')$.
Similarly, if $x_n \uparrow x$, we have $J_x = \bigcap_n J_{x_n}$[a] so $[f(x),\infty) = \bigcap_n [f(x_n),\infty)$ so $f(x_n) \to f(x)$ as $x_n \to x$. □

[a] As $y \in \bigcap_n J_{x_n} \iff g(y) \geq x_n \;\forall\, n \iff g(y) \geq x \iff y \in J_x$.

Theorem 2.2

Let $g\colon \mathbb{R} \to \mathbb{R}$ as in the previous lemma. Then $\exists$ a unique Radon measure μ_g on $\mathbb{R}$ such that $\mu_g((a,b]) = g(b) - g(a)$ for all $a < b$. Further, all Radon measures on $\mathbb{R}$ can be obtained in this way.

Proof. Define I, f as in the previous lemma and λ the Lebesgue measure on I.

f is Borel measurable since $f^{-1}((-\infty, z]) = \{x \in I\colon f(x) \leq z\} = \{x \in I\colon x \leq g(z)\} = (-g(\infty), g(z)] \in \mathcal{B}$. As $\{(-\infty, z] : z \in \mathbb{R}\}$ generate $\mathcal{B}$, f measurable.

Therefore, the image measure $\mu_g = \lambda \circ f^{-1}$ exists on $\mathcal{B}$. Then for any $-\infty < a < b < \infty$, we have

$$\begin{aligned}
\mu_g((a,b]) &= \lambda\left(f^{-1}\left((a,b]\right)\right) \\
&= \lambda\left(\{x\colon a < f(x) \leq f(b)\}\right) \\
&= \lambda\left(\{x\colon g(a) < x \leq g(b)\}\right) \\
&= g(b) - g(a)
\end{aligned}$$

By the Uniqueness of Extension for σ-finite measures, μ_g is uniquely defined.

Conversely, let ν be a Radon measure on $\mathbb{R}$. Define $g : \mathbb{R} \to \mathbb{R}$ as

$$g(y) = \begin{cases} \nu((0,y]) & \text{if } y \geq 0 \\ -\nu((y,0]) & \text{if } y < 0 \end{cases}$$

ν Radon tells us that g is finite. Easy to check g is right-continuous[a]. This is an increasing function in y, since ν is a measure. Finally, $\nu((a,b]) = g(b) - g(a)$ which can be seen by case analysis and additivity of the measure ν. By uniqueness as before, this characterises ν in its entirety. $\square$

[a] For $y_n \downarrow y$ where $y \geq 0$, $(0, y_n] \downarrow (0, y]$ and then $\nu((0, y_n]) \downarrow \nu((0, y])$ by countably additivity. Similarly for $y < 0$.

Remark 17. Such image measures μ_g are called **Lebesgue–Stieltjes measures** associated with g, where g is the **Stieltjes distribution**.

Example 2.3

Fix $x \in \mathbb{R}$ and take $g = 1_{[x,\infty)}$. Then $\mu_g = \delta_x$ the *dirac measure at* x defined for all $A \in \mathcal{B}$ by

$$\delta_x(A) = \begin{cases} 1 & \text{if } x \in A \\ 0 & \text{otherwise} \end{cases}$$

§2.4 Random variables

Definition 2.4 (Random Variable)

Let $(\Omega, \mathcal{F}, \mathbb{P})$ be a probability space, and $(E, \mathcal{E})$ be a measurable space. If $X : \Omega \to E$ a measurable function then X is a **random variable** in E.

When $E = \mathbb{R}$ or $\mathbb{R}^d$ with the Borel σ-algebra, we simply call X a random variable or random vector.

Example 2.4

X models a "random" outcome of an experiment, e.g. when tossing a coin $\Omega = \{H, T\}, X = \#\text{ heads} : \Omega \to \{0, 1\}$.

Definition 2.5 (Distribution)

The **law** or **distribution** μ_X of a random variable X is given by the image measure $\mu_X = \mathbb{P} \circ X^{-1}$. It is a measure on $(E, \mathcal{E})$.

When $(E, \mathcal{E}) = (\mathbb{R}, \mathcal{B})$, μ_X is uniquely determined by its values on any π-system, we shall take $\{(-\infty, x] : x \in \mathbb{R}\}$ and

$$F_X(z) = \mu_X((-\infty, z]) = \mathbb{P}(X^{-1}(-\infty, z]) = \mathbb{P}\left(\{\omega \in \Omega : X(\omega) \leq z\}\right) = \mathbb{P}\left(X \leq z\right)$$

The function F_x is called the **distribution function** of X, because it uniquely determines the distribution of X.

Using the properties of measures, we can show that any distribution function satisfies:

1. F_X is increasing;
2. F_X is right-continuous[3];
3. $F_X(-\infty) = \lim_{z \to -\infty} F_X(z) = \mu_X(\varnothing) = 0$;
4. $F_X(\infty) = \lim_{z \to \infty} F_X(z) = \mu_X(\mathbb{R}) = \mathbb{P}(\Omega) = 1$.

Proposition 2.2

Given any function F satisfying the previous properties, $\exists$ a random variable X s.t. $F = F_X$.

Proof. Let $\Omega = (0, 1)$, $\mathcal{F} = \mathcal{B}(0, 1)$, $\mathbb{P}$ the Lebesgue measure $\lambda|_{(0,1)}$.
Let F be any function satisfying the properties, then F is increasing and right con-

[3] $x_n \downarrow x \implies (-\infty, x_n] \downarrow (-\infty, x]$ hence by countable additivity of $\mathbb{P} \circ X^{-1}$.

tinuous so we can define the generalised inverse

$$X(\omega) = \inf\{x : \omega \leq F(x)\} : (0,1) \to \mathbb{R}$$

Hence X is a measurable function and thus a random variable.

$$\begin{aligned} F_X(x) &= \mathbb{P}(X \leq x) = \mathbb{P}(\{\omega \in \Omega : X(\omega) \leq x\}) = \mathbb{P}(\{\omega \in \Omega : \omega \leq F(x)\}) \\ &= \mathbb{P}(\{\omega \in (0,1) : \omega \leq F(x)\}) \\ &= \mathbb{P}((0, F(x)]) \\ &= F(x) - 0 \end{aligned}$$

□

Remark 18. This is similar to what we saw in IB Probability, if we have F then r.v. $F^{-1}(U)$ where $U \sim U(0,1)$ has the distribution function F, where F^{-1} is the generalised inverse.

Definition 2.6 (Independent)
Consider a countable collection $(X_i \colon (\Omega, \mathcal{F}, \mathbb{P}) \to (E, \mathcal{E}))$ for $i \in I$. This collection of random variables is called **independent** if the σ-algebras $\sigma(X_i)$ are independent, recall $\sigma(X_i)$ is generated by $\left\{X_i^{-1}(A) \colon A \in \mathcal{E}\right\}$, the smallest σ-algebra s.t. X_i measurable.

For $(E, \mathcal{E}) = (\mathbb{R}, \mathcal{B})$ we show on an Sheet 1 that this is equivalent to the condition

$$\mathbb{P}(X_1 \leq x_1, \dots, X_n \leq x_n) = \mathbb{P}(X_1 \leq x_1) \dots \mathbb{P}(X_n \leq x_n)$$

for all finite subsets $\{X_1, \dots, X_n\}$ of the X_i.

§2.5 Constructing independent random variables

Question
Given a distribution function F, we know $\exists$ a r.v. X corresponding to it. But given an infinite sequence of distribution functions $F_1, F_2, \dots$ does $\exists$ independent r.v. $(X_1, X_2, \dots)$ corresponding to them?

Let $(\Omega, \mathcal{F}, \mathbb{P}) = ((0,1), \mathcal{B}(0,1), \lambda|_{(0,1)})$. We start with Bernoulli random variables.

Any $\omega \in (0,1)$ has a binary representation given by $(\omega_i) \in \{0,1\}^{\mathbb{N}}$ where $\omega = \sum_{i=1}^{\infty} 2^{-i}\omega_i$, which is unique if we exclude infinitely long tails of zeroes from the binary representation (same reasoning as $1.00000\dots = 0.99999\dots$).

Definition 2.7 (nth Rademacher function)

The **nth Rademacher function** $R_n : \Omega \to 0, 1$ is given by $R_n(\omega) = \omega_n$, it extracts the nth bit from the binary expansion.

Observe that $R_1 = 1_{(1/2,1]}$, $R_2 = 1_{(1/4,1/2]} + 1_{(3/4,1]}$ and so on. Since each R_n can be given as the sum of finite (2^{n-1}) indicator functions on measurable sets, they are measurable functions and are hence random variables.

Claim 2.1

R_i are iid $\mathrm{Ber}(\frac{1}{2})$.

Proof. $\mathbb{P}(R_n = 1) = \frac{1}{2} = \mathbb{P}(R_n = 0)$ can be checked by induction.

We now show they are independent. For a finite set $(x_i)_{i=1}^n$, by considering the size of the intervals that ω can lie in,

$$\mathbb{P}(R_1 = x_1, \ldots, R_n = x_n) = 2^{-n} = \mathbb{P}(R_1 = x_1) \ldots \mathbb{P}(R_n = x_n)$$

□

Therefore, the R_n are all independent, so countable sequences of independent random variables indeed exist.
The next step is to construct a sequence of iid r.v.s on $\mathrm{U}(0, 1)$.

Now, take a bijection $m \colon \mathbb{N}^2 \to \mathbb{N}$ and define $Y_{k,n} = R_{m(k,n)}$, the Rademacher functions. We now define $Y_n = \sum_{k=1}^{\infty} 2^{-k} Y_{k,n}$[4].

Claim 2.2

Y_n are iid $\mathrm{U}(0, 1)$, i.e. $\mu_{Y_n} = \lambda|_{(0,1)}$ and Y_n independent.

Lemma 2.2

Any measurable functions of independent random variables are independent.

Proof. They are independent because the Y_i are measurable functions of independent random variables, e.g. Y_1 is a measurable function of $Y_{1,1}, Y_{2,1}, \ldots$; Y_2 of $Y_{1,2}, Y_{2,2}, \ldots$

The π-system of intervals $\left(\frac{i}{2^m}, \frac{i+1}{2^m}\right]$ for $i = 0, \ldots, 2^m - 1$ for $m \in \mathbb{N}$ generates $\mathcal{B}(0, 1)$

[4] This converges for all $\omega \in \Omega$ since $|Y_{k,n}| \leq 1$.

as $\mathbb{Q}$ dense in $\mathbb{R}$. So by theorem 1.2 the distribution of Y_n is identified on the intervals.

$$\begin{aligned}
\mathbb{P}\left(Y_n \in \left(\frac{i}{2^m}, \frac{i+1}{2^m}\right]\right) &= \mathbb{P}\left(\frac{i}{2^m} < \sum_{k=1}^{\infty} 2^{-k} Y_{k,n} \leq \frac{i+1}{2^n}\right)^{a} \\
&= \mathbb{P}(Y_{1,n} = y_1, \ldots, Y_{m,n} = y_m) \text{ where } \frac{i}{2^m} = 0.y_1 y_2 \ldots y_m \\
&= \prod_{i=1}^{m} \mathbb{P}(Y_{m,n} = y_m) \text{ by independence.} \\
&= 2^{-m} = \lambda\left(\frac{i}{2^m}, \frac{i+1}{2^m}\right]
\end{aligned}$$

Hence $\mu_{Y_n} = \lambda|_{(0,1)}$ on the π-system and so on $\mathcal{B}(0,1)$. □

[a] This specifies the first m digits in the binary expansion of Y_n.

As before, set $G_n(x) = F_n^{-1}(x)$ which is the generalised inverse. Then G_n are Borel functions, set $X_n = G_n(Y_n)$ for $n \in \mathbb{N}$, then as before $F_{X_n} = F_n$ and X_n are independent as Y_n are.

§2.6 Convergence of measurable functions

Let $(E, \mathcal{E}, \mu)$ be a measure space. Let $A \in \mathcal{E}$ be defined by some property.

Definition 2.8 (Almost everywhere)
We say that a property defining a set $A \in \mathcal{E}$ holds **μ-almost everywhere** if $\mu(A^c) = 0$.

Definition 2.9 (Almost surely)
If μ is a $\mathbb{P}$- measure, we say a property holds **$\mathbb{P}$-almost surely** or **with probability one**, if $\mathbb{P}(A^c) = 0$, i.e. if $\mathbb{P}(A) = 1$.

Definition 2.10 (Convergence almost everywhere)
If f_n and f are measurable functions on $(E, \mathcal{E}, \mu) \to (\mathbb{R}, \mathcal{B})$, we say **$f_n$ converges to f μ-almost everywhere** if $\mu(\{x \in E : f_n(x) \nrightarrow f(x)\}) = 0$.

For r.v.s, we say $X_n \to X$ **$\mathbb{P}$-almost surely** if $\mathbb{P}(\{\omega \in \Omega : X_n(\omega) \to X(\omega)\}) = 1$.

Definition 2.11 (Convergence in Measure)

We say f_n **converges to f in μ-measure** if for all $\varepsilon > 0$

$$\mu(\{x \in E : |f_n(x) - f(x)| > \varepsilon\}) \to 0,$$

as $n \to \infty$.

We say $X_n \to X$ in **$\mathbb{P}$-probability** if $\forall\, \varepsilon > 0$

$$\mathbb{P}(|X_n - X| > \varepsilon) \to 0$$

as $n \to \infty$.

Theorem 2.3

Let $f_n \colon (E, \mathcal{E}, \mu) \to \mathbb{R}$ be measurable functions.

1. If $\mu(E) < \infty$, then $f_n \to 0$ a.e. $\implies$ $f_n \to 0$ in measure;
2. If $f_n \to 0$ in measure, $\exists$ subsequence n_k s.t. $f_{n_k} \to 0$ a.e.

Example 2.5

Let $f_n = 1_{(n,\infty)}$ and the Lebesgue measure, then $f_n \to 0$ a.e. but $\mu(|f_n| > \varepsilon) = \infty\ \forall\, n$.

Proof. Fix $\varepsilon > 0$. Suppose $f_n \to 0$ a.e., then for every n,

$$\mu(E) \geq \mu(|f_n| \leq \varepsilon) \geq \mu\left(\bigcap_{m \geq n} \{|f_m| \leq \varepsilon\}\right)$$

Let $A_n = \bigcap_{m \geq n} \{|f_m| \leq \varepsilon\}$ which is increasing to $\bigcup_n \bigcap_{m \geq n} \{|f_m| \leq \varepsilon\}$. So by the countable additivity of μ,

$$\begin{aligned}
\mu\left(\bigcap_{m \geq n} \{|f_m| \leq \varepsilon\}\right) &\to \mu\left(\bigcup_n \bigcap_{m \geq n} \{|f_m| \leq \varepsilon\}\right) \\
&= \mu(|f_n| \leq \varepsilon \text{ eventually}) \\
&\geq \mu(|f_n| \to 0) \\
&= \mu(E) \text{ as } f_n \to 0 \text{ a.e. and } \mu \text{ finite.}
\end{aligned}$$

Hence,

$$\liminf_{n \to \infty} \mu(|f_n| \leq \varepsilon) = \mu(E) \implies \limsup_{n \to \infty} \mu(|f_n| > \varepsilon) \leq 0 \implies \mu(|f_n| > \varepsilon) \to 0$$

$\square$

Proof. Suppose $f_n \to 0$ in measure, choosing $\varepsilon = \frac{1}{k}$ we have

$$\mu\left(|f_n| > \frac{1}{k}\right) \to 0.$$

So we can choose n_k s.t. $\mu\left(|f_n| > \frac{1}{k}\right) \leq \frac{1}{k^2}$. We can choose n_{k+1} in the same way s.t. $n_{k+1} > n_k$. So we get a subsequence n_k s.t. $\mu\left(|f_{n_k}| > \frac{1}{k}\right) < \frac{1}{k^2}$. Also $\sum_k \frac{1}{k^2} < \infty$, so $\sum_k \mu\left(|f_{n_k}| > \frac{1}{k}\right) < \infty$. So by the first Borel–Cantelli lemma, we have

$$\mu\left(\underbrace{|f_{n_k}| > \frac{1}{k} \text{ infinitely often}}_{f_{n_k} \not\to 0}\right) = 0$$

so $f_{n_k} \to 0$ a.e. □

Remark 19. The first statement is false if $\mu(E)$ is infinite: consider $f_n = 1_{(n,\infty)}$ on $(\mathbb{R}, \mathcal{B}, \mu)$, since $f_n \to 0$ almost everywhere but $\mu(f_n) = \infty$.

The second statement is false if we do not restrict to subsequences: consider independent events A_n such that $\mathbb{P}(A_n) = \frac{1}{n}$, then $1_{A_n} \to 0$ in probability since $\mathbb{P}(1_{A_n} > \varepsilon) = \mathbb{P}(A_n) = \frac{1}{n} \to 0$, but $\sum_n \mathbb{P}(A_n) = \infty$, and by the second Borel–Cantelli lemma, $\mathbb{P}(1_{A_n} > \varepsilon \text{ infinitely often}) = 1$, so $1_{A_n} \not\to 0$ almost surely.

Definition 2.12 (Convergence in Distribution)

For X and X_n a sequence of r.v.s, we say $X_n \xrightarrow{d} X$[a] if $F_{X_n}(t) \to F_X(t)$ as $n \to \infty$ for all $t \in \mathbb{R}$ which are continuity points of F_X.

[a] X_n converges to X in distribution

Remark 20. This definition does not require X_n to be defined on the same probability space.

Remark 21. If $X_n \to X$ in probability, then $X_n \xrightarrow{d} X$, see Sheet 2 for proof.

Example 2.6

Let $(X_n)_{n\in\mathbb{N}}$ be iid $\mathrm{Exp}(1)$, i.e. $\mathbb{P}(X_n > x) = e^{-x}$ for $x \geq 0$.

Question

Find a deterministic fcn $g : \mathbb{N} \to \mathbb{R}$ s.t. a.s. $\limsup \frac{X_n}{g(n)} = 1$.

Define $A_n = \{X_n \geq \alpha \log n\}$ where $\alpha > 0$, so $\mathbb{P}(A_n) = n^{-\alpha}$, and in particular, $\sum_n \mathbb{P}(A_n) < \infty$ if and only if $\alpha > 1$. By the Borel–Cantelli lemmas, we have for all $\varepsilon > 0$,

$$\mathbb{P}\left(\frac{X_n}{\log n} \geq 1 \text{ infinitely often}\right) = 1; \quad \mathbb{P}\left(\frac{X_n}{\log n} \geq 1 + \varepsilon \text{ infinitely often}\right) = 0$$

In other words, $\mathbb{P}(\limsup_n \frac{X_n}{\log n} = 1) = 1$.

§2.7 Kolmogorov's zero-one law

Let $(X_n)_{n\in\mathbb{N}}$ be a sequence of r.v.s. We can define $\mathcal{T}_n = \sigma(X_{n+1}, X_{n+2}, \dots)$[5]. Let $\mathcal{T} = \bigcap_{n\in\mathbb{N}} \mathcal{T}_n$ be the **tail σ-algebra**, which contains all events in $\mathcal{F}$ that depend only on the 'limiting behaviour' of (X_n).

Theorem 2.4 (Kolmogorov 0-1 Law)
Let $(X_n)_{n\in\mathbb{N}}$ be a sequence of independent r.v.s. Let $A \in \mathcal{T}$ be an event in the tail σ-algebra. Then $\mathbb{P}(A) = 1$ or $\mathbb{P}(A) = 0$.
If $Y : (\Omega, \mathcal{T}) \to (\mathbb{R}, \mathcal{B})$ is measurable, it is constant almost surely.

Proof. Let $\mathcal{F}_n = \sigma(X_1, \dots, X_n)$. Then $\mathcal{F}_n$ is generated by the π-system of sets $A = (X_1 \leq x_1, \dots, X_n \leq x_n)$ for any $x_i \in \mathbb{R}$.
Note that the π-system of sets $B = (X_{n+1} \leq x_{n+1}, \dots, X_{n+k} \leq x_{n+k})$, for arbitrary $k \in \mathbb{N}$ and $x_i \in \mathbb{R}$, generates $\mathcal{T}_n$.
By independence of the sequence, we see that $\mathbb{P}(A \cap B) = \mathbb{P}(A)\mathbb{P}(B)$ for all such sets A, B, and so the σ-algebras $\mathcal{T}_n, \mathcal{F}_n$ generated by these π-systems are independent. As $\mathcal{T} \subseteq \mathcal{T}_n$, $\mathcal{F}_n$ and $\mathcal{T}$ are independent $\forall\, n$.

Let $\mathcal{F}_\infty = \sigma(X_1, X_2, \dots)$. Then, $\bigcup_n \mathcal{F}_n$ is a π-system that generates $\mathcal{F}_\infty$. As $\mathcal{F}_n$ and $\mathcal{T}$ are independent $\forall\, n$, $\bigcup_n \mathcal{F}_n$ independent of $\mathcal{T}$. So $\mathcal{F}_\infty, \mathcal{T}$ are independent.

Since $\mathcal{T} \subseteq \mathcal{F}_\infty$, if $A \in \mathcal{T}$, A is independent from $A \in \mathcal{F}_\infty$. So $\mathbb{P}(A) = \mathbb{P}(A \cap A) = \mathbb{P}(A)\mathbb{P}(A)$, so $\mathbb{P}(A)^2 - \mathbb{P}(A) = 0$ as required.

