

PROBABILITY AND MEASURE

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SCHEDULE

Measure spaces, σ -algebras, π -systems and uniqueness of extension, statement *and proof* of Carathéodory's extension theorem. Construction of Lebesgue measure on \mathbb{R} , Borel σ -algebra of \mathbb{R} , existence of a non-measurable subset of \mathbb{R} . Lebesgue–Stieltjes measures and probability distribution functions. Independence of events, independence of σ -algebras. Borel–Cantelli lemmas. Kolmogorov's zero–one law.

Measurable functions, random variables, independence of random variables. Construction of the integral, expectation. Convergence in measure and convergence almost everywhere. Fatou's lemma, monotone and dominated convergence, uniform integrability, differentiation under the integral sign. Discussion of product measure and statement of Fubini's theorem.

Chebyshev's inequality, tail estimates. Jensen's inequality. Completeness of L^p for $1 \leq p \leq \infty$. Hölder's and Minkowski's inequalities, uniform integrability.

L^2 as a Hilbert space. Orthogonal projection, relation with elementary conditional probability. Variance and covariance. Gaussian random variables, the multivariate normal distribution.

The strong law of large numbers, proof for independent random variables with bounded fourth moments. Measure preserving transformations, Bernoulli shifts. Statements *and proofs* of the maximal ergodic theorem and Birkhoff's almost everywhere ergodic theorem, proof of the strong law.

The Fourier transform of a finite measure, characteristic functions, uniqueness and inversion. Weak convergence, statement of Lévy's continuity theorem for characteristic functions. The central limit theorem.

Appropriate books

P. Billingsley *Probability and Measure*. Wiley 1995 (£71.50 hardback).

R.M. Dudley *Real Analysis and Probability*. Cambridge University Press 2002 (£35.00 paperback).

R.T. Durrett *Probability: Theory and Examples*. (£49.99 hardback).

D. Williams *Probability with Martingales*. Cambridge University Press 1991 (£26.99 paperback).

1. MEASURES

1.1. Definitions. Let E be a set. A σ -algebra \mathcal{E} on E is a set of subsets of E , containing the empty set \emptyset and such that, for all $A \in \mathcal{E}$ and all sequences $(A_n : n \in \mathbb{N})$ in \mathcal{E} ,

$$A^c \in \mathcal{E}, \quad \bigcup_n A_n \in \mathcal{E}.$$

The pair (E, \mathcal{E}) is called a *measurable space*. Given (E, \mathcal{E}) , each $A \in \mathcal{E}$ is called a *measurable set*.

A *measure* μ on (E, \mathcal{E}) is a function $\mu : \mathcal{E} \rightarrow [0, \infty]$, with $\mu(\emptyset) = 0$, such that, for any sequence $(A_n : n \in \mathbb{N})$ of disjoint elements of \mathcal{E} ,

$$\mu\left(\bigcup_n A_n\right) = \sum_n \mu(A_n).$$

The triple (E, \mathcal{E}, μ) is called a *measure space*. If $\mu(E) = 1$ then μ is a *probability measure* and (E, \mathcal{E}, μ) is a *probability space*. The notation $(\Omega, \mathcal{F}, \mathbb{P})$ is often used to denote a probability space.

1.2. Discrete measure theory. Let E be a countable set and let $\mathcal{E} = \mathcal{P}(E)$. A *mass function* is any function $m : E \rightarrow [0, \infty]$. If μ is a measure on (E, \mathcal{E}) , then, by countable additivity,

$$\mu(A) = \sum_{x \in A} \mu(\{x\}), \quad A \subseteq E.$$

So there is a one-to-one correspondence between measures and mass functions, given by

$$m(x) = \mu(\{x\}), \quad \mu(A) = \sum_{x \in A} m(x).$$

This sort of measure space provides a ‘toy’ version of the general theory, where each of the results we prove for general measure spaces reduces to some straightforward fact about the convergence of series. This is all one needs to do elementary discrete probability and discrete-time Markov chains, so these topics are usually introduced without discussing measure theory.

Discrete measure theory is essentially the only context where one can define a measure explicitly, because, in general, σ -algebras are not amenable to an explicit presentation which would allow us to make such a definition. Instead one specifies the values to be taken on some smaller set of subsets, which generates the σ -algebra. This gives rise to two problems: first to know that there is a measure extending the given set function, second to know that there is not more than one. The first problem, which is one of construction, is often dealt with by Carathéodory’s extension theorem. The second problem, that of uniqueness, is often dealt with by Dynkin’s π -system lemma.

1.3. Generated σ -algebras. Let \mathcal{A} be a set of subsets of E . Define

$$\sigma(\mathcal{A}) = \{A \subseteq E : A \in \mathcal{E} \text{ for all } \sigma\text{-algebras } \mathcal{E} \text{ containing } \mathcal{A}\}.$$

Then $\sigma(\mathcal{A})$ is a σ -algebra, which is called the σ -algebra generated by \mathcal{A} . It is the smallest σ -algebra containing \mathcal{A} .

1.4. π -systems and d -systems. Let \mathcal{A} be a set of subsets of E . Say that \mathcal{A} is a π -system if $\emptyset \in \mathcal{A}$ and, for all $A, B \in \mathcal{A}$,

$$A \cap B \in \mathcal{A}.$$

Say that \mathcal{A} is a d -system if $E \in \mathcal{A}$ and, for all $A, B \in \mathcal{A}$ with $A \subseteq B$ and all increasing sequences $(A_n : n \in \mathbb{N})$ in \mathcal{A} ,

$$B \setminus A \in \mathcal{A}, \quad \bigcup_n A_n \in \mathcal{A}.$$

Note that, if \mathcal{A} is both a π -system and a d -system, then \mathcal{A} is a σ -algebra.

Lemma 1.4.1 (Dynkin's lemma). *Let \mathcal{A} be a π -system. Then any d -system containing \mathcal{A} contains also $\sigma(\mathcal{A})$, the σ -algebra generated by \mathcal{A} .*

Proof. Denote by \mathcal{D} the intersection of all d -systems containing \mathcal{A} . Then \mathcal{D} is itself a d -system. We shall show that \mathcal{D} is also a π -system and hence a σ -algebra, thus proving the lemma. We need to show that for every $A, B \in \mathcal{D}$, then $A \cap B \in \mathcal{D}$. This is certainly true if $A, B \in \mathcal{A}$, as \mathcal{A} is a π -system. We proceed to extend this in two steps.

Fix $B \in \mathcal{A}$ and consider $\mathcal{A}_B = \{A \subset E, A \cap B \in \mathcal{D}\}$. The \mathcal{A}_B is a d -system containing \mathcal{A} , and hence contains \mathcal{D} . Thus $A \cap B \in \mathcal{D}$ for all $A \in \mathcal{D}, B \in \mathcal{A}$. Now fix $A \in \mathcal{D}$ and consider $\mathcal{B}_A = \{B \subset E, B \cap A \in \mathcal{D}\}$. Then \mathcal{B}_A is a d -system, and by the above contains \mathcal{A} . Thus \mathcal{B}_A contains \mathcal{D} , and we are done. \square

1.5. Uniqueness of measures.

Theorem 1.5.1 (Uniqueness of extension). *Let μ_1, μ_2 be measures on (E, \mathcal{E}) with $\mu_1(E) = \mu_2(E) < \infty$. Suppose that $\mu_1 = \mu_2$ on \mathcal{A} , for some π -system \mathcal{A} generating \mathcal{E} . Then $\mu_1 = \mu_2$ on \mathcal{E} .*

Proof. We first need the following important lemma.

Lemma 1.5.2. *Let μ be a measure on (E, \mathcal{E}) . Let A_n be an increasing sequence of \mathcal{E} , and let $A = \bigcup_n A_n$. Then $\mu(A) = \lim_n \mu(A_n)$.*

Proof. Let $B_1 = A_1, B_2 = A_2 \setminus A_1, \dots, B_n = A_n \setminus A_{n-1} \in \mathcal{E}$. Moreover, the B_n are disjoint, and $\bigcup_{i=1}^n B_i = \bigcup_{i=1}^n A_i$ for all $n \geq 1$. Thus by σ -additivity, we get

$$\mu(A) = \mu\left(\bigcup_{i=1}^{\infty} B_i\right) = \sum_{i=1}^{\infty} \mu(B_i) = \lim_n \sum_{i=1}^n \mu(B_i) = \lim_n \mu(A_n)$$

\square

Consider $\mathcal{D} = \{A \in \mathcal{E} : \mu_1(A) = \mu_2(A)\}$. By hypothesis, $E \in \mathcal{D}$; for $A, B \in \mathcal{E}$ with $A \subseteq B$, we have

$$\mu_1(A) + \mu_1(B \setminus A) = \mu_1(B) < \infty, \quad \mu_2(A) + \mu_2(B \setminus A) = \mu_2(B) < \infty$$

so, if $A, B \in \mathcal{D}$, then also $B \setminus A \in \mathcal{D}$; if $A_n \in \mathcal{D}, n \in \mathbb{N}$, with $A_n \uparrow A$, then

$$\mu_1(A) = \lim_n \mu_1(A_n) = \lim_n \mu_2(A_n) = \mu_2(A)$$

so $A \in \mathcal{D}$. Thus \mathcal{D} is a d -system containing the π -system \mathcal{A} , so $\mathcal{D} = \mathcal{E}$ by Dynkin's lemma. \square

1.6. Set functions and properties. Let \mathcal{A} be any set of subsets of E containing the empty set \emptyset . A *set function* is a function $\mu : \mathcal{A} \rightarrow [0, \infty]$ with $\mu(\emptyset) = 0$. Let μ be a set function. Say that μ is *increasing* if, for all $A, B \in \mathcal{A}$ with $A \subseteq B$,

$$\mu(A) \leq \mu(B).$$

Say that μ is *additive* if, for all disjoint sets $A, B \in \mathcal{A}$ with $A \cup B \in \mathcal{A}$,

$$\mu(A \cup B) = \mu(A) + \mu(B).$$

Say that μ is *countably additive* if, for all sequences of disjoint sets $(A_n : n \in \mathbb{N})$ in \mathcal{A} with $\bigcup_n A_n \in \mathcal{A}$,

$$\mu\left(\bigcup_n A_n\right) = \sum_n \mu(A_n).$$

Say that μ is *countably subadditive* if, for all sequences $(A_n : n \in \mathbb{N})$ in \mathcal{A} with $\bigcup_n A_n \in \mathcal{A}$,

$$\mu\left(\bigcup_n A_n\right) \leq \sum_n \mu(A_n).$$

1.7. Construction of measures. Let \mathcal{A} be a set of subsets of E . Say that \mathcal{A} is a *ring* on E if $\emptyset \in \mathcal{A}$ and, for all $A, B \in \mathcal{A}$,

$$B \setminus A \in \mathcal{A}, \quad A \cup B \in \mathcal{A}.$$

Let \mathcal{A} be a ring. Since the symmetric difference may be written

$$A \Delta B = (A \setminus B) \cup (B \setminus A),$$

then for all $A, B \in \mathcal{A}$, we get $A \Delta B \in \mathcal{A}$. Thus, $A \cap B = (A \cup B) \setminus (A \Delta B) \in \mathcal{A}$, so any ring is also stable by intersection.

Theorem 1.7.1 (Carathéodory's extension theorem). *Let \mathcal{A} be a ring of subsets of E and let $\mu : \mathcal{A} \rightarrow [0, \infty]$ be a countably additive set function. Then μ extends to a measure on the σ -algebra generated by \mathcal{A} .*

Proof. For any $B \subseteq E$, define the *outer measure*

$$\mu^*(B) = \inf \sum_n \mu(A_n)$$

where the infimum is taken over all sequences $(A_n : n \in \mathbb{N})$ in \mathcal{A} such that $B \subseteq \bigcup_n A_n$ and is taken to be ∞ if there is no such sequence. Note that μ^* is increasing and $\mu^*(\emptyset) = 0$. Let us say that $A \subseteq E$ is μ^* -*measurable* if, for all $B \subseteq E$,

$$\mu^*(B) = \mu^*(B \cap A) + \mu^*(B \cap A^c).$$

Write \mathcal{M} for the set of all μ^* -measurable sets. We shall show that \mathcal{M} is a σ -algebra containing \mathcal{A} and that μ^* is a measure on \mathcal{M} , extending μ . This will prove the theorem.

Step I. We show that μ^* is countably subadditive. Suppose that $B \in E$, $B_n \in E$ and $B \subseteq \bigcup_n B_n$. If $\mu^*(B_n) < \infty$ for all n , then, given $\varepsilon > 0$, there exist sequences $(A_{nm} : m \in \mathbb{N})$ in \mathcal{A} , with

$$B_n \subseteq \bigcup_m A_{nm}, \quad \mu^*(B_n) + \varepsilon/2^n \geq \sum_m \mu(A_{nm}).$$

Then

$$B \subseteq \bigcup_n \bigcup_m A_{nm}$$

so

$$\mu^*(B) \leq \sum_n \sum_m \mu(A_{nm}) \leq \sum_n \mu^*(B_n) + \varepsilon.$$

Hence, in any case,

$$\mu^*(B) \leq \sum_n \mu^*(B_n).$$

Step II. We show that μ^* extends μ . We need the following important lemma.

Lemma 1.7.2. *Let μ be a countably additive set function on a ring \mathcal{A} . Then μ is countably subadditive and increasing.*

Proof. μ is increasing since μ is additive and \mathcal{A} is a ring. Now assume that $A_n \in \mathcal{A}$ and $A = \bigcup_n A_n \in \mathcal{A}$. Let $B_1 = A_1$, $B_2 = A_2 \setminus A_1, \dots, B_n = A_n \setminus B_{n-1}$. By induction, $B_n \in \mathcal{A}$ since \mathcal{A} is a ring. Moreover, the B_n are disjoint, and $\bigcup_{i=1}^n B_i = \bigcup_{i=1}^n A_i$. Thus by countable additivity, we get

$$\mu(A) = \mu\left(\bigcup_{i=1}^{\infty} B_i\right) = \sum_{i=1}^{\infty} \mu(B_i) \leq \sum_{i=1}^{\infty} \mu(A_i)$$

since μ is increasing. □

Hence, for $A \in \mathcal{A}$ and any sequence $(A_n : n \in \mathbb{N})$ in \mathcal{A} with $A \subseteq \bigcup_n A_n$, then $A \subset \bigcup_n (A_n \cap A) = A \in \mathcal{A}$, so

$$\mu(A) \leq \sum_n \mu(A \cap A_n) \leq \sum_n \mu(A_n).$$

On taking the infimum over all such sequences, we see that $\mu(A) \leq \mu^*(A)$. On the other hand, it is obvious that $\mu^*(A) \leq \mu(A)$ for $A \in \mathcal{A}$.

Step III. We show that \mathcal{M} contains \mathcal{A} . Let $A \in \mathcal{A}$ and $B \subseteq E$. We have to show that

$$\mu^*(B) = \mu^*(B \cap A) + \mu^*(B \cap A^c).$$

By subadditivity of μ^* , it is enough to show that

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

If $\mu^*(B) = \infty$, this is clearly true, so let us assume that $\mu^*(B) < \infty$. Then, given $\varepsilon > 0$, we can find a sequence $(A_n : n \in \mathbb{N})$ in \mathcal{A} such that

$$B \subseteq \bigcup_n A_n, \quad \mu^*(B) + \varepsilon \geq \sum_n \mu(A_n).$$

Then

$$B \cap A \subseteq \bigcup_n (A_n \cap A), \quad B \cap A^c \subseteq \bigcup_n (A_n \cap A^c)$$

so, since $A_n \cap A \in \mathcal{A}$ and $A_n \cap A^c = A_n \setminus A \in \mathcal{A}$,

$$\mu^*(B \cap A) + \mu^*(B \cap A^c) \leq \sum_n \mu(A_n \cap A) + \sum_n \mu(A_n \cap A^c) = \sum_n \mu(A_n) \leq \mu^*(B) + \varepsilon.$$

Since $\varepsilon > 0$ was arbitrary, we are done.

Step IV. We show that \mathcal{M} is a σ -algebra and that μ^ is a measure on \mathcal{M} .* Clearly $E \in \mathcal{M}$ and $A^c \in \mathcal{M}$ whenever $A \in \mathcal{M}$. Suppose that $A_1, A_2 \in \mathcal{M}$ and $B \subseteq E$. Then

$$\begin{aligned} \mu^*(B) &= \mu^*(B \cap A_1) + \mu^*(B \cap A_1^c) \\ &= \mu^*(B \cap A_1 \cap A_2) + \mu^*(B \cap A_1 \cap A_2^c) + \mu^*(B \cap A_1^c) \\ &= \mu^*(B \cap A_1 \cap A_2) + \mu^*(B \cap (A_1 \cap A_2)^c \cap A_1) + \mu^*(B \cap (A_1 \cap A_2)^c \cap A_1^c) \\ &= \mu^*(B \cap (A_1 \cap A_2)) + \mu^*(B \cap (A_1 \cap A_2)^c). \end{aligned}$$

Hence $A_1 \cap A_2 \in \mathcal{M}$. Since \mathcal{M} is stable by complement and by taking (finite) intersections, it is also stable by taking (finite) unions. We now show that, for any sequence of disjoint sets $(A_n : n \in \mathbb{N})$ in \mathcal{M} , for $A = \bigcup_n A_n$ we have

$$A \in \mathcal{M}, \quad \mu^*(A) = \sum_n \mu^*(A_n).$$

So, take any $B \subseteq E$, then

$$\begin{aligned}\mu^*(B) &= \mu^*(B \cap A_1) + \mu^*(B \cap A_1^c) \\ &= \mu^*(B \cap A_1) + \mu^*(B \cap A_2) + \mu^*(B \cap A_1^c \cap A_2^c) \\ &= \cdots = \sum_{i=1}^n \mu^*(B \cap A_i) + \mu^*(B \cap A_1^c \cap \cdots \cap A_n^c).\end{aligned}$$

Note that $\mu^*(B \cap A_1^c \cap \cdots \cap A_n^c) \geq \mu^*(B \cap A^c)$ for all n . Hence, on letting $n \rightarrow \infty$ we get

$$\begin{aligned}\mu^*(B) &\geq \sum_{n=1}^{\infty} \mu^*(B \cap A_n) + \mu^*(B \cap A^c) \\ &\geq \mu^*(B \cap A) + \mu^*(B \cap A^c),\end{aligned}$$

by using countable subadditivity. The reverse inequality always holds (by subadditivity), so we have equality. Hence $A \in \mathcal{M}$ and, setting $B = A$, we get

$$\mu^*(A) = \sum_{n=1}^{\infty} \mu^*(A_n).$$

□

1.8. Borel sets and measures. Let E be a topological space. The σ -algebra generated by the set of open sets in E is called the *Borel σ -algebra* of E and is denoted $\mathcal{B}(E)$. The Borel σ -algebra of \mathbb{R} is denoted simply by \mathcal{B} . A measure μ on $(E, \mathcal{B}(E))$ is called a *Borel measure* on E . If moreover $\mu(K) < \infty$ for all compact sets K , then μ is called a *Radon measure* on E .

1.9. Lebesgue measure.

Theorem 1.9.1. *There exists a unique Borel measure μ on \mathbb{R} such that, for all $a, b \in \mathbb{R}$ with $a < b$,*

$$\mu((a, b]) = b - a.$$

The measure μ is called *Lebesgue measure* on \mathbb{R} .

Proof. (Existence.) Consider the ring \mathcal{A} of finite unions of disjoint intervals of the form

$$A = (a_1, b_1] \cup \cdots \cup (a_n, b_n].$$

We note that \mathcal{A} generates \mathcal{B} . Define for such $A \in \mathcal{A}$

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

Note that the presentation of A is not unique, as $(a, b] \cup (b, c] = (a, c]$ whenever $a < b < c$. Nevertheless, it is easy to check that μ is well-defined and additive. We aim to show that μ is countably additive on \mathcal{A} , which then proves existence by Carathéodory's extension theorem.

By additivity, it suffices to show that, if $A \in \mathcal{A}$ and if $(A_n : n \in \mathbb{N})$ is an increasing sequence in \mathcal{A} with $A_n \uparrow A$, then $\mu(A_n) \rightarrow \mu(A)$. Set $B_n = A \setminus A_n$ then $B_n \in \mathcal{A}$ and $B_n \downarrow \emptyset$. By additivity again, it suffices to show that $\mu(B_n) \rightarrow 0$. Suppose, in fact, that for some $\varepsilon > 0$, we have $\mu(B_n) \geq 2\varepsilon$ for all n . For each n we can find $C_n \in \mathcal{A}$ with $\bar{C}_n \subseteq B_n$ and $\mu(B_n \setminus C_n) \leq \varepsilon 2^{-n}$. Then

$$\mu(B_n \setminus (C_1 \cap \cdots \cap C_n)) \leq \mu((B_1 \setminus C_1) \cup \cdots \cup (B_n \setminus C_n)) \leq \sum_{n \in \mathbb{N}} \varepsilon 2^{-n} = \varepsilon.$$

Since $\mu(B_n) \geq 2\varepsilon$, we must have $\mu(C_1 \cap \cdots \cap C_n) \geq \varepsilon$, so $C_1 \cap \cdots \cap C_n \neq \emptyset$, and so $K_n = \bar{C}_1 \cap \cdots \cap \bar{C}_n \neq \emptyset$. Now $(K_n : n \in \mathbb{N})$ is a decreasing sequence of bounded non-empty closed sets in \mathbb{R} , so $\emptyset \neq \bigcap_n K_n \subseteq \bigcap_n B_n$, which is a contradiction.

(Uniqueness.) Let λ be any measure on \mathcal{B} with $\mu((a, b]) = b - a$ for all $a < b$. Fix n and consider

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

Then μ_n and λ_n are probability measures on \mathcal{B} and $\mu_n = \lambda_n$ on the π -system of intervals of the form $(a, b]$, which generates \mathcal{B} . So, by Theorem 1.5.1, $\mu_n = \lambda_n$ on \mathcal{B} . Hence, for all $A \in \mathcal{B}$, we have

$$\mu(A) = \sum_n \mu_n(A) = \sum_n \lambda_n(A) = \lambda(A).$$

□

1.10. Existence of a non-Lebesgue-measurable subset of \mathbb{R} . For $x, y \in [0, 1)$, let us write $x \sim y$ if $x - y \in \mathbb{Q}$. Then \sim is an equivalence relation. Using the Axiom of Choice, we can find a subset S of $[0, 1)$ containing exactly one representative of each equivalence class. Set $Q = \mathbb{Q} \cap [0, 1)$ and, for each $q \in Q$, define $S_q = S + q = \{s + q \pmod{1} : s \in S\}$. It is an easy exercise to check that the sets S_q are all disjoint and their union is $[0, 1)$.

Now, Lebesgue measure μ on $\mathcal{B} = \mathcal{B}([0, 1))$ is translation invariant. That is to say, $\mu(B) = \mu(B + x)$ for all $B \in \mathcal{B}$ and all $x \in [0, 1)$. If S were a Borel set, then we would have

$$1 = \mu([0, 1)) = \sum_{q \in Q} \mu(S + q) = \sum_{q \in Q} \mu(S)$$

which is impossible. Hence $S \notin \mathcal{B}$.

A *Lebesgue measurable* set in \mathbb{R} is any set of the form $A \cup N$, with A Borel and $N \subseteq B$ for some Borel set B with $\mu(B) = 0$. Thus the set of Lebesgue measurable sets is the *completion* of the Borel σ -algebra with respect to μ . See Exercise 1.8. The same argument shows that S cannot be Lebesgue measurable either.

1.11. **Independence.** A probability space $(\Omega, \mathcal{F}, \mathbb{P})$ provides a model for an experiment whose outcome is subject to chance, according to the following interpretation:

- Ω is the set of possible outcomes
- \mathcal{F} is the set of observable sets of outcomes, or *events*
- $\mathbb{P}(A)$ is the probability of the event A .

Recall that $\mathbb{P}(\Omega) = 1$. In addition, we say that an event $A \in \mathcal{F}$ occurs *almost surely* (often abbreviated *a.s.*) if $\mathbb{P}(A) = 1$.

Example. Consider $\Omega = [0, 2\pi)$, \mathcal{F} the Borel σ -algebra, and $\mathbb{P}(A) = \mu(A)/2\pi$ where μ is the Lebesgue measure. Then $(\Omega, \mathcal{F}, \mathbb{P})$ may serve as a probability space for choosing a point uniformly at random on the unit circle. E.g., if $A = [0, 2\pi) \setminus \mathbb{Q}$ then $\mathbb{P}(A) = 1$ so almost surely, the angle of the point is irrational.

Example. Infinite coin-tossing. Let

$$\Omega = \{0, 1\}^{\mathbb{N}} = \{\omega = (\omega_1, \dots) \text{ with } \omega_n \in \{0, 1\} \text{ for all } n \geq 1\}.$$

We interpret the event $\{\omega_n = 1\}$ as the event that the n th toss results in a heads. We take for the σ -algebra \mathcal{F} the σ -algebra generated by events of the form

$$A = \{\omega_1 = \varepsilon_1, \dots, \omega_n = \varepsilon_n\}$$

for any fixed sequence $(\varepsilon_1, \dots) \in \{0, 1\}^{\mathbb{N}}$. We will soon define a probability measure \mathbb{P} on (Ω, \mathcal{F}) under which the outcomes of the tosses are *independent*.

Relative to measure theory, probability theory is enriched by the significance attached to the notion of independence. Let I be a countable set. Say that events $A_i, i \in I$, are *independent* if, for all finite subsets $J \subseteq I$,

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

Say that σ -algebras $\mathcal{A}_i \subseteq \mathcal{F}, i \in I$, are *independent* if $A_i, i \in I$, are independent whenever $A_i \in \mathcal{A}_i$ for all i . Here is a useful way to establish the independence of two σ -algebras.

Theorem 1.11.1. Let \mathcal{A}_1 and \mathcal{A}_2 be π -systems contained in \mathcal{F} and suppose that

$$\mathbb{P}(A_1 \cap A_2) = \mathbb{P}(A_1)\mathbb{P}(A_2)$$

whenever $A_1 \in \mathcal{A}_1$ and $A_2 \in \mathcal{A}_2$. Then $\sigma(\mathcal{A}_1)$ and $\sigma(\mathcal{A}_2)$ are independent.

Proof. Fix $A_1 \in \mathcal{A}_1$ and define for $A \in \mathcal{F}$

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

Then μ and ν are measures which agree on the π -system \mathcal{A}_2 , with $\mu(\Omega) = \nu(\Omega) = \mathbb{P}(A_1) < \infty$. So, by uniqueness of extension, for all $A_2 \in \sigma(\mathcal{A}_2)$,

$$\mathbb{P}(A_1 \cap A_2) = \mu(A_2) = \nu(A_2) = \mathbb{P}(A_1)\mathbb{P}(A_2).$$

Now fix $A_2 \in \sigma(\mathcal{A}_2)$ and repeat the argument with

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

to show that, for all $A_1 \in \sigma(\mathcal{A}_1)$, $\mathbb{P}(A_1 \cap A_2) = \mathbb{P}(A_1)\mathbb{P}(A_2)$. \square

1.12. Borel-Cantelli lemmas. Given a sequence of events $(A_n : n \in \mathbb{N})$, we may ask for the probability that infinitely many occur. Set

$$\limsup A_n = \bigcap_n \bigcup_{m \geq n} A_m, \quad \liminf A_n = \bigcup_n \bigcap_{m \geq n} A_m.$$

We sometimes write $\{A_n \text{ infinitely often}\}$ as an alternative for $\limsup A_n$, because $\omega \in \limsup A_n$ if and only if $\omega \in A_n$ for infinitely many n . Similarly, we write $\{A_n \text{ eventually}\}$ for $\liminf A_n$. The abbreviations i.o. and ev. are often used. Note that

$$\{A_n \text{ i.o.}\}^c = \{A_n^c \text{ ev.}\}.$$

Lemma 1.12.1 (First Borel-Cantelli lemma). *If $\sum_n \mathbb{P}(A_n) < \infty$, then $\mathbb{P}(A_n \text{ i.o.}) = 0$. Hence A_n^c holds eventually, almost surely.*

Proof. As $n \rightarrow \infty$ we have

$$\mathbb{P}(A_n \text{ i.o.}) \leq \mathbb{P}\left(\bigcup_{m \geq n} A_m\right) \leq \sum_{m \geq n} \mathbb{P}(A_m) \rightarrow 0.$$

\square

We note that this argument is valid whether or not \mathbb{P} is a probability measure.

Lemma 1.12.2 (Second Borel-Cantelli lemma). *Assume that the events $(A_n : n \in \mathbb{N})$ are independent. If $\sum_n \mathbb{P}(A_n) = \infty$, then $\mathbb{P}(A_n \text{ i.o.}) = 1$. Hence, almost surely, A_n holds infinitely often.*

Proof. We use the inequality $1 - a \leq e^{-a}$. Set $a_n = \mathbb{P}(A_n)$. Then, for all n we have

$$\mathbb{P}\left(\bigcap_{m \geq n} A_m^c\right) = \prod_{m \geq n} (1 - a_m) \leq \exp\left\{-\sum_{m \geq n} a_m\right\} = 0.$$

Hence $\mathbb{P}(A_n \text{ i.o.}) = 1 - \mathbb{P}\left(\bigcup_n \bigcap_{m \geq n} A_m^c\right) = 1$. \square

Example. Infinite coin-toss. Let $0 < p < 1$. We will soon construct a probability measure \mathbb{P} on the infinite coin-toss measurable space (Ω, \mathcal{F}) such that the events $\{\omega_n = 1\}, n \in \mathbb{N}$ are independent and for all $n \geq 1$, $\mathbb{P}(\omega_n = 1) = p$. Then

$$\sum_n \mathbb{P}(\omega_n = 1) = \infty,$$

hence, since these events are independent, we may apply the second Borel-Cantelli lemma: almost surely, there are infinitely many heads in the coin-tossing experiment (no matter how small p is).

2. MEASURABLE FUNCTIONS AND RANDOM VARIABLES

2.1. Measurable functions. Let (E, \mathcal{E}) and (G, \mathcal{G}) be measurable spaces. A function $f : E \rightarrow G$ is *measurable* if $f^{-1}(A) \in \mathcal{E}$ whenever $A \in \mathcal{G}$. Here $f^{-1}(A)$ denotes the *inverse image* of A by f

$$f^{-1}(A) = \{x \in E : f(x) \in A\}.$$

Usually $G = \mathbb{R}$ or $G = [-\infty, \infty]$, in which case \mathcal{G} is always taken to be the Borel σ -algebra. If E is a topological space and $\mathcal{E} = \mathcal{B}(E)$, then a measurable function on E is called a *Borel function*. For any function $f : E \rightarrow G$, the inverse image preserves set operations

$$f^{-1}\left(\bigcup_i A_i\right) = \bigcup_i f^{-1}(A_i), \quad f^{-1}(B \setminus A) = f^{-1}(B) \setminus f^{-1}(A).$$

Therefore, the set $\{f^{-1}(A) : A \in \mathcal{G}\}$ is a σ -algebra on E and $\{A \subseteq G : f^{-1}(A) \in \mathcal{E}\}$ is a σ -algebra on G .

In particular, if $\mathcal{G} = \sigma(\mathcal{A})$ and $f^{-1}(A) \in \mathcal{E}$ whenever $A \in \mathcal{A}$, then $\{A : f^{-1}(A) \in \mathcal{E}\}$ is a σ -algebra containing \mathcal{A} and hence \mathcal{G} , so f is measurable. In the case $G = \mathbb{R}$, the Borel σ -algebra is generated by intervals of the form $(-\infty, y]$, $y \in \mathbb{R}$, so, to show that $f : E \rightarrow \mathbb{R}$ is Borel measurable, it suffices to show that $\{x \in E : f(x) \leq y\} \in \mathcal{E}$ for all y .

If E is any topological space and $f : E \rightarrow \mathbb{R}$ is continuous, then $f^{-1}(U)$ is open in E and hence measurable, whenever U is open in \mathbb{R} ; the open sets U generate \mathcal{B} , so *any continuous function is measurable*.

For $A \subseteq E$, the *indicator function* 1_A of A is the function $1_A : E \rightarrow \{0, 1\}$ which takes the value 1 on A and 0 otherwise. Note that *the indicator function of any measurable set is a measurable function*. Also, *the composition of measurable functions is measurable*.

Given any family of functions $f_i : E \rightarrow G, i \in I$, we can make them all measurable by taking

$$\mathcal{E} = \sigma(f_i^{-1}(A) : A \in \mathcal{G}, i \in I).$$

Then \mathcal{E} is the σ -algebra generated by $(f_i : i \in I)$, and is often denoted by $\sigma(F_i, i \in I)$.

Proposition 2.1.1. *Let $f_n : E \rightarrow \mathbb{R}, n \in \mathbb{N}$, be measurable functions. Then so are $f_1 + f_2, f_1 f_2$ and each of the following:*

$$\inf_n f_n, \quad \sup_n f_n, \quad \liminf_n f_n, \quad \limsup_n f_n.$$

Proof. Note that $\{f_1 + f_2 < y\} = \bigcup_{q \in \mathbb{Q}} \{f_1 < q\} \cap \{f_2 < y - q\}$, hence $f_1 + f_2$ is measurable. Using the identity $ab = (1/4)((a + b)^2 - (a - b)^2)$, it suffices to prove that f^2 is measurable whenever f is, which is easy to see. The rest is left as an exercise. \square

Theorem 2.1.2 (Monotone class theorem). *Let (E, \mathcal{E}) be a measurable space and let \mathcal{A} be a π -system generating \mathcal{E} . Let \mathcal{V} be a vector space of bounded functions $f : E \rightarrow \mathbb{R}$ such that:*

- (i) $1 \in \mathcal{V}$ and $1_A \in \mathcal{V}$ for all $A \in \mathcal{A}$;
- (ii) if $f_n \in \mathcal{V}$ for all n and f is bounded with $0 \leq f_n \uparrow f$, then $f \in \mathcal{V}$.

Then \mathcal{V} contains every bounded measurable function.

Proof. Consider $\mathcal{D} = \{A \in \mathcal{E} : 1_A \in \mathcal{V}\}$. Then \mathcal{D} is a d -system containing \mathcal{A} , so $\mathcal{D} = \mathcal{E}$. Since \mathcal{V} is a vector space, it thus contains all finite linear combinations of indicator functions of measurable sets. If f is a bounded and non-negative measurable function, then the functions $f_n = 2^{-n} \lfloor 2^n f \rfloor \wedge n, n \in \mathbb{N}$, belong to \mathcal{V} and $0 \leq f_n \uparrow f$, so $f \in \mathcal{V}$. Finally, any bounded measurable function is the difference of two non-negative such functions, hence in \mathcal{V} . \square

2.2. Image measures. Let (E, \mathcal{E}) and (G, \mathcal{G}) be measurable spaces and let μ be a measure on \mathcal{E} . Then any measurable function $f : E \rightarrow G$ induces an *image measure* $\nu = \mu \circ f^{-1}$ on \mathcal{G} , given by

$$\nu(A) = \mu(f^{-1}(A)).$$

We shall construct some new measures from Lebesgue measure in this way.

Lemma 2.2.1. *Let $g : \mathbb{R} \rightarrow \mathbb{R}$ be non-constant, right-continuous and non-decreasing. Let $I = (g(-\infty), g(\infty))$ and define $f : I \rightarrow \mathbb{R}$ by $f(x) = \inf\{y \in \mathbb{R} : x \leq g(y)\}$. Then f is left-continuous and non-decreasing. Moreover, for $x \in I$ and $y \in \mathbb{R}$,*

$$f(x) \leq y \quad \text{if and only if} \quad x \leq g(y).$$

Proof. Fix $x \in I$ and consider the set $J_x = \{y \in \mathbb{R} : x \leq g(y)\}$. Note that J_x is non-empty and is not the whole of \mathbb{R} . Since g is non-decreasing, if $y \in J_x$ and $y' \geq y$, then $y' \in J_x$. Since g is right-continuous, if $y_n \in J_x$ and $y_n \downarrow y$, then $y \in J_x$. Hence $J_x = [f(x), \infty)$ and $x \leq g(y)$ if and only if $f(x) \leq y$. For $x \leq x'$, we have $J_x \supseteq J_{x'}$ and so $f(x) \leq f(x')$. For $x_n \uparrow x$, we have $J_x = \bigcap_n J_{x_n}$, so $f(x_n) \rightarrow f(x)$. (Indeed, $\ell = \lim f(x_n)$ exists by monotonicity, and then $\bigcap_n [f(x_n), \infty) = [\ell, \infty)$.) So f is left-continuous and non-decreasing, as claimed. \square

Theorem 2.2.2. *Let $g : \mathbb{R} \rightarrow \mathbb{R}$ be non-constant, right-continuous and non-decreasing. Then there exists a unique Radon measure dg on \mathbb{R} such that, for all $a, b \in \mathbb{R}$ with $a < b$,*

$$dg((a, b]) = g(b) - g(a).$$

Moreover, we obtain in this way all non-zero Radon measures on \mathbb{R} .

The measure dg is called the *Lebesgue-Stieltjes measure* associated with g .