Finally, if $Y : (\Omega, \mathcal{T}) \to (\mathbb{R}, \mathcal{B})$ measurable, the preimages of $\{Y \leq y\}$ lie in $\mathcal{T}$, which give probability one or zero. Let $c = \inf\{y : F_Y(y) = 1\}$, so $Y = c$ almost surely. □

Remark 22. This tells us that for X_i iid with finite expectation, $\liminf_{n\to\infty} \frac{1}{n}\sum_{i=1}^n X_i$, $\limsup_{n\to\infty} \frac{1}{n}\sum_{i=1}^n X_i$ are constants a.s.

[5] The smallest σ-algebra s.t. $X_{n+1}, \dots$ are measurable.

§3 Integration

§3.1 Notation

Let $f\colon (E, \mathcal{E}, \mu) \to \mathbb{R}$ be measurable and $f \geq 0$[6].

Notation. We will then define the integral with respect to μ, either written $\mu(f)$ or $\int_E f \,\mathrm{d}\mu = \int_E f(x) \,\mathrm{d}\mu(x)$.

When $(E, \mathcal{E}, \mu) = (\mathbb{R}, \mathcal{B}, \lambda)$, we write it as $\int f(x)dx$.

Notation. If X is a random variable, we will define its expectation $\mathbb{E}[X] = \int_\Omega X \,\mathrm{d}\mathbb{P} = \int_\Omega X(\omega) \,\mathrm{d}\mathbb{P}(\omega)$.

§3.2 Definition

Definition 3.1 (Simple)
We say that a function $f\colon (E, \mathcal{E}, \mu) \to \mathbb{R}$ is **simple** if it is of the form

$$f = \sum_{k=1}^{m} a_k 1_{A_k}; \quad a_k \geq 0; \quad A_k \in \mathcal{E}; \quad m \in \mathbb{N}$$

Definition 3.2 (μ-integral)
The **μ-integral** of a simple function f defined as above is

$$\mu(f) = \sum_{k=1}^{m} a_k \mu(A_k)^a$$

which is independent of the choice of representation of the simple function, i.e. well-defined.

[a] Note we take $0 \cdot \infty = 0$.

Remark 23.

- We have $\mu(\alpha f + \beta g) = \alpha\mu(f) + \beta\mu(g)$ for all nonnegative coefficients α, β and simple functions f, g.
- If $g \leq f$, $\mu(g) \leq \mu(f)$, so μ is increasing.
- $f = 0$ a.e. $\iff \mu(f) = 0$.

[6] f is measurable when mapped to $\mathbb{R}$ and $f \geq 0$, this is different from saying f non-negative, measurable.

Definition 3.3 (μ-integral)
For a general non-negative function $f\colon (E, \mathcal{E}, \mu) \to \mathbb{R}$, we define its **$\mu$-integral** to be

$$\mu(f) = \sup\{\mu(g) : g \leq f, g \text{ simple}\}$$

which agrees with the above definition for simple functions.

Clearly if $0 \leq f_1 \leq f_2$ then $\mu(f_1) \leq \mu(f_2)$.

Now, for $f\colon (E, \mathcal{E}, \mu) \to \mathbb{R}$ measurable but not necessarily non-negative, we define $f^+ = \max(f, 0)$ and $f^- = \max(-f, 0)$, so that $f = f^+ - f^-$ and $|f| = f^+ + f^-$.

Definition 3.4 (μ-integrable)
A measurable function $f\colon (E, \mathcal{E}, \mu) \to \mathbb{R}$ is **μ-integrable** if $\mu(|f|) < \infty$. In this case, we define its integral to be

$$\mu(f) = \mu(f^+) - \mu(f^-)$$

which is a well-defined real number.

Later we shall prove that $\mu(|f|) = \mu(f^+) + \mu(f^-)$ hence $|\mu(f)| \leq \mu(|f|)$.

If one of $\mu(f^+)$ or $\mu(f^-)$ is ∞ and the other finite, we define $\mu(f)$ to be ∞ or $-\infty$ respectively (though f is not integrable).

§3.3 Monotone Convergence Theorem

Notation.

- We say $x_n \uparrow x$ to mean $x_n \leq x_{n+1}\ \forall\, n$ and $x_n \to x$.
- We say $f_n \uparrow f$ to mean $f_n(x) \leq f_{n+1}(x)\ \forall\, n$ and $f_n(x) \to f$.

Theorem 3.1 (Monotone Convergence Theorem)
Let $f_n, f\colon (E, \mathcal{E}, \mu) \to \mathbb{R}$ be measurable and non-negative s.t. $f_n \uparrow f$. Then, $\mu(f_n) \uparrow \mu(f)$.

Remark 24. This is a theorem that allows us to interchange a pair of limits, $\mu(f) = \mu(\lim_n f_n) = \lim_n \mu(f_n)$, i.e. $\lim_n \int f_n \,\mathrm{d}\mu = \int \lim_n f_n \,\mathrm{d}\mu$ for $f_n \geq 0$ and $f_n \uparrow f$.
If $g_n \geq 0$, letting $f_n = \sum_{k=1}^n g_k$ and $f_n \uparrow f = \sum_{k=1}^\infty g_k$ we get $\lim_n \int \sum_{k=1}^n g_k \,\mathrm{d}\mu = \int \sum_{k=1}^\infty g_k \,\mathrm{d}\mu \implies \sum_{k=\infty}^n \int g_k \,\mathrm{d}\mu = \int \sum_k g_k \,\mathrm{d}\mu$ or equivalently $\mu(\sum_k g_k) = \sum_k \mu(g_k)$.
This generalises the countable additivity of μ to integrals of non-negative functions.

If we consider the approximating sequence $\widetilde{f}_n = 2^{-n}\lfloor 2^n f \rfloor$, as defined in the monotone class theorem, then this is a non-negative sequence converging to f. So in particular, $\mu(f)$ is equal to the limit of the integrals of these simple functions.

It suffices to require convergence of $f_n \to f$ a.e., the general argument does not need to change. The non-negativity constraint is not required if the first term in the sequence f_0 is integrable, by subtracting f_0 from every term.

Proof. Recall that $\mu(f) = \sup\{\mu(g) : g \leq f, g \text{ simple}\}$. Let $M = \sup_n \mu(f_n)$, then $\mu(f_n) \uparrow M$.

We now show $M = \mu(f)$.

Since $f_n \leq f$, $\mu(f_n) \leq \mu(f)$, so taking suprema, $M \leq \mu(f)$.

Now, we need to show $\mu(f) \leq M$, or equivalently, $\mu(g) \leq M$ for all simple g s.t. $g \leq f$, so by taking suprema, $\mu(f) = \sup_g \mu(g) \leq M$.
Now let $g = \sum_{k=1}^m a_k 1_{A_k}$ where $a_k \geq 0$ and wlog the $A_k \in \mathcal{E}$ are disjoint. We define $g_n = \min(\overline{f}_n, g)$, where $\overline{f}_n$ is the nth approximation of f_n by simple functions as in the Monotone Class Theorem. So g_n is simple, $g_n \leq \overline{f}_n \leq f_n \uparrow f$, so $g_n \uparrow \min(f, g) = g$. I.e. $g \uparrow g$ and g_n simple with $g_n \leq f$.

Fix $\varepsilon \in (0, 1)$, and define sets $A_k(n) = \{x \in A_k : g_n(x) \geq (1-\varepsilon)a_k\}$. Since $g = a_k$ on A_k, and since $g_n \uparrow g$, $A_k(n) \uparrow A_k$ for all k. Since μ is a measure, $\mu(A_k(n)) \uparrow \mu(A_k)$ by countable additivity.

Also, we have $g_n 1_{A_k} \geq g_n 1_{A_k(n)} \geq (1-\varepsilon)a_k 1_{A_k(n)}$ as $A_k(n) \subseteq A_k$. So as $\mu(f)$ is increasing, we have $\mu(g_n 1_{A_k}) \geq \mu\Big((1-\varepsilon)a_k 1_{A_k(n)}\Big)$ and so $\mu(g_n 1_{A_k}) \geq (1-\varepsilon)a_k \mu(1_{A_k(n)})$ as they are simple functions.

Finally, $g_n = \sum_{k=1}^n g_n 1_{A_k}$ as $g_n \leq g$ and g supported on $\bigcup_{k=1}^n A_k$ and A_k disjoint. So So as $g_n 1_{A_k}$ is simple,

$$\begin{aligned}
\mu(g_n) &= \mu\left(\sum_{k=1}^n g_n 1_{A_k}\right) \\
&= \sum_{k=1}^n \mu(g_n 1_{A_k}) \\
&\geq \sum_{k=1}^n (1-\varepsilon)a_k \mu(A_k(n)) \\
&\uparrow \sum_{k=1}^n (1-\varepsilon)a_k \mu(A_k) \\
&= (1-\varepsilon)\mu(g).
\end{aligned}$$

Then,

$$(1-\varepsilon)\mu(g) \leq \lim_n \mu(g_n) \leq^{a} \lim_n \mu(f_n) \leq M$$

so $\mu(g) \leq \frac{M}{1-\varepsilon}$ $\forall\ \varepsilon \in (0,1)$ hence $\mu(g) \leq M$. Since ε was arbitrary, this completes the proof. □

[a] As $g_n \leq f_n$

§3.4 Linearity of Integral

Theorem 3.2 (Linearity of Integral)
Let $f, g\colon (E, \mathcal{E}, \mu) \to \mathbb{R}$ be nonnegative measurable functions. Then $\forall\ \alpha, \beta \geq 0$,

- $\mu(\alpha f + \beta g) = \alpha\mu(f) + \beta\mu(g)$;
- $f \leq g \implies \mu(f) \leq \mu(g)$;
- $f = 0$ a.e. $\iff \mu(f) = 0$.

Proof. If $\widetilde{f}_n, \widetilde{g}_n$ are the approximations of f and g by simple functions from the Monotone Class Theorem let $f_n = \min(\widetilde{f}_n, n)$ and $g_n = \min(\widetilde{g}_n, n)$. Then f_n, g_n are simple and $f_n \uparrow f$ and $g_n \uparrow g$. Then $\alpha f_n + \beta g_n \uparrow \alpha f + \beta g$, so by MCT[a], $\mu(f_n) \uparrow \mu(f)$, $\mu(g_n) \uparrow \mu(g)$ and $\mu(\alpha f_n + \beta g_n) \uparrow \mu(\alpha f + \beta g)$. As f_n, g_n simple $\mu(\alpha f_n + \beta g_n) = \alpha\mu(f_n) + \beta\mu(g_n) \uparrow \alpha\mu(f) + \beta\mu(g)$. So $\alpha\mu(f) + \beta\mu(g) = \mu(\alpha f + \beta g)$.

The second part is obvious from definition.

If $f = 0$ a.e, then $0 \leq f_n \leq f$, so $f_n = 0$ a.e. but f_n simple $\implies \mu(f_n) = 0$. As $\mu(f_n) \uparrow \mu(f)$ so $\mu(f) = 0$.
Conversely, if $\mu(f) = 0$, then $0 \leq \mu(f_n) \uparrow \mu(f)$ so $\mu(f_n) = 0\ \forall\ n \implies f_n = 0$ a.e.
But $f_n \uparrow f \implies f = 0$ a.e. □

[a] Monotone Convergence Theorem

Remark 25. Functions such as $1_{\mathbb{Q}}$ are integrable and have integral zero. They are 'identified' with the zero element in the theory of integration.

Theorem 3.3 (Linearity of Integral)
Let $f, g\colon (E, \mathcal{E}, \mu) \to \mathbb{R}$ be integrable. Then $\forall\ \alpha, \beta \in \mathbb{R}$,

- $\mu(\alpha f + \beta g) = \alpha\mu(f) + \beta\mu(g)$;
- $f \leq g \implies \mu(f) \leq \mu(g)$;

- $f = 0$ a.e. $\implies \mu(f) = 0$.

Proof. Left as an exercise, just use $f = f^+ - f^-$ and use definitions and $\mu(f) = \mu(f^+) - \mu(f^-)$ etc. $\square$

§3.5 Fatou's lemma

Example 3.1

Let $f_n = 1_{(n,n+1)}$, $f_n \geq 0$ with $f_n \to 0$ as $n \to \infty$. $\lambda(f_n) = 1$ but $\lambda(0) = 0$.

Lemma 3.1 (Fatou's lemma)

Let $f_n \colon (E, \mathcal{E}, \mu) \to \mathbb{R}$ be measurable, non-negative functions. Then $\mu(\liminf_n f_n) \leq \liminf_n \mu(f_n)$.

Remark 26. Recall that $\liminf_n x_n = \sup_n \inf_{m \geq n} x_m$ and $\limsup_n x_n = \inf_n \sup_{m \geq n} x_m$. In particular, $\limsup_n x_n = \liminf_n x_n$ implies that $\lim_n x_n$ exists and is equal to $\limsup_n x_n$ and $\liminf_n x_n$. Hence, if the f_n converge to some measurable function f, we must have $\mu(f) \leq \liminf_n \mu(f_n)$.

Proof. We have $\inf_{m \geq n} f_m \leq f_k$ for all $k \geq n$, so by taking integrals, $\mu(\inf_{m \geq n} f_m) \leq \mu(f_k)$. Thus,

$$\mu\left(\inf_{m \geq n} f_m\right) \leq \inf_{k \geq n} \mu(f_k) \leq \sup_n \inf_{k \geq n} \mu(f_k) = \liminf \mu(f_k) \tag{$\dagger$}$$

Note that $\inf_{m \geq n} f_m$ increases to $\sup_n \inf_{m \geq n} f_m = \liminf_n f_n$.

Let $g_n = \inf_{m \geq n} f_n$, then $g_n \geq 0$ and $g_n \uparrow \sup_n g_n = \sup_n \inf_{m \geq n} f_m = \liminf_n f_n$. By MCT $\mu(g_n) \uparrow \mu(\liminf_n f_n)$ so by taking limits in ($\dagger$), $\mu(\liminf_n f_n) \leq \liminf \mu(f_n)$. $\square$

§3.6 Dominated Convergence Theorem

Theorem 3.4 (Dominated Convergence Theorem)

Let $f_n, f \colon (E, \mathcal{E}, \mu)$ be measurable functions s.t. $|f_n| \leq g$ a.e., for some integrable fcn g, so $\mu(g) < \infty$[a], and $f_n \to f$ pointwise (or a.e.) on E.
Then f_n and f are also integrable, and $\mu(f_n) \to \mu(f)$.

[a] Note $g \geq |f_n| \geq 0$.

Proof. Clearly $\mu(|f_n|) \leq \mu(g) < \infty$, so the f_n are integrable. Taking limits in $|f_n| \leq g$, we have $|f| \leq g$, so f is also integrable by the same argument and as the limit of measurable fcns is measurable.

Now, $g \pm f_n \geq 0$, and converges pointwise to $g \pm f$. Since limits are equal to the limit inferior when they exist, by Fatou's lemma, we have

$$\mu(g)+\mu(f) = \mu(g+f) = \mu\Big(\liminf_n (g+f_n)\Big) \leq \liminf_n \mu(g+f_n) = \mu(g)+\liminf_n \mu(f_n)$$

Hence $\mu(f) \leq \liminf_n \mu(f_n)$ as $\mu(g)$ finite. Likewise, $\mu(g) - \mu(f) \leq \mu(g) - \limsup_n \mu(f_n)$, so $\mu(f) \geq \limsup_n \mu(f_n)$, so

$$\limsup_n \mu(f_n) \leq \mu(f) \leq \liminf_n \mu(f_n)$$

But since $\liminf_n \mu(f_n) \leq \limsup_n \mu(f_n)$, the result follows. □

Remark 27. In fact, $\mu(|f_n - f|) \to 0$ as $|f_n - f| \leq |f_n| + |f| \leq g + g = 2g$ and $2g$ is integrable so by DCT (Dominated Convergence Theorem) proved.

If $X_n \to X$ $\mathbb{P}$ a.s., and $|X_n| \leq Y$ and $\mathbb{E}[Y] < \infty$ then $\mathbb{E}[X_n] \to \mathbb{E}[X]$ and $\mathbb{E}[|X_n - X|] \to 0$. In particular, if $|X_n| \leq M$ $\forall$ n, for some $M > 0, M \in \mathbb{R}$ then $\mathbb{E}[|X_n - X|] \to 0$ (Bounded Convergence Theorem)[7].

Remark 28. DCT also holds for convergence in $\mathbb{P}$-prob, where if $X_n \to X$ in $\mathbb{P}$ probability then we get $\mathbb{E}[X_n] \to \mathbb{E}[X]$ and $\mathbb{E}[|X_n - X|] \to 0$.

Proof. Suppose $\mathbb{E}[|X_n - X|] \not\to 0$. Then $\exists$ a subsequence n_k s.t. $\mathbb{E}[|X_{n_k} - X|] > \varepsilon$ $\forall$ k for some $\varepsilon > 0$. Now $X_n \to X$ in $\mathbb{P}$ prob then $X_{n_k} \to X$ in $\mathbb{P}$ prob by definition. By theorem 2.3, $\exists$ n_{k_l} s.t. $X_{n_{k_l}} \to X$ a.s. But then by DCT, $\mathbb{E}\Big[\Big|X_{n_{k_l}} - X\Big|\Big] \to 0$ ↯. □

Theorem 3.5 (Bounded Convergence Theorem)

If $X_n \to X$ in $\mathbb{P}$ prob and $|X_n| \leq M$ for some constant $M > 0$, $\forall$ $n \geq 0$. Then $\mathbb{E}[|X_n - X|] \to 0$.

This is quite useful in probability.

Example 3.2

Let $E = [0, 1]$ with the Lebesgue measure. Let $f_n \to f$ pointwise and the f_n are uniformly bounded, so $\sup_n \|f_n\|_\infty \leq g$ for some $g \in \mathbb{R}$. Then since $\mu(g) = g < \infty$, the DCT implies that f_n, f are integrable and $\mu(f_n) \to \mu(f)$ as $n \to \infty$.

[7] This works as with finite measure then $\mathbb{E}[M]$ is finite

In particular, no notion of uniform convergence of the f_n is required as in Riemann Integrals.

Remark 29. FTC states that

1. Let $f : [a, b] \to \mathbb{R}$ be continuous and set $F(t) = \int_a^t f(x)\,\mathrm{d}x$[8]. Then F is differentiable on $[a, b]$ with $F' = f$.
2. Let $F : [a, b] \to \mathbb{R}$ be differentiable and F' is continuous, then $\int_a^b F'(x)\,\mathrm{d}x = F(b) - F(a)$.

The proof of the fundamental theorem of calculus requires only the fact that

$$\int_x^{x+h} \mathrm{d}t = h$$

This is a fact which is obviously true of the Riemann integral and also of the Lebesgue integral.

Therefore, for any continuous function $f\colon [0, 1] \to \mathbb{R}$, we have

$$\underbrace{\int_0^x f(t)\,\mathrm{d}t}_{\text{Riemann integral}} = F(x) = \underbrace{\int_0^x f(t)\,\mathrm{d}\mu(t)}_{\text{Lebesgue integral}}$$

So these integrals coincide for continuous functions.

Remark 30. We can generalise the FTC for Lebesgue integrals:
If $f : [a, b] \to \mathbb{R}$ is Lebesgue integrable and $F(t) = \int_a^t f(x)\,\mathrm{d}x$. Then,

$$\lim_{h\to 0} \frac{F(t+h) - F(t)}{h} = \lim_{h\to 0} \frac{\int_t^{t+h} f(x)\,\mathrm{d}x}{h} = f(t) \text{ a.e.}$$

This is the Lebesgue differentiation theorem, studied in Analysis of Functions.

Remark 31. We can show that all Riemann integrable functions are $\mu^\star$-measurable, where $\mu^\star$ is the outer measure of the Lebesgue measure, as defined in the proof of Carathéodory's theorem. However, there exist certain Riemann integrable functions that are not Borel measurable. We can modify such an f on a Lebesgue measure 0 set to make it Borel measurable, i.e. $\exists\ \widetilde{f}$ s.t. $\widetilde{f} = f$ on A and $\lambda(A^c) = 0$ and $\int \widetilde{f}\,\mathrm{d}x = \int f\,\mathrm{d}x$.

A (bounded) Riemann integrable fcn $f : [a, b] \to \mathbb{R}$ is Lebesgue integrable in the following sense. If f is bounded on $[a, b]$, f is R-integrable iff

$$\lambda(\mu(\{x \in [0, 1] : f \text{ is discontinuous at } x\})) = 0,$$

i.e. f is continuous a.e.

The standard techniques of Riemann integration, such as substitution and integration by parts, extend to all bounded measurable functions by the monotone class theorem.

[8]This is a Lebesgue integral

Example 3.3

$1_{\mathbb{Q}}$ on $[0,1]$ is a bounded function on a bounded interval. The set of discontinuity points is $[0,1]$ which is not measure 0, thus not Riemann Integrable. But this is Lebesgue integrable and $1_{\mathbb{Q}} = 0$ λ a.s., so $\lambda(1_{\mathbb{Q}}) = 0$.