Proof. Define I and f as in the lemma and let μ denote Lebesgue measure on I . Then f is Borel measurable, since for all $y \in \mathbb{R}$,

$$\{x \in I : f(x) \leq y\} = \{x \in I : x \leq g(y)\} = I \cap (-\infty, g(y)].$$

The induced measure $dg = \mu \circ f^{-1}$ on \mathbb{R} satisfies

$$dg((a, b]) = \mu(\{x : f(x) > a \text{ and } f(x) \leq b\}) = \mu((g(a), g(b)]) = g(b) - g(a).$$

The argument used for uniqueness of Lebesgue measure shows that there is at most one Borel measure with this property. Finally, if ν is any Radon measure on \mathbb{R} , we can define $g : \mathbb{R} \rightarrow \mathbb{R}$, by

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

Then g is nondecreasing and nonconstant, since ν is not the zero measure. Moreover it is right-continuous (see exercise 1.4). Then $\nu((a, b]) = g(b) - g(a)$ whenever $a < b$, so $\nu = dg$ by uniqueness. \square

Example. Let $\alpha > 0$, and let $g(x) = \alpha x$. Then g satisfies the conditions of Theorem 2.2.2, and the Lebesgue-Stieltjes measure associated with it is

$$dg = \alpha \mu,$$

where μ is Lebesgue measure.

Example. Dirac mass. Let $x \in \mathbb{R}$, and $k > 0$. The Dirac mass at the point x , with mass k is the Borel measure on \mathbb{R} defined by

$$\mu(A) = k \text{ if } x \in A$$

and $\mu(A) = 0$ else, for all Borel set A . Then μ is a Radon measure. Moreover, it is the Radon measure dg associated with the function:

$$g(t) = k1_{[x, \infty)}(t).$$

If $k = 1$, then μ is simply denoted by $\delta_{\{x\}}$, and hence in general μ is written $k\delta_{\{x\}}$. We thus have

$$\delta_{\{x\}} = dg,$$

where g is the Heavyside function at x .

2.3. Random variables. Let $(\Omega, \mathcal{F}, \mathbb{P})$ be a probability space and let (E, \mathcal{E}) be a measurable space. A measurable function $X : \Omega \rightarrow E$ is called a *random variable in E* . It has the interpretation of a quantity, or state, determined by chance. Where no space E is mentioned, it is assumed that X takes values in \mathbb{R} . The image measure $\mu_X = \mathbb{P} \circ X^{-1}$ is called the *law* or *distribution* of X . For real-valued random variables, μ_X is uniquely determined by its values on the π -system of intervals $(-\infty, x], x \in \mathbb{R}$, given by

$$F_X(x) = \mu_X((-\infty, x]) = \mathbb{P}(X \leq x).$$

The function F_X is called the *distribution function of X* .

Note that $F = F_X$ is increasing and right-continuous, with

$$\lim_{x \rightarrow -\infty} F(x) = 0, \quad \lim_{x \rightarrow \infty} F(x) = 1.$$

Let us call any function $F : \mathbb{R} \rightarrow [0, 1]$ satisfying these conditions a *distribution function*.

Set $\Omega = (0, 1]$ and $\mathcal{F} = \mathcal{B}((0, 1])$. Let \mathbb{P} denote the restriction of Lebesgue measure to \mathcal{F} . Then $(\Omega, \mathcal{F}, \mathbb{P})$ is a probability space. Let F be any distribution function. Define $X : \Omega \rightarrow \mathbb{R}$ by

$$X(\omega) = \inf\{x : \omega \leq F(x)\}.$$

Then, by Lemma 2.2.1, X is a random variable and $X(\omega) \leq x$ if and only if $\omega \leq F(x)$. So

$$F_X(x) = \mathbb{P}(X \leq x) = \mathbb{P}((0, F(x)]) = F(x).$$

Thus every distribution function is the distribution function of a random variable.

A countable family of random variables $(X_i : i \in I)$ is said to be *independent* if the σ -algebras $(\sigma(X_i) : i \in I)$ are independent. For a sequence $(X_n : n \in \mathbb{N})$ of real valued random variables, 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 $x_1, \dots, x_n \in \mathbb{R}$ and all n . A sequence of random variables $(X_n : n \geq 0)$ is often regarded as a *process* evolving in time. The σ -algebra generated by X_0, \dots, X_n

$$\mathcal{F}_n = \sigma(X_0, \dots, X_n)$$

contains those events depending (measurably) on X_0, \dots, X_n and represents what is known about the process by time n .

2.4. Rademacher functions. We continue with the particular choice of probability space $(\Omega, \mathcal{F}, \mathbb{P})$ made in the preceding section. Provided that we forbid infinite sequences of 0's, each $\omega \in \Omega$ has a unique binary expansion

$$\omega = 0.\omega_1\omega_2\omega_3\dots$$

Define random variables $R_n : \Omega \rightarrow \{0, 1\}$ by $R_n(\omega) = \omega_n$. Then

$$R_1 = 1_{(\frac{1}{2}, 1]}, \quad R_2 = 1_{(\frac{1}{4}, \frac{1}{2}]} + 1_{(\frac{3}{4}, 1]}, \quad R_3 = 1_{(\frac{1}{8}, \frac{1}{4}]} + 1_{(\frac{3}{8}, \frac{1}{2}]} + 1_{(\frac{5}{8}, \frac{3}{4}]} + 1_{(\frac{7}{8}, 1]}.$$

These are called the *Rademacher functions*. The random variables R_1, R_2, \dots are independent and *Bernoulli*, that is to say

$$\mathbb{P}(R_n = 0) = \mathbb{P}(R_n = 1) = 1/2.$$

Let $E = \{0, 1\}^{\mathbb{N}}$ and let \mathcal{E} denote the σ -algebra generated by events of the form $A = \{\omega \in E : \omega_1 = y_1, \dots, \omega_n = y_n\}, n \in \mathbb{N}, y_i \in \{0, 1\}$. Let $R = (R_1, \dots) \in E$. Then R is a random variable in (E, \mathcal{E}) . (*check it!*) The distribution μ of R is then a probability measure on (E, \mathcal{E}) , under which the events $\{\omega_n = y_n\}, n \in \mathbb{N}$, are independent. This is the *infinite (fair) coin-toss measure*.

We now use a trick involving the Rademacher functions to construct on $\Omega = (0, 1]$, not just one random variable, but an infinite sequence of independent random variables with given distribution functions.

Proposition 2.4.1. *Let $(\Omega, \mathcal{F}, \mathbb{P})$ be the probability space of Lebesgue measure on the Borel subsets of $(0, 1]$. Let $(F_n : n \in \mathbb{N})$ be a sequence of distribution functions. Then there exists a sequence $(X_n : n \in \mathbb{N})$ of independent random variables on $(\Omega, \mathcal{F}, \mathbb{P})$ such that X_n has distribution function $F_{X_n} = F_n$ for all n .*

Proof. Choose a bijection $m : \mathbb{N}^2 \rightarrow \mathbb{N}$ and set $Y_{k,n} = R_{m(k,n)}$, where R_m is the m th Rademacher function. Set

$$Y_n = \sum_{k=1}^{\infty} 2^{-k} Y_{k,n}.$$

Then Y_1, Y_2, \dots are independent and, for all n , for $i2^{-k} = 0.y_1 \dots y_k$, we have

$$\mathbb{P}(i2^{-k} < Y_n \leq (i+1)2^{-k}) = \mathbb{P}(Y_{1,n} = y_1, \dots, Y_{k,n} = y_k) = 2^{-k}$$

so $\mathbb{P}(Y_n \leq x) = x$ for all $x \in (0, 1]$. Set

$$G_n(y) = \inf\{x : y \leq F_n(x)\}$$

then, by Lemma 2.2.1, G_n is Borel and $G_n(y) \leq x$ if and only if $y \leq F_n(x)$. So, if we set $X_n = G_n(Y_n)$, then X_1, X_2, \dots are independent random variables on Ω and

$$\mathbb{P}(X_n \leq x) = \mathbb{P}(G_n(Y_n) \leq x) = \mathbb{P}(Y_n \leq F_n(x)) = F_n(x).$$

□

2.5. Tail events. Let $(X_n : n \in \mathbb{N})$ be a sequence of random variables. Define

$$\mathcal{T}_n = \sigma(X_{n+1}, X_{n+2}, \dots), \quad \mathcal{T} = \bigcap_n \mathcal{T}_n.$$

Then \mathcal{T} is a σ -algebra, called the *tail σ -algebra* of $(X_n : n \in \mathbb{N})$. It contains the events which depend only on the limiting behaviour of the sequence.

Theorem 2.5.1 (Kolmogorov's zero-one law). *Suppose that $(X_n : n \in \mathbb{N})$ is a sequence of independent random variables. Then the tail σ -algebra \mathcal{T} of $(X_n : n \in \mathbb{N})$ contains only events of probability 0 or 1. Moreover, any \mathcal{T} -measurable random variable is almost surely constant.*

Proof. Set $\mathcal{F}_n = \sigma(X_1, \dots, X_n)$. Then \mathcal{F}_n is generated by the π -system of events

$$A = \{X_1 \leq x_1, \dots, X_n \leq x_n\}$$

whereas \mathcal{T}_n is generated by the π -system of events

$$B = \{X_{n+1} \leq x_{n+1}, \dots, X_{n+k} \leq x_{n+k}\}, \quad k \in \mathbb{N}.$$

We have $\mathbb{P}(A \cap B) = \mathbb{P}(A)\mathbb{P}(B)$ for all such A and B , by independence. Hence \mathcal{F}_n and \mathcal{T}_n are independent, by Theorem 1.11.1. It follows that \mathcal{F}_n and \mathcal{T} are independent. Now $\bigcup_n \mathcal{F}_n$ is a π -system which generates the σ -algebra $\mathcal{F}_\infty = \sigma(X_n : n \in \mathbb{N})$. So by Theorem 1.11.1 again, \mathcal{F}_∞ and \mathcal{T} are independent. But $\mathcal{T} \subseteq \mathcal{F}_\infty$. So, if $A \in \mathcal{T}$,

$$\mathbb{P}(A) = \mathbb{P}(A \cap A) = \mathbb{P}(A)\mathbb{P}(A)$$

so $\mathbb{P}(A) \in \{0, 1\}$.

Finally, if Y is any \mathcal{T} -measurable random variable, then $F_Y(y) = \mathbb{P}(Y \leq y)$ takes values in $\{0, 1\}$, so $\mathbb{P}(Y = c) = 1$, where $c = \inf\{y : F_Y(y) = 1\}$. \square

2.6. Convergence in measure and convergence almost everywhere.

Let (E, \mathcal{E}, μ) be a measure space. A set $A \in \mathcal{E}$ is sometimes defined by a property shared by its elements. If $\mu(A^c) = 0$, then we say that property holds *almost everywhere* (or *a.e.*). The alternative *almost surely* (or *a.s.*) is often used in a probabilistic context. Thus, for a sequence of measurable functions $(f_n : n \in \mathbb{N})$, we say f_n *converges to f a.e.* to mean that

$$\mu(\{x \in E : f_n(x) \not\rightarrow f(x)\}) = 0.$$

If, on the other hand, we have that

$$\mu(\{x \in E : |f_n(x) - f(x)| > \varepsilon\}) \rightarrow 0, \quad \text{for all } \varepsilon > 0,$$

then we say f_n *converges to f in measure* or, in a probabilistic context, *in probability*.

Example. Let U be a uniform random variable on $(0, 1)$ and let R_n be its binary expansion. Then R_n are independent and identically distributed (i.i.d.). The strong law of large numbers (proved later in §10) applies here to show that $n^{-1} \sum_{i=1}^n R_i$

converges to $1/2$, almost surely. Reformulating this using the canonical probability space $(\Omega, \mathcal{F}, \mathbb{P})$,

$$\mathbb{P} \left(\left\{ \omega \in (0, 1] : \frac{|\{k \leq n : \omega_k = 1\}|}{n} \rightarrow \frac{1}{2} \right\} \right) = \mathbb{P} \left(\frac{R_1 + \cdots + R_n}{n} \rightarrow \frac{1}{2} \right) = 1.$$

This is called Borel's normal number theorem: *almost every point in $(0, 1]$ is normal, that is, has 'equal' proportions of 0's and 1's in its binary expansion.*

Theorem 2.6.1. *Let $(f_n : n \in \mathbb{N})$ be a sequence of measurable functions.*

- (a) *Assume that $\mu(E) < \infty$. If $f_n \rightarrow 0$ a.e. then $f_n \rightarrow 0$ in measure.*
- (b) *If $f_n \rightarrow 0$ in measure then $f_{n_k} \rightarrow 0$ a.e. for some subsequence (n_k) .*

Proof. (a) Suppose $f_n \rightarrow 0$ a.e.. For each $\varepsilon > 0$,

$$\mu(|f_n| \leq \varepsilon) \geq \mu \left(\bigcap_{m \geq n} \{|f_m| \leq \varepsilon\} \right) \uparrow \mu(|f_n| \leq \varepsilon \text{ ev.}) \geq \mu(f_n \rightarrow 0) = \mu(E).$$

Hence $\mu(|f_n| > \varepsilon) \rightarrow 0$ and $f_n \rightarrow 0$ in measure.

(b) Suppose $f_n \rightarrow 0$ in measure, then we can find a subsequence (n_k) such that

$$\sum_k \mu(|f_{n_k}| > 1/k) < \infty.$$

So, by the first Borel–Cantelli lemma,

$$\mu(|f_{n_k}| > 1/k \text{ i.o.}) = 0$$

so $f_{n_k} \rightarrow 0$ a.e.. □

2.7. Large values in sequences of IID random variables. Consider a sequence $(X_n : n \in \mathbb{N})$ of independent random variables, all having the same distribution function F . Assume that $F(x) < 1$ for all $x \in \mathbb{R}$. Then, almost surely, the sequence $(X_n : n \in \mathbb{N})$ is unbounded above, so $\limsup_n X_n = \infty$. A way to describe the occurrence of large values in the sequence is to find a function $g : \mathbb{N} \rightarrow (0, \infty)$ such that, almost surely,

$$\limsup_n X_n/g(n) = 1.$$

We now show that $g(n) = \log n$ is the right choice when $F(x) = 1 - e^{-x}$. The same method adapts to other distributions.

Fix $\alpha > 0$ and consider the event $A_n = \{X_n \geq \alpha \log n\}$. Then $\mathbb{P}(A_n) = e^{-\alpha \log n} = n^{-\alpha}$, so the series $\sum_n \mathbb{P}(A_n)$ converges if and only if $\alpha > 1$. By the Borel–Cantelli Lemmas, we deduce that, for all $\varepsilon > 0$,

$$\mathbb{P}(X_n/\log n \geq 1 \text{ i.o.}) = 1, \quad \mathbb{P}(X_n/\log n \geq 1 + \varepsilon \text{ i.o.}) = 0.$$

Hence, almost surely,

$$\limsup_n X_n/\log n = 1.$$

3. INTEGRATION

3.1. Definition of the integral and basic properties. Let (E, \mathcal{E}, μ) be a measure space. We shall define, where possible, for a measurable function $f : E \rightarrow [-\infty, \infty]$, the *integral* of f , to be denoted

$$\mu(f) = \int_E f d\mu = \int_E f(x) \mu(dx).$$

For a random variable X on a probability space $(\Omega, \mathcal{F}, \mathbb{P})$, the integral is usually called instead the *expectation* of X and written $\mathbb{E}(X)$ instead of $\int X d\mathbb{P}$.

A *simple* function is one of the form

$$f = \sum_{k=1}^m a_k 1_{A_k}$$

where $0 \leq a_k < \infty$ and $A_k \in \mathcal{E}$ for all k , and where $m \in \mathbb{N}$. For simple functions f , we define

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

where we adopt the convention $0 \cdot \infty = 0$. Although the representation of f is not unique, it is straightforward to check that $\mu(f)$ is well defined and, for simple functions f, g and constants $\alpha, \beta \geq 0$, we have

- (a) $\mu(\alpha f + \beta g) = \alpha \mu(f) + \beta \mu(g)$,
- (b) $f \leq g$ implies $\mu(f) \leq \mu(g)$,
- (c) $f = 0$ a.e. if and only if $\mu(f) = 0$.

In particular, for simple functions f , we have

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

We define the integral $\mu(f)$ of a non-negative measurable function f by

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

We have seen that this is consistent with our definition for simple functions. Note that, for all non-negative measurable functions f, g with $f \leq g$, we have $\mu(f) \leq \mu(g)$. For any measurable function f , set $f^+ = f \vee 0$ and $f^- = (-f) \vee 0$. Then $f = f^+ - f^-$ and $|f| = f^+ + f^-$. If $\mu(|f|) < \infty$, then we say that f is *integrable* and define

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

Note that $|\mu(f)| \leq \mu(|f|)$ for all integrable functions f . We sometimes define the integral $\mu(f)$ by the same formula, even when f is not integrable, but when either $\mu(f^-)$ or $\mu(f^+)$ is finite. In such cases the integral take the value ∞ or $-\infty$.

Here is the key result for the theory of integration. For $x \in [0, \infty]$ and a sequence $(x_n : n \in \mathbb{N})$ in $[0, \infty]$, we write $x_n \uparrow x$ to mean that $x_n \leq x_{n+1}$ for all n and $x_n \rightarrow x$

as $n \rightarrow \infty$. For a non-negative function f on E and a sequence of such functions $(f_n : n \in \mathbb{N})$, we write $f_n \uparrow f$ to mean that $f_n(x) \uparrow f(x)$ for all $x \in E$.

Theorem 3.1.1 (Monotone convergence). *Let f be a non-negative measurable function and let $(f_n : n \in \mathbb{N})$ be a sequence of such functions. Suppose that $f_n \uparrow f$. Then $\mu(f_n) \uparrow \mu(f)$.*

Proof. Case 1: $f_n = 1_{A_n}$, $f = 1_A$.

The result is a simple consequence of countable additivity.

Case 2: f_n simple, $f = 1_A$.

Fix $\varepsilon > 0$ and set $A_n = \{f_n > 1 - \varepsilon\}$. Then $A_n \uparrow A$ and

$$(1 - \varepsilon)1_{A_n} \leq f_n \leq 1_A$$

so

$$(1 - \varepsilon)\mu(A_n) \leq \mu(f_n) \leq \mu(A).$$

But $\mu(A_n) \uparrow \mu(A)$ by Case 1 and $\varepsilon > 0$ was arbitrary, so the result follows.

Case 3: f_n simple, f simple.

We can write f in the form

$$f = \sum_{k=1}^m a_k 1_{A_k}$$

with $a_k > 0$ for all k and the sets A_k disjoint. Then $f_n \uparrow f$ implies

$$a_k^{-1} 1_{A_k} f_n \uparrow 1_{A_k}$$

so, by Case 2,

$$\mu(f_n) = \sum_k \mu(1_{A_k} f_n) \uparrow \sum_k a_k \mu(A_k) = \mu(f).$$

Case 4: f_n simple, $f \geq 0$ measurable.

Let g be simple with $g \leq f$. Then $f_n \uparrow f$ implies $f_n \wedge g \uparrow g$ so, by Case 3,

$$\mu(f_n) \geq \mu(f_n \wedge g) \uparrow \mu(g).$$

Since g was arbitrary, the result follows.

Case 5: $f_n \geq 0$ measurable, $f \geq 0$ measurable.

Set $g_n = (2^{-n} \lfloor 2^n f_n \rfloor) \wedge n$ then g_n is simple and $g_n \leq f_n \leq f$, so

$$\mu(g_n) \leq \mu(f_n) \leq \mu(f).$$

But $f_n \uparrow f$ forces $g_n \uparrow f$, so $\mu(g_n) \uparrow \mu(f)$, by Case 4, and so $\mu(f_n) \uparrow \mu(f)$. □

Theorem 3.1.2. For all non-negative measurable functions f, g and all constants $\alpha, \beta \geq 0$,

- (a) $\mu(\alpha f + \beta g) = \alpha\mu(f) + \beta\mu(g)$,
- (b) $f \leq g$ implies $\mu(f) \leq \mu(g)$,
- (c) $f = 0$ a.e. if and only if $\mu(f) = 0$.

Proof. Define simple functions f_n, g_n by

$$f_n = (2^{-n} \lfloor 2^n f \rfloor) \wedge n, \quad g_n = (2^{-n} \lfloor 2^n g \rfloor) \wedge n.$$

Then $f_n \uparrow f$ and $g_n \uparrow g$, so $\alpha f_n + \beta g_n \uparrow \alpha f + \beta g$. Hence, by monotone convergence,

$$\mu(f_n) \uparrow \mu(f), \quad \mu(g_n) \uparrow \mu(g), \quad \mu(\alpha f_n + \beta g_n) \uparrow \mu(\alpha f + \beta g).$$

We know that $\mu(\alpha f_n + \beta g_n) = \alpha\mu(f_n) + \beta\mu(g_n)$, so we obtain (a) on letting $n \rightarrow \infty$. As we noted above, (b) is obvious from the definition of the integral. If $f = 0$ a.e., then $f_n = 0$ a.e., for all n , so $\mu(f_n) = 0$ and $\mu(f) = 0$. On the other hand, if $\mu(f) = 0$, then $\mu(f_n) = 0$ for all n , so $f_n = 0$ a.e. and $f = 0$ a.e. \square

Theorem 3.1.3. For all integrable functions f, g and all constants $\alpha, \beta \in \mathbb{R}$,

- (a) $\mu(\alpha f + \beta g) = \alpha\mu(f) + \beta\mu(g)$,
- (b) $f \leq g$ implies $\mu(f) \leq \mu(g)$,
- (c) $f = 0$ a.e. implies $\mu(f) = 0$.

Proof. We note that $\mu(-f) = -\mu(f)$. For $\alpha \geq 0$, we have

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

If $h = f + g$ then $h^+ + f^- + g^- = h^- + f^+ + g^+$, so

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

and so $\mu(h) = \mu(f) + \mu(g)$. That proves (a). If $f \leq g$ then $\mu(g) - \mu(f) = \mu(g - f) \geq 0$, by (a). Finally, if $f = 0$ a.e., then $f^\pm = 0$ a.e., so $\mu(f^\pm) = 0$ and so $\mu(f) = 0$. \square

Note that in Theorem 3.1.3(c) we lose the reverse implication. The following result is sometimes useful:

Proposition 3.1.4. Let \mathcal{A} be a π -system containing E and generating \mathcal{E} . Then, for any integrable function f ,

$$\mu(f1_A) = 0 \text{ for all } A \in \mathcal{A} \text{ implies } f = 0 \text{ a.e..}$$

Here are some minor variants on the monotone convergence theorem.

Proposition 3.1.5. Let $(f_n : n \in \mathbb{N})$ be a sequence of measurable functions, with $f_n \geq 0$ a.e.. Then

$$f_n \uparrow f \text{ a.e.} \implies \mu(f_n) \uparrow \mu(f).$$

Thus the pointwise hypotheses of non-negativity and monotone convergence can be relaxed to hold almost everywhere.

Proposition 3.1.6. *Let $(g_n : n \in \mathbb{N})$ be a sequence of non-negative measurable functions. Then*

$$\sum_{n=1}^{\infty} \mu(g_n) = \mu \left(\sum_{n=1}^{\infty} g_n \right).$$

This reformulation of monotone convergence makes it clear that it is the counterpart for the integration of functions of the countable additivity property of the measure on sets.

3.2. Integrals and limits. In the monotone convergence theorem, the hypothesis that the given sequence of functions is non-decreasing is essential. In this section we obtain some results on the integrals of limits of functions without such a hypothesis.

Lemma 3.2.1 (Fatou's lemma). *Let $(f_n : n \in \mathbb{N})$ be a sequence of non-negative measurable functions. Then*

$$\mu(\liminf f_n) \leq \liminf \mu(f_n).$$

Proof. For $k \geq n$, we have

$$\inf_{m \geq n} f_m \leq f_k$$

so

$$\mu(\inf_{m \geq n} f_m) \leq \inf_{k \geq n} \mu(f_k) \leq \liminf \mu(f_n).$$

But, as $n \rightarrow \infty$,

$$\inf_{m \geq n} f_m \uparrow \sup_n \left(\inf_{m \geq n} f_m \right) = \liminf f_n$$

so, by monotone convergence,

$$\mu(\inf_{m \geq n} f_m) \uparrow \mu(\liminf f_n).$$

□

Theorem 3.2.2 (Dominated convergence). *Let f be a measurable function and let $(f_n : n \in \mathbb{N})$ be a sequence of such functions. Suppose that $f_n(x) \rightarrow f(x)$ for all $x \in E$ and that $|f_n| \leq g$ for all n , for some integrable function g . Then f and f_n are integrable, for all n , and $\mu(f_n) \rightarrow \mu(f)$.*

Proof. The limit f is measurable and $|f| \leq g$, so $\mu(|f|) \leq \mu(g) < \infty$, so f is integrable. We have $0 \leq g \pm f_n \rightarrow g \pm f$ so certainly $\liminf(g \pm f_n) = g \pm f$. By Fatou's lemma,

$$\mu(g) + \mu(f) = \mu(\liminf(g + f_n)) \leq \liminf \mu(g + f_n) = \mu(g) + \liminf \mu(f_n),$$

$$\mu(g) - \mu(f) = \mu(\liminf(g - f_n)) \leq \liminf \mu(g - f_n) = \mu(g) - \limsup \mu(f_n).$$

Since $\mu(g) < \infty$, we can deduce that

$$\mu(f) \leq \liminf \mu(f_n) \leq \limsup \mu(f_n) \leq \mu(f).$$

This proves that $\mu(f_n) \rightarrow \mu(f)$ as $n \rightarrow \infty$.

□

3.3. Transformations of integrals.

Proposition 3.3.1. *Let (E, \mathcal{E}, μ) be a measure space and let $A \in \mathcal{E}$. Then the set \mathcal{E}_A of measurable subsets of A is a σ -algebra and the restriction μ_A of μ to \mathcal{E}_A is a measure. Moreover, for any non-negative measurable function f on E , we have*

$$\mu(f1_A) = \mu_A(f|_A).$$

In the case of Lebesgue measure on \mathbb{R} , we write, for any interval I with $\inf I = a$ and $\sup I = b$,

$$\int_{\mathbb{R}} f1_I(x)dx = \int_I f(x)dx = \int_a^b f(x)dx.$$

Note that the sets $\{a\}$ and $\{b\}$ have measure zero, so we do not need to specify whether they are included in I or not.

Proposition 3.3.2. *Let (E, \mathcal{E}) and (G, \mathcal{G}) be measure spaces and let $f : E \rightarrow G$ be a measurable function. Given a measure μ on (E, \mathcal{E}) , define $\nu = \mu \circ f^{-1}$, the image measure on (G, \mathcal{G}) . Then, for all non-negative measurable functions g on G ,*

$$\nu(g) = \mu(g \circ f).$$

In particular, for a G -valued random variable X on a probability space $(\Omega, \mathcal{F}, \mathbb{P})$, for any non-negative measurable function g on G , we have

$$\mathbb{E}(g(X)) = \mu_X(g).$$

(Recall that for a random variable Y , by definition $\mathbb{E}(Y) = \int Y d\mathbb{P}$.)

Proposition 3.3.3. *Let (E, \mathcal{E}, μ) be a measure space and let f be a non-negative measurable function on E . Define $\nu(A) = \mu(f1_A)$, $A \in \mathcal{E}$. Then ν is a measure on E and, for all non-negative measurable functions g on E ,*

$$\nu(g) = \mu(fg).$$

f is called the density of ν with respect to μ . If f, g are two densities of ν with respect to μ , then $f = g$, μ -a.e., and ν -a.e. as well.

In particular, to each non-negative Borel function f on \mathbb{R} , there corresponds a Borel measure μ on \mathbb{R} given by $\mu(A) = \int_A f(x)dx$. Then, for all non-negative Borel functions g ,

$$\mu(g) = \int_{\mathbb{R}^n} g(x)f(x)dx.$$

We say that μ has density f (with respect to Lebesgue measure).

If the law μ_X of a real-valued random variable X has a density f_X , then we call f_X a density function for X . Then $\mathbb{P}(X \in A) = \int_A f_X(x)dx$, for all Borel sets A , and, for all non-negative Borel functions g on \mathbb{R} ,

$$\mathbb{E}(g(X)) = \mu_X(g) = \int_{\mathbb{R}} g(x)f_X(x)dx.$$

Example. Let X be an exponential random variable. That is, its distribution function satisfies $F(t) = 1 - e^{-t}, t \geq 0$. Therefore, by the fundamental theorem of calculus (proved below),

$$\mathbb{P}(X \leq t) = 1 - e^{-t} = \int_0^t e^{-x} dx = \int_0^\infty e^{-x} \mathbf{1}_{[0,t]}(x) dx.$$

Then by Dynkin's lemma, it follows that for all Borel set A ,

$$\mathbb{P}(X \in A) = \int f_X(x) \mathbf{1}_A(x) dx,$$

where $f_X(x) = e^{-x} \mathbf{1}_{[0,\infty)}(x)$. Thus μ_X has a density function f_X . Hence, by Proposition 3.3.2 and 3.3.3,

$$\mathbb{E}(X^2) = \int x^2 f_X(x) dx = \int_0^\infty x^2 e^{-x} dx$$

which we will compute as being equal to 2.

3.4. Relation to Riemann integral. Let $I = [a, b]$ and let $f : I \rightarrow \mathbb{R}$. Let μ denote the Lebesgue measure. Recall that f is said Riemann integrable with integral R if for all $\varepsilon > 0$, there exists $\delta > 0$ such that for all finite partitions of I into subintervals I_1, \dots, I_k of mesh $\max_{1 \leq j \leq k} \mu(I_j) \leq \delta$,

$$(3.1) \quad \left| R - \sum_{j=1}^k f(x_j) \mu(I_j) \right| \leq \varepsilon$$

whenever $x_j \in I_j, 1 \leq j \leq k$. We will show that if f is bounded on I and Riemann-integrable then it is Lebesgue-integrable (hence measurable) $\int_I f d\mu = R$, so the two definitions coincide.

We may assume without loss of generality that $I = [0, 1]$. Consider the dyadic subdivision of I , that is, for $n \geq 1$ and $0 \leq k \leq 2^n - 1$, let $I_{k,n} = [k2^{-n}, (k+1)2^{-n}]$. Consider the function

$$\bar{f}_n(x) = \sum_{k=0}^{2^n-1} \sup_{I_{k,n}}(f) \mathbf{1}_{I_{k,n}}(x)$$

and

$$\underline{f}_n(x) = \sum_{k=0}^{2^n-1} \inf_{I_{k,n}}(f) \mathbf{1}_{I_{k,n}}(x)$$

Then note that \bar{f}_n is simple and is a decreasing sequence of functions, while \underline{f}_n is nondecreasing. Moreover

$$\underline{f}_n \leq f \leq \bar{f}_n.$$

Since f is Riemann integrable, we know

$$\lim_{n \rightarrow \infty} \int_I \bar{f}_n d\mu = \lim_{n \rightarrow \infty} \int_I \underline{f}_n d\mu = R.$$

To see this, use (3.1) with x_k such that $f(x_k)$ is arbitrarily close to $\sup_{I_{k,n}}(f)$ and $\inf_{I_{k,n}}(f)$ respectively, $0 \leq k \leq 2^n - 1$. Let $\bar{f} = \lim \bar{f}_n$ and $\underline{f} = \lim \underline{f}_n$. Then \bar{f} and \underline{f} are measurable as pointwise limits of simple functions, and moreover by the dominated convergence theorem, $R = \int_I \bar{f} d\mu = \int_I \underline{f} d\mu = R$. Since $\underline{f} \leq \bar{f}$, we deduce that $\underline{f} = \bar{f}$ a.e. Thus, since $\underline{f} \leq f \leq \bar{f}$, we deduce that $f = \bar{f}$ a.e. and thus f is Lebesgue-measurable. Moreover $\int_I f d\mu = \int_I \bar{f} d\mu = R$, which proves the result.

Note that any Riemann-integrable function on $I = [a, b]$ is necessarily bounded, so the result holds for any Riemann-integrable function on I . If $I = (a, b]$ and $|f|$ has finite Riemann integral, then it follows that f is integrable and the two integrals coincide.

Example. The function $\mathbf{1}_{\mathbb{Q} \cap [0,1]}$ is Lebesgue-integrable with integral 0 (since $\mu(\mathbb{Q}) = 0$) but is not Riemann-integrable.

3.5. Fundamental theorem of calculus. We show that integration with respect to Lebesgue measure on \mathbb{R} acts as an inverse to differentiation. Since we restrict here to the integration of continuous functions, the proof is the same as for the Riemann integral.

Theorem 3.5.1 (Fundamental theorem of calculus).

(i) Let $f : [a, b] \rightarrow \mathbb{R}$ be a continuous function and set

$$F_a(t) = \int_a^t f(x)dx.$$

Then F_a is differentiable on $[a, b]$, with $F'_a = f$.

(ii) Let $F : [a, b] \rightarrow \mathbb{R}$ be differentiable with continuous derivative f . Then

$$\int_a^b f(x)dx = F(b) - F(a).$$

Proof. Fix $t \in [a, b]$. Given $\varepsilon > 0$, there exists $\delta > 0$ such that $|f(x) - f(t)| \leq \varepsilon$ whenever $|x - t| \leq \delta$. So, for $0 < h \leq \delta$,

$$\begin{aligned} \left| \frac{F_a(t+h) - F_a(t)}{h} - f(t) \right| &= \frac{1}{h} \left| \int_t^{t+h} (f(x) - f(t))dx \right| \\ &\leq \frac{1}{h} \int_t^{t+h} |f(x) - f(t)|dx \leq \frac{\varepsilon}{h} \int_t^{t+h} dx = \varepsilon. \end{aligned}$$

Here we have used that if f is integrable then $|\int f d\mu| \leq \int |f| d\mu$, by the triangle inequality. Hence F_a is differentiable on the right at t with derivative $f(t)$. Similarly, for all $t \in (a, b]$, F_a is differentiable on the left at t with derivative $f(t)$. Finally, $(F - F_a)'(t) = 0$ for all $t \in (a, b)$ so $F - F_a$ is constant (by the mean value theorem), and so

$$F(b) - F(a) = F_a(b) - F_a(a) = \int_a^b f(x)dx.$$

□

Proposition 3.5.2. *Let $\phi : [a, b] \rightarrow \mathbb{R}$ be continuously differentiable and strictly increasing. Then, for all non-negative Borel functions g on $[\phi(a), \phi(b)]$,*

$$\int_{\phi(a)}^{\phi(b)} g(y)dy = \int_a^b g(\phi(x))\phi'(x)dx.$$

The proposition can be proved as follows. First, the case where g is the indicator function of an interval follows from the Fundamental Theorem of Calculus. Next, show that the set of Borel sets B such that the conclusion holds for $g = 1_B$ is a d -system, which must then be the whole Borel σ -algebra by Dynkin's lemma. The identity extends to simple functions by linearity and then to all non-negative measurable functions g by monotone convergence, using approximating simple functions $(2^{-n} \lfloor 2^n g \rfloor) \wedge n$.

An general formulation of this procedure, which is often used, is given in the monotone class theorem Theorem 2.1.2.

3.6. Differentiation under the integral sign. Integration in one variable and differentiation in another can be interchanged subject to some regularity conditions.

Theorem 3.6.1 (Differentiation under the integral sign). *Let $U \subseteq \mathbb{R}$ be open and suppose that $f : U \times E \rightarrow \mathbb{R}$ satisfies:*

- (i) $x \mapsto f(t, x)$ is integrable for all t ,
- (ii) $t \mapsto f(t, x)$ is differentiable for all x ,
- (iii) for some integrable function g , for all $x \in E$ and all $t \in U$,

$$\left| \frac{\partial f}{\partial t}(t, x) \right| \leq g(x).$$

Then the function $x \mapsto (\partial f / \partial t)(t, x)$ is integrable for all t . Moreover, the function $F : U \rightarrow \mathbb{R}$, defined by

$$F(t) = \int_E f(t, x) \mu(dx),$$

is differentiable and

$$\frac{d}{dt} F(t) = \int_E \frac{\partial f}{\partial t}(t, x) \mu(dx).$$

Proof. Take any sequence $h_n \rightarrow 0$ and set

$$g_n(x) = \frac{f(t + h_n, x) - f(t, x)}{h_n} - \frac{\partial f}{\partial t}(t, x).$$

Then $g_n(x) \rightarrow 0$ for all $x \in E$ and, by the mean value theorem, $|g_n| \leq 2g$ for all n . In particular, for all t , the function $x \mapsto (\partial f / \partial t)(t, x)$ is the limit of measurable functions, hence measurable, and hence integrable, by (iii). Then, by dominated convergence,

$$\frac{F(t + h_n) - F(t)}{h_n} - \int_E \frac{\partial f}{\partial t}(t, x) \mu(dx) = \int_E g_n(x) \mu(dx) \rightarrow 0.$$

□

3.7. Product measure and Fubini's theorem. Let $(E_1, \mathcal{E}_1, \mu_1)$ and $(E_2, \mathcal{E}_2, \mu_2)$ be *finite* measure spaces. The set

$$\mathcal{A} = \{A_1 \times A_2 : A_1 \in \mathcal{E}_1, A_2 \in \mathcal{E}_2\}$$

is a π -system of subsets of $E = E_1 \times E_2$. Define the *product σ -algebra*

$$\mathcal{E}_1 \otimes \mathcal{E}_2 = \sigma(\mathcal{A}).$$

Set $\mathcal{E} = \mathcal{E}_1 \otimes \mathcal{E}_2$.