Theorem 3.6 (Substitution Formula)

Let $\varphi : [a,b] \to \mathbb{R}$, φ strictly increasing and continuously differentiable. Then $\forall\ g$ Borel fcns, $g \geq 0$ on $[\varphi(a), \varphi(b)]$, $\int_{\varphi(a)}^{\varphi(b)} g(y)\,\mathrm{d}y = \int_a^b g(\varphi(x))\varphi'(x)\,\mathrm{d}x\ (\star)$.

Proof. Let $\mathcal{V}$ be the set of all measurable fcns g for which $(\star)$ holds. Then by linearity of integral, $\mathcal{V}$ is a vector space.

- $1 \in \mathcal{V}$ by FTC (2), $1_{(c,d]} \in \mathcal{V}$ by FTC (2).
- If $f_n \in \mathcal{V}$, $f_n \uparrow \mathcal{V}$, $f_n \geq 0$ then by Monotone Convergence Theorem.

Hence by Monotone Class Theorem, $(\star)$ holds $\forall\ g \geq 0$ measurable. □

Theorem 3.7 (Differentiation Under The Integral Sign)

Let $U \subseteq \mathbb{R}$ be an open set and $(E, \mathcal{E}, \mu)$ be a measure space. Let $f\colon U \times E \to \mathbb{R}$ be s.t.

1. $x \mapsto f(t,x)$ is integrable $\forall\ t$;
2. $t \mapsto f(t,x)$ is differentiable $\forall\ x \in E$;
3. $\exists\ g : E \to \mathbb{R}$ integrable s.t. $\left|\frac{\partial f}{\partial t}\right| < g(x)\ \forall\ t \in U, x \in E$;

Then $x \mapsto \frac{\partial f}{\partial t}$ is integrable $\forall\ t$ and,

$$F(t) = \int_E f(t,x)\,\mathrm{d}\mu(x) \implies F'(t) = \int_E \frac{\partial f}{\partial t}(t,x)\,\mathrm{d}\mu(x)$$

Proof. Fix t. By the mean value theorem,

$$g_h(x) = \frac{f(t+h,x) - f(t,x)}{h} - \frac{\partial f}{\partial t}(t,x) \implies |g_h(x)| = \left|\frac{\partial f}{\partial t}\left(\tilde{t}, x\right) - \frac{\partial f}{\partial t}(t,x)\right| \leq 2g(x)$$

Note that g is integrable. By differentiability of f, we have $g_h \to 0$ as $h \to 0$, so by

DCT, $\mu(g_h) \to \mu(0) = 0$. By linearity of the integral,

$$\mu(g_h) = \frac{\int_E f(t+h,x) - f(t,x)\,\mathrm{d}\mu(x)}{h} - \int_E \frac{\partial f}{\partial t}(t,x)\,\mathrm{d}\mu(x)$$

Hence, $\frac{F(t+h)-F(t)}{h} - F'(t) \to 0$. □

Example 3.4 (Integrals and Image Measures.)
For a measurable function $f\colon (E,\mathcal{E},\mu) \to (G,\mathcal{G})$ with image measure $\nu = \mu \circ f^{-1}$ on $(G,\mathcal{G})$. If $g\colon G \to \mathbb{R}$ is a measurable, non-negative function then,

$$\mu \circ f^{-1}(g) = \nu(g) = \int_G g\,\mathrm{d}\nu = \int_G g\,\mathrm{d}\mu \circ f^{-1} = \int_E g(f(x))\,\mathrm{d}\mu(x) = \mu(g \circ f)$$

Proof on Sheet 2, use monotone class theorem and first prove for g indicator fcns and then simple functions.

In particular, for $X : (\Omega,\mathcal{F},\mathbb{P}) \to \mathbb{R}$ and $X \geq 0$ measurable, we have,

$$\mathbb{E}\left[g(X)\right] = \int_\Omega g(X(\omega))\,\mathrm{d}\mathbb{P}(\omega) = \int g\,\mathrm{d}\mu_X\,,$$

where $\mu_X = \mathbb{P} \circ X^{-1}$ is the distribution of X.

Example 3.5 (Densities of Measures)
If $f\colon (E,\mathcal{E},\mu) \to \mathbb{R}$ is a measurable non-negative function, we can define $\nu(A) = \mu(f1_A)$ for any measurable set A, which is again a measure on $(E,\mathcal{E})$ by the Monotone Convergence Theorem. For A_n disjoint,

$$\begin{aligned}
\nu\left(\bigcup_{i=1}^\infty A_i\right) &= \mu(f1_{\cup A_i}) \\
&= \mu(f\sum_{i=1}^\infty 1_{A_i}) \\
&= \mu(\sum_{i=1}^\infty f1_{A_i}) \\
&= \sum_{i=1}^\infty \mu(f1_{A_i}) \text{ by MCT} \\
&= \sum_{i=1}^\infty \nu(A_i).
\end{aligned}$$

In particular, if $g\colon (E, \mathcal{E}) \to \mathbb{R}$ is measurable, $\nu(g) = \mu(fg) = \int_E g(x)f(x)\,\mathrm{d}\mu(x) = \int_E g\,\mathrm{d}\nu(f)$. This follows by definition for g indicator functions, by additivity extends to simple functions and by Monotone Convergence Theorem to all measurable non-negative functions.

We call f the **density** of ν with respect to μ. This is unique as $\mu(f1_A) = \mu(g1_A)\ \forall\ A \in \mathcal{E} \implies f = g\ \mu$ a.e. (proved on Sheet 2).

In particular, for $\mu = \lambda$, $\forall\ f$ Borel $\exists$ a Borel measure ν on $\mathbb{R}$ given by $\nu(A) = \int_A f(x)\,\mathrm{d}x$ and then $\forall\ g$ Borel, $g \geq 0$ $\nu(g) = \int f(x)g(x)\,\mathrm{d}x$. We say ν has **density** f. This ν is a prob measure on $(\mathbb{R}, \mathcal{B})$ iff $\int f(x)\,\mathrm{d}x = 1$.

For $\lambda : (\Omega, \mathcal{F}, \mathbb{P}) \to \mathbb{R}$, if the law $\mu_X = \mathbb{P} \circ X^{-1}$ has the density f_X (wrt λ), we call f_X the **probability density function** of X. Then $\mathbb{P}(X \in A) = \mathbb{P} \circ X^{-1}(A) = \mu_X(A) = \int_A f_X(x)\,\mathrm{d}x\ \ \forall\ A \in \mathcal{B}$ and $\forall\ g$ Borel, $g \geq 0$. Taking $A = (-\infty, x]$, we get $\mathbb{P}(X \leq x) = F_X(x) = \int_{-\infty}^{x} f_X(x)\,\mathrm{d}x$.
$\mathbb{E}[g(x)] = \int g(x)\,\mathrm{d}\mu_X(x)$ from previous example and $\int g(x)\,\mathrm{d}\mu_X(x) = \nu(g) = \int g(x)f_X(x)\,\mathrm{d}x$.

§4 Product Measures

§4.1 Integration in product spaces

Let $(E_1, \mathcal{E}_1, \mu_1), (E_2, \mathcal{E}_2, \mu_2)$ be finite measure spaces. On $E = E_1 \times E_2$, we can consider the π-system of 'rectangles' $\mathcal{A} = \{A_1 \times A_2 : A_1 \in \mathcal{E}_1, A_2 \in \mathcal{E}_2\}$. Then we define the σ-algebra $\mathcal{E} = \mathcal{E}_1 \otimes \mathcal{E}_2 = \sigma(\mathcal{A})$ on the product space.

If the E_i are topological spaces with a countable basis, then $\mathcal{B}(E_1 \times E_2) = \mathcal{B}(E_1) \otimes \mathcal{B}(E_2)$ where we take the product topology.

Lemma 4.1

Let $f\colon (E, \mathcal{E}) \to \mathbb{R}$ be measurable. Then $\forall\, x_1 \in E_1$, the fcn $(x_2 \mapsto f(x_1, x_2))\colon (E_2, \mathcal{E}_2) \to \mathbb{R}$ is $\mathcal{E}_2$-measurable.

Proof. Let

$$\mathcal{V} = \{f\colon (E, \mathcal{E}) \to \mathbb{R} : f \text{ bounded, measurable, conclusion of the lemma holds}\}$$

This is a $\mathbb{R}$-vector space, and $1_E, 1_A \in \mathcal{V}\ \forall\ A = A_1 \times A_2 \in \mathcal{A}$, since $1_A(x_1, x_2) = 1_{A_1}(x_1) 1_{A_2}(x_2)$ thus fixing x_1 gives 0 or 1_{A_2}.

Now, let $0 \leq f_n$ increase to f, $f_n \in \mathcal{V}$. Then $(x_2 \mapsto f(x_1, x_2)) = \lim_n (x_2 \mapsto f_n(x_1, x_2))$, so it is $\mathcal{E}_2$-measurable as it's a limit of a sequence of measurable functions. Then by the Monotone Class Theorem, $\mathcal{V}$ contains all bounded measurable functions. This extends to all measurable functions by truncating the absolute value of f to $n \in \mathbb{N}$, then the sequence of such bounded truncations converges pointwise to f. □

Lemma 4.2

Let $f\colon (E, \mathcal{E}) \to \mathbb{R}$ be measurable s.t.

1. f is bounded; or
2. f is nonnegative.

Then the map $x_1 \mapsto \int_{E_2} f(x_1, x_2)\, \mathrm{d}\mu_2(x_2)$ is μ_1-measurable and is bounded[a] or nonnegative respectively.

[a] As μ_2 is a finite measure.

Remark 32. In case (ii), the map on x_1 may evaluate to infinity, but the set of values

$$\left\{x_1 \in E_1 : \int_{E_2} f(x_1, x_2)\, \mathrm{d}\mu_2(x_2) = \infty\right\}$$

lies in $\mathcal{E}_1$.

Generally, a fcn f taking values in $[0, \infty]$ is measurable means $f^{-1}(\{\infty\}) \in \mathcal{E}_1$ and $f^{-1}(A) \in \mathcal{E}_1 \ \forall \ A \in \mathcal{B}$.

Proof. Let

$$\mathcal{V} = \{f \colon (E, \mathcal{E}) \to \mathbb{R} : f \text{ bounded, measurable, conclusion of the lemma holds}\}$$

This is a vector space by linearity of the integral. $1_E \in \mathcal{V}$, since $\int_{E_2} 1_E(x_1, x_2) \, \mathrm{d}\mu_2(x_2) = 1_{E_1}\mu_2(E_2)$ is non-negative and bounded. $1_A \in \mathcal{V} \ \forall \ A \in \mathcal{A}$, because $1_{A_1}(x_1)\mu_2(A_2)$ is $\mathcal{E}_1$-measurable, non-negative, and bounded since it is at most $\mu_2(E_2) < \infty$.

Now let f_n be a sequence of non-negative functions that increase to f, where $f_n \in \mathcal{V}$. Then by the Monotone Convergence Theorem,

$$\int_{E_2} \lim_{n\to\infty} f_n(x_1, x_2) \, \mathrm{d}\mu_2(x_2) = \lim_{n\to\infty} \int_{E_2} f_n(x_1, x_2) \, \mathrm{d}\mu_2(x_2)$$

is an increasing limit of $\mathcal{E}_1$-measurable functions, so is $\mathcal{E}_1$-measurable. It is bounded by $\mu_2(E_2)\|f\|_\infty$, or non-negative as required. So $f \in \mathcal{V}$. By the Monotone Class Theorem, the result for bounded functions holds.

□

Theorem 4.1 (Product Measure)
There $\exists$ a unique measure $\mu = \mu_1 \otimes \mu_2$ on $(E, \mathcal{E})$ such that $\mu(A_1 \times A_2) = \mu_1(A_1)\mu_2(A_2)$ for all $A_1 \in \mathcal{E}_1$, $A_2 \in \mathcal{E}_2$.

Proof. $\mathcal{A}$ is a π-system generating $\mathcal{E}$ and μ a finite measure, so by the Uniqueness of Extension, μ unique.

We define for $A \in \mathcal{E}$,

$$\mu(A) = \int_{E_1} \left(\int_{E_2} 1_A(x_1, x_2) \, \mathrm{d}\mu_2(x_2) \right) \mathrm{d}\mu_1(x_1) \, .$$

This is well-defined by the two previous lemmas.

We have

$$\begin{aligned} \mu(A_1 \times A_2) &= \int_{E_1} \left(\int_{E_2} 1_{A_1}(x_1) 1_{A_2}(x_2) \, \mathrm{d}\mu_2(x_2) \right) \mathrm{d}\mu_1(x_1) \\ &= \int_{E_1} 1_{A_1}(x_1)\mu_2(A_2) \, \mathrm{d}\mu_1(x_1) \\ &= \mu_1(A_1)\mu_2(A_2) \end{aligned}$$

Clearly $\mu(\varnothing) = 0$, so it suffices to show countable additivity. Let A_n be disjoint sets in $\mathcal{E}$. Then

$$1_{(\bigcup_n A_n)} = \sum_n 1_{A_n} = \lim_{n\to\infty} \sum_{i=1}^n 1_{A_n}$$

Then by the Monotone Convergence Theorem and the previous lemmas,

$$\begin{aligned}
\mu\left(\bigcup_n A_n\right) &= \int_{E_1} \left(\int_{E_2} \lim_{n\to\infty} \sum_{i=1}^n 1_{A_i} \,\mathrm{d}\mu_2(x_2)\right) \mathrm{d}\mu_1(x_1) \\
&= \int_{E_1} \left(\lim_{n\to\infty} \int_{E_2} \sum_{i=1}^n 1_{A_i} \,\mathrm{d}\mu_2(x_2)\right) \mathrm{d}\mu_1(x_1) \\
&= \lim_{n\to\infty} \int_{E_1} \left(\int_{E_2} \sum_{i=1}^n 1_{A_i} \,\mathrm{d}\mu_2(x_2)\right) \mathrm{d}\mu_1(x_1) \\
&= \lim_{n\to\infty} \sum_{i=1}^n \int_{E_1} \left(\int_{E_2} 1_{A_i} \,\mathrm{d}\mu_2(x_2)\right) \mathrm{d}\mu_1(x_1) \\
&= \lim_{n\to\infty} \sum_{i=1}^n \mu(A_i) \\
&= \sum_{n=1}^\infty \mu(A_n)
\end{aligned}$$

□

Remark 33. Note $\mu(A) = \int_{E_2} \left(\int_{E_1} 1_A(x_1, x_2) \,\mathrm{d}\mu_1(x_1)\right) \mathrm{d}\mu_2(x_2)$ by just swapping the order of integration in the previous lemmas and proofs and then by Dynkin's Lemma/π-λ/π-d theorem.

§4.2 Fubini's theorem

Theorem 4.2 (Fubini-Tonelli)
Let $(E, \mathcal{E}, \mu) = (E_1 \times E_2, \mathcal{E}_1 \otimes \mathcal{E}_2, \mu_1 \otimes \mu_2)$ be a finite measure space.

1. Let $f\colon E \to \mathbb{R}$ be measurable, non-negative. Then

$$\begin{aligned}
\mu(f) &= \int_E f \,\mathrm{d}\mu \\
&= \int_{E_1} \left(\int_{E_2} f(x_1, x_2) \,\mathrm{d}\mu_2(x_2)\right) \mathrm{d}\mu_1(x_1) \\
&= \int_{E_2} \left(\int_{E_1} f(x_1, x_2) \,\mathrm{d}\mu_1(x_1)\right) \mathrm{d}\mu_2(x_2)
\end{aligned}$$

2. Let $f\colon E \to \mathbb{R}$ be a μ-integrable function (on the product measure). Let

$$A_1 = \left\{ x_1 \in E_1 : \int_{E_2} |f(x_1, x_2)| \, \mathrm{d}\mu_2(x_2) < \infty \right\}.$$

Define $f_1 : E_1 \to \mathbb{R}$ by $f_1(x_1) = \int_{E_2} f(x_1, x_2) \, \mathrm{d}\mu_2(x_2)$ on A_1 and 0 elsewhere. Then $\mu_1(A_1^c) = 0$, f_1 is μ_1-integrable and $\mu(f) = \mu_1(f_1) = \mu_1(f_1 1_{A_1})$, and defining A_2 symmetrically, $\mu(f) = \mu_2(f_2) = \mu_2(f_2 1_{A_2})$.

Remark 34. If f is bounded, $A_1 = E_1$. Note, for $f(x_1, x_2) = \frac{x_1^2 - x_2^2}{(x_1^2 + x_2^2)^2}$ on $(0,1)^2$, we have $\mu_1(f_1) \neq \mu_2(f_2)$, but f is not Lebesgue integrable on $(0,1)^2$.

Proof. By the definition of the product measure, first statement is true for $f = 1_A$ for $A \in \mathcal{E}$. Then, by linearity of the integral, this extends to simple functions. For general fcn $f \geq 0$ by Monotone Convergence Theorem and the standard approximation by simple fcns $f_n = \min(2^{-n}\lfloor 2^n f \rfloor, n)$, the first statement follows.

Now let f be μ-integrable. Define $h : E_1 \to [0, \infty]$ as $h(x_1) = \int_{E_2} |f(x_1, x_2)| \, \mathrm{d}\mu_2(x_2)$. By Lemma 4.2, h is measurable (as $|f| \geq 0$), is non-negative, so $A_1 \in \mathcal{E}_1$[a].
Then by the first part, $\mu_1(h) \leq \mu(|f|) < \infty$. So f_1 is μ_1-integrable. We have $\mu_1(A_1^c) = 0$, otherwise $\mu_1(h) \geq \mu_1(h 1_{A_1^c}) = \infty$ ↯.

Setting, $f_1^{\pm} = \int_{E_2} f^{\pm}(x_1, x_2) \, \mathrm{d}\mu_2(x_2)$ we see than $f_1 = (f_1^+ - f_1^-) 1_{A_1}$. Also by the first part, $\mu_1(f_1^+) = \mu(f^+) < \infty$ and $\mu_1(f_1^-) = \mu(f^-) < \infty$. Hence, $\mu(f) =$ [b]$\mu(f^+) - \mu(f^-) = \mu_1(f_1^+) - \mu_1(f_1^-) =$ [c]$\mu_1(f_1)$ as required. □

[a] h measurable $\implies h^{-1}(\{\infty\}) \in \mathcal{E}_1$. $A_1 = h^{-1}(\{\infty\})^c$ thus in $\mathcal{E}_1$.
[b] As f integrable
[c] As f_1 integrable due to $\mu_1(A_1^c) = 0$.

Remark 35. The proofs above extend to σ-finite measures μ.

Let $(E_i, \mathcal{E}_i, \mu_i)$ be measure spaces with σ-finite measures. Note that $(\mathcal{E}_1 \otimes \mathcal{E}_2) \otimes \mathcal{E}_3 = \mathcal{E}_1 \otimes (\mathcal{E}_2 \otimes \mathcal{E}_3)$, by a π-system argument using Dynkin's lemma. So we can iterate the construction of the product measure to obtain a measure $\mu_1 \otimes \cdots \otimes \mu_n$[9], which is a unique measure on $(\prod_{i=1}^n E_i, \bigotimes_{i=1}^n \mathcal{E}_i)$ with the property that the measure of a hypercube $\mu(A_1 \times A_n)$ is the product of the measures of its sides $\mu_i(A_i)$.

In particular, we have constructed the Lebesgue measure $\mu^n = \bigotimes_{i=1}^n \mu$ on $\mathbb{R}^n$. Applying Fubini's theorem, for functions f that are either non-negative and measurable or μ^n-integrable, we have

$$\int_{\mathbb{R}^n} f \, \mathrm{d}\mu^n = \int \cdots \int_{\mathbb{R} \ldots \mathbb{R}} f(x_1, \ldots, x_n) \, \mathrm{d}\mu(x_1) \ldots \mathrm{d}\mu(x_n)$$

[9] This is associative.

§4.3 Product probability spaces and independence

Proposition 4.1

Let $X_1, \dots, X_n$ be r.v.s, $X_i : (\Omega, \mathcal{F}, \mathbb{P}) \to (E_i, \mathcal{E}_i)$. Set $(E, \mathcal{E}) = (\prod_{i=1}^n E_i, \bigotimes_{i=1}^n \mathcal{E}_i)$. Consider $X \colon (\Omega, \mathcal{F}, \mathbb{P}) \to (E, \mathcal{E})$ given by $X(\omega) = (X_1(\omega), X_2(\omega), \dots, X_n(\omega))$. Then X is $\mathcal{E}$-measurable and the following are equivalent.

1. $X_1, \dots, X_n$ are independent random variables;
2. $\mu_X = \bigotimes_{i=1}^n \mu_{X_i}$;
3. for all bounded and measurable $f_i \colon E_i \to \mathbb{R}$, $\mathbb{E}\left[\prod_{i=1}^n f_i(X_i)\right] = \prod_{i=1}^n \mathbb{E}\left[f_i(X_i)\right]$.

Proof. To show X measurable suffices to check $X^{-1}(A_1 \times \cdots \times A_n) \in \mathcal{F}$, where $A_i \in \mathcal{E}_i \; \forall \, i$ as this is a π-system generating $\mathcal{E}$.

$$\begin{aligned} X^{-1}(A_1 \times \cdots \times A_n) &= \{\omega : X_1(\omega) \in A_1, \dots, X_n(\omega) \in A_n\} \\ &= \bigcap_{i=1}^n X_i^{-1}(A_i). \end{aligned}$$

X_i measurable so $X_i^{-1}(A_i) \in \mathcal{F}$ and so the intersection is in $\mathcal{F}$.

(1) $\implies$ (2): Consider the π-system $\mathcal{A}$ of rectangles $A = \prod_{i=1}^n A_i$ for $A_i \in \mathcal{E}_i$, as this generates $\mathcal{E}$ suffices to check equality on it.

Since μ_X is an image measure, then

$$\begin{aligned} \mu_X(A_1 \times \cdots \times A_n) = \mathbb{P}\left(X_1 \in A_1, \dots, X_n \in A_n\right) = \mathbb{P}\left(X_1\right) \dots \mathbb{P}\left(A_n\right) &= \prod_{i=1}^n \mu_{X_i}(A_i) \\ &= \left(\bigotimes_{i=1}^n \mu_{X_i}\right)(A). \end{aligned}$$

(2) $\implies$ (3): By Fubini's theorem,

$$\begin{aligned} \mathbb{E}\left[\prod_{i=1}^n f_i(X_i)\right] &= \mu_X\left(\prod_{i=1}^n f_i(x_i)\right) \\ &= \int_E f(x) \, \mathrm{d}\mu_X(x) \\ &= \int \cdots \int_{E_i} \left(\prod_{i=1}^n f_i(x_i)\right) \mathrm{d}\mu_{X_1}(x_1) \dots \mathrm{d}\mu_{X_2}(x_2) \\ &= \prod_{i=1}^n \int_{E_i} f_i(x_i) \, \mathrm{d}\mu_{X_i}(x_i) \end{aligned}$$

$$= \prod_{i=1}^{n} \mathbb{E}\left[f_i(X_i)\right]$$

(3) $\implies$ (1): Let $f_i = 1_{A_i}$ for any $A_i \in \mathcal{E}_i$. These are bounded and measurable functions. Then

$$\mathbb{P}\left(X_1 \in A_1, \ldots, X_n \in A_n\right) = \mathbb{E}\left[\prod_{i=1}^{n} 1_{A_i}(X_i)\right] = \prod_{i=1}^{n} \mathbb{E}\left[1_{A_i}(X_i)\right] = \prod_{i=1}^{n} \mathbb{P}\left(X_i \in A_i\right)$$

So the σ-algebras generated by the X_i are independent as required. $\square$

§5 L^p Spaces, Norms and Inequalities

§5.1 Norms

Definition 5.1 (Norm)
A **norm** on a real vector space is a map $\|\cdot\|_V : V \to [0, \infty)$ s.t.

1. $\|\lambda v\| = |\lambda| \cdot \|v\|$;
2. $\|u + v\| \leq \|u\| + \|v\|$;
3. $\|v\| = 0 \iff v = 0$.

Definition 5.2
Let $(E, \mathcal{E}, \mu)$ be a measure space. We define $L^p(E, \mathcal{E}, \mu) = L^p(\mu) = L^p$ for the space of measurable functions $f : E \to \mathbb{R}$ s.t. $\|f\|_p$ is finite, where

$$\|f\|_p = \begin{cases} \left(\int_E |f(x)|^p \, \mathrm{d}\mu(x)\right)^{\frac{1}{p}} & 1 \leq p < \infty \\ \text{ess sup}\, |f| = \inf\{\lambda \geq 0 : |f| \leq \lambda \ \mu\text{-a.e.}\} & p = \infty \end{cases}$$

We must check that $\|\cdot\|_p$ as defined is a norm.
Clearly (1) holds for all $1 \leq p < \infty$ by linearity of integral and for $p = \infty$ its obvious.
Property (2) holds for $p = 1$ and $p = \infty$, and we will prove later that this holds for other values of p by Minkowski inequality.
The last property does not hold: $f = 0$ implies $\|f\|_p = 0$, but $\|f\|_p = 0$ implies only that $|f|^p = 0$ a.e., so f is zero a.e. on E.

Therefore, to rigorously define the norm, we must construct the quotient space $\mathcal{L}^p$ of functions that coincide a.e.. We write $[f]$ for the equivalence class of functions that are equal a.e. The functional $\|\cdot\|_p$ is then a norm on $\mathcal{L}^p = \{[f] : f \in L^p\}$.

Proposition 5.1 (Chebyshev's inequality, Markov's inequality)
Let $f : E \to \mathbb{R}$ be non-negative and measurable. Then $\forall\, \lambda > 0$,

$$\mu(\{x \in E : f(x) \geq \lambda\}) = \mu(f \geq \lambda) \leq \frac{\mu(f)}{\lambda}$$

Proof. Integrate the inequality $\lambda 1_{\{f \geq \lambda\}} \leq f$, which holds on E. □

In particular if $g \in L^p$, $p < \infty$ and $\lambda > 0$ then $\mu(|g| \geq \lambda) = \mu(|g|^p \geq \lambda^p) \leq \frac{\mu(|g|^p)}{\lambda^p} \leq \infty$ this gives the tail estimates as $\lambda \to \infty$.

Definition 5.3 (Convex Function)

Let $I \subseteq R$ be an interval. Then we say a map $c\colon I \to \mathbb{R}$ is **convex** if for all $x, y \in I$ and $t \in [0,1]$, we have $c(tx + (1-t)y) \leq tc(x) + (1-t)c(y)$. Equivalently, for all $x < t < y$ and $x, y \in I$, we have $\frac{c(t)-c(x)}{t-x} \leq \frac{c(y)-c(t)}{y-t}$.

Thus a convex function is continuous on the interior of the interval and so is Borel measurable.

Lemma 5.1

Let $I \subseteq R$ be an interval and $c : I \to \mathbb{R}$, and let $m \in$ the interior of I. If c is convex on I, $\exists\, a, b$ s.t. $c(x) \geq ax + b \;\forall\, x \in I$, and $c(m) = am + b$.

Proof. Define $a = \sup\left\{\frac{c(m)-c(x)}{m-x} : x < m, x \in I\right\}$. This exists in $\mathbb{R}$ by the second definition of convexity. Let $x, y \in I$, and $y > m > x$. Then $\frac{c(m)-c(x)}{m-x} \leq a \leq \frac{c(y)-c(m)}{y-m}$, so $c(y) \geq ay - am + c(m) = ay + b$ where we define $b = c(m) - am$. Similarly, for x, we have $c(x) \geq ax + b$. □

Theorem 5.1 (Jensen's inequality)

Let X be a integrable[a] r.v. taking values in an interval $I \subseteq \mathbb{R}$. Let $c\colon I \to \mathbb{R}$ be a convex function. Then $\mathbb{E}[c(X)]$ well-defined and

$$c(\mathbb{E}[X]) \leq \mathbb{E}[c(X)].$$

[a] $\mathbb{E}[|X|] < \infty$

Proof. If X is a constant a.s., then done.

Otherwise, then $m = \mathbb{E}[X] \in \operatorname{int}{}^{a} I$

Using the previous lemma, $\exists\; a, b$ s.t. $c(X) \geq aX + b$. In particular, $(c(X))^- \leq |a||X| + |b|$[b]. Hence, $\mathbb{E}[c^-(X)] \leq |a|\mathbb{E}[|X|] + |b| < \infty$, and $\mathbb{E}[c(X)] = \mathbb{E}[c^+(X)] - \mathbb{E}[c^-(X)]$ is well-defined in $(-\infty, \infty]$.

Integrating $c(X) \geq aX + b$[c],

$$\mathbb{E}[c(X)] \geq a\mathbb{E}[X] + b^{d} = am + b = c(m) = c(\mathbb{E}[X])$$

□

[a] Interior of
[b] $f \geq g$ gives $-f \leq -g$ so $f^- \leq g^- \leq |g|$.
[c] If $\mathbb{E}[c(X)] = \infty$ done.

[d]Expectation of Ω is 1, this is why we need a probability measure.

Example 5.1

$(\Omega, \mathcal{F}, \mathbb{P})$ and $1 \leq p \leq \infty$. If $X \in L^\infty(\mathbb{P})$ then $X \in L^p(\mathbb{P})$ as $\|X\|_p \leq \|X\|_\infty$ as $\mathbb{P}(\Omega)$ finite.

Example 5.2

If $1 \leq p < q < \infty$, $c(x) = |x|^{\frac{q}{p}}$ is a convex function. If X is a bounded r.v. (so lies in $L^\infty(\mathbb{P})$), we then have

$$\|X\|_p = \mathbb{E}\left[|X^p|\right]^{\frac{1}{p}} = c(\mathbb{E}\left[|X|^p\right])^{\frac{1}{q}} \leq {}^a\mathbb{E}\left[c(|X|^p)\right]^{\frac{1}{q}} = \|X\|_q$$

Using the monotone convergence theorem, this extends to all $X \in L^q(\mathbb{P})$ when $\|X\|_q$ is finite. So $X \in L^q \implies X \in L^p$ so $L^q(\mathbb{P}) \subseteq L^p(\mathbb{P})$ for all $1 \leq p \leq q \leq \infty$.

[a]By Jensen

Theorem 5.2 (Hölder's inequality)

Let f, g be measurable functions on $(E, \mathcal{E}, \mu)$. If p, q are **conjugate**, so $\frac{1}{p} + \frac{1}{q} = 1$ and $1 \leq p \leq q \leq \infty$, we have

$$\mu(|fg|) = \int_E |f(x)g(x)| \, \mathrm{d}\mu \leq \|f\|_p \cdot \|g\|_q$$

Remark 36. For $p = q = 2$, this is exactly the Cauchy–Schwarz inequality on L^2 (Simpler proof on Sheet 3 by considering $\int (f+g)^2 \geq 0$.).

Proof. The cases $p = 1$ or $p = \infty$ are obvious. We can assume $f \in L^p$ and $g \in L^q$ wlog since the right hand side would otherwise be infinite. We can also assume f is not equal to zero a.e., otherwise this reduces to $0 \leq 0$.

Hence, $\|f\|_p > 0$. Then, we can divide both sides by $\|f\|_p$ and then assume $\|f\|_p = 1$.

Define a prob measure $\mathbb{P}$ on $\mathcal{E}$ by $\mathbb{P}(A) = \int_A |f|^p \, \mathrm{d}\mu$ ($\mathbb{P}$ has prob density $|f|^p$ wrt μ). Note, for $h \geq 0$ $\int h \, \mathrm{d}\mathbb{P} = \int h|f|^p \, \mathrm{d}\mu$.

The

$$\begin{aligned}\mu(|fg|) &= \mu(|fg|1_{|f|>0}) \\ &= \int |f||g|1_{|f|>0} \, \mathrm{d}\mu\end{aligned}$$

$$
\begin{aligned}
&= \int \frac{|f|^p}{|f|^{p-1}} |g| 1_{|f|>0} \, \mathrm{d}\mu \\
&= \int \frac{|g|}{|f|^{p-1}} 1_{|f|>0} |f|^p \, \mathrm{d}\mu \\
&= \int \frac{|g|}{|f|^{p-1}} 1_{|f|>0} \, \mathrm{d}\mathbb{P} \\
&= \mathbb{E}\left[\frac{|g|}{|f|^{p-1}} 1_{|f|>0}\right] \\
&\leq \mathbb{E}\left[\left(\frac{|g|}{|f|^{p-1}} 1_{|f|>0}\right)^q\right]^{\frac{1}{q}} {}^{a} \\
&= \mathbb{E}\left[\frac{|g|^q}{|f|^p} 1_{|f|>0}\right]^{\frac{1}{q}} \\
&= \left(\int \frac{|g|^q}{|f|^p} 1_{|f|>0} \, \mathrm{d}\mathbb{P}\right)^{\frac{1}{q}} \\
&= \left(\int |g|^q 1_{|f|>0} \, \mathrm{d}\mu\right)^{\frac{1}{q}} \\
&\leq \left(\int |g|^q \, \mathrm{d}\mu\right)^{\frac{1}{q}} \\
&= \|g\|_q
\end{aligned}
$$

□

[a] Proven earlier by Jensen's that $\|X\|_p \leq \|X\|_q$ for $1 \leq p \leq q$.

Theorem 5.3 (Minkowski's inequality)

Let $f, g \colon (E, \mathcal{E}, \mu) \to \mathbb{R}$ be measurable functions. Then for all $1 \leq p \leq \infty$, we have $\|f+g\|_p \leq \|f\|_p + \|g\|_p$.

Proof. The results for $p = 1, \infty$ are clear. Suppose $1 < p < \infty$. We can assume wlog that $f, g \in L^p$.

We can integrate the pointwise inequality $|f+g|^p \leq 2^p(|f|^p + |g|^p)$ to deduce that $\mu(|f+g|^p) \leq 2^p[\mu(|f|^p) + \mu(|g|^p)] < \infty$ so $f + g \in L^p$. We assume that $0 < \|f+g\|_p$, otherwise the result is trivial. Now, using Hölder's inequality with q conjugate to p,

$$
\begin{aligned}
\|f+g\|_p^p = \int |f+g|^p \, \mathrm{d}\mu &= \int |f+g|^{p-1} |f+g| \, \mathrm{d}\mu \\
&\leq \int |f+g|^{p-1} |f| \, \mathrm{d}\mu + \int |f+g|^{p-1} |g| \, \mathrm{d}\mu
\end{aligned}
$$

$$
\begin{aligned}
&\leq {}^{a}\|f\|_p \left\|(f+g)^{p-1}\right\|_q + \|g\|_p \left\|(f+g)^{p-1}\right\|_q \\
&\leq \left(\int |f+g|^{q(p-1)}\,\mathrm{d}\mu\right)^{\frac{1}{q}} \left(\|f\|_p + \|g\|_p\right) \\
&\leq \left(\int |f+g|^{p}\,\mathrm{d}\mu\right)^{\frac{1}{q}} \left(\|f\|_p + \|g\|_p\right) \\
&\leq \|f+g\|_p^{\frac{p}{q}} \left(\|f\|_p + \|g\|_p\right)
\end{aligned}
$$

Dividing both sides by $\|f+g\|_p^{\frac{p}{q}}$ noting $\frac{p}{q} = p-1$, we obtain $\|f+g\|_p \leq \|f\|_p + \|g\|_p$. □

[a] By Holder's Inequality

So the L^p spaces are indeed normed spaces.

§5.2 Banach spaces

Definition 5.4 (Banach Space)
A **Banach space** is a complete normed vector space.

Theorem 5.4 ($\mathcal{L}^p$ is a Banach space)
Let $1 \leq p \leq \infty$, and let $f_n \in L^p$ be a Cauchy sequence, so $\forall\, \varepsilon > 0\, \exists\, N$ s.t. $\forall\, m, n \geq N$, we have $\|f_m - f_n\|_p < \varepsilon$. Then $\exists\, f \in L^p$ s.t. $f_n \to f$ in L^p, so $\|f_n - f\|_p \to 0$ as $n \to \infty$.

Proof. For this proof, we assume $p < \infty$; the other case is already proven in IB Analysis and Topology.

Since f_n is Cauchy, using $\varepsilon = 2^{-k}$ we extract a subsequence f_{N_k} of L^p functions s.t.

$$S = \sum_{k=1}^{\infty} \left\|f_{N_{k+1}} - f_{N_k}\right\|_p \leq \sum_{k=1}^{\infty} 2^{-k} < \infty$$

By Minkowski's inequality, for any K, we have

$$\left\|\sum_{k=1}^{K} \left|f_{N_{k+1}} - f_{N_k}\right|\right\|_p \leq \sum_{k=1}^{K} \left\|f_{N_{k+1}} - f_{N_k}\right\|_p \leq S < \infty.$$

So $\int \left|\sum_{k=1}^{K} \left|f_{N_{k+1}} - f_{N_k}\right|\right|^p \mathrm{d}\mu \leq S^p < \infty$.

By the monotone convergence theorem applied to $\left|\sum_{k=1}^{K}\left|f_{N_{k+1}}-f_{N_k}\right|\right|^p$ which increases to $\left|\sum_{k=1}^{\infty}\left|f_{N_{k+1}}-f_{N_k}\right|\right|^p$, we find

$$\left\|\sum_{k=1}^{\infty}\left|f_{N_{k+1}}-f_{N_k}\right|\right\|_p \leq S < \infty$$

Since the integral is finite, we see that $\sum_{k=1}^{\infty}\left|f_{N_{k+1}}-f_{N_k}\right|$ is finite μ-a.e.. Let A be the set where this sum is finite, then $\mu(A^c)=0$. For any $x \in A$, $(f_{N_k}(x))$ is Cauchy as sum finite, and since $\mathbb{R}$ complete it converges. Define,

$$f(x) = \begin{cases} \lim_{k\to\infty} f_{N_k}(x) & x \in A \\ 0 & x \in A^c \end{cases}$$

so $f_{N_k} \to f$ as $k \to \infty$ μ-a.e. and f measurable as the limit of measurable fcns.

Now, by Fatou's lemma,

$$\begin{aligned} \|f_n - f\|_p^p &= \mu(|f_n - f|^p) \\ &= \mu(\lim_k |f_n - f_{N_k}|^p) \\ &= \mu(\liminf_k |f_n - f_{N_k}|^p) \\ &\leq \liminf_k \mu(|f_n - f_{N_k}|^p) \\ &\leq \varepsilon^p \quad \forall\, n \geq N. \end{aligned}$$

Since the f_n are Cauchy,

$$\|f\|_p \leq \underbrace{\|f - f_N\|_p}_{\leq\varepsilon} + \underbrace{\|f_N\|_p}_{<\infty} < \infty$$

so $f \in L^p$ and so $f_n \to f$ in L^p. □

Remark 37. If V is any of the spaces

$$C([0,1]); \quad \{f \text{ simple}\}; \quad \{f \text{ a finite linear combination of indicators of intervals}\}$$

then V is dense in $L^p((0,1),\mathcal{B},\lambda)$. So the completion $\overline{(V,\|\cdot\|)}$ is exactly $L^1(\lambda)$ (Proof on Sheet 3, first prove for finite linear combinations, use monotone class theorem, approximate continuous fcns by indicators of intervals so done).

§5.3 Hilbert spaces

Definition 5.5 (Inner Product)

A symmetric bilinear form $\langle \cdot\,,\,\cdot \rangle : V \times V \to \mathbb{R}$ on a real vector space V is called an **inner product** if $\langle v, v \rangle \geq 0$ and $\langle v, v \rangle = 0 \iff v = 0$.
In this case, we can define a norm[a] $\|v\| = \sqrt{\langle v, v \rangle}$.

[a] Cauchy-Schwarz gives triangle inequality

Definition 5.6 (Hilbert Space)

If $(V, \langle \cdot\,,\,\cdot \rangle)$ is complete, we say that it is a **Hilbert space**.

Corollary 5.1

The space $\mathcal{L}^2$ is a Hilbert space for the inner product $\langle f, g \rangle = \int_E fg \,\mathrm{d}\mu$.

Example 5.3

An analog of the Pythagorean theorem holds. Let $f, g \in L^2$, then $\|f+g\|_2^2 = \|f\|_2^2 + 2\langle f, g \rangle + \|g\|_2^2$.

Example 5.4

The parallelogram identity holds: $\|f+g\|_2^2 + \|f-g\|_2^2 = 2\left(\|f\|_2^2 + \|g\|_2^2\right)$

Definition 5.7 (Orthogonal)

We say f is **orthogonal** to g if $\langle f, g \rangle = 0$.

Remark 38. f and g are orthogonal iff $\|f+g\|_2^2 = \|f\|_2^2 + \|g\|_2^2$.

For centred (mean zero) r.v.s X, Y, we have $\langle X, Y \rangle = \mathbb{E}[XY] = \mathbb{E}[(X - \mathbb{E}[X])(Y - \mathbb{E}[Y])] = \operatorname{Cov} X, Y$ which vanishes when X and Y are orthogonal.

Definition 5.8 (Orthogonal Complement)

Let $V \subseteq L^2(\mu)$. We define its **orthogonal complement** to be

$$V^\perp = \left\{ f \in L^2(\mu) : \langle f, g \rangle = 0 \quad \forall g \in V \right\}$$

Definition 5.9 (Closed Set)
We say that a subset V of $\mathcal{L}^2$ is **closed** if any sequence $f_n \in V$ that converges in $\mathcal{L}^2$, its limit f coincides a.e. with some $v \in V$.

Theorem 5.5 (Orthogonal Projection)
Let V be a closed linear subspace of $\mathcal{L}^2(\mu)$. Then $\forall\ f \in \mathcal{L}^2$, $\exists$ an orthogonal decomposition $f = v + u$ where $v \in V$ and $u \in V^\perp$.
Moreover, $\|f - v\|_2 \leq \|f - g\|_2\ \forall\ g \in V$ with equality iff $v = g$ a.e..

Definition 5.10 (Projection)
We call v the **projection** of f onto V.