Lemma 3.7.1. *Let $f : E \rightarrow \mathbb{R}$ be \mathcal{E} -measurable. Then, for all $x_1 \in E_1$, the function $x_2 \mapsto f(x_1, x_2) : E_2 \rightarrow \mathbb{R}$ is \mathcal{E}_2 -measurable.*

Proof. Denote by \mathcal{V} the set of bounded \mathcal{E} -measurable functions for which the conclusion holds. Then \mathcal{V} is a vector space, containing the indicator function 1_A of every set $A \in \mathcal{A}$. Moreover, if $f_n \in \mathcal{V}$ for all n and if f is bounded with $0 \leq f_n \uparrow f$, then also $f \in \mathcal{V}$. So, by the monotone class theorem, \mathcal{V} contains all bounded \mathcal{E} -measurable functions. The rest is easy. \square

Lemma 3.7.2. *For all bounded \mathcal{E} -measurable functions f , the function*

$$x_1 \mapsto f_1(x_1) = \int_{E_2} f(x_1, x_2) \mu_2(dx_2) : E_1 \rightarrow \mathbb{R}$$

is bounded and \mathcal{E}_1 -measurable.

Proof. Apply the monotone class theorem, as in the preceding lemma. \square

Theorem 3.7.3 (Product measure). *There exists a unique measure $\mu = \mu_1 \otimes \mu_2$ on \mathcal{E} such that*

$$\mu(A_1 \times A_2) = \mu_1(A_1) \mu_2(A_2)$$

for all $A_1 \in \mathcal{E}_1$ and $A_2 \in \mathcal{E}_2$.

Proof. Uniqueness holds because \mathcal{A} is a π -system generating \mathcal{E} . For existence, by the lemmas, we can define

$$\mu(A) = \int_{E_1} \left(\int_{E_2} 1_A(x_1, x_2) \mu_2(dx_2) \right) \mu_1(dx_1)$$

and use monotone convergence to see that μ is countably additive. \square

Note that if μ_1 or μ_2 is not finite, then one cannot define the measure μ as we may be faced with taking the product $0 \times \infty$. This cannot be done without breaking the symmetry between E_1 and E_2 .

Proposition 3.7.4. *Let $\hat{\mathcal{E}} = \mathcal{E}_2 \otimes \mathcal{E}_1$ and $\hat{\mu} = \mu_2 \otimes \mu_1$. For a function f on $E_1 \times E_2$, write \hat{f} for the function on $E_2 \times E_1$ given by $\hat{f}(x_2, x_1) = f(x_1, x_2)$. Suppose that f is \mathcal{E} -measurable. Then \hat{f} is $\hat{\mathcal{E}}$ -measurable, and if f is also non-negative, then $\hat{\mu}(\hat{f}) = \mu(f)$.*

Theorem 3.7.5 (Fubini's theorem).

(a) Let f be \mathcal{E} -measurable and non-negative. Then

$$\mu(f) = \int_{E_1} \left(\int_{E_2} f(x_1, x_2) \mu_2(dx_2) \right) \mu_1(dx_1).$$

(b) Let f be μ -integrable. Then

- (i) $x_2 \mapsto f(x_1, x_2)$ is μ_2 -integrable for μ_1 -almost all x_1 ,
 - (ii) $x_1 \mapsto \int_{E_2} f(x_1, x_2) \mu_2(dx_2)$ is μ_1 -integrable
- and the formula for $\mu(f)$ in (a) holds.

Note that the *iterated integral* in (a) is well defined, for all bounded or non-negative measurable functions f , by Lemmas 3.7.1 and 3.7.2. Note also that, in combination with Proposition 3.7.4, Fubini's theorem allows us to interchange the order of integration in multiple integrals, whenever the integrand is non-negative or μ -integrable.

Proof. Denote by \mathcal{V} the set of all bounded \mathcal{E} -measurable functions f for which the formula holds. Then \mathcal{V} contains the indicator function of every \mathcal{E} -measurable set so, by the monotone class theorem, \mathcal{V} contains all bounded \mathcal{E} -measurable functions. Hence, for all \mathcal{E} -measurable functions f , we have

$$\mu(f_n) = \int_{E_1} \left(\int_{E_2} f_n(x_1, x_2) \mu_2(dx_2) \right) \mu_1(dx_1)$$

where $f_n = (-n) \vee f \wedge n$.

For f non-negative, we can pass to the limit as $n \rightarrow \infty$ by monotone convergence to extend the formula to f . That proves (a).

If f is μ -integrable, then, by (a)

$$\int_{E_1} \left(\int_{E_2} |f(x_1, x_2)| \mu_2(dx_2) \right) \mu_1(dx_1) = \mu(|f|) < \infty.$$

Hence we obtain (i) and (ii). Then, by dominated convergence, we can pass to the limit as $n \rightarrow \infty$ in the formula for $\mu(f_n)$ to obtain the desired formula for $\mu(f)$. \square

Example. Let f_n be a sequence of measurable functions on some finite measure space (E, \mathcal{E}, μ) . If $\sum_n \int |f_n| d\mu < \infty$, then $\sum_n \int f_n d\mu = \int \sum_n f_n d\mu$. If $f_n \geq 0$, then $\int \sum_n f_n d\mu = \infty$ if and only if $\sum_n \int f_n d\mu = \infty$. This follows from considering the function $f(n, x) = f_n(x)$ on the product space $\mathbb{N} \times E$ where \mathbb{N} is equipped with the σ -algebra $\mathcal{P}(\mathbb{N})$ and the counting measure, $\nu(A) = \#A$.

The existence of product measure and Fubini's theorem extend easily to σ -finite measure spaces, i.e., measure spaces such that there exists $E_n \uparrow E$ with $E_n \in \mathcal{E}$ and $\mu(E_n) < \infty$.

Remark. The assumption that the spaces $(E_1, \mathcal{E}_1, \mu_1)$ and $(E_2, \mathcal{E}_2, \mu_2)$ are σ -finite is essential. Consider the following counter-examples. Let $(E_1, \mathcal{E}_1) = (E_2, \mathcal{E}_2) = ([0, 1], \mathcal{B})$. Let μ_1 be the Lebesgue measure on $[0, 1]$ and let μ_2 be the counting

measure on $[0, 1]$. That is, for $A \in \mathcal{B}$, $\mu_2(A) = \#A$. Consider the set $\Delta \in \mathcal{E}_1 \otimes \mathcal{E}_2$ given by $\{(x, y) : x = y\}$. Then

$$\int_{E_1} \left(\int_{E_2} f(x, y) \mu_2(dx) \right) \mu_1(dx) = \int_{E_1} 1 \cdot \mu_1(dx) = 1,$$

but

$$\int_{E_2} \left(\int_{E_1} \mu_1(dx) \right) \mu_2(dx) = \int_{E_2} 0 \cdot \mu_2(dx) = 0.$$

The operation of taking the product of two measure spaces is associative, by a π -system uniqueness argument. So we can, by induction, take the product of a finite number, without specifying the order. The measure obtained by taking the n -fold product of Lebesgue measure on \mathbb{R} is called *Lebesgue measure on \mathbb{R}^n* . The corresponding integral is written

$$\int_{\mathbb{R}^n} f(x) dx = \int \cdots \int f(x_1, \dots, x_n) dx_1 \cdots dx_n.$$

Proposition 3.7.6. *Let $\mathcal{B}(\mathbb{R}^n)$ denote the Borel σ -algebra on \mathbb{R}^n , and let \mathcal{E} denote the product $\mathcal{B}(\mathbb{R}) \otimes \cdots \otimes \mathcal{B}(\mathbb{R})$. Then $\mathcal{E} = \mathcal{B}(\mathbb{R}^n)$.*

Proof. We do the proof when $n = 2$. Let $A = A_1 \times A_2$, where $A_1, A_2 \in \mathcal{B}(\mathbb{R})$. Let us show $A \in \mathcal{B}(\mathbb{R}^2)$. If I is an open interval then $\{B : I \times B \in \mathcal{B}(\mathbb{R}^2)\}$ is a d -system containing finite union of open intervals and thus contains $\mathcal{B}(\mathbb{R})$, and hence A_2 in particular. As a consequence, $\{A : A \times A_2 \in \mathcal{B}(\mathbb{R}^2)\}$ is a d -system containing open intervals, and (likewise) hence contains A_1 . Thus $\mathcal{A} \subset \mathcal{B}(\mathbb{R}^2)$, and since \mathcal{A} is a π -system, $\mathcal{E} \subset \mathcal{B}(\mathbb{R}^2)$.

For the converse direction, it suffices to prove that $\mathcal{B}(\mathbb{R}^2)$ is generated by rectangles $(a, b] \times (c, d]$. To see this, let G be open in \mathbb{R}^2 . Then for $y \in G$, we can find a rectangle $A_y = (a, b] \times (c, d]$ where $a, b, c, d \in \mathbb{Q}$ and $A_y \subset G$. Then $G = \bigcup_{y \in G} A_y$, and this union is countable as there are only countably many rational rectangles. Thus G is in the σ -algebra generated by such rectangles, and the result is proved. \square

4. NORMS AND INEQUALITIES

4.1. L^p -norms. Let (E, \mathcal{E}, μ) be a measure space. For $1 \leq p < \infty$, we denote by $L^p = L^p(E, \mathcal{E}, \mu)$ the set of measurable functions f with finite L^p -norm:

$$\|f\|_p = \left(\int_E |f|^p d\mu \right)^{1/p} < \infty.$$

We denote by $L^\infty = L^\infty(E, \mathcal{E}, \mu)$ the set of measurable functions f with finite L^∞ -norm:

$$\|f\|_\infty = \inf\{\lambda : |f| \leq \lambda \text{ a.e.}\}.$$

Note that $\|f\|_p \leq \mu(E)^{1/p} \|f\|_\infty$ for all $1 \leq p < \infty$. For $1 \leq p \leq \infty$ and $f_n \in L^p$, we say that f_n converges to f in L^p if $\|f_n - f\|_p \rightarrow 0$.

4.2. Chebyshev's inequality. Let f be a non-negative measurable function and let $\lambda \geq 0$. We use the notation $\{f \geq \lambda\}$ for the set $\{x \in E : f(x) \geq \lambda\}$. Note that

$$\lambda 1_{\{f \geq \lambda\}} \leq f$$

so on integrating we obtain *Chebyshev's inequality*

$$\lambda \mu(f \geq \lambda) \leq \mu(f).$$

Now let g be any measurable function. We can deduce inequalities for g by choosing some non-negative measurable function ϕ and applying Chebyshev's inequality to $f = \phi \circ g$. For example, if $g \in L^p, p < \infty$ and $\lambda > 0$, then

$$\mu(|g| \geq \lambda) = \mu(|g|^p \geq \lambda^p) \leq \lambda^{-p} \mu(|g|^p) < \infty.$$

So we obtain the *tail estimate*

$$\mu(|g| \geq \lambda) = O(\lambda^{-p}), \quad \text{as } \lambda \rightarrow \infty.$$

It follows that if $f_n \rightarrow f$ in L^p then $f_n \rightarrow f$ also in measure. See also (5.1) for another inequality of the same flavour which is in practice very useful.

4.3. Jensen's inequality. Let $I \subseteq \mathbb{R}$ be an interval. A function $c : I \rightarrow \mathbb{R}$ is *convex* if, for all $x, y \in I$ and $t \in [0, 1]$,

$$c(tx + (1-t)y) \leq tc(x) + (1-t)c(y).$$

Lemma 4.3.1. *Let $c : I \rightarrow \mathbb{R}$ be convex and let m be a point in the interior of I . Then there exist $a, b \in \mathbb{R}$ such $c(x) \geq ax + b$ for all x , with equality at $x = m$.*

Proof. By convexity, for $m, x, y \in I$ with $x < m < y$, we have

$$\frac{c(m) - c(x)}{m - x} \leq \frac{c(y) - c(m)}{y - m}.$$

So, fixing an interior point m , there exists $a \in \mathbb{R}$ such that, for all $x < m$ and all $y > m$

$$\frac{c(m) - c(x)}{m - x} \leq a \leq \frac{c(y) - c(m)}{y - m}.$$

Then $c(x) \geq a(x - m) + c(m)$, for all $x \in I$. \square

Theorem 4.3.2 (Jensen's inequality). *Let X be an integrable random variable with values in I and let $c : I \rightarrow \mathbb{R}$ be convex. Then $\mathbb{E}(c(X))$ is well defined and*

$$\mathbb{E}(c(X)) \geq c(\mathbb{E}(X)).$$

Proof. The case where X is almost surely constant is easy. We exclude it. Then $m = \mathbb{E}(X)$ must lie in the interior of I . Choose $a, b \in \mathbb{R}$ as in the lemma. Then $c(X) \geq aX + b$. In particular $\mathbb{E}(c(X)^-) \leq |a|\mathbb{E}(|X|) + |b| < \infty$, so $\mathbb{E}(c(X))$ is well defined. Moreover

$$\mathbb{E}(c(X)) \geq a\mathbb{E}(X) + b = am + b = c(m) = c(\mathbb{E}(X)).$$

\square

We deduce from Jensen's inequality *the monotonicity of L^p -norms with respect to a probability measure*. Let $1 \leq p < q < \infty$. Set $c(x) = |x|^{q/p}$, then c is convex. So, for any $X \in L^p(\mathbb{P})$,

$$\|X\|_p = (\mathbb{E}|X|^p)^{1/p} = (c(\mathbb{E}|X|^p))^{1/q} \leq (\mathbb{E}c(|X|^p))^{1/q} = (\mathbb{E}|X|^q)^{1/q} = \|X\|_q.$$

In particular, $L^p(\mathbb{P}) \supseteq L^q(\mathbb{P})$.

4.4. Hölder's inequality and Minkowski's inequality. For $p, q \in [1, \infty]$, we say that p and q are *conjugate indices* if

$$\frac{1}{p} + \frac{1}{q} = 1.$$

Theorem 4.4.1 (Hölder's inequality). *Let $p, q \in (1, \infty)$ be conjugate indices. Then, for all measurable functions f and g , we have*

$$\mu(|fg|) \leq \|f\|_p \|g\|_q.$$

Proof. The cases where $\|f\|_p = 0$ or $\|f\|_p = \infty$ are obvious. We exclude them. Then, by multiplying f by an appropriate constant, we are reduced to the case where $\|f\|_p = 1$. So we can define a probability measure \mathbb{P} on \mathcal{E} by

$$\mathbb{P}(A) = \int_A |f|^p d\mu.$$

For measurable functions $X \geq 0$,

$$\mathbb{E}(X) = \mu(X|f|^p), \quad \mathbb{E}(X) \leq \mathbb{E}(X^q)^{1/q}.$$

Note that $q(p-1) = p$. Then

$$\mu(|fg|) = \mu\left(\frac{|g|}{|f|^{p-1}}|f|^p\right) = \mathbb{E}\left(\frac{|g|}{|f|^{p-1}}\right) \leq \mathbb{E}\left(\frac{|g|^q}{|f|^{q(p-1)}}\right)^{1/q} = \mu(|g|^q)^{1/q} = \|f\|_p \|g\|_q.$$

□

Theorem 4.4.2 (Minkowski's inequality). *For $p \in [1, \infty)$ and measurable functions f and g , we have*

$$\|f + g\|_p \leq \|f\|_p + \|g\|_p.$$

Proof. The cases where $p = 1$ or where $\|f\|_p = \infty$ or $\|g\|_p = \infty$ are easy. We exclude them. Then, since $|f + g|^p \leq 2^p(|f|^p + |g|^p)$, we have

$$\mu(|f + g|^p) \leq 2^p\{\mu(|f|^p) + \mu(|g|^p)\} < \infty.$$

The case where $\|f + g\|_p = 0$ is clear, so let us assume $\|f + g\|_p > 0$. Observe that

$$\| |f + g|^{p-1} \|_q = \mu(|f + g|^{(p-1)q})^{1/q} = \mu(|f + g|^p)^{1-1/p}.$$

So, by Hölder's inequality,

$$\begin{aligned} \mu(|f + g|^p) &\leq \mu(|f||f + g|^{p-1}) + \mu(|g||f + g|^{p-1}) \\ &\leq (\|f\|_p + \|g\|_p) \| |f + g|^{p-1} \|_q. \end{aligned}$$

The result follows on dividing both sides by $\| |f + g|^{p-1} \|_q$. □

5. COMPLETENESS OF L^p AND ORTHOGONAL PROJECTION

5.1. \mathcal{L}^p as a Banach space. Let V be a vector space. A map $v \mapsto \|v\| : V \rightarrow [0, \infty)$ is a *norm* if

- (i) $\|u + v\| \leq \|u\| + \|v\|$ for all $u, v \in V$,
- (ii) $\|\alpha v\| = |\alpha| \|v\|$ for all $v \in V$ and $\alpha \in \mathbb{R}$,
- (iii) $\|v\| = 0$ implies $v = 0$.

We note that, for any norm, if $\|v_n - v\| \rightarrow 0$ then $\|v_n\| \rightarrow \|v\|$.

A symmetric bilinear map $(u, v) \mapsto \langle u, v \rangle : V \times V \rightarrow \mathbb{R}$ is an *inner product* if $\langle v, v \rangle \geq 0$, with equality only if $v = 0$. For any inner product, $\langle \cdot, \cdot \rangle$, the map $v \mapsto \sqrt{\langle v, v \rangle}$ is a norm, by the Cauchy–Schwarz inequality.

Minkowski's inequality shows that each L^p space is a vector space and that the L^p -norms satisfy condition (i) above. Condition (ii) also holds. Condition (iii) fails, because $\|f\|_p = 0$ does not imply that $f = 0$, only that $f = 0$ a.e.. However, it is possible to make the L^p -norms into true norms by quotienting out by the subspace of measurable functions vanishing a.e.. This quotient will be denoted \mathcal{L}^p . Note that, for $f \in L^2$, we have $\|f\|_2^2 = \langle f, f \rangle$, where $\langle \cdot, \cdot \rangle$ is the symmetric bilinear form on L^2 given by

$$\langle f, g \rangle = \int_E f g d\mu.$$

Thus \mathcal{L}^2 is an inner product space. The notion of convergence in L^p defined in §4.1 is the usual notion of convergence in a normed space.