Proof. In this proof, we use $p = 2$ for all norms. We define $d(f, V) = \inf_{g \in V} \|g - f\|$, and let $g_n \in V$ be a sequence of functions s.t. $\|g_n - f\| \to d(f, V)$.
By the parallelogram law,

$$\begin{aligned}
2\|f - g_n\|^2 + 2\|f - g_m\|^2 &= \|2f - (g_n + g_m)\|^2 + \|g_n - g_m\|^2 \\
&= 4\left\| f - \underbrace{\frac{g_n + g_m}{2}}_{\in V} \right\|^2 + \|g_n - g_m\|^2 \\
&\geq 4d(f, V)^2 + \|g_n - g_m\|^2
\end{aligned}$$

Taking the limit superior as $n, m \to \infty$, $\limsup_{m,n} 0$. So the sequence g_n is Cauchy in L^2, so by completeness, it converges to some $v \in L^2$. Since V is closed, $v \in V$. In particular, $d(f, V) = \inf_{g \in V} \|g - f\| = \|v - f\|$.

Note that $d(f, V)^2 \leq F(t) = \|f - (v + th)\|^2 > d(f, V)^2 - 2t\langle f - v, h\rangle + t^2\|h\|^2$ where $t \in \mathbb{R}$ and $h \in V$. Letting $t \downarrow 0$ and $t \uparrow 0$, we obtain $\langle f - v, h\rangle = 0$ for all h. Defining $f - v = u$, we have $f = u + v$ and $u \in V^\perp$ since h was arbitrary.

For any $g \in V$, $f - g = \underbrace{f - v}_{\in V^\perp} + \underbrace{v - g}_{\in V}$. So $\|f - g\|^2 = \|f - v\|^2 + \|v - g\|^2$ hence $\|f - g\| \geq \|f - v\|$ with equality iff $\|v - g\| = 0$, i.e. $v = g$ a.e..

$\square$

§5.4 Conditional Expectation

If $\mathcal{G}$ a sub-σ algebra of $\mathcal{F}$ (i.e. $\mathcal{G} \subseteq \mathcal{F}$), then $L^2(\Omega, \mathcal{G}, \mathbb{P})$ is a closed subspace of $L^2(\Omega, \mathcal{F}, \mathbb{P})$.

Definition 5.11 (Conditional Expectation)
For $X \in L^2(\Omega, \mathcal{F}, \mathbb{P})$ s.t. X measurable wrt $\mathcal{G}$, $\|X - Y\|_2 \geq \|X - \mathbb{E}[X \mid \mathcal{G}]\|_2 \quad \forall\, Y$ that are $\mathcal{G}$ measurable.
The[a] **conditional expectation of X given $\mathcal{G}$**, $\mathbb{E}[X \mid \mathcal{G}]$ is defined as the orthogonal projection of X on $L^2(\Omega, \mathcal{G}, \mathbb{P})$.

[a] This is actually only a variant of the conditional expectation.

Question
How to define $\mathcal{E}[X \mid \mathcal{G}]$ if $X \in L^1(\Omega, \mathcal{F}, \mathbb{P})$, see Advanced Probability.

Example 5.5
Let $(G_i)_{i \in I}$ be a countable family of disjoint events whose union is Ω and set $\mathcal{G} = \sigma(G_i : i \in I)$. Let X be integrable. Then the conditional expectation of X given $\mathcal{G}$ is given by:

$$\mathbb{E}[X \mid G_i] = \frac{\mathbb{E}[X 1_{G_i}]}{\mathbb{P}(G_i)} \quad \forall\, i \in I.$$

Let $Y = \sum_i \mathbb{E}[X \mid G_i] 1_{G_i}$ (i.e. if $\omega \in G_i$, $Y(\omega) = \mathbb{E}[X \mid G_i]$). Check that Y is $\mathcal{G}$-measurable; $Y \in L^2(\Omega, \mathcal{G}, \mathbb{P})$; and that Y is "the" orthogonal projection of X onto $L^2(\Omega, \mathcal{G}, \mathbb{P})$ if $X \in L^2(\Omega, \mathcal{F}, \mathbb{P})$.

§5.5 L^p Convergence and Uniform Integrability

For $(\Omega, \mathcal{F}, \mathbb{P})$, what are the implications between convergence: a.s., in L^p for $1 \leq p < \infty$, in $\mathbb{P}$ and in distribution.

Let $f_n = n 1_{(0,1/n)}$ on $((0,1), \mathcal{B}, \lambda)$. $f_n \to 0$ a.s. but $\mathbb{E}|f_n| = \mathbb{E}[f_n] = 1 \;\forall\, n$ so a.s. $\not\Longrightarrow$ L^p convergence.

$\mathbb{P}(|X_n - X| > \varepsilon) \leq \frac{\mathbb{E}|X_n - X|^p}{\varepsilon^p}$ by Markov's Inequality, so convergence in L^p for $1 \leq p < \infty \implies$ convergence in $\mathbb{P}$.

Theorem 5.6 (Dominated Convergence Theorem)
Let X_n be r.v.s on $(\Omega, \mathcal{F}, \mathbb{P})$ s.t. $|X_n| \leq Y$ for integrable r.v. Y and they converge in $\mathbb{P}$ to X. Then $X_n \to X$ in $L^1(\mathbb{P})$, i.e. $\mathbb{E}|X_n - X| \to 0$.

Question
What is the "minimum condition" on (X_n) under which $X_n \to X$ in $\mathbb{P}$ implies $X_n \to X$ in $L^1(\mathbb{P})$.

Answer
Uniformly Integrable

For $X \in L^1(\mathbb{P})$, then as $\delta \to 0$,

$$I_X(\delta) = \sup\left\{\mathbb{E}\left[|X|1_A\right] : \mathbb{P}(A) \leq \delta, A \in \mathcal{F}\right\} \to 0$$

If not, $\exists\, \varepsilon > 0$ and $A_n \in \mathcal{F}$ s.t. $\mathbb{P}(A_n) \leq 2^{-n}$ but $\mathbb{E}\left[|X|1_{A_n}\right] \geq \varepsilon$. Since $\sum_n \mathbb{P}(A_n) < \infty$, by the first Borel–Cantelli lemma, we have $\mathbb{P}\left(\bigcap_n \bigcup_{m \geq n} A_m\right) = 0$. But $\mathbb{E}\left[|X|1_{A_n}\right] \leq \mathbb{E}\left[|X|1_{\bigcup_{m\geq n} A_m}\right]$. Note that $1_{\bigcup_{m\geq n} A_m} \to 1_{\bigcap_n \bigcup_{m\geq n} A_n} = 0$ a.s., so $\mathbb{E}\left[|X|1_{\bigcup_{m\geq n} A_m}\right] \to \mathbb{E}\left[|X|1_{\bigcap_n \bigcup_{m\geq n}}\right] = 0$ by DCT ↯.

Definition 5.12 (Uniformly Integrable)
For a collection $\mathcal{X} \subseteq L^1(\mathbb{P})$ of r.v.s, we say $\mathcal{X}$ is **uniformly integrable (UI)** if it is bounded in $L^1(\mathbb{P})$[a], and

$$I_{\mathcal{X}}(\delta) = \sup\left\{\mathbb{E}\left[|X|1_A\right] : \mathbb{P}(A) \leq \delta, X \in \mathcal{X}\right\} \to 0 \text{ as } \delta \to 0.$$

[a] I.e. $\sup_{x\in\mathcal{X}} \|X\|_1 = \sup_{x\in\mathcal{X}} \mathbb{E}[|X|] = I_{\mathcal{X}}(1) < \infty$.

Remark 39.
1. Any single integrable r.v. is UI. Also, true for any finite collection of integrable r.v.s. Also, if $\mathcal{X} = \{X : X \text{ a r.v. s.t. } |X| \leq Y \text{ for some } Y \in L^1\}$ as $\sup_{X\in\mathcal{X}} \mathbb{E}[|X|1_A] \leq \mathbb{E}[Y1_A]$ implies $I_{\mathcal{X}}(\delta) \leq I_Y(\lambda) \to 0$ as $\delta \to 0$.
2. If $\mathcal{X}$ is bounded in $L^p(\mathbb{P})$ for $p > 1$, then by Hölder's inequality,

$$\mathbb{E}\left[|X|1_A\right] \leq \underbrace{\|X\|_p}_{\text{bounded}} \cdot \underbrace{\mathbb{P}(A)^{\frac{1}{q}}}_{\leq \delta^{\frac{1}{q}} \to 0}$$

Lemma 5.2

$\mathcal{X} \subseteq L^1(\mathbb{P})$ is UI $\iff$ $\sup_{X \in \mathcal{X}} \mathbb{E}\left[|X| 1_{\{|X|>K\}}\right] \to 0$ as $K \to \infty$.

Proof. ($\implies$): Applying Markov's inequality, as $K \to \infty$,

$$\mathbb{P}(|X| > K) \leq \frac{\mathbb{E}[|X|]}{K} = \frac{\mathbb{E}[|X| 1_\Omega]}{K} \leq \frac{I_{\mathcal{X}}(1)}{K} \to 0$$

Using the uniform integrability property using $A = \{|X| > K\}$, we obtain the required limit.

($\Longleftarrow$):

$$\mathbb{E}[|X|] = \mathbb{E}\left[|X|\left(1_{\{|X| \leq K\}} + 1_{\{|X|>K\}}\right)\right] \leq K + \frac{\varepsilon}{2}$$

for sufficiently large K. So $\mathcal{X}$ is bounded in $L^1(\mathbb{P})$ as required. Then for A s.t. $\mathbb{P}(A) \leq \delta$,

$$\mathbb{E}\left[|X| 1_A\left(1_{\{|X| \leq K\}} + 1_{\{|X|>K\}}\right)\right] \leq K\mathbb{P}(A) + \mathbb{E}\left[|X| 1_{\{|X|>K\}}\right] \leq K\delta + \frac{\varepsilon}{2} < \varepsilon$$

for sufficiently small δ. $\square$

Theorem 5.7

Let X_n, X be r.v.s on $(\Omega, \mathcal{F}, \mathbb{P})$. Then the following are equivalent.

1. $X_n, X \in L^1(\mathbb{P})$ and $X_n \to X$ in $L^1(\mathbb{P})$.
2. $\{X_n : n \in \mathbb{N}\}$ is uniformly integrable, and $X_n \to X$ in $\mathbb{P}$.

Proof. (1) $\implies$ (2i): Using Markov's inequality,

$$\mathbb{P}(|X_n - X| > \varepsilon) \leq \frac{\mathbb{E}[|X_n - X|]}{\varepsilon} \to 0$$

so $X_n \to X$ in $\mathbb{P}$.

Choose N s.t. $\mathbb{E}|X_n - X| < \frac{\varepsilon}{2}\ \forall\, n \geq N$. Choose δ s.t. $\mathbb{E}[|X| 1_A] \leq \frac{\varepsilon}{2}$ and $\mathbb{E}[|X_n| 1_A] \leq \varepsilon \quad \forall\, n = 1, \ldots, N-1$ when $\mathbb{P}(A) < \delta$.

$$\mathbb{E}[|X_n| 1_A] \leq \mathbb{E}[|X_n - X| 1_A] + \mathbb{E}[|X| 1_A] \leq \frac{\varepsilon}{2} + \frac{\varepsilon}{2}$$

$\{X_1, \ldots, X_{N-1}, X\}$ is finite so UI. So $\mathcal{X}$ is UI.

(2) $\implies$ (1): $X_n \to X$ in $\mathbb{P}$, so take a subsequence n_k s.t. $X_{n_k} \to X$ a.s.. Then,

$$\mathbb{E}\left[|X|\right] = \mathbb{E}\left[\liminf_k |X_{n_k}|\right] \leq {}^{a} \liminf_k \mathbb{E}\left[|X_{n_k}|\right] \leq I_{\mathcal{X}}(1) < {}^{b}\infty,$$

so $X \in L^1(\mathbb{P})$.

Next, we define truncated r.v.s $X_n^K = \max(-K, \min(K, X_n))$ and $X^K = \max(-K, \min(K, X))$. Then $X_n^K \to X^K$ in $\mathbb{P}$ (as $\mathbb{P}(|X_n^K - X^K| > \varepsilon) \leq \mathbb{P}(|X_n - X| < \varepsilon)$)[c]. And $|X_n^K| \leq K \quad \forall\, n$ so by BCT, $X_n^K \to X^K$ in L^1. Now,

$$\begin{aligned}
\mathbb{E}\left[|X_n - X|\right] &\leq \mathbb{E}\left[\left|X_n - X_n^K\right|\right] + \mathbb{E}\left[\left|X_n^K - X^K\right|\right] + \mathbb{E}\left[\left|X^K - X\right|\right] \\
&= \mathbb{E}\left[|X_n| 1_{\{|X_n| > k\}}\right] + \mathbb{E}\left[\left|X_n^K - X^K\right|\right] + \mathbb{E}\left[|X| 1_{\{|X| > K\}}\right] \\
&< \varepsilon
\end{aligned}$$

by choosing sufficiently large K (by UI) and n. $\square$

[a] Fatou's lemma
[b] As $\mathcal{X}$ is UI, hence L^1 bounded.
[c] Aside: If $X_n \to X$ in $\mathbb{P}$ and f cts, then $f(X_n) \to f(X)$ in $\mathbb{P}$.

§6 Fourier transforms

§6.1 Fourier transforms

In this section, we will write $L^p = L^p(\mathbb{R}^d)$ for the space of complex valued Borel measurable fcns on $\mathbb{R}^d$, i.e. $f\colon \mathbb{R}^d \to \mathbb{C}$ s.t. $\|f\|_p = (\int_{\mathbb{R}^d} |f(x)|^p \,\mathrm{d}\mu(x))^{\frac{1}{p}} < \infty$ for $1 \le p < \infty$.

Remark 40. For g measurable s.t. $\int |g| < \infty$, define $\int g(x)\,\mathrm{d}\mu(x) = \int \mathrm{Re}(g(x))\,\mathrm{d}\mu(x) + i\int \mathrm{Im}(g(x))\,\mathrm{d}\mu(x)$.

For $f, g \in L^2$, $\langle f, g\rangle = \int f(x)\overline{g(x)}\,\mathrm{d}\mu(x)$ is an inner product on $L^2(\mu)$.

For any $y \in \mathbb{R}^d$,

$$\begin{aligned}\int f(x-y)\,\mathrm{d}x &= \int f(y-x)\,\mathrm{d}x - \int f(x)\,\mathrm{d}x \\ &= \int f(-x)\,\mathrm{d}x\,.\end{aligned}$$

This is by the translation invariance and $x \mapsto -x$ symmetry of λ, proved in Sheet 3. Also, for $a \in \mathbb{R}$ with $a \neq 0$, $\int f(ax)\,\mathrm{d}x = \frac{1}{a^d}\int f(x)\,\mathrm{d}x$.

Definition 6.1 (Fourier Transform)
Let $f \in L^1(\mathbb{R}^d)$. We define the **Fourier transform** $\hat{f}$ by

$$\hat{f}(u) = \int_{\mathbb{R}^d} f(x)e^{i\langle u,x\rangle}\,\mathrm{d}x$$

where $u \in \mathbb{R}^d$ and $\langle u, x\rangle = \sum_{i=1}^d u_i x_i$.

Remark 41. Note that $\left|\hat{f}(u)\right| \le \|f\|_1 \quad \forall\, u \in \mathbb{R}^d$, i.e. $\hat{f} \in L^\infty$.

Also, if $u_n \to u$, then $e^{i\langle u_n,x\rangle} \to e^{i\langle u,x\rangle}$ so $f(x)e^{i\langle u_n,x\rangle} \to f(x)e^{i\langle u,x\rangle}$; $|f(x)e^{i\langle u_n,x\rangle}| \le |f(x)|$ and $f \in L^1$. By the DCT $\hat{f}(u_n) \to \hat{f}(u)$. Moreover, $\lim_{\|u\|\to\infty} \hat{f}(u) = 0$ (Riemann-Lebesgue Lemma, Sheet 3). Thus $\hat{f} \in C_0(\mathbb{R}^d) = \{f$ bounded cts and vanishing at $\pm\infty\}$.

The map is $1-1$ (but not onto), its injective but not surjective.

Definition 6.2 (Fourier Transform)
Let μ be a finite Borel measure on $\mathbb{R}^d$. We define the **Fourier transform** of the measure for $u \in \mathbb{R}^d$ by

$$\hat{\mu}(u) = \int_{\mathbb{R}^d} e^{i\langle u,x\rangle}\,\mathrm{d}\mu(x)$$

Note that $|\hat{\mu}(u)| \leq \mu(\mathbb{R}^d)$, and $\hat{\mu}$ a bounded cts fcn on $\mathbb{R}^d$. If μ has a density f (wrt λ), $\hat{\mu} = \int_{\mathbb{R}^d} e^{i\langle u,x\rangle} f(x)\,\mathrm{d}x = \hat{f}$.

Definition 6.3 (Characteristic Function)
Let X be an $\mathbb{R}^d$-valued r.v.. The **characteristic function (c.f.)** φ_X of X is the Fourier transform of its law $\mu_X = \mathbb{P} \circ X^{-1}$. So,

$$\varphi_X(u) = \hat{\mu}_X(u) = \int e^{i\langle u,x\rangle} \underbrace{\mathrm{d}\mu_X(x)}_{\mathrm{d}\mathbb{P}\circ X^{-1}(x)} \overset{a}{=} \int e^{i\langle u,x\rangle}\,\mathrm{d}\mathbb{P} = \mathbb{E}\left[e^{i\langle u,X\rangle}\right].$$

[a] Note that $\nu \circ f^{-1}(g) = \nu(f \circ g)$.

In particular if X has pdf f, then $\varphi_X(u) = \hat{f}(u)$.

Definition 6.4 (Fourier Inversion Formula)
Let $f \in L^1(\mathbb{R}^d)$ s.t. $\hat{f} \in L^1(\mathbb{R}^d)$. Then we say that the **Fourier inversion formula** holds for f if

$$f(x) = \frac{1}{(2\pi)^d} \int_{\mathbb{R}^d} \hat{f}(u) e^{-i\langle u,x\rangle}\,\mathrm{d}u$$

a.e. in $\mathbb{R}^d$.

Remark 42. The RHS is cts by DCT, so for f cts the equality is everywhere.

Remark 43. The map from $L^1 \to C_0$ by $f \mapsto \hat{f}$ is $1-1$ (for $f, g \in L^1$ with $\hat{f} = \hat{g}$, then $f - g \in L^1$ and $\widehat{f-g} = \hat{f} - \hat{g} = 0$. So by Fourier Inversion $f - g = 0$ a.e.)

§6.2 Convolutions

A key concept in Fourier analysis is convolutions.

Definition 6.5 (Convolution)
Let $f \in L^p(\mathbb{R}^d)$, $1 \leq p < \infty$ and ν be a probability measure on $\mathbb{R}^d$. We define their **convolution** $f * \nu$ by

$$(f * \nu)(x) = \begin{cases} \int_{\mathbb{R}^d} f(x-y)\,\mathrm{d}\nu(y) & \text{if integral exists;} \\ 0 & \text{else.} \end{cases}$$

Remark 44. If $1 \leq p < \infty$, by Jensen's inequality,

$$\int_{\mathbb{R}^d} |f * \nu(x)|^p\,\mathrm{d}x \leq \int_{\mathbb{R}^d} \left(\int_{\mathbb{R}^d} |f(x-y)|\,\mathrm{d}\nu(y) \right)^p \mathrm{d}x$$

$$
\begin{aligned}
&\leq \int_{\mathbb{R}^d} \int_{\mathbb{R}^d} |f(x-y)|^p \, \mathrm{d}\nu(y) \, \mathrm{d}x \quad \text{as } p \geq 1 \\
&= \int_{\mathbb{R}^d} \int_{\mathbb{R}^d} |f(x-y)|^p \, \mathrm{d}x \, \mathrm{d}\nu(y) \\
&= \int_{\mathbb{R}^d} \int_{\mathbb{R}^d} |f(x)| \, \mathrm{d}\nu(y) \, \mathrm{d}x \\
&= \int_{\mathbb{R}^d} |f(x)| \, \mathrm{d}x \\
&= \|f\|_p^p < \infty.
\end{aligned}
$$

Hence $f * v$ is defined a.e., and $\|f * v\|_p \leq \|f\|_p < \infty$. When ν has pdf $g \in L^1$[10], $f * v(x) = \int f(x-y)g(y) \, \mathrm{d}y = f * g(x)$.

In the case where ν has a density g with respect to the Lebesgue measure, we write $f * g = f * \nu$.

Definition 6.6 (Convolution)

For probability measures μ, ν on $\mathbb{R}^d$, their **convolution** $\mu * \nu$ is a probability measure on $\mathbb{R}^d$ given by the law of $X + Y$ where X, Y are independent r.v.s with laws μ and ν, so

$$
\begin{aligned}
(\mu * \nu)(A) &= \mathbb{P}(X + Y \in A) \\
&= \int_{\mathbb{R}^d \times \mathbb{R}^d} 1_A(x+y) \, \mathrm{d}(\mu \otimes \nu)(x, y) \\
&= \int_{\mathbb{R}^d} \int_{\mathbb{R}^d} 1_A(x+y) \, \mathrm{d}\nu(y) \, \mathrm{d}\mu(x)
\end{aligned}
$$

If μ has density $f \in L^1$ wrt the Lebesgue measure, $\mu * \nu$ has density $f * \nu$ wrt the Lebesgue measure. Indeed,

$$
\begin{aligned}
(\mu * \nu)(A) &= \int_{\mathbb{R}^d} \int_{\mathbb{R}^d} 1_A(x+y) f(x) \, \mathrm{d}x \, \mathrm{d}\nu(y) \\
&= \int_{\mathbb{R}^d} \int_{\mathbb{R}^d} 1_A(x) f(x-y) \, \mathrm{d}x \, \mathrm{d}\nu(y) \\
&= \int_{\mathbb{R}^d} 1_A(x) \int_{\mathbb{R}^d} f(x-y) \, \mathrm{d}\nu(y) \, \mathrm{d}x \\
&= \int_{\mathbb{R}^d} 1_A(x) (f * \nu)(x) \, \mathrm{d}x
\end{aligned}
$$

Proposition 6.1

[10] wrt Lebesgue measure

$\widehat{f * \nu}(u) = \hat{f}(u)\hat{\nu}(u)$ for all $f \in L^1$ and ν a prob measure.

Proposition 6.2
$\widehat{\mu * \nu}(u) = \mathbb{E}\left[e^{i\langle u, X+Y\rangle}\right] = \mathbb{E}\left[e^{i\langle u,X\rangle}e^{i\langle u,Y\rangle}\right] = \hat{\mu}(u)\hat{\nu}(u)$ for all μ, ν prob measures.

§6.3 Fourier transforms of Gaussians

Definition 6.7 (Normal Distribution)
The **normal distribution** $N(0,t)$ is given by the probability density function

$$g_t(x) = \frac{1}{\sqrt{2\pi t}} e^{-\frac{x^2}{2t}}$$

If φ_Z is the characteristic function of $Z \sim N(0,1)$, i.e. $\varphi_Z(u) = \int \frac{1}{\sqrt{2\pi}} e^{-z^2/2} e^{iuz}\,\mathrm{d}z$ then by a previous thm, φ_z is differentiable and can use DUTIS.

$$\begin{aligned}
\frac{\mathrm{d}}{\mathrm{d}u}\varphi_Z(u) &= \frac{1}{\sqrt{2\pi}} \int \frac{\mathrm{d}}{\mathrm{d}u}\left(e^{-z^2/2}e^{iuz}\right)\mathrm{d}z \\
&= \frac{1}{\sqrt{2\pi}} \int ize^{-z^2/2}e^{iuz}\,\mathrm{d}z \\
&= \frac{i}{\sqrt{2\pi}} \int \underbrace{e^{iuz}}_{v}\underbrace{ze^{-\frac{z^2}{2}}}_{w'}\,\mathrm{d}z \\
&= \frac{i}{\sqrt{2\pi}} \int iue^{iux}e^{-\frac{z^2}{2}}\,\mathrm{d}x \\
&= -u\varphi_Z(u)
\end{aligned}$$

Hence,

$$\frac{\mathrm{d}}{\mathrm{d}u}\left(e^{\frac{u^2}{2}}\varphi_Z(u)\right) = ue^{\frac{u^2}{2}}\varphi_Z(u) - e^{\frac{u^2}{2}}u\varphi_Z(u) = 0$$

In particular, $\varphi_Z(u) = \varphi_Z(0)e^{-\frac{u^2}{2}} = e^{-\frac{u^2}{2}}$. In other words, $\hat{g}_1(u) = \sqrt{2\pi}g_1(u)$.

In $\mathbb{R}^d$, consider a Gaussian r.v. $Z = (Z_1, \ldots, Z_d)$ with iid entries $N(0,1)$. Then, the joint pdf (wrt λ^d) of $\sqrt{t}Z$ is

$$g_t(x) = \prod_{j=1}^{d} \frac{1}{\sqrt{2\pi t}} e^{-\frac{x_j^2}{2t}} = (2\pi t)^{-\frac{d}{2}} e^{-\frac{\|x\|^2}{2t}}$$

The Fourier transform of g_t is

$$\hat{g}_t(u) = \mathbb{E}\left[e^{i\langle u, \sqrt{t}Z\rangle}\right] = \mathbb{E}\left[\prod_{j=1}^{d} e^{iu_j\sqrt{t}z_j}\right] = \prod_{j=1}^{d} \underbrace{\mathbb{E}\left[e^{iu_j\sqrt{t}z_j}\right]}_{\varphi_{Z_i}(\sqrt{t}u_i)} = \prod_{j=1}^{d} e^{-u_j^2\frac{t}{2}} = e^{-\frac{\|u\|^2 t}{2}}$$

which implies that in general, $\hat{g}_t(u) = \frac{(2\pi)^{\frac{d}{2}}}{t^{\frac{d}{2}}}\frac{t^{\frac{d}{2}}}{(2\pi)^{\frac{d}{2}}}e^{-\frac{\|u\|^2 t}{2}} = (2\pi)^{\frac{d}{2}}t^{-\frac{d}{2}}g_{\frac{1}{t}}(u)$. Taking the Fourier transform with respect to u, $\hat{\hat{g}}_t = (2\pi)^{\frac{d}{2}}t^{-\frac{d}{2}}\hat{g}_{\frac{1}{t}}(u) = (2\pi)^d g_t$. Since $g_t(-x) = g_t(x)$ and the Lebesgue measure is translation invariant, we have

$$g_t(x) = \frac{1}{(2\pi)^d}\hat{\hat{g}}_t(x) = \frac{1}{(2\pi)^d}\int_{\mathbb{R}^d} e^{-i\langle u,x\rangle}\hat{g}_t(u)\,\mathrm{d}u$$

so the Fourier inversion theorem holds for g_t.

Definition 6.8 (Gaussian Convolution)
We say that a function on $\mathbb{R}^d$ is a **Gaussian convolution** if it is of the form

$$f * g_t(x) = \int_{\mathbb{R}^d} f(x-y)g_t(y)\,\mathrm{d}y$$

where $x \in \mathbb{R}^d, t > 0, f \in L^1(\mathbb{R}^d)$.

1. $f * g_t \in L^1$ as $f \in L^1$ (proved earlier) and $\|f * g_t\|_1 \leq \|f\|_1 < \infty$.
2. $f * g_t$ is continuous on $\mathbb{R}^d$ by noting $f * g_t(x) = \int_{\mathbb{R}^d} f(y)g_t(x-y)\,\mathrm{d}y$ by translation invariance, then using DCT noting g bounded as cts.
3. $f * g_t$ is bdd.
4. $\widehat{f * g_t}(u) = \hat{f}(u)\hat{g}_t(u) = \hat{f}(u)e^{-\frac{\|u\|^2 t}{2}}$.
5. $\widehat{f * g_t}$ is bdd cts as $f * g_t \in L_1$.
6. $\left\|\widehat{f * g_t}\right\| \leq c_t\left\|\hat{f}\right\|_\infty \leq c_t\left\|\hat{f}\right\|_1$.
7. For μ a prob measure and any $t > 0$, $\mu * g_t$ is a Gaussian convolution. As, $g_t = g_{t/2} * g_{t/2}$ as g_t is the density of a $N(0,t)$ r.v.. Then $\mu * g_t = \underbrace{(\mu * g_{t/2})}_{L_1} * g_{t/2}$.

Lemma 6.1
The Fourier inversion theorem holds for all Gaussian convolutions.

Proof. Let $f \in L_1$ and $t > 0$. We can use the Fourier inversion theorem for $g_t(y)$ to see that

$$\begin{aligned}
(2\pi)^d f * g_t(x) &= (2\pi)^d \int_{\mathbb{R}^d} f(x-y) g_t(y)\,\mathrm{d}y \\
&= \int_{\mathbb{R}^d} f(x-y) \int_{\mathbb{R}^d} e^{-i\langle u,y\rangle} \hat{g}_t(u)\,\mathrm{d}u\,\mathrm{d}y \\
&= \int_{\mathbb{R}^d} e^{-i\langle u,x\rangle} \int_{\mathbb{R}^d} f(x-y) e^{i\langle u,x-y\rangle}\,\mathrm{d}y\,\hat{g}_t(u)\,\mathrm{d}u \\
&= \int_{\mathbb{R}^d} e^{-i\langle u,x\rangle} \int_{\mathbb{R}^d} f(z) e^{i\langle u,z\rangle}\,\mathrm{d}z\,\hat{g}_t(u)\,\mathrm{d}u \\
&= \int_{\mathbb{R}^d} e^{-i\langle u,x\rangle} \hat{f}(u)\hat{g}_t(u)\,\mathrm{d}u \\
&= \int_{\mathbb{R}^d} e^{-i\langle u,x\rangle} \widehat{f * g_t}(u)\,\mathrm{d}u
\end{aligned}$$

□

Lemma 6.2 (Gaussian convolutions are dense in L^p)
Let $f \in L^p(\mathbb{R}^d)$ where $1 \le p < \infty$. Then $\|f * g_t - f\|_p \to 0$ as $t \to 0$.

Proof. One can easily show that the space $C_c(\mathbb{R}^d)$ of continuous functions of compact support is dense in L^p. Hence, given $\varepsilon > 0$, $\exists\ h \in C_c(\mathbb{R}^d)$ s.t. $\|f - h\|_p < \frac{\varepsilon}{3}$. Then by linearity of convolution,

$$\|f * g_t - h * g_t\|_p = \|(f-h) * g_t\|_p \le \|f - h\|_p < \frac{\varepsilon}{3}$$

So by Minkowski's inequality,

$$\|f * g_t - f\|_p \le \underbrace{\|f * g_t - h * g_t\|_p}_{\le \varepsilon/3} + \underbrace{\|h - f\|_p}_{\le \varepsilon/3} + \|h * g_t + h\|_p \le \frac{2\varepsilon}{3} + \|h * g_t - h\|_p$$

so it suffices to prove the result for $f = h \in C_c(\mathbb{R}^d)$. We define a new map

$$e(y) = \int_{\mathbb{R}^d} |h(x-y) - h(x)|^p\,\mathrm{d}x$$

Since h is bdd (cts on compact support) and supported on $[-M, M]$ say, for some $M > 0$. As $y \to 0$, $|h(x-y) - h(x)|^p \to 0$ as h cts. Also for $|y| < 1$, $|h(x-y) - h(x)|^p \le 2^p \|h(x)\|_\infty^p 1_{|x| \le M+1}$, with the RHS being integrable. Hence by DCT, $e(y) \to 0$ as $y \to 0$.

Hence, by Jensen's inequality,

$$\begin{aligned}
\|h * g_t - h\|_p^p &= \int_{\mathbb{R}^d} \left| \int_{\mathbb{R}^d} h(x-y) g_t(y)\,\mathrm{d}y - h(x) \right|^p \mathrm{d}x \\
&= \int_{\mathbb{R}^d} \left| \int_{\mathbb{R}^d} (h(x-y) - h(x)) g_t(y)\,\mathrm{d}y \right|^p \mathrm{d}x \\
&\leq \int_{\mathbb{R}^d} \int_{\mathbb{R}^d} |h(x-y) - h(x)|^p \,\mathrm{d}x\, g_t(y)\,\mathrm{d}y \\
&= \int_{\mathbb{R}^d} e(y) g_t(y)\,\mathrm{d}y \\
&= \int_{\mathbb{R}^d} e(y) \frac{1}{t^{d/2}} g_1\left(\frac{y}{\sqrt{t}}\right)^{a} \mathrm{d}y \\
&= \int_{\mathbb{R}^d} \underbrace{e(\sqrt{t}z)}_{\to e(0)=0 \text{ as } t\to 0} g_1(z)\,\mathrm{d}z \\
&\to 0 \text{ by DCT.}
\end{aligned}$$

□

[a] Note that $g_t(u) = \frac{1}{t^{d/2}} g_1\left(\frac{u}{\sqrt{t}}\right)$

Theorem 6.1 (Fourier Inversion)
Let $f \in L^1(\mathbb{R}^d)$ be s.t. $\hat{f} \in L^1(\mathbb{R}^d)$. Then a.e. in $\mathbb{R}^d$,

$$f(x) = \frac{1}{(2\pi)^d} \int_{\mathbb{R}^d} e^{-i\langle u,x\rangle} \hat{f}(u)\,\mathrm{d}u$$

Remark 45. This proves that the Fourier transform is injective; $\hat{f} = \hat{g}$ implies $\widehat{f-g} = 0$ so by Fourier inversion, $f = g$ almost everywhere. The identity holds everywhere on $\mathbb{R}^d$ for the (unique) continuous representative f in its equivalence class.

Proof. Consider $f * g_t$ and

$$f_t(x) = \frac{1}{(2\pi)^d} \int_{\mathbb{R}^d} e^{-i\langle u,x\rangle} \hat{f}(u) \underbrace{e^{\frac{-|u|^2 t}{2}}}_{\hat{g}_t(u)} \mathrm{d}u$$

As FI holds for $f * g_t$, $f * g_t = f_t$.

So, $\|f_t - f\|_1 \xrightarrow{t\to 0} 0$ by density of Gaussian convolutions and as $f \in L^1$. So $f_t \to f$ in $\mathbb{P}$ and thus $\exists$ a subsequence s.t. $f_{t_n} \to f_t$ a.e.

Also, $e^{-i\langle u,x\rangle} \hat{f}(u) e^{\frac{-|u|^2 t}{2}}$ is bounded by $\left|\hat{f}(u)\right|$, which is integrable, and $\to$

$e^{-i\langle u,x\rangle}\hat{f}(u)$ as $t \to 0$. So by DCT, $f_t(x) \to \frac{1}{(2\pi)^d}\int e^{-i\langle u,x\rangle}\hat{f}(u)\,\mathrm{d}u$ as $t \to 0$ a.e.

Hence $f = \frac{1}{(2\pi)^d}\int e^{-i\langle u,x\rangle}\hat{f}(u)\,\mathrm{d}u$ a.e. as f_t converges to it a.e. and $f_{t_n} \to f$. □

Theorem 6.2 (Plancherel)
Let $f \in L^1(\mathbb{R}^d) \cap L^2(\mathbb{R}^d)$. Then $\|f\|_2 = (2\pi)^{-\frac{d}{2}}\left\|\hat{f}\right\|_2$.

Remark 46. By the Pythagorean identity, $\langle f, g\rangle = (2\pi)^{-d}\left\langle \hat{f}, \hat{g}\right\rangle$.

Proof. Initially, we assume $f, \hat{f} \in L^1$. In this case, $f, \hat{f} \in L^\infty$, and $(x,u) \mapsto f(x)\hat{f}(u)$ is integrable for the product Lebesgue measure $\mathrm{d}x \otimes \mathrm{d}u$ on $\mathbb{R}^d \times \mathbb{R}^d$, so Fubini's theorem for bounded functions applies.

$$\begin{aligned}
(2\pi)^d\|f\|_2^2 &= (2\pi)^d\int_{\mathbb{R}^d} f(x)\overline{f(x)}\,\mathrm{d}x \\
&= \int_{\mathbb{R}^d}\left(\int_{\mathbb{R}^d} e^{-i\langle u,x\rangle}\hat{f}(u)\,\mathrm{d}u\right)\overline{f(x)}\,\mathrm{d}x < \infty (\text{ and } f \in L^2) \\
&= \int_{\mathbb{R}^d}\hat{f}(u)\overline{\int_{\mathbb{R}^d} e^{i\langle u,x\rangle}f(x)\,\mathrm{d}x}\,\mathrm{d}u \\
&= \int_{\mathbb{R}^d}\hat{f}(u)\overline{\hat{f}(u)}\,\mathrm{d}u \\
&= \left\|\hat{f}\right\|_2^2
\end{aligned}$$

Now, let $f \in L^1 \cap L^2$. For $t > 0$, take $f_t = f * g_t \underset{t\to 0}{\longrightarrow} f$ in L^2 and so $\|f_t\|_2 \underset{t\to 0}{\longrightarrow} \|f\|_2$ continuity of the norm. Also, $\hat{f}_t(u) = \hat{f}(u)\hat{g}_t(u) = \hat{f}(u)e^{-\frac{|u|^2 t}{2}}$. So $\left|\hat{f}_t(u)\right| \uparrow \left|\hat{f}(u)\right|$ as $t \to 0$. Thus $\left\|\hat{f}_t(u)\right\|_2^2 = \int |\hat{f}_t(u)|^2\,\mathrm{d}u \underset{t\to 0}{\longrightarrow} \int |\hat{f}(u)|^2\,\mathrm{d}u = \left\|\hat{f}\right\|_2^2$ by MCT.

But, $f_t = f * g_t \in L^1$, and $\hat{f}_t \in L^1$. So by the first part of the proof, $(2\pi)^d\|f_t\|_2^2 = \left\|\hat{f}_t\right\|_2^2$. Letting $t \to 0$, we get $(2\pi)^d\|f\|_2^2 = \left\|\hat{f}\right\|_2^2$. □

Remark 47. Since $L_1 \cap L_2$ is dense in L^2, we can extend the linear operator $F_0(f) = (2\pi)^{-\frac{d}{2}}\hat{f}$ to L^2 by continuity to a linear isometry $F\colon L^2 \to L^2$ known as the *Fourier–Plancherel transform*. One can show that F is surjective with inverse $F^{-1}\colon L^2 \to L^2$.

§6.4 Characteristic fcns, Weak Convergence and the CLT

Definition 6.9 (Characteristic Function)
For a r.v. X, its **characteristic function** is

$$\varphi_X(t) = \mathbb{E}[e^{itx}] = \widehat{\mu_X} = \int e^{i\langle t,x\rangle}\,\mathrm{d}\mu_X(x)$$

Example 6.1
Consider the Dirac measure δ_0 on $\mathbb{R}$, so $\hat{\delta}_0(u) = \int_{\mathbb{R}} e^{iux}\,\mathrm{d}\delta_0(x) = 1$. But the inverse Fourier transform would be $\frac{1}{2\pi}\int_{\mathbb{R}} e^{iux}\,\mathrm{d}u$ which is not a Lebesgue integrable function.

To circumvent this, we test 'μ' on nice test fcns f.

Remark 48. 2 p.m.s μ, ν on $\mathbb{R}^d$ coincide $\iff \int f\,\mathrm{d}\mu = \int f\,\mathrm{d}\nu \quad \forall\, f : \mathbb{R}^d \to \mathbb{R}$ bdd cts[11]. In fact, enough to have condition holding $\forall\ f \in C_c^\infty$ (space of infinitely differentiable fcns with compact support).

Aside

($\mu : f \mapsto \mu(f)$ mapping from $C_c^\infty \to \mathbb{R}$ is a linear, cts (on L_f topology), hence μ is "Schwartz distribution").

Definition 6.10 (Weak Convergence)
Let $(\mu_n), \mu$ be Borel prob measures on $\mathbb{R}^d$. Then **μ_n converges to μ weakly** if $\int f\,\mathrm{d}\mu_n \to \int f\,\mathrm{d}\mu$ as $n \to \infty$ for all $f : \mathbb{R}^d \to \mathbb{R}$ bdd cts.

Remark 49.

1. For a sequence of r.v.s X_n and X another r.v., $X_n \to X$ weakly if $\mu_{X_n} \to \mu_X$ weakly.
2. A sequence of prob measures, μ_n, can have at most one weak limit by previous remark (the one about 2 p.m.s. coinciding).
3. If $X_n \to X$ weakly, and $h : \mathbb{R}^d \to \mathbb{R}^k$ cts, then $h(X_n) \to h(X)$ weakly (as r.v. in $\mathbb{R}^k$). (Continuous Mapping Theorem) (from definition as $f \circ h$ bdd cts if f bdd cts).

[11] RHS implies LHS, as if true for $f = 1_{[a,b]}$ then true on a π-system so done. We can approximate such f with bdd cts fcns and so done.

4. Suffices to check condition in definition for all $f \in C_c^\infty$. ("tightness" argument, i.e. $\exists\, K$ compact s.t. $\mu_n(K^c) < \varepsilon \quad \forall\, n$ if $\mu_n \to \mu$ weakly. Sheet 4)
5. When $d = 1$, this is equivalent to $X_n \to X$ in distribution (i.e., $F_{X_n}(x) \to F_X(x)$ at all points where $x \mapsto F_X(x)$ is cts). Sheet 4 Q1, approximate indictators.

Theorem 6.3

Let X be a r.v. in $\mathbb{R}^d$ with law μ_X. Then the characteristic function $\varphi_X = \hat{\mu}_X$ uniquely determines μ_X. In addition, if $\varphi_X \in L^1$, then μ_X has a bdd cts pdf $f_X(x) = \frac{1}{(2\pi)^d} \int_{\mathbb{R}^d} e^{-i\langle u,x\rangle} \varphi_X(u)\,\mathrm{d}u$ a.e..

Proof. Let $Z = (Z_1, \ldots, Z_d)$ be a vector of independent and identically distributed r.v.s, independent of X, with $Z_j \sim N(0,1)$. Then $\sqrt{t}Z$ has pdf g_t and $X + \sqrt{t}Z$ has pdf $f_t = \mu_X * g_t$. Then, $\hat{f}_t(u) = \hat{\mu_X}(u)\hat{g}_t(u) = \varphi_X(u)e^{-\frac{|u|^2 t}{2}}$. So by F.I. of Gaussian convolutions,

$$f_t(x) = \frac{1}{(2\pi)^d} \int_{\mathbb{R}^d} e^{-i\langle u,x\rangle} \varphi_X(u) e^{-\frac{|u|^2 t}{2}}\,\mathrm{d}u$$

which is uniquely determined by φ_X.