A normed vector space V is *complete* if every Cauchy sequence in V converges, that is to say, given any sequence $(v_n : n \in \mathbb{N})$ in V such that $\|v_n - v_m\| \rightarrow 0$ as $n, m \rightarrow \infty$, there exists $v \in V$ such that $\|v_n - v\| \rightarrow 0$ as $n \rightarrow \infty$. A complete normed vector space is called a *Banach space*. A complete inner product space is called a *Hilbert space*. Such spaces have many useful properties, which makes the following result important.

Theorem 5.1.1 (Completeness of L^p). *Let $p \in [1, \infty]$. Let $(f_n : n \in \mathbb{N})$ be a sequence in L^p such that*

$$\|f_n - f_m\|_p \rightarrow 0 \quad \text{as } n, m \rightarrow \infty.$$

Then there exists $f \in L^p$ such that

$$\|f_n - f\|_p \rightarrow 0 \quad \text{as } n \rightarrow \infty.$$

Proof. Some modifications of the following argument are necessary in the case $p = \infty$, which are left as an exercise. We assume from now on that $p < \infty$. Choose a subsequence (n_k) such that

$$S = \sum_{k=1}^{\infty} \|f_{n_{k+1}} - f_{n_k}\|_p < \infty.$$

By Minkowski's inequality, for any $K \in \mathbb{N}$,

$$\left\| \sum_{k=1}^K |f_{n_{k+1}} - f_{n_k}| \right\|_p \leq S.$$

By monotone convergence this bound holds also for $K = \infty$, so

$$\sum_{k=1}^{\infty} |f_{n_{k+1}} - f_{n_k}| < \infty \quad \text{a.e..}$$

Hence, by completeness of \mathbb{R} , f_{n_k} converges a.e.. We define a measurable function f by

$$f(x) = \begin{cases} \lim f_{n_k}(x) & \text{if the limit exists,} \\ 0 & \text{otherwise.} \end{cases}$$

Given $\varepsilon > 0$, we can find N so that $n \geq N$ implies

$$\mu(|f_n - f_m|^p) \leq \varepsilon, \quad \text{for all } m \geq n,$$

in particular $\mu(|f_n - f_{n_k}|^p) \leq \varepsilon$ for all sufficiently large k . Hence, by Fatou's lemma, for $n \geq N$,

$$\mu(|f_n - f|^p) = \mu(\liminf_k |f_n - f_{n_k}|^p) \leq \liminf_k \mu(|f_n - f_{n_k}|^p) \leq \varepsilon.$$

Hence $f \in L^p$ and, since $\varepsilon > 0$ was arbitrary, $\|f_n - f\|_p \rightarrow 0$. □

Corollary 5.1.2. *We have*

- (a) \mathcal{L}^p is a Banach space, for all $1 \leq p \leq \infty$,
- (b) \mathcal{L}^2 is a Hilbert space.

5.2. \mathcal{L}^2 as a Hilbert space. We shall apply some general Hilbert space arguments to L^2 . First, we note *Pythagoras' rule*

$$\|f + g\|_2^2 = \|f\|_2^2 + 2\langle f, g \rangle + \|g\|_2^2$$

and the *parallelogram law*

$$\|f + g\|_2^2 + \|f - g\|_2^2 = 2(\|f\|_2^2 + \|g\|_2^2).$$

If $\langle f, g \rangle = 0$, then we say that f and g are *orthogonal*. For any subset $V \subseteq L^2$, we define

$$V^\perp = \{f \in L^2 : \langle f, v \rangle = 0 \text{ for all } v \in V\}.$$

A subset $V \subseteq L^2$ is *closed* if, for every sequence $(f_n : n \in \mathbb{N})$ in V , with $f_n \rightarrow f$ in L^2 , we have $f = v$ a.e., for some $v \in V$.

Theorem 5.2.1 (Orthogonal projection). *Let V be a closed subspace of L^2 . Then each $f \in L^2$ has a decomposition $f = v + u$ a.e., with $v \in V$ and $u \in V^\perp$. Moreover, $\|f - v\|_2 \leq \|f - g\|_2$ for all $g \in V$, with equality only if $g = v$ a.e..*

The function v is called (a version of) the *orthogonal projection of f on V* .

Proof. Choose a sequence $g_n \in V$ such that

$$\|f - g_n\|_2 \rightarrow d(f, V) = \inf\{\|f - g\|_2 : g \in V\}.$$

By the parallelogram law,

$$\|2(f - (g_n + g_m)/2)\|_2^2 + \|g_n - g_m\|_2^2 = 2(\|f - g_n\|_2^2 + \|f - g_m\|_2^2).$$

But $\|2(f - (g_n + g_m)/2)\|_2^2 \geq 4d(f, V)^2$, so we must have $\|g_n - g_m\|_2 \rightarrow 0$ as $n, m \rightarrow \infty$. By completeness, $\|g_n - g\|_2 \rightarrow 0$, for some $g \in L^2$. By closure, $g = v$ a.e., for some $v \in V$. Hence

$$\|f - v\|_2 = \lim_n \|f - g_n\|_2 = d(f, V).$$

Now, for any $h \in V$ and $t \in \mathbb{R}$, we have

$$d(f, V)^2 \leq \|f - (v + th)\|_2^2 = d(f, V)^2 - 2t\langle f - v, h \rangle + t^2\|h\|_2^2.$$

So we must have $\langle f - v, h \rangle = 0$. Hence $u = f - v \in V^\perp$, as required. \square

5.3. Variance, covariance and conditional expectation. In this section we look at some L^2 notions relevant to probability. For $X, Y \in L^2(\mathbb{P})$, with means $m_X = \mathbb{E}(X)$, $m_Y = \mathbb{E}(Y)$, we define *variance*, *covariance* and *correlation* by

$$\begin{aligned}\text{var}(X) &= \mathbb{E}[(X - m_X)^2], \\ \text{cov}(X, Y) &= \mathbb{E}[(X - m_X)(Y - m_Y)], \\ \text{corr}(X, Y) &= \text{cov}(X, Y) / \sqrt{\text{var}(X) \text{var}(Y)}.\end{aligned}$$

Note that $\text{var}(X) < \infty$ by Minkowski's inequality, and $\text{var}(X) = 0$ if and only if $X = m_X$ a.s.. Note also that, if X and Y are independent, then $\text{cov}(X, Y) = 0$. The converse is generally false. Finally, note that $\text{var}(X) = \text{cov}(X, X)$ and $\text{var}(X) = \mathbb{E}(X^2) - \mathbb{E}(X)^2$, a formula which is easier to use in practice for computing the variance.

Let $X \in L^2(\mathbb{P})$, and let $m = \mathbb{E}(X)$, which exists since necessarily $X \in L^1(\mathbb{P})$. Then we obtain another and very useful form of Chebyshev's inequality: for all $A > 0$,

$$(5.1) \quad \mathbb{P}(|X - m| > A) \leq \frac{\text{var}(X)}{A^2}.$$

To obtain (5.1), apply Chebyshev's inequality to the random variable $(X - m)^2$. The number $\sigma = \sqrt{\text{var}(X)}$ is called the standard deviation of X . Then (5.1) says that if $X \in L^2(\mathbb{P})$, then X has fluctuations of the order of the standard deviation: for instance $\mathbb{P}(|X - m| \geq k\sigma) \leq 1/k^2$ for all $k > 0$.

For a random variable $X = (X_1, \dots, X_n)$ in \mathbb{R}^n , we define its *covariance matrix*

$$\text{var}(X) = (\text{cov}(X_i, X_j))_{i,j=1}^n.$$

Proposition 5.3.1. *Every covariance matrix is symmetric non-negative definite.*

Suppose now we are given a countable family of disjoint events $(G_i : i \in I)$, whose union is Ω . Set $\mathcal{G} = \sigma(G_i : i \in I)$. Let X be an integrable random variable. The *conditional expectation* of X given \mathcal{G} is given by

$$Y = \sum_i \mathbb{E}(X|G_i) 1_{G_i},$$

where we set $\mathbb{E}(X|G_i) = \mathbb{E}(X 1_{G_i}) / \mathbb{P}(G_i)$ when $\mathbb{P}(G_i) > 0$, and define $\mathbb{E}(X|G_i)$ in some arbitrary way when $\mathbb{P}(G_i) = 0$. Set $V = L^2(\mathcal{G}, \mathbb{P})$ and note that $Y \in V$. Then V is a subspace of $L^2(\mathcal{F}, \mathbb{P})$, and V is complete and therefore closed.

Proposition 5.3.2. *If $X \in L^2$, then Y is a version of the orthogonal projection of X on V . In particular, Y depends only on \mathcal{G} .*

6. CONVERGENCE IN $L^1(\mathbb{P})$

6.1. Bounded convergence. We begin with a basic, but easy to use, condition for convergence in $L^1(\mathbb{P})$.

Theorem 6.1.1 (Bounded convergence). *Let $(X_n : n \in \mathbb{N})$ be a sequence of random variables, with $X_n \rightarrow X$ in probability and $|X_n| \leq C$ for all n , for some constant $C < \infty$. Then $X_n \rightarrow X$ in L^1 .*

Proof. By Theorem 2.6.1, X is the almost sure limit of a subsequence, so $|X| \leq C$ a.s.. For $\varepsilon > 0$, there exists N such that $n \geq N$ implies

$$\mathbb{P}(|X_n - X| > \varepsilon/2) \leq \varepsilon/(4C).$$

Then

$$\mathbb{E}|X_n - X| = \mathbb{E}(|X_n - X|1_{|X_n - X| > \varepsilon/2}) + \mathbb{E}(|X_n - X|1_{|X_n - X| \leq \varepsilon/2}) \leq 2C(\varepsilon/4C) + \varepsilon/2 = \varepsilon.$$

□

6.2. Uniform integrability.

Lemma 6.2.1. *Let X be an integrable random variable and set*

$$I_X(\delta) = \sup\{\mathbb{E}(|X|1_A) : A \in \mathcal{F}, \mathbb{P}(A) \leq \delta\}.$$

Then $I_X(\delta) \downarrow 0$ as $\delta \downarrow 0$.

Proof. Suppose not. Then, for some $\varepsilon > 0$, there exist $A_n \in \mathcal{F}$, with $\mathbb{P}(A_n) \leq 2^{-n}$ and $\mathbb{E}(|X|1_{A_n}) \geq \varepsilon$ for all n . By the first Borel–Cantelli lemma, $\mathbb{P}(A_n \text{ i.o.}) = 0$. But then, by dominated convergence,

$$\varepsilon \leq \mathbb{E}(|X|1_{\bigcup_{m \geq n} A_m}) \rightarrow \mathbb{E}(|X|1_{\{A_n \text{ i.o.}\}}) = 0$$

which is a contradiction. □

Let \mathcal{X} be a family of random variables. For $1 \leq p \leq \infty$, we say that \mathcal{X} is *bounded in L^p* if $\sup_{X \in \mathcal{X}} \|X\|_p < \infty$. Let us define

$$I_{\mathcal{X}}(\delta) = \sup\{\mathbb{E}(|X|1_A) : X \in \mathcal{X}, A \in \mathcal{F}, \mathbb{P}(A) \leq \delta\}.$$

Obviously, \mathcal{X} is bounded in L^1 if and only if $I_{\mathcal{X}}(1) < \infty$. We say that \mathcal{X} is *uniformly integrable* or *UI* if \mathcal{X} is bounded in L^1 and

$$I_{\mathcal{X}}(\delta) \downarrow 0, \quad \text{as } \delta \downarrow 0.$$

Note that, by Hölder’s inequality, for conjugate indices $p, q \in (1, \infty)$,

$$\mathbb{E}(|X|1_A) \leq \|X\|_p (\mathbb{P}(A))^{1/q}.$$

Hence, if \mathcal{X} is bounded in L^p , for some $p \in (1, \infty)$, then \mathcal{X} is UI. It is important to assume $p > 1$: the sequence $X_n = n1_{(0,1/n]}$ is bounded in L^1 for Lebesgue measure on $(0, 1]$, but not uniformly integrable.

Lemma 6.2.1 shows that any single integrable random variable is uniformly integrable. This extends easily to any finite collection of integrable random variables. Moreover, for any integrable random variable Y , the set

$$\mathcal{X} = \{X : X \text{ a random variable, } |X| \leq Y\}$$

is uniformly integrable, because $\mathbb{E}(|X|1_A) \leq \mathbb{E}(Y1_A)$ for all A .

The following result gives an alternative characterization of uniform integrability, which is often more easy to check in practice.

Lemma 6.2.2. *Let \mathcal{X} be a family of random variables. Then \mathcal{X} is UI if and only if*

$$\sup\{\mathbb{E}(|X|1_{|X| \geq K}) : X \in \mathcal{X}\} \rightarrow 0, \quad \text{as } K \rightarrow \infty.$$

Proof. Suppose \mathcal{X} is UI. Given $\varepsilon > 0$, choose $\delta > 0$ so that $I_{\mathcal{X}}(\delta) < \varepsilon$, then choose $K < \infty$ so that $I_{\mathcal{X}}(1) \leq K\delta$. Then, for $X \in \mathcal{X}$ and $A = \{|X| \geq K\}$, we have $\mathbb{P}(A) \leq \delta$ so $\mathbb{E}(|X|1_A) < \varepsilon$. Hence, as $K \rightarrow \infty$,

$$\sup\{\mathbb{E}(|X|1_{|X| \geq K}) : X \in \mathcal{X}\} \rightarrow 0.$$

On the other hand, if this condition holds, then, since

$$\mathbb{E}(|X|) \leq K + \mathbb{E}(|X|1_{|X| \geq K}),$$

we have $I_{\mathcal{X}}(1) < \infty$. Given $\varepsilon > 0$, choose $K < \infty$ so that $\mathbb{E}(|X|1_{|X| \geq K}) < \varepsilon/2$ for all $X \in \mathcal{X}$. Then choose $\delta > 0$ so that $K\delta < \varepsilon/2$. For all $X \in \mathcal{X}$ and $A \in \mathcal{F}$ with $\mathbb{P}(A) < \delta$, we have

$$\mathbb{E}(|X|1_A) \leq \mathbb{E}(|X|1_{|X| \geq K}) + K\mathbb{P}(A) < \varepsilon.$$

Hence \mathcal{X} is UI. □

Here is the definitive result on L^1 -convergence of random variables.

Theorem 6.2.3. *Let X be a random variable and let $(X_n : n \in \mathbb{N})$ be a sequence of random variables. The following are equivalent:*

- (a) $X_n \in L^1$ for all n , $X \in L^1$ and $X_n \rightarrow X$ in L^1 ,
- (b) $\{X_n : n \in \mathbb{N}\}$ is UI and $X_n \rightarrow X$ in probability.

Proof. Suppose (a) holds. By Chebyshev's inequality, for $\varepsilon > 0$,

$$\mathbb{P}(|X_n - X| > \varepsilon) \leq \varepsilon^{-1} \mathbb{E}(|X_n - X|) \rightarrow 0$$

so $X_n \rightarrow X$ in probability. Moreover, given $\varepsilon > 0$, there exists N such that $\mathbb{E}(|X_n - X|) < \varepsilon/2$ whenever $n \geq N$. Then we can find $\delta > 0$ so that $\mathbb{P}(A) \leq \delta$ implies

$$\mathbb{E}(|X|1_A) \leq \varepsilon/2, \quad \mathbb{E}(|X_n|1_A) \leq \varepsilon, \quad n = 1, \dots, N.$$

Then, for $n \geq N$ and $\mathbb{P}(A) \leq \delta$,

$$\mathbb{E}(|X_n|1_A) \leq \mathbb{E}(|X_n - X|) + \mathbb{E}(|X|1_A) \leq \varepsilon.$$

Hence $\{X_n : n \in \mathbb{N}\}$ is UI. We have shown that (a) implies (b).

Suppose, on the other hand, that (b) holds. Then there is a subsequence (n_k) such that $X_{n_k} \rightarrow X$ a.s.. So, by Fatou's lemma, $\mathbb{E}(|X|) \leq \liminf_k \mathbb{E}(|X_{n_k}|) < \infty$. Now, given $\varepsilon > 0$, there exists $K < \infty$ such that, for all n ,

$$\mathbb{E}(|X_n|1_{|X_n| \geq K}) < \varepsilon/3, \quad \mathbb{E}(|X|1_{|X| \geq K}) < \varepsilon/3.$$

Consider the uniformly bounded sequence $X_n^K = (-K) \vee X_n \wedge K$ and set $X^K = (-K) \vee X \wedge K$. Then $X_n^K \rightarrow X^K$ in probability, so, by bounded convergence, there exists N such that, for all $n \geq N$,

$$\mathbb{E}|X_n^K - X^K| < \varepsilon/3.$$

But then, for all $n \geq N$,

$$\mathbb{E}|X_n - X| \leq \mathbb{E}(|X_n|1_{|X_n| \geq K}) + \mathbb{E}|X_n^K - X^K| + \mathbb{E}(|X|1_{|X| \geq K}) < \varepsilon.$$

Since $\varepsilon > 0$ was arbitrary, we have shown that (b) implies (a). \square

7. CHARACTERISTIC FUNCTIONS AND WEAK CONVERGENCE

7.1. Definitions. For a finite Borel measure μ on \mathbb{R}^n , we define the *Fourier transform*

$$\hat{\mu}(u) = \int_{\mathbb{R}^n} e^{i\langle u, x \rangle} \mu(dx), \quad u \in \mathbb{R}^n.$$

Here, $\langle \cdot, \cdot \rangle$ denotes the usual inner product on \mathbb{R}^n . Note that $\hat{\mu}$ is in general complex-valued, with $\overline{\hat{\mu}(u)} = \hat{\mu}(-u)$, and $\|\hat{\mu}\|_\infty = \hat{\mu}(0) = \mu(\mathbb{R}^n)$. Moreover, by bounded convergence, $\hat{\mu}$ is continuous on \mathbb{R}^n .

For a random variable X in \mathbb{R}^n , we define the *characteristic function*

$$\phi_X(u) = \mathbb{E}(e^{i\langle u, X \rangle}), \quad u \in \mathbb{R}^n.$$

Thus $\phi_X = \hat{\mu}_X$, where μ_X is the law of X .