We show on an example sheet that two Borel prob measures μ, ν on $\mathbb{R}^d$ coincide iff $\mu(g) = \nu(g) \quad \forall\, g\colon \mathbb{R}^d \to \mathbb{R}$ bdd, cts with compact support. Now,

$$\int_{\mathbb{R}^d} g(x) f_t(x)\,\mathrm{d}x = \mathbb{E}\left[\underbrace{g(X + \sqrt{t}Z)}_{\to X \text{ a.s.}}\right]$$

Since $\left|g(X + \sqrt{t}Z)\right| \le \|g\|_\infty < \infty$, by BCT, this converges as $t \to 0$ to $\mathbb{E}\left[g(X)\right] = \int_{\mathbb{R}^d} g(x)\,\mathrm{d}\mu_X(x)\ (*)$. So by uniqueness of limits, φ_X determines μ_X.

If $\varphi_X \in L^1$, then $e^{-i\langle u,x\rangle}\varphi_X(u)e^{-\frac{|u|^2 t}{2}} \underset{t\to 0}{\longrightarrow} e^{-i\langle u,x\rangle}\varphi_X(u)$. By DCT, $f_t(x) \to f_X(x)$ as $t \to 0\ \forall\, x$.

In particular, since $\mu_X * g_t \ge 0$, $f_X \ge$ on $\mathbb{R}^d$. Then, for any bdd cts g with compact support,

$$\int \underbrace{g(x) f_t(x)}_{\text{As } t\to 0, \to g(x) f_X(x)} \mathrm{d}x \to \int g(x) f_X(x)\,\mathrm{d}x \text{ by DCT as } |f_t(x)| \le \frac{1}{(2\pi)^d}\|\varphi_X\|_1$$

By $(*)$, $\int g(x)\,\mathrm{d}\mu_X(x) = \int g(x) f_X(x)\,\mathrm{d}x \quad \forall\, g$ bdd cts with compact support.

Thus μ_X has density f_X. $\square$

Theorem 6.4 (Lévy's continuity theorem)
Let X_n, X be r.v.s on $\mathbb{R}^d$, s.t. $\varphi_{X_n}(u) \to \varphi_X(u)\ \forall\ u$, as $n \to \infty$. Then $X_n \to X$ weakly.

Remark 50.

1. A stronger version of this theorem is that if $\varphi_{X_n}(u) \to \varphi(u)\ \forall\ u$ for some fcn φ that is cts in a nbd of 0, then φ is the c.f. of some r.v. X and $X_n \to X$ weakly.
2. <u>Cramer' Wold device</u>: Let $(X_n), X$ be r.v.s on $\mathbb{R}^d$, then $X_n \to X$ weakly iff $\langle u, X_n\rangle \to \langle u, X\rangle\ \forall\ u \in \mathbb{R}^d$ weakly or in distribution in $\mathbb{R}$. LHS $\implies$ RHS by continuous mapping theorem, the converse holds as $e^{i\langle u,x\rangle}$ is bdd cts so $\mathbb{E}[f(X_n)] = \varphi_{X_n}(u) \to \mathbb{E}[f(x)] = \varphi_X(u)\ \forall\ u$. So done by Lévy's.
3. The converse holds by definition of weak convergence, testing against the complex exponentials in the Fourier transform.

Proof. Let $g : \mathbb{R}^d \to \mathbb{R}$ be compactly suppored and Lipschitz cts, i.e. $|g(x) - g(y)| \leq C_g|x - y|\ \forall\ x, y \in \mathbb{R}^d$. This includes all $g \in C_c^\infty$ as any fcn with bounded derivative is Lipschitz. Enough to show, $\mathbb{E}[g(X_n)] \to \mathbb{E}[g(X)]$.

Let $Z \sim N(0, I_d)$ indep of $(X_n), X$. Then for fixed $\varepsilon > 0$, choose $t > 0$ small enough s.t. $C_g\sqrt{t}\mathbb{E}\left[|Z|\right] \leq \frac{\varepsilon}{3}$. Then,

$$\begin{aligned}
|\mu_{X_n}(g) - \mu_X(g)| &= |\mathbb{E}\left[g(X_n)\right] - \mathbb{E}\left[g(X)\right]| \\
&\leq \left|\mathbb{E}\left[g(X_n) - g(X_n + \sqrt{t}Z)\right]\right| + \left|\mathbb{E}\left[g(X) - g(X + \sqrt{t}Z)\right]\right| \\
&+ \left|\mathbb{E}\left[g(X_n + \sqrt{t}Z) - g(X + \sqrt{t}Z)\right]\right| \\
&\leq \mathbb{E}\left[\left|g(X_n) - g(X_n + \sqrt{t}Z)\right|\right] + \mathbb{E}\left[\left|g(X) - g(X + \sqrt{t}Z)\right|\right] \\
&+ \left|\mathbb{E}\left[g(X_n + \sqrt{t}Z) - g(X + \sqrt{t}Z)\right]\right| \\
&\leq 2C_g\sqrt{t}\mathbb{E}\left[|Z|\right] + \left|\mathbb{E}\left[g(X_n + \sqrt{t}Z) - g(X + \sqrt{t}Z)\right]\right| \\
&\leq \frac{2\varepsilon}{3} + \left|\mathbb{E}\left[g(X_n + \sqrt{t}Z) - g(X + \sqrt{t}Z)\right]\right|
\end{aligned}$$

We show that the remaining term can be made less than $\frac{\varepsilon}{3}$ as $n \to \infty$. Let $f_{t,n}(x) = g_t * \mu_{X_n}$. Then, by Fourier inversion for Gaussian convolutions,

$$\begin{aligned}
\mathbb{E}\left[g(X_n + \sqrt{t}Z)\right] &= \int_{\mathbb{R}^d} g(x) f_{t,n}(x)\,\mathrm{d}x \\
&= \frac{1}{(2\pi)^d}\int_{\mathbb{R}^d} g(x) \int_{\mathbb{R}^d} e^{-i\langle u,x\rangle}\varphi_{X_n}(u)e^{-\frac{|u|^2 t}{2}}\,\mathrm{d}u\,\mathrm{d}x
\end{aligned}$$

Since characteristic functions are bounded by 1, we can use DCT with dominating

function $|g(x)|e^{-\frac{|u|^2 t}{2}}$ to find

$$\begin{aligned}\mathbb{E}\left[g(X_n+\sqrt{t}Z)\right] &\to \frac{1}{(2\pi)^d}\int_{\mathbb{R}^d} g(x)\int_{\mathbb{R}^d} e^{-i\langle u,x\rangle}\varphi_X(u)e^{-\frac{|u|^2 t}{2}}\,\mathrm{d}u\,\mathrm{d}x \\ &= \int_{\mathbb{R}^d} g(x) f_t(x)\,\mathrm{d}x \\ &= \mathbb{E}\left[g(X+\sqrt{t}Z)\right]\end{aligned}$$

where $f_t = g_t * \mu_X$. So as $n \to \infty$, the difference between these two terms can be made less than $\frac{\varepsilon}{3}$ as required. □

Note. We like adding Gaussians, as pdf of $X+Z$ exists due to Z having a pdf. Also, pdf is a Gaussian convolution, which is nice.

Theorem 6.5 (Central Limit Theorem)
Let $X_1, \ldots, X_n$ be iid r.v.s on $\mathbb{R}$ with $\mathbb{E}[X_i] = 0$ and $\operatorname{Var} X_i = 1$. Let $S_n = \sum_{i=1}^n X_n$. Then

$$\frac{1}{\sqrt{n}} S_n \xrightarrow{\text{weakly}}{}^{a} Z \sim N(0,1)$$

In particular,

$$\mathbb{P}\left(\frac{1}{\sqrt{n}} S_n \le x\right) \to \mathbb{P}(Z \le x)$$

[a] As $d = 1$ equiv to in distribution convergence

Proof. Let $X = X_1$. The characteristic function $\varphi(u) = \varphi_X(u) = \mathbb{E}\left[e^{iuX}\right]$ satisfies $\varphi(0) = 1$. As $\mathbb{E}[X^2] < \infty$ by DUTIS $\varphi'(u) = i\mathbb{E}\left[Xe^{iuX}\right]$, $\varphi''(u) = i^2\mathbb{E}\left[X^2 e^{iuX}\right]$ (Sheet 3). We can find $\varphi'(0) = i\mathbb{E}[X] = 0$ and $\varphi''(0) = -\mathbb{E}[X^2] = -\operatorname{Var} X = -1$. By Taylor's theorem, $\varphi(v) = 1 - \frac{v^2}{2} + o(v^2)$ as $v \to 0$. Now, denoting $\varphi_n(u) = \varphi_{\frac{1}{\sqrt{n}}S_n}(u)$, we can write

$$\begin{aligned}\varphi_n(u) &= \mathbb{E}\left[e^{i\frac{u}{\sqrt{n}}(X_1+\cdots+X_n)}\right] \\ &= \prod_{j=1}^n \mathbb{E}\left[e^{i\frac{u}{\sqrt{n}}X_j}\right] \\ &= \left[\varphi\left(\frac{u}{\sqrt{n}}\right)\right]^n \\ &= \left[1 - \frac{u^2}{2n} + o\left(\frac{u^2}{n}\right)\right]^n\end{aligned}$$

$$= \left[1 - \frac{u^2}{2n} + o\left(\frac{1}{n}\right)\right]^n \text{ fixing } u \text{ and letting } n \to \infty$$

The complex logarithm satisfies $\log(1+z) = z + o(z)$, so by taking logarithms, we find

$$\log \varphi_n(u) = n \log\left(1 - \frac{u^2}{2n} + o\left(\frac{1}{n}\right)\right) = -\frac{u^2}{2}$$

Hence, $\varphi_n(u) \to e^{-\frac{|u|^2}{2}} = \varphi_Z(u)$. So by Lévy's continuity theorem, the result follows. □

Remark 51. The CLT in $\mathbb{R}^d$ can be proved similarly using the Cramer-Wold device and properties of multi-variate Guassians.

Definition 6.11 (Gaussian)

A r.v. on $\mathbb{R}$ is **Gaussian** ($N(\mu, \sigma)$), if it has density

$$\frac{1}{\sqrt{2\pi}\sigma} e^{-\frac{(x-\mu)^2}{2\sigma^2}}$$

for $\mu \in \mathbb{R}$, $\sigma > 0$.

Definition 6.12 (Gaussian)

A r.v. X in $\mathbb{R}^d$ is **Gaussian** if $\langle X_n, v\rangle$ are Gaussian for each $v \in \mathbb{R}^d$.

Example 6.2

If $X_1, \ldots, X_n \overset{\text{iid}}{\sim} N(0,1)$, then $X = (X_1, \ldots, X_n)$ is Gaussian in $\mathbb{R}^n$. Check the c.f. of $\langle X, v\rangle$

Proposition 6.3

Let X be Gaussian in $\mathbb{R}^n$, A is an $m \times n$ matrix and $b \in \mathbb{R}^m$. Then

1. $AX + b$ is Guassian in $\mathbb{R}^m$.
2. $X \in L^2(\mathbb{R}^d)$, and $\mu = \mathbb{E}[X]$ and $V = \mathrm{Cov}(X_i, X_j)$ exist and determine μ_X.
3. $\varphi_X(u) = e^{i\langle \mu, u\rangle - \frac{\langle u, Vu\rangle}{2}}$

4. If V is invertible, then μ_X has pdf

$$f_X(x) = (2\pi)^{-\frac{d}{2}}(\det V)^{-\frac{1}{2}} \exp\left\{-\frac{1}{2}\left\langle x-\mu, V^{-1}(x-\mu)\right\rangle\right\}$$

5. Subvectors $X_{(1)}, X_{(2)}$[a] of X are indep iff $\operatorname{Cov}(X_{(1)}, X_{(2)}) = 0$.

[a] $X_{(1)}, X_{(2)}$ disjoint X

Proof. Proofs are easy, and in examples sheets and James Norris' notes. □

§7 Ergodic theory

§7.1 Laws of Large Numbers

Proposition 7.1 (Weak Law of Large Numbers)
Let $(X_n)_{n\in\mathbb{N}}$ be iid s.t. $\mathbb{E}[X_n] = \mu$ and $\operatorname{Var} X_n = \sigma^2 < \infty$. Then $\frac{1}{n}\sum_{i=1}^n X_i \to \mu$ in probability as $n \to \infty$.

Proof. By Chebyshev's inequality,

$$\mathbb{P}\left(\left|\frac{1}{n}\sum_{i=1}^n X_i - \mu\right| > \varepsilon\right) \le \frac{1}{n^2\varepsilon^2}\operatorname{Var}\sum_{i=1}^n X_i = \frac{\sigma^2}{n\varepsilon^2} \xrightarrow{n\to\infty} 0$$

So $\frac{1}{n}\sum_{i=1}^n X_i \to \mu$ in probability. □

This result has several weaknesses, and we can provide stronger results.

Proposition 7.2 (Strong Law of Large Numbers)
Let $(X_n)_{n\in\mathbb{N}}$ be iid s.t. $\mathbb{E}[X_n] = \mu < \infty$. Then $\frac{1}{n}\sum_{i=1}^n X_i \to \mu$ a.s. as $n \to \infty$.

Proposition 7.3
Let $(X_n)_{n\in\mathbb{N}}$ be indep with $\mathbb{E}[X_n] = \mu$ and $\mathbb{E}[X_n^4] \le M \ \forall\, n$. Then $\frac{1}{n}\sum_{i=1}^n X_i \to \mu$ a.s. as $n \to \infty$.

Proof. Let $Y_n = X_n - \mu$. Then $\mathbb{E}[Y_n] = 0$, and $\mathbb{E}[Y_n^4] \le 2^4(\mathbb{E}[X_n^4] + \mu^4) < \infty$. So we can assume $\mu = 0$. For distinct indices i, j, k, ℓ, by independence and the Cauchy–Schwarz inequality, we have

$$0 = \mathbb{E}[X_iX_jX_kX_\ell] = \mathbb{E}\left[X_i^2X_jX_k\right] = \mathbb{E}\left[X_i^3X_j\right]; \quad \mathbb{E}\left[X_i^2X_j^2\right] \le \sqrt{\mathbb{E}[X_i^4]}\sqrt{\mathbb{E}\left[X_j^4\right]} \le M$$

So we can compute

$$\mathbb{E}\left[\left(\sum_{i=1}^n X_i\right)^4\right] = \mathbb{E}\left[\sum_{i=1}^n X_i^4\right] + 6\mathbb{E}\left[\sum_{i<j} X_i^2X_j^2\right] \le nM + 3n(n-1)M \le 3n^2M$$

Let $S_n = \sum_{i=1}^n X_i$. Then,

$$\begin{aligned}
\mathbb{E}\left[\left(\frac{S_n}{n}\right)^4\right] &\leq \frac{1}{n^4} 3n^2 M \\
\mathbb{E}\left[\sum_{n=1}^{\infty}\left(\frac{S_n}{n}\right)^4\right] &= \sum_{i=1}^{\infty} \mathbb{E}\left[\left(\frac{S_n}{n}\right)\right] \text{ by MCT} \\
&\leq \sum_{n=1}^{\infty} \frac{3M}{n^2} < \infty
\end{aligned}$$

Hence $\sum_{n=1}^{\infty}\left(\frac{S_n}{n}\right)^4 < \infty$ almost surely. But then $\left(\frac{S_n}{n}\right)^4 \to 0$ almost surely, so $\frac{S_n}{n} \to 0$ almost surely. $\square$

§7.2 Invariants - Measure Preserving Transformations

Let $(E, \mathcal{E}, \mu)$ be a σ-finite measure space.

Definition 7.1 (Measure Preserving)
A measurable map $\Theta\colon E \to E$ is **measure-preserving** (m.p.) if $\mu \circ \Theta^{-1} = \mu$, i.e. $\mu(\Theta^{-1}(A)) = \mu(A)\ \forall\ A \in \mathcal{E}$.

In this case, for any integrable function $f \in L^1(\mu)$, we have $\int_E f \,\mathrm{d}\mu = \int_E f \circ \Theta \,\mathrm{d}\mu$ as $\int_E f \circ \Theta \,\mathrm{d}\mu = \int_E f \,\mathrm{d}\mu \circ \Theta^{-1}$.

Definition 7.2 (Θ-Invariant)
A measurable map $f\colon E \to \mathbb{R}$ is called **Θ-invariant** if $f \circ \Theta = f$.
A set $A \in \mathcal{E}$ is **Θ-invariant** if $\Theta^{-1}(A) = A$, or equivalently, 1_A is Θ-invariant.

The collection $\mathcal{E}_\Theta$ of Θ-invariant sets forms a σ-algebra over E. A function $f\colon E \to \mathbb{R}$ is invariant iff f is $\mathcal{E}_\Theta$-measurable (Sheet 4).

Definition 7.3 (Ergodic)
Θ is called **ergodic** if the $\mathcal{E}_\Theta$ is μ-trivial, i.e. $\forall\ A \in \mathcal{E}_\Theta\ \mu(A) = 0$ or $\mu(A^c) = 0$.

The point is an ergodic transformation mixes the sets well. Boltzman (1880), Ergodic hypothesis - over long times, a gas particle in some space will "fill the whole space", it will be arbitrarily close to any point in the space. You could imagine a trajectory, $x, \Theta(x), \Theta^2(x), \ldots$

For Markov Chains, ergodicity $\iff$ irreducibility.

Fact 7.1

If f is Θ-invariant, Θ is ergodic iff f is constant a.s. on E (Sheet 4).

Proof. See Sheet 4. We want to show distribution of f is a step function. We know $f \in \mathcal{E}_\Theta$ so $\mu(f^{-1}(-\infty, x)) = 0$ or $\mu(f^{-1}[x, \infty)) = 0$. Hence f is constant. $\square$

Example 7.1

Consider $(E, \mathcal{E}) = ((0, 1], \mathcal{B})$ with the Lebesgue measure μ. The maps $\Theta_a(x) = x + a$ modulo 1 and $\Theta(x) = 2x$ modulo 1 are both m.p., and ergodic unless $a \in \mathbb{Q}$ (Sheet 4).

§7.3 Ergodic theorems

Lemma 7.1 (Maximal Ergodic Lemma)

Let $(E, \mathcal{E}, \mu)$ be a σ-finite measure space. Let $\Theta \colon E \to E$ be m.p.. For $f \in L^1(\mu)$, we define $S_0(f) = 0$ and $S_n(f) = \sum_{k=0}^{n-1} f \circ \Theta^k$. Let $S^\star = S^\star(f) = \sup_{n \geq 0} S_n(f)$. Then $\int_{\{S^\star > 0\}} f \, \mathrm{d}\mu \geq 0$.

Proof (non-examinable). Define $S_n^\star = \max_{k \leq n} S_k$. Then clearly $S_n^\star \uparrow S^\star$, and $S_k \leq S_n^\star$ for all $k \leq n$. Note that for $m \leq n + 1$, $S_m = S_{m-1} \circ \Theta + f \leq S_n^\star \circ \Theta + f$.

Define $A_n = \{S_n^\star > 0\}$. On A_n, we have

$$S_n^\star = \max_{1 \leq k \leq n} S_k \leq \max_{1 \leq k \leq n+1} S_k \leq S_n^\star \circ \Theta + f$$

since $S_0 = 0$. We can integrate this inequality to find

$$\int_{A_n} S_n^\star \, \mathrm{d}\mu \leq \int_{A_n} S_n^\star \circ \Theta \, \mathrm{d}\mu + \int_{A_n} f \, \mathrm{d}\mu$$

On A_n^c, we must have $S_n^\star = 0 \leq S_n^\star \circ \Theta$. Hence,

$$\int_E S_n^\star \, \mathrm{d}\mu \leq \int_E S_n^\star \circ \Theta \, \mathrm{d}\mu + \int_{A_n} f \, \mathrm{d}\mu$$

Since Θ is m.p.,

$$\int_E S_n^\star \, \mathrm{d}\mu \leq \int_E S_n^\star \, \mathrm{d}\mu + \int_{A_n} f \, \mathrm{d}\mu$$

so we obtain (as $S_n^\star \in L_1$)

$$\int_{A_n} f \, \mathrm{d}\mu \geq 0 \quad \forall \, n.$$

$A_n = \{S_n^\star > 0\} = \{\max_{0 \leq m \leq n} S_m > 0\} = \bigcup_{m=0}^n \{S_m > 0\} \uparrow \bigcup_{m=0}^\infty \{S_m > 0\} = \{\sup S_m > 0\} = S^\star$.
Hence, $1_{A_n} \to 1_{S^\star > 0}$ and so $f 1_{A_n} \to f 1_{\{S^\star > 0\}}$ pointwise, and $|f 1_{A_n}| \leq |f| \in L^1(\mu)$, by DCT

$$\int_{\{S^\star > 0\}} f \, \mathrm{d}\mu = \lim_{n \to \infty} \int_{A_n} f \, \mathrm{d}\mu \geq 0$$

as required. □

Remark 52. Let μ be a finite measure. Then for $f \in L^1$ and any $\alpha > 0$, define $\overline{S}_k = \frac{S_k(f)}{k}$ and $\overline{S}^\star = \sup_{k \geq 0} \overline{S}_k$, then

$$\mu\left(\overline{S}^\star > \alpha\right) \leq \frac{1}{\alpha} \int_{\overline{S}^\star > \alpha} f \, \mathrm{d}\mu \leq \frac{1}{\alpha} \int_E |f| \, \mathrm{d}\mu \, .$$

Proof. Proof is left as an exercise, follows from shifting f by α and then applying the maximal ergodic theorem. □

Exercise 7.1. For μ a prob measure and $f \in L^1(\mu)$, show that $\left\{\frac{S_n(f)}{n} : n \in \mathbb{N}\right\}$ is UI. Hence $\frac{S_n(f)}{n} \to \overline{f}$ in L^1 by Birkhoff's. If Θ is ergodic, then $\overline{f} = \int f \, \mathrm{d}\mu$ a.s..