A random variable X in \mathbb{R}^n is *standard Gaussian* if

$$\mathbb{P}(X \in A) = \int_A \frac{1}{(2\pi)^{n/2}} e^{-|x|^2/2} dx, \quad A \in \mathcal{B}.$$

Let us compute the characteristic function of a standard Gaussian random variable X in \mathbb{R} . We have

$$\phi_X(u) = \int_{\mathbb{R}} e^{iux} \frac{1}{\sqrt{2\pi}} e^{-x^2/2} dx = e^{-u^2/2} I$$

where

$$I = \int_{\mathbb{R}} \frac{1}{\sqrt{2\pi}} e^{-(x-iu)^2/2} dx.$$

The integral I can be evaluated by considering the integral of the analytic function $e^{-z^2/2}$ around the rectangular contour with corners $R, R - iu, -R - iu, -R$: by Cauchy's theorem, the integral round the contour vanishes, as do, in the limit $R \rightarrow \infty$, the contributions from the vertical sides of the rectangle. We deduce that

$$I = \int_{\mathbb{R}} \frac{1}{\sqrt{2\pi}} e^{-x^2/2} dx = 1.$$

Hence $\phi_X(u) = e^{-u^2/2}$.

7.2. Uniqueness and inversion. We now show that a finite Borel measure is determined uniquely by its Fourier transform and obtain, where possible, an *inversion* formula by which to compute the measure from its transform.

Define, for $t > 0$ and $x, y \in \mathbb{R}^n$, the *heat kernel*

$$p(t, x, y) = \frac{1}{(2\pi t)^{n/2}} e^{-|y-x|^2/2t} = \prod_{k=1}^n \frac{1}{\sqrt{2\pi t}} e^{-|y_k - x_k|^2/2t}.$$

(This is the fundamental solution for the heat equation $(\partial/\partial t - (1/2)\Delta)p = 0$ in \mathbb{R}^n , but we shall not pursue this property here.)

Lemma 7.2.1. *Let Z be a standard Gaussian random variable in \mathbb{R}^n . Let $x \in \mathbb{R}^n$ and $t \in (0, \infty)$. Then*

- (a) *the random variable $x + \sqrt{t}Z$ has density function $p(t, x, \cdot)$ on \mathbb{R}^n ,*
- (b) *for all $y \in \mathbb{R}^n$, we have*

$$p(t, x, y) = \frac{1}{(2\pi)^n} \int_{\mathbb{R}^n} e^{i\langle u, x \rangle} e^{-|u|^2 t/2} e^{-i\langle u, y \rangle} du.$$

Proof. The component random variables $Y_k = x_k + \sqrt{t}Z_k$ are independent Gaussians with mean x_k and variance t (see Subsection 8.1). So Y_k has density

$$\frac{1}{\sqrt{2\pi t}} e^{-|y_k - x_k|^2/2t}$$

and we obtain the claimed density function for Y as the product of the marginal densities.

For $u \in \mathbb{R}$ and $t \in (0, \infty)$, we know that

$$\int_{\mathbb{R}} e^{iux} \frac{1}{\sqrt{2\pi t}} e^{-|x|^2/2t} dx = \mathbb{E}(e^{iu\sqrt{t}Z_1}) = e^{u^2 t/2}.$$

By relabelling the variables we obtain, for $x_k, y_k, u_k \in \mathbb{R}$ and $t \in (0, \infty)$,

$$\int_{\mathbb{R}} e^{iu_k(x_k - y_k)} \frac{\sqrt{t}}{\sqrt{2\pi}} e^{-t|u_k|^2/2} du_k = e^{(x_k - y_k)^2/2t},$$

so

$$\frac{1}{\sqrt{2\pi t}} e^{-|y_k - x_k|^2/2t} = \frac{1}{2\pi} \int_{\mathbb{R}} e^{iu_k x_k} e^{-u_k^2 t/2} e^{-i\langle u_k, y_k \rangle} du_k.$$

On taking the product over $k \in \{1, \dots, n\}$, we obtain the claimed formula for $p(t, x, y)$. \square

Theorem 7.2.2. *Let X be a random variable in \mathbb{R}^n . The law μ_X of X is uniquely determined by its characteristic function ϕ_X . Moreover, if ϕ_X is integrable, and we define*

$$f_X(x) = \frac{1}{(2\pi)^n} \int_{\mathbb{R}^n} \phi_X(u) e^{-i\langle u, x \rangle} du,$$

then f_X is a continuous, bounded and non-negative function, which is a density function for X .

Proof. Let Z be a standard Gaussian random variable in \mathbb{R}^n , independent of X , and let g be a continuous function on \mathbb{R}^n of compact support. Then, for any $t \in (0, \infty)$, by Fubini's theorem,

$$\mathbb{E}(g(X + \sqrt{t}Z)) = \int_{\mathbb{R}^n} \int_{\mathbb{R}^n} g(x + \sqrt{t}z) (2\pi)^{-n/2} e^{-|z|^2/2} dz \mu_X(dx).$$

By the lemma, we have

$$\begin{aligned} \int_{\mathbb{R}^n} g(x + \sqrt{t}z) (2\pi)^{-n/2} e^{-|z|^2/2} dz &= \mathbb{E}(g(x + \sqrt{t}Z)) \\ &= \int_{\mathbb{R}^n} g(y) p(t, x, y) dy = \int_{\mathbb{R}^n} g(y) \frac{1}{(2\pi)^n} \int_{\mathbb{R}^n} e^{i\langle u, x \rangle} e^{-|u|^2 t/2} e^{-i\langle u, y \rangle} du dy, \end{aligned}$$

so, by Fubini again,

$$\mathbb{E}(g(X + \sqrt{t}Z)) = \int_{\mathbb{R}^n} \left(\frac{1}{(2\pi)^n} \int_{\mathbb{R}^n} \phi_X(u) e^{-|u|^2 t/2} e^{-i\langle u, y \rangle} du \right) g(y) dy.$$

By this formula, ϕ_X determines $\mathbb{E}(g(X + \sqrt{t}Z))$. For any such function g , by bounded convergence, we have

$$\mathbb{E}(g(X + \sqrt{t}Z)) \rightarrow \mathbb{E}(g(X))$$

as $t \downarrow 0$, so ϕ_X determines $\mathbb{E}(g(X))$. Hence ϕ_X determines μ_X .

Suppose now that ϕ_X is integrable. Then

$$|\phi_X(u)| |g(y)| \in L^1(du \otimes dy).$$

So, by Fubini's theorem, $g \cdot f_X \in L^1$ and, by dominated convergence, as $t \downarrow 0$,

$$\begin{aligned} \int_{\mathbb{R}^n} \left(\frac{1}{(2\pi)^n} \int_{\mathbb{R}^n} \phi_X(u) e^{-|u|^2 t/2} e^{-i\langle u, y \rangle} du \right) g(y) dy \\ \rightarrow \int_{\mathbb{R}^n} \left(\frac{1}{(2\pi)^n} \int_{\mathbb{R}^n} \phi_X(u) e^{-i\langle u, y \rangle} du \right) g(y) dy = \int_{\mathbb{R}^n} g(x) f_X(x) dx. \end{aligned}$$

Hence we obtain the identity

$$\mathbb{E}(g(X)) = \int_{\mathbb{R}^n} g(x) f_X(x) dx.$$

Since $\overline{\phi_X(u)} = \phi_X(-u)$, we have $\bar{f}_X = f_X$, so f_X is real-valued. Moreover, f_X is continuous, by bounded convergence, and $\|f_X\|_\infty \leq (2\pi)^{-n} \|\phi_X\|_1$. Since f_X is

continuous, if it took a negative value anywhere, it would do so on an open interval of positive length, I say. There would exist a continuous function g , positive on I and vanishing outside I . Then we would have $\mathbb{E}(g(X)) \geq 0$ and $\int_{\mathbb{R}^n} g(x)f_X(x)dx < 0$, a contradiction. Hence, f_X is non-negative. It is now straightforward to extend the identity to all bounded Borel functions g , by a monotone class argument. In particular f_X is a density function for X . \square

7.3. Characteristic functions and independence.

Theorem 7.3.1. *Let $X = (X_1, \dots, X_n)$ be a random variable in \mathbb{R}^n . Then the following are equivalent:*

- (a) X_1, \dots, X_n are independent,
- (b) $\mu_X = \mu_{X_1} \otimes \dots \otimes \mu_{X_n}$,
- (c) $\mathbb{E}(\prod_k f_k(X_k)) = \prod_k \mathbb{E}(f_k(X_k))$, for all bounded Borel functions f_1, \dots, f_n ,
- (d) $\phi_X(u) = \prod_k \phi_{X_k}(u_k)$, for all $u = (u_1, \dots, u_n) \in \mathbb{R}^n$.

Proof. If (a) holds, then

$$\mu_X(A_1 \times \dots \times A_n) = \prod_k \mu_{X_k}(A_k)$$

for all Borel sets A_1, \dots, A_n , so (b) holds, since this formula characterizes the product measure.

If (b) holds, then, for f_1, \dots, f_n bounded Borel, by Fubini's theorem,

$$\mathbb{E}\left(\prod_k f_k(X_k)\right) = \int_{\mathbb{R}^n} \prod_k f_k(x_k) \mu_X(dx) = \prod_k \int_{\mathbb{R}} f_k(x_k) \mu_{X_k}(dx_k) = \prod_k \mathbb{E}(f_k(X_k)),$$

so (c) holds. Statement (d) is a special case of (c). Suppose, finally, that (d) holds and take independent random variables $\tilde{X}_1, \dots, \tilde{X}_n$ with $\mu_{\tilde{X}_k} = \mu_{X_k}$ for all k . We know that (a) implies (d), so

$$\phi_{\tilde{X}}(u) = \prod_k \phi_{\tilde{X}_k}(u_k) = \prod_k \phi_{X_k}(u_k) = \phi_X(u)$$

so $\mu_{\tilde{X}} = \mu_X$ by uniqueness of characteristic functions. Hence (a) holds. \square

7.4. Weak convergence. Let F be a distribution function on \mathbb{R} and let $(F_n : n \in \mathbb{N})$ be a sequence of such distribution functions. We say that $F_n \rightarrow F$ *as distribution functions* if $F_{X_n}(x) \rightarrow F_X(x)$ at every point $x \in \mathbb{R}$ where F_X is continuous.

Let μ be a Borel probability measure on \mathbb{R} and let $(\mu_n : n \in \mathbb{N})$ be a sequence of such measures. We say that $\mu_n \rightarrow \mu$ *weakly* if $\mu_n(f) \rightarrow \mu(f)$ for all continuous bounded functions f on \mathbb{R} .

Let X be a random variable and let $(X_n : n \in \mathbb{N})$ be a sequence of random variables. We say that $X_n \rightarrow X$ *in distribution* if $F_{X_n} \rightarrow F_X$.

7.5. Equivalent modes of convergence. Suppose we are given a sequence of Borel probability measures $(\mu_n : n \in \mathbb{N})$ on \mathbb{R} and a further such measure μ . Write F_n and F for the corresponding distribution functions, and write ϕ_n and ϕ for the corresponding Fourier transforms.

Theorem 7.5.1. *The following are equivalent:*

- (i) $\mu_n \rightarrow \mu$ *weakly*,
- (ii) $F_n \rightarrow F$ *as distribution functions*,
- (iii) *There exists a probability space $(\Omega, \mathcal{F}, \mathbb{P})$ and random variables X_n, X on this space such that $\mu_{X_n} = \mu_n$ (resp. $\mu_X = \mu$) and $X_n \rightarrow X$ almost surely,*
- (iv) $\phi_{X_n}(u) \rightarrow \phi_X(u)$ *for all $u \in \mathbb{R}$,*

We do not prove this result in full, but note how to see certain of the implications. We first check (i) \Rightarrow (ii). Let x be a continuity point of F and let $\varepsilon > 0$. Then by continuity of F , there exists $\delta > 0$ such that $\mathbb{P}(X \in (x - 2\delta, x + 2\delta]) \leq \varepsilon$. Find f_δ a continuous bounded function such that: $f_\delta \geq 0$, $f_\delta(t) = 1$ if $t \leq x - 2\delta$, $f_\delta(t) = 0$ if $t \geq x - \delta$. Then

$$F_n(x) = \mathbb{P}(X_n \leq x) \geq \mathbb{E}(f_\delta(X_n)) \rightarrow \mathbb{E}(f_\delta(X))$$

by (a). Hence $\liminf F_n(x) \geq \mathbb{E}(f_\delta(X)) \geq F(x - 2\delta) \geq F(x) - \varepsilon$. Since $\varepsilon > 0$ was arbitrary, it follows $\liminf F_n(x) \geq F(x)$. A similar argument for the limsup concludes (ii).

(ii) \Rightarrow (iii). The idea is to consider the canonical space: let $(\Omega, \mathcal{F}, \mathbb{P}) = ([0, 1], \mathcal{B}([0, 1]), dx)$, define random variables $X_n(\omega) = \inf\{x \in \mathbb{R} : \omega \leq F_n(x)\}$ and $X(\omega) = \inf\{x \in \mathbb{R} : \omega \leq F(x)\}$, as in Subsection 2.3. Recall that $X_n \sim \mu_n$ and $X \sim \mu$. Since X_n and X are essentially inverse functions to F_n and F , the almost sure convergence follows with some care; the key is that the set of discontinuity points of a nondecreasing function is countable and hence has zero Lebesgue measure. Recall that $\omega \leq F_n(x)$ if and only if $X_n(\omega) \leq x$; and that $\omega \leq F(X(\omega))$. Let $\omega \in (0, 1)$. Let $\varepsilon > 0$. Let $x \in \mathbb{R}$ such that $X(\omega) - \varepsilon < x < X(\omega)$ and $\mathbb{P}(X = x) = 0$. Since $X(\omega) > x$ we see that $\omega > F(x)$. Since $F_n(x) \rightarrow F(x)$ we have that for n large enough $F_n(x) < \omega$ as well. Thus $x < X_n(\omega)$ for n large enough, and hence $\liminf X_n(\omega) \geq x \geq X(\omega) - \varepsilon$. Thus $\liminf X_n(\omega) \geq X(\omega)$. Conversely, assume that $X : \Omega \rightarrow \mathbb{R}$ is continuous at ω . Let $\omega' > \omega$ and fix $\varepsilon > 0$. Fix y such that $X(\omega') < y < X(\omega') + \varepsilon$ and $\mathbb{P}(X = y) = 0$.

Then $\omega < \omega' \leq F(X(\omega')) \leq F(y)$. Thus $\omega < F(y)$ and hence for n large enough $F_n(y) > \omega$ as well, i.e., $X_n(\omega) < y < X(\omega') + \varepsilon$. Hence $\limsup X_n(\omega) \leq X(\omega') + \varepsilon$. Since ε was arbitrary and X is continuous at ω , we get that $\lim X_n(\omega) = X(\omega)$. Since X is nondecreasing, the set of discontinuities of X has zero Lebesgue measure and we have almost sure convergence.

(iii) \Rightarrow (i) follows directly by dominated convergence.

At this stage (i), (ii) and (iii) are equivalent. It suffices to prove any of those is equivalent to (iv).

(i) \Rightarrow (iv) is trivial since e^{ix} is continuous and bounded.

The fact that (iv) implies (ii) is a consequence of the following famous result, which we do not prove here.

Theorem 7.5.2 (Lévy's continuity theorem for characteristic functions). *Let $(X_n : n \in \mathbb{N})$ be a sequence of random variables and suppose that $\phi_{X_n}(u)$ converges as $n \rightarrow \infty$, with limit $\phi(u)$ say, for all $u \in \mathbb{R}$. If ϕ is continuous in a neighbourhood of 0, then it is the characteristic function of some random variable X , and $X_n \rightarrow X$ in distribution.*

8. GAUSSIAN RANDOM VARIABLES

8.1. Gaussian random variables in \mathbb{R} . A random variable X in \mathbb{R} is *Gaussian* if, for some $\mu \in \mathbb{R}$ and some $\sigma^2 \in (0, \infty)$, X has density function

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

We also admit as Gaussian any random variable X with $X = \mu$ a.s., this degenerate case corresponding to taking $\sigma^2 = 0$. We write $X \sim N(\mu, \sigma^2)$.

Proposition 8.1.1. *Suppose $X \sim N(\mu, \sigma^2)$ and $a, b \in \mathbb{R}$. Then (a) $\mathbb{E}(X) = \mu$, (b) $\text{var}(X) = \sigma^2$, (c) $aX + b \sim N(a\mu + b, a^2\sigma^2)$, (d) $\phi_X(u) = e^{iu\mu - u^2\sigma^2/2}$.*

8.2. Gaussian random variables in \mathbb{R}^n . A random variable X in \mathbb{R}^n is *Gaussian* if $\langle u, X \rangle$ is Gaussian, for all $u \in \mathbb{R}^n$. An example of such a random variable is provided by $X = (X_1, \dots, X_n)$, where X_1, \dots, X_n are independent $N(0, 1)$ random variables. To see this, we note that

$$\mathbb{E} e^{i\langle u, X \rangle} = \mathbb{E} \prod_k e^{iu_k X_k} = e^{-|u|^2/2}$$

so $\langle u, X \rangle$ is $N(0, |u|^2)$ for all $u \in \mathbb{R}^n$.

Theorem 8.2.1. *Let X be a Gaussian random variable in \mathbb{R}^n . Then*

- (a) $AX + b$ is a Gaussian random variable in \mathbb{R}^m , for all $m \times n$ matrices A and all $b \in \mathbb{R}^m$,
- (b) $X \in L^2$ and its distribution is determined by its mean μ and its covariance matrix V ,

- (c) $\phi_X(u) = e^{i\langle u, \mu \rangle - \langle u, Vu \rangle / 2}$,
(d) if V is invertible, then X has a density function on \mathbb{R}^n , given by

$$f_X(x) = (2\pi)^{-n/2} (\det V)^{-1/2} \exp\{-\langle x - \mu, V^{-1}(x - \mu) \rangle / 2\},$$

- (e) suppose $X = (X_1, X_2)$, with X_1 in \mathbb{R}^{n_1} and X_2 in \mathbb{R}^{n_2} , then

$$\text{cov}(X_1, X_2) = 0 \quad \text{implies} \quad X_1, X_2 \text{ independent.}$$

Proof. For $u \in \mathbb{R}^n$, we have $\langle u, AX + b \rangle = \langle A^T u, X \rangle + \langle u, b \rangle$ so $\langle u, AX + b \rangle$ is Gaussian, by Proposition 8.1.1. This proves (a).

Each component X_k is Gaussian, so $X \in L^2$. Set $\mu = \mathbb{E}(X)$ and $V = \text{var}(X)$. For $u \in \mathbb{R}^n$ we have $\mathbb{E}(\langle u, X \rangle) = \langle u, \mu \rangle$ and $\text{var}(\langle u, X \rangle) = \text{cov}(\langle u, X \rangle, \langle u, X \rangle) = \langle u, Vu \rangle$. Since $\langle u, X \rangle$ is Gaussian, by Proposition 8.1.1, we must have $\langle u, X \rangle \sim N(\langle u, \mu \rangle, \langle u, Vu \rangle)$ and $\phi_X(u) = \mathbb{E} e^{i\langle u, X \rangle} = e^{i\langle u, \mu \rangle - \langle u, Vu \rangle / 2}$. This is (c) and (b) follows by uniqueness of characteristic functions.