Theorem 7.2 (Birkhoff's Ergodic Theorem)
Let $(E, \mathcal{E}, \mu)$ be a σ-finite measure space. Let $\Theta \colon E \to E$ be m.p.. For $f \in L^1(E, \mathcal{E}, \mu)$, we define $S_0 = 0$ and $S_n = S_n(f) = \sum_{k=0}^{n-1} f \circ \Theta^k$. Then $\exists$ a Θ-invariant integrable fcn $\overline{f} \in L^1(E, \mathcal{E}, \mu)$ with $\mu\left(\left|\overline{f}\right|\right) \leq \mu(|f|)$ s.t. $\frac{S_n(f)}{n} \to \overline{f}$ a.e. as $n \to \infty$.

Remark 53. If Θ ergodic, $\hat{f}$ is constant a.e..

Relating back to the gas example, $x, \Theta(x), \ldots$ is the trajectory of a gas particle. Then $\frac{S_n(f)}{n}$ is the average of f along the trajectory (time average). Then $\overline{f}$ is the average of f over the whole space.

Proof (non-examinable). Since $\mu(|f \circ \Theta^{n-1}|) = \mu(|f|)$, we have $\mu(|S_n|) \leq n\mu(|f|)$ and

thus by Fatou's

$$\mu\left(\left|\overline{f}\right|\right) = \mu\left(\liminf_n \left|\frac{S_n}{n}\right|\right) \leq \liminf_n \mu\left(\left|\frac{S_n}{n}\right|\right) \leq \mu(|f|).$$

Note that

$$\frac{S_n \circ \Theta}{n} = \frac{S_{n+1} - f}{n+1}\frac{n+1}{n}$$

$$\text{So, } \limsup_n \frac{S_n \circ \Theta}{n} = \limsup_n \frac{S_{n+1}}{n+1} = \limsup_n \frac{S_n}{n}$$

and the same holds for $\liminf_n$. Hence $\limsup_n \frac{S_n}{n}$ and $\liminf_n \frac{S_n}{n}$ are invariant functions. So they are $\mathcal{E}_\Theta$-measurable. Hence

$$D = D_{a,b} = \left\{\liminf_n \frac{S_n}{n} < a < b < \limsup_n \frac{S_n}{n}\right\}$$

are measurable and invariant sets.

It suffices to show that $\mu(D) = 0$. As letting $\Delta = \left\{\liminf\left(\frac{S_n}{n}\right) < \limsup\left(\frac{S_n}{n}\right)\right\} = \bigcup_{a<b\in\mathbb{Q}} D_{a,b}$. Hence if $\mu(D) = 0 \implies \mu\left(\bigcup_{a<b\in\mathbb{Q}} D_{a,b}\right) = 0 \implies \mu(\Delta) = 0$. Define,

$$\overline{f} = \begin{cases} \liminf \frac{S_n}{n} = \limsup \frac{S_n}{n} & x \in \Delta^c \\ 0 & x \in \Delta \end{cases}$$

then $S_n/n \to \overline{f}$ μ-a.e. and $\overline{f}$ is Θ invariant (as $\liminf \frac{S_n}{n}$ and Δ are Θ-invariant).

Fix $a < b$. Note that $\Theta : D \to D$ by invariance and Θ is $\mu\,|_D$-measure preserving as

$$\begin{aligned} \mu\Big|_D(A) &= \mu(A \cap D) = \mu(\Theta^{-1}(A \cap D)) = \mu(\Theta^{-1}(A) \cap \Theta^{-1}(D)) \\ &= \mu(\Theta^{-1}(A) \cap D) = \mu\Big|_D(\Theta^{-1}(A)) \end{aligned}$$

Also either $b > 0$ or $a < 0$ (if $a < 0$ change f to $-f$ and a to $-b$, then $b = -a > 0$). So assume $b > 0$ wlog.

We will apply the Maximal Ergodic Lemma with $E = D$ and $\mu = \mu|_D$. Let $B \in \mathcal{E}$, where $B \subseteq D$ s.t. $\mu(B) < \infty$. Let $g = f - b1_B \in L^1(\mu)$. Then,

$$S_n(g) = S_n(f) - bS_n(1_B) \geq S_n - bn$$

which is positive on D for some n by the definition of $\limsup_n$. Hence, $S^\star(g) = \sup_{n\geq 0} S_n(g) > 0$ on D.

So $D = \{S^\star > 0\} \cap D$. Then by the Maximal Ergodic Lemma on D,

$$0 \leq \int_{D \cap \{S^\star > 0\}} g \, \mathrm{d}\mu = \int_D g \, \mathrm{d}\mu = \int_D f \, \mathrm{d}\mu - b\mu(B)$$

Hence, $b\mu(B) \leq \int_D f \, \mathrm{d}\mu$.

By σ-finiteness, this inequality extends to D; $\exists$ measurable sets $B_n \uparrow D$ where $\mu(B_n) < \infty$. Hence taking limits, $b\mu(D) = b \lim_n \mu(B_n) \leq \int_D f \, \mathrm{d}\mu$. Thus $\mu(D) < \infty$ as $f \in L^1$.
Repeating a similar arguments to above[a] for $-f$ and $-a$, we obtain $-a\mu(D) \leq \int_D -f \, \mathrm{d}\mu$. Combining these two inequalities gives

$$b\mu(D) \leq \int_D f \, \mathrm{d}\mu \leq a\mu(D)$$

But $a < b$ and $\mu(D) < \infty$, so $\mu(D) = 0$. □

[a] Now however we take D instead of B and $g' = -f - (-a)1_D$. On D, $S_n(g') = S_n(-f) - (-a)S_(1_D) = S_n(-f) - (-a)n1_D = S_n(-f) - (-a)n$ as D is Θ invariant, i.e. $S_n(1_D) = 1_D$.

Theorem 7.3 (von Neumann's L^p ergodic theorem)
Let $\mu(E) < \infty$ and $1 \leq p < \infty$. Then for $f \in L^p(\mu)$, $\frac{S_n(f)}{n} \to \overline{f}$ in L^p as $n \to \infty$.

Proof. Since Θ is m.p., we have

$$\|f \circ \Theta^n\|_p^p = \int_E |f|^p \circ \Theta^n \, \mathrm{d}\mu = \int_E |f|^p \, \mathrm{d}\mu = \int_E |f|^p \, \mathrm{d}\mu = \|f\|_p^p$$

Thus, by Minkowski's inequality, for all $f \in L^p$ we have

$$\left\| \frac{S_n(f)}{n} \right\|_p \leq \frac{1}{n} \sum_{i=0}^{n-1} \left\| f \circ \Theta^i \right\|_p = \|f\|_p$$

So $\frac{S_n(f)}{n}$ is a contraction in L^p. For each $K > 0$, we define $f_K = \max(\min(f, K), -K)$. Then

$$\|f - f_K\|_p^p = \int_E |f - f_K|^p 1_{|f|>K} \, \mathrm{d}\mu$$

Since $1_{|f|>K} \to 0$ pointwise, and $|f - f_K| \leq 2|f|^p \in L^1$, we find $\|f - f_K\|_p < \frac{\varepsilon}{3}$ by DCT, for sufficiently large $K = K_\varepsilon$.

As $|f_K| \leq K$, we have $\left|\frac{S_n(f_K)}{n}\right| \leq K$. Since μ is finite, $f_K \in L^1(\mu)$, so by Birkhoff's ergodic theorem, $\frac{S_n(f_K)}{n} \to \overline{f}_K$ a.e. for some invariant $\overline{f}_K \in L^1$. Note that $\overline{f}_k$ is

bounded by K as $\frac{S_n(f_K)}{n}$ is bounded by K. By the bounded convergence theorem, we deduce that $\left\|\frac{S_n(f_K)}{n} - \overline{f}_K\right\| \to 0$ as $n \to \infty$. Further, this holds in L^p since

$$\left\|\frac{S_n(f_K)}{n} - \overline{f}_K\right\|_p \leq (2K)^{\frac{p-1}{p}} \left\|\frac{S_n(f_K)}{n} - \overline{f}_K\right\|_1 < \frac{\varepsilon}{3}$$

where the last inequality holds for sufficiently large n. Since μ is a finite measure, $L^p(\mu) \subseteq L^1(\mu)$, hence by Birkhoff's ergodic theorem, $\frac{S_n(f)}{n} \to \overline{f}$ a.e. as $f \to \infty$. Then, by the contraction property applied to $f - f_K$,

$$\begin{aligned}
\left\|\overline{f} - \overline{f}_K\right\|_p^p &= \int_E \left|\overline{f} - \overline{f}_K\right|^p \mathrm{d}\mu \\
&= \int_E \liminf_n \left|\frac{S_n(f) - S_n(f_K)}{n}\right|^p \mathrm{d}\mu \\
&\leq \liminf_n \int_E \left|\frac{S_n(f) - S_n(f_K)}{n}\right|^p \mathrm{d}\mu \quad \text{by Fatou} \\
&= \liminf_n \int_E \left|\frac{S_n(f - f_K)}{n}\right|^p \mathrm{d}\mu \\
&= \liminf_n \left\|\frac{S_n(f - f_K)}{n}\right\|_p^p \\
&\leq \liminf_n \|f - f_K\|_p^p \quad \text{shown earlier by Minkowski} \\
&= \|f - f_K\|_p^p < \left(\frac{\varepsilon}{3}\right)^p
\end{aligned}$$

So in particular, $\overline{f} \in L^p$. Then by Minkowski,

$$\begin{aligned}
\left\|\frac{S_n(f)}{n} - \overline{f}\right\|_p &\leq \left\|\frac{S_n(f) - S_n(f_K)}{n}\right\|_p + \left\|\frac{S_n(f_K)}{n} - \overline{f}_K\right\|_p + \left\|\overline{f} - \overline{f}_K\right\|_p \\
&< \left\|\frac{S_n(f) - S_n(f_K)}{n}\right\|_p + \frac{2\varepsilon}{3} \\
&\leq \|f - f_K\|_p + \frac{2\varepsilon}{3} = \varepsilon
\end{aligned}$$

for sufficiently large n. □

Remark 54.

1. If μ a prob measure and Θ ergodic, then $\bar{f}$ is a constant a.s., so $\bar{f} = \int \bar{f}\,\mathrm{d}\mu$. Also, $\int f\,\mathrm{d}\mu = \int \frac{S_n(f)}{n}\,\mathrm{d}\mu \to \int \bar{f}\,\mathrm{d}\mu \ \forall\, n$. Hence $\int \bar{f}\,\mathrm{d}\mu = \int f\,\mathrm{d}\mu$.
   Then, $\frac{S_n(f)}{n} \xrightarrow{n\to\infty} \mathbb{E}[f]$ μ a.s. and in L^1.

2. For Θ m.p. and $f \in L^1$, $\frac{S_n(f)}{n} \xrightarrow{\mu \text{ a.s and } L^1} \mathbb{E}[f \mid \mathcal{E}_\Theta]$. For $f \in L^2$, $\mathbb{E}[f \mid \mathcal{E}_\Theta]$ is (a version of) the projection of f on $L^2(\mathcal{E}_\Theta)$.

§7.4 Infinite product spaces

Let $E = \mathbb{R}^{\mathbb{N}} = \{x = (x_n)_{n\in\mathbb{N}}\}$ be the space of real sequences. Consider

$$\mathcal{C} = \left\{ A = \prod_{n=1}^{\infty} A_n : A_n \in \mathcal{B}, \exists N \in \mathbb{N}, \forall n > N, A_n = \mathbb{R} \right\}$$

This forms a π-system, which generates the **cylindrical σ-algebra** $\sigma(\mathcal{C})$. We can show that $\sigma(\mathcal{C}) = \sigma(\{f_n : n \in \mathbb{N}\})$ where $f_n(x) = x_n$ are the coordinate projection functions on E. We can also show $\sigma(\mathcal{C}) = \mathcal{B}(\mathbb{R}^{\mathbb{N}})$ for the product topology or the topology of pointwise convergence.

Let $(X_n)_{n\in\mathbb{N}}$ be a sequence of iid r.v.s defined on $(\Omega, \mathcal{F}, \mathbb{P})$ with common law $\mu_{X_n} = m$ for all n; this exists by an earlier theorem. We define a map $X \colon \Omega \to E$ by $X(\omega) = (X_1(\omega), X_2(\omega), \dots)$. This is $\mathcal{F}$–$\sigma(\mathcal{C})$ measurable, since for all $A \in \mathcal{C}$, we have

$$X^{-1}(A) = \{\omega : X_1(\omega) \in A_1, \dots, X_N(\omega) \in A_N\} = \bigcap_{n=1}^{N} X_n^{-1}(A_n) \in \mathcal{F}$$

We denote $\mu = \mathbb{P} \circ X^{-1}$, which is the unique product prob measure in E satisfying

$$\begin{aligned}
\mu\left(\prod_{n=1}^{\infty} A_n\right) &= \lim_{N\to\infty} \mu\left(\prod_{n=1}^{N} A_n\right) \\
&= \lim_{N\to\infty} \mathbb{P}\left(X_1 \in A_1, \dots, X_N \in A_N\right) \\
&= \lim_{N\to\infty} \mathbb{P}\left(X_1 \in A_1\right) \cdots \mathbb{P}\left(X_N \in A_N\right) \\
&= \prod_{n=1}^{\infty} \mathbb{P}\left(X_n \in A_n\right) \\
&= \prod_{n=1}^{\infty} m(A_n)
\end{aligned}$$

Note that we need to use finiteness of N to exploit independence of the X_i. We call $(E, \mathcal{E}, \mu) = (\mathbb{R}^{\mathbb{N}}, \sigma(\mathcal{C}), m^{\mathbb{N}})$ the **canonical model** for an infinite sequence of iid r.v.s of law m.

Definition 7.4 (Shift Map)

The shift map $\Theta \colon E \to E$ is defined by $\Theta(x_1, x_2, \dots) = (x_2, x_3, \dots)$.

Theorem 7.4

On $(E, \mathcal{E}, \mu)$, the shift map Θ is measurable, m.p. and ergodic.

Proof. Measurability is obvious.

For $A \in \mathcal{C}$,

$$\begin{aligned}
\mu(A) &= \mathbb{P}\left(X_1 \in A_1, \dots, X_N \in A_N\right) \\
&= \mathbb{P}\left(X_1 \in A_1\right) \cdots \mathbb{P}\left(X_N \in A_N\right) \\
&= \prod_{n=1}^{N} m(A_n) \\
&= \mathbb{P}\left(X_2 \in A_1\right) \cdots \mathbb{P}\left(X_{N+1} \in A_N\right) \\
&= \mu(\Theta^{-1}(A))
\end{aligned}$$

so Θ is m.p. as m.p. on π-system $\mathcal{C}$.

Recall that the tail σ-algebra is defined by $\mathcal{T} = \bigcap_n \mathcal{T}_n$ where $\mathcal{T}_n = \sigma(X_{n+1}, X_{n+2}, \dots) = \sigma(f_{n+1}, f_{n+2}, \dots)$. Note that for all $A \in \mathcal{C}$, we have

$$\begin{aligned}
\Theta^{-n}(A) &= \mathbb{R}^n \times A_1 \times A_2 \times \dots \\
&= \left\{x \in \mathbb{R}^{\mathbb{N}} : (x_{n+1}, x_{n+2}, \dots) \in A\right\} \in \mathcal{T}_n
\end{aligned}$$

Now, if A is invariant, $A = \Theta^{-n}(A) \in \mathcal{T}_n \ \forall\, n$, so $A \in \mathcal{T}$. By Kolmogorov's zero-one law as (X_i) iid, $\mu(A) = 0$ or $\mu(A) = 1$ as required for ergodicity. $\square$

We can apply Birkhoff's ergodic theorem to Θ. If $f \in L^1(\mu)$, then $\frac{S_n(f)}{n} \to \overline{f} \in L^1(\mu)$ almost surely. Since $\overline{f}$ is invariant and Θ is ergodic, $\overline{f}$ is almost surely constant. By von Neumann's L^p-ergodic theorem, convergence holds in fact in L^1.

§7.5 Strong law of large numbers

Theorem 7.5

Let m be a prob measure on $\mathbb{R}$ s.t. $\int_{\mathbb{R}} |x| \, \mathrm{d}m(x) < \infty$ and $\int_{\mathbb{R}} x \, \mathrm{d}m(x) = \nu$. Let $(E, \mathcal{E}, \mu)$ be the canonical model where the coordinate maps $f_n(x) = x_n$ are iid with law m. Then

$$\mu\left(\left\{x \in \mathbb{R}^{\mathbb{N}} : \frac{x_1 + x_2 + \dots + x_n}{n} \to \nu\right\}\right) = 1$$

Proof. Let $\Theta : E \to E$ be the shift map $\Theta(x_1, x_2, \dots) = (x_2, x_3, \dots)$. It is m.p. and ergodic by previous thm. Consider $f : E \to \mathbb{R}$ as $f(x) = x_1$. Then $f \in L^1(\mu)$, since $\int_E |f| \, \mathrm{d}\mu = \int_{\mathbb{R}} \mathrm{d}\mu \circ f^{-1}(x) = \int_{\mathbb{R}} |x_1| \, \mathrm{d}m(x_1) < \infty$. So by Birkhoff and von-Neumann, as Θ ergodic, by remark 54, $\frac{S_n(f)}{n} = \frac{x_1 + \dots + x_n}{n} \to \overline{f} = \int f \, \mathrm{d}\mu = \int_{\mathbb{R}} x_1 \, \mathrm{d}m(x_1) = \nu$ μ a.s. $\square$

Theorem 7.6 (Kolmogorov SLLN (1930))
Let $(X_n)_{n\in\mathbb{N}}$ be iid r.v.s s.t. $\mathbb{E}\left[|X_1|\right] < \infty$. Then $\frac{1}{n}\sum_{i=1}^n X_i \to \mathbb{E}\left[X_1\right]$ almost surely.

Proof. Let m be the law of X_n, $\nu = \mathbb{E}[X_1]$ and $\mu = \mathbb{P} \circ X^{-1}$ where $X : \Omega \to E = \mathbb{R}^{\mathbb{N}}$ is $X(\omega) = (X_1(\omega), X_2(\omega), \dots)$. Then apply theorem 7.5,

$$\mathbb{P}\left(\frac{1}{n}\sum_{i=1}^n X_i \to \mathbb{E}\left[X_1\right]\right) = \mu\left(\left\{x : \frac{x_1 + \cdots + x_n}{n} \to \nu\right\}\right) = 1$$

□

Remark 55. The hypothesis $\mathbb{E}\left[|X|\right] < \infty$ cannot be weakened; we see on an example sheet that $\frac{1}{n}\sum_{i=1}^n X_i$ can exhibit various behaviours. Note that this notion of convergence is stronger than the weak convergence seen in the central limit theorem. The law of the iterated logarithm is that

$$\limsup_n \frac{X_1 + \cdots + X_n}{\sqrt{2n \log\log n}} = 1$$

almost surely, and -1 for the limit inferior. In particular, the central limit theorem does not hold almost surely.

Corollary 7.1
By von Neumann's ergodic theorem, in the strong law of large numbers, we have $\mathbb{E}\left[\left|\frac{1}{n}\sum_{i=1}^n X_i - \mathbb{E}\left[X\right]\right|\right] \to 0$ as $n \to \infty$.

Aside

1. If (μ_n) is a sequence in $\mathbb{R}^n$ of prob measures that converges weakly to μ, then (μ_n) is "tight", i.e. $\forall\, \varepsilon > 0$, $\exists$ a compact set K s.t. $\mu_n(K^c) < \varepsilon\ \forall\, n$.
2. If (μ_n) a sequence of prob measures that are tight, then $\exists$ a subsequence (n_k) and a prob measure μ s.t. $(\mu_{n_k}) \to \mu$ weakly (Prokhorov's thm or Banach-Alaoglu thm).
3. If distributions $F_n \xrightarrow{d} F$, then $\exists$ a prob space s.t. μ_{X_n} is F_n and $X_n \to X$ a.s..
4. If $X_n \xrightarrow{d} X$ and $Y_n \xrightarrow{d} Y$ then $X_n + Y_n \not\to X + Y$, in fact X, Y are not necessarily even defined in the same prob space so $X + Y$ doesn't even make sense.
   However, if $X_n \xrightarrow{d} X$ and $Y_n \xrightarrow{\mathbb{P}} c$, where c is a constant. Then $(X_n, Y_n) \to (X, c)$ so by continuous mapping thm $X_n + Y_n \xrightarrow{d} X + c$ (Slutsky's thm).

This is quite useful in stats where with CLT if variance is unknown we can replace it with std deviation which converges to the true value a.s. and hence in $\mathbb{P}$.