Let Y_1, \dots, Y_n be independent $N(0, 1)$ random variables. Then $Y = (Y_1, \dots, Y_n)$ has density

$$f_Y(y) = (2\pi)^{-n/2} \exp\{-|y|^2/2\}.$$

Set $\tilde{X} = V^{1/2}Y + \mu$, then \tilde{X} is Gaussian, with $\mathbb{E}(\tilde{X}) = \mu$ and $\text{var}(\tilde{X}) = V$, so $\tilde{X} \sim X$. If V is invertible, then \tilde{X} and hence X has the density claimed in (d), by a linear change of variables in \mathbb{R}^n .

Finally, if $X = (X_1, X_2)$ with $\text{cov}(X_1, X_2) = 0$, then, for $u = (u_1, u_2)$, we have

$$\langle u, Vu \rangle = \langle u_1, V_{11}u_1 \rangle + \langle u_2, V_{22}u_2 \rangle,$$

where $V_{11} = \text{var}(X_1)$ and $V_{22} = \text{var}(X_2)$. Then $\phi_X(u) = \phi_{X_1}(u_1)\phi_{X_2}(u_2)$ so X_1 and X_2 are independent by Theorem 7.3.1. \square

9. ERGODIC THEORY

9.1. Measure-preserving transformations. Let (E, \mathcal{E}, μ) be a measure space. A measurable function $\theta : E \rightarrow E$ is called a *measure-preserving transformation* if

$$\mu(\theta^{-1}(A)) = \mu(A), \quad \text{for all } A \in \mathcal{E}.$$

A set $A \in \mathcal{E}$ is *invariant* if $\theta^{-1}(A) = A$. A measurable function f is *invariant* if $f = f \circ \theta$. The class of all invariant sets forms a σ -algebra, which we denote by \mathcal{E}_θ . Then f is invariant if and only if f is \mathcal{E}_θ -measurable. We say that θ is *ergodic* if \mathcal{E}_θ contains only sets of measure zero and their complements.

Here are two simple examples of measure preserving transformations.

- (i) *Translation map on the torus.* Take $E = [0, 1)^n$ with Lebesgue measure on its Borel σ -algebra, and consider addition modulo 1 in each coordinate. For $a \in E$ set

$$\theta_a(x_1, \dots, x_n) = (x_1 + a_1, \dots, x_n + a_n).$$

- (ii) *Bakers' map.* Take $E = [0, 1)$ with Lebesgue measure. Set

$$\theta(x) = 2x - \lfloor 2x \rfloor.$$

Proposition 9.1.1. *If f is integrable and θ is measure-preserving, then $f \circ \theta$ is integrable and*

$$\int_E f d\mu = \int_E f \circ \theta d\mu.$$

Proposition 9.1.2. *If θ is ergodic and f is invariant, then $f = c$ a.e., for some constant c .*

9.2. Bernoulli shifts. Let m be a probability measure on \mathbb{R} . In §2.4, we constructed a probability space $(\Omega, \mathcal{F}, \mathbb{P})$ on which there exists a sequence of independent random variables $(Y_n : n \in \mathbb{N})$, all having distribution m . Consider now the infinite product space

$$E = \mathbb{R}^{\mathbb{N}} = \{x = (x_n : n \in \mathbb{N}) : x_n \in \mathbb{R} \text{ for all } n\}$$

and the σ -algebra \mathcal{E} on E generated by the coordinate maps $X_n(x) = x_n$

$$\mathcal{E} = \sigma(X_n : n \in \mathbb{N}).$$

Note that \mathcal{E} is also generated by the π -system

$$\mathcal{A} = \left\{ \prod_{n \in \mathbb{N}} A_n : A_n \in \mathcal{B} \text{ for all } n, A_n = \mathbb{R} \text{ for sufficiently large } n \right\}.$$

Define $Y : \Omega \rightarrow E$ by $Y(\omega) = (Y_n(\omega) : n \in \mathbb{N})$. Then Y is measurable and the image measure $\mu = \mathbb{P} \circ Y^{-1}$ satisfies, for $A = \prod_{n \in \mathbb{N}} A_n \in \mathcal{A}$,

$$\mu(A) = \prod_{n \in \mathbb{N}} m(A_n).$$

By uniqueness of extension, μ is the unique measure on \mathcal{E} having this property. Note that, under the probability measure μ , the coordinate maps $(X_n : n \in \mathbb{N})$ are themselves a sequence of independent random variables with law m . The probability space (E, \mathcal{E}, μ) is called the *canonical model* for such sequences. Define the *shift map* $\theta : E \rightarrow E$ by

$$\theta(x_1, x_2, \dots) = (x_2, x_3, \dots).$$

Theorem 9.2.1. *The shift map is an ergodic measure-preserving transformation.*

Proof. The details of showing that θ is measurable and measure-preserving are left as an exercise. To see that θ is ergodic, we recall the definition of the tail σ -algebras

$$\mathcal{T}_n = \sigma(X_m : m \geq n+1), \quad \mathcal{T} = \bigcap_n \mathcal{T}_n.$$

For $A = \prod_{n \in \mathbb{N}} A_n \in \mathcal{A}$ we have

$$\theta^{-n}(A) = \{X_{n+k} \in A_k \text{ for all } k\} \in \mathcal{T}_n.$$

Since \mathcal{T}_n is a σ -algebra, it follows that $\theta^{-n}(A) \in \mathcal{T}_n$ for all $A \in \mathcal{E}$, so $\mathcal{E}_\theta \subseteq \mathcal{T}$. Hence θ is ergodic by Kolmogorov's zero-one law. \square

9.3. Birkhoff's and von Neumann's ergodic theorems. Throughout this section, (E, \mathcal{E}, μ) will denote a measure space, on which is given a measure-preserving transformation θ . Given an measurable function f , set $S_0 = 0$ and define, for $n \geq 1$,

$$S_n = S_n(f) = f + f \circ \theta + \dots + f \circ \theta^{n-1}.$$

Lemma 9.3.1 (Maximal ergodic lemma). *Let f be an integrable function on E . Set $S^* = \sup_{n \geq 0} S_n(f)$. Then*

$$\int_{\{S^* > 0\}} f d\mu \geq 0.$$

Proof. Set $S_n^* = \max_{0 \leq m \leq n} S_m$ and $A_n = \{S_n^* > 0\}$. Then, for $m = 1, \dots, n$,

$$S_m = f + S_{m-1} \circ \theta \leq f + S_n^* \circ \theta.$$

On A_n , we have $S_n^* = \max_{1 \leq m \leq n} S_m$, so

$$S_n^* \leq f + S_n^* \circ \theta.$$

On A_n^c , we have

$$S_n^* = 0 \leq S_n^* \circ \theta.$$

So, integrating and adding, we obtain

$$\int_E S_n^* d\mu \leq \int_{A_n} f d\mu + \int_E S_n^* \circ \theta d\mu.$$

But S_n^* is integrable, so

$$\int_E S_n^* \circ \theta d\mu = \int_E S_n^* d\mu < \infty$$

which forces

$$\int_{A_n} f d\mu \geq 0.$$

As $n \rightarrow \infty$, $A_n \uparrow \{S^* > 0\}$ so, by dominated convergence, with dominating function $|f|$,

$$\int_{\{S^* > 0\}} f d\mu = \lim_{n \rightarrow \infty} \int_{A_n} f d\mu \geq 0.$$

□

Theorem 9.3.2 (Birkhoff's almost everywhere ergodic theorem). *Assume that (E, \mathcal{E}, μ) is σ -finite and that f is an integrable function on E . Then there exists an invariant function \bar{f} , with $\mu(|\bar{f}|) \leq \mu(|f|)$, such that $S_n(f)/n \rightarrow \bar{f}$ a.e. as $n \rightarrow \infty$.*

Proof. The functions $\liminf_n (S_n/n)$ and $\limsup_n (S_n/n)$ are invariant. Therefore, for $a < b$, so is the following set

$$D = D(a, b) = \{\liminf_n (S_n/n) < a < b < \limsup_n (S_n/n)\}.$$

We shall show that $\mu(D) = 0$. First, by invariance, we can restrict everything to D and thereby reduce to the case $D = E$. Note that either $b > 0$ or $a < 0$. We can interchange the two cases by replacing f by $-f$. Let us assume then that $b > 0$.

Let $B \in \mathcal{E}$ with $\mu(B) < \infty$, then $g = f - b1_B$ is integrable and, for each $x \in D$, for some n ,

$$S_n(g)(x) \geq S_n(f)(x) - nb > 0.$$

Hence $S^*(g) > 0$ everywhere and, by the maximal ergodic lemma,

$$0 \leq \int_D (f - b1_B) d\mu = \int_D f d\mu - b\mu(B).$$

Since μ is σ -finite, there is a sequence of sets $B_n \in \mathcal{E}$, with $\mu(B_n) < \infty$ for all n and $B_n \uparrow D$. Hence,

$$b\mu(D) = \lim_{n \rightarrow \infty} b\mu(B_n) \leq \int_D f d\mu.$$

In particular, we see that $\mu(D) < \infty$. A similar argument applied to $-f$ and $-a$, this time with $B = D$, shows that

$$(-a)\mu(D) \leq \int_D (-f) d\mu.$$

Hence

$$b\mu(D) \leq \int_D f d\mu \leq a\mu(D).$$

Since $a < b$ and the integral is finite, this forces $\mu(D) = 0$. Set

$$\Delta = \{\liminf_n (S_n/n) < \limsup_n (S_n/n)\}$$

then Δ is invariant. Also, $\Delta = \bigcup_{a,b \in \mathbb{Q}, a < b} D(a, b)$, so $\mu(\Delta) = 0$. On the complement of Δ , S_n/n converges in $[-\infty, \infty]$, so we can define an invariant function \bar{f} by

$$\bar{f} = \begin{cases} \lim_n (S_n/n) & \text{on } \Delta^c \\ 0 & \text{on } \Delta. \end{cases}$$

Finally, $\mu(|f \circ \theta^n|) = \mu(|f|)$, so $\mu(|S_n|) \leq n\mu(|f|)$ for all n . Hence, by Fatou's lemma,

$$\mu(|\bar{f}|) = \mu(\liminf_n |S_n/n|) \leq \liminf_n \mu(|S_n/n|) \leq \mu(|f|).$$

□

Theorem 9.3.3 (von Neumann's L^p ergodic theorem). *Assume that $\mu(E) < \infty$. Let $p \in [1, \infty)$. Then, for all $f \in L^p(\mu)$, $S_n(f)/n \rightarrow \bar{f}$ in L^p .*

Proof. We have

$$\|f \circ \theta^n\|_p = \left(\int_E |f|^p \circ \theta^n d\mu \right)^{1/p} = \|f\|_p.$$

So, by Minkowski's inequality,

$$\|S_n(f)/n\|_p \leq \|f\|_p.$$

Given $\varepsilon > 0$, choose $K < \infty$ so that $\|f - g\|_p < \varepsilon/3$, where $g = (-K) \vee f \wedge K$. By Birkhoff's theorem, $S_n(g)/n \rightarrow \bar{g}$ a.e.. We have $|S_n(g)/n| \leq K$ for all n so, by bounded convergence, there exists N such that, for $n \geq N$,

$$\|S_n(g)/n - \bar{g}\|_p < \varepsilon/3.$$

By Fatou's lemma,

$$\begin{aligned} \|\bar{f} - \bar{g}\|_p^p &= \int_E \liminf_n |S_n(f - g)/n|^p d\mu \\ &\leq \liminf_n \int_E |S_n(f - g)/n|^p d\mu \leq \|f - g\|_p^p. \end{aligned}$$

Hence, for $n \geq N$,

$$\begin{aligned} \|S_n(f)/n - \bar{f}\|_p &\leq \|S_n(f - g)/n\|_p + \|S_n(g)/n - \bar{g}\|_p + \|\bar{g} - \bar{f}\|_p \\ &< \varepsilon/3 + \varepsilon/3 + \varepsilon/3 = \varepsilon. \end{aligned}$$

□

10. SUMS OF INDEPENDENT RANDOM VARIABLES

10.1. Strong law of large numbers for finite fourth moment. The result we obtain in this section will be largely superseded in the next. We include it because its proof is much more elementary than that needed for the definitive version of the strong law which follows.

Theorem 10.1.1. *Let $(X_n : n \in \mathbb{N})$ be a sequence of independent random variables such that, for some constants $\mu \in \mathbb{R}$ and $M < \infty$,*

$$\mathbb{E}(X_n) = \mu, \quad \mathbb{E}(X_n^4) \leq M \quad \text{for all } n.$$

Set $S_n = X_1 + \cdots + X_n$. Then

$$S_n/n \rightarrow \mu \quad \text{a.s., as } n \rightarrow \infty.$$

Proof. Consider $Y_n = X_n - \mu$. Then $Y_n^4 \leq 2^4(X_n^4 + \mu^4)$, so

$$\mathbb{E}(Y_n^4) \leq 16(M + \mu^4)$$

and it suffices to show that $(Y_1 + \cdots + Y_n)/n \rightarrow 0$ a.s.. So we are reduced to the case where $\mu = 0$.

Note that X_n, X_n^2, X_n^3 are all integrable since X_n^4 is. Since $\mu = 0$, by independence,

$$\mathbb{E}(X_i X_j^3) = \mathbb{E}(X_i X_j X_k^2) = \mathbb{E}(X_i X_j X_k X_l) = 0$$

for distinct indices i, j, k, l . Hence

$$\mathbb{E}(S_n^4) = \mathbb{E} \left(\sum_{1 \leq i \leq n} X_i^4 + 6 \sum_{1 \leq i < j \leq n} X_i^2 X_j^2 \right).$$

Now for $i < j$, by independence and the Cauchy–Schwarz inequality

$$\mathbb{E}(X_i^2 X_j^2) = \mathbb{E}(X_i^2) \mathbb{E}(X_j^2) \leq \mathbb{E}(X_i^4)^{1/2} \mathbb{E}(X_j^4)^{1/2} \leq M.$$

So we get the bound

$$\mathbb{E}(S_n^4) \leq nM + 3n(n-1)M \leq 3n^2M.$$

Thus

$$\mathbb{E} \sum_n (S_n/n)^4 \leq 3M \sum_n 1/n^2 < \infty$$

which implies

$$\sum_n (S_n/n)^4 < \infty \quad \text{a.s.}$$

and hence $S_n/n \rightarrow 0$ a.s.. □

10.2. Strong law of large numbers.

Theorem 10.2.1. *Let m be a probability measure on \mathbb{R} , with*

$$\int_{\mathbb{R}} |x| m(dx) < \infty, \quad \int_{\mathbb{R}} x m(dx) = \nu.$$

Let (E, \mathcal{E}, μ) be the canonical model for a sequence of independent random variables with law m . Then

$$\mu(\{x : (x_1 + \cdots + x_n)/n \rightarrow \nu \text{ as } n \rightarrow \infty\}) = 1.$$

Proof. By Theorem 9.2.1, the shift map θ on E is measure-preserving and ergodic. The coordinate function $f = X_1$ is integrable and $S_n(f) = f + f \circ \theta + \cdots + f \circ \theta^{n-1} = X_1 + \cdots + X_n$. So $(X_1 + \cdots + X_n)/n \rightarrow \bar{f}$ a.e. and in L^1 , for some invariant function \bar{f} , by Birkhoff's theorem. Since θ is ergodic, $\bar{f} = c$ a.e., for some constant c and then $c = \mu(\bar{f}) = \lim_n \mu(S_n/n) = \nu$. \square

Theorem 10.2.2 (Strong law of large numbers). *Let $(Y_n : n \in \mathbb{N})$ be a sequence of independent, identically distributed, integrable random variables with mean ν . Set $S_n = Y_1 + \cdots + Y_n$. Then*

$$S_n/n \rightarrow \nu \quad \text{a.s., as } n \rightarrow \infty.$$

Proof. In the notation of Theorem 10.2.1, take m to be the law of the random variables Y_n . Then $\mu = \mathbb{P} \circ Y^{-1}$, where $Y : \Omega \rightarrow E$ is given by $Y(\omega) = (Y_n(\omega) : n \in \mathbb{N})$. Hence

$$\mathbb{P}(S_n/n \rightarrow \nu \text{ as } n \rightarrow \infty) = \mu(\{x : (x_1 + \cdots + x_n)/n \rightarrow \nu \text{ as } n \rightarrow \infty\}) = 1.$$

\square

10.3. Central limit theorem.

Theorem 10.3.1 (Central limit theorem). *Let $(X_n : n \in \mathbb{N})$ be a sequence of independent, identically distributed, random variables with mean 0 and variance 1. Set $S_n = X_1 + \cdots + X_n$. Then, for all $a < b$, as $n \rightarrow \infty$,*

$$\mathbb{P}\left(\frac{S_n}{\sqrt{n}} \in [a, b]\right) \rightarrow \int_a^b \frac{1}{\sqrt{2\pi}} e^{-y^2/2} dy.$$

Proof. Set $\phi(u) = \mathbb{E}(e^{iuX_1})$. Since $\mathbb{E}(X_1^2) < \infty$, we can differentiate $\mathbb{E}(e^{iuX_1})$ twice under the expectation, to show that

$$\phi(0) = 1, \quad \phi'(0) = 0, \quad \phi''(0) = -1.$$

Hence, by Taylor's theorem, as $u \rightarrow 0$,

$$\phi(u) = 1 - u^2/2 + o(u^2).$$

So, for the characteristic function ϕ_n of S_n/\sqrt{n} ,

$$\phi_n(u) = \mathbb{E}(e^{iu(X_1 + \cdots + X_n)/\sqrt{n}}) = \{\mathbb{E}(e^{i(u/\sqrt{n})X_1})\}^n = (1 - u^2/2n + o(u^2/n))^n.$$

The complex logarithm satisfies, as $z \rightarrow 0$,

$$\log(1 + z) = z + o(|z|)$$

so, for each $u \in \mathbb{R}$, as $n \rightarrow \infty$,

$$\log \phi_n(u) = n \log(1 - u^2/2n + o(u^2/n)) = -u^2/2 + o(1).$$

Hence $\phi_n(u) \rightarrow e^{-u^2/2}$ for all u . But $e^{-u^2/2}$ is the characteristic function of the $N(0, 1)$ distribution, so $S_n/\sqrt{n} \rightarrow N(0, 1)$ in distribution by Theorem 7.5.1, as required. \square

Here is an alternative argument, which does not rely on Lévy's continuity theorem. Take a random variable $Y \sim N(0, 1)$, independent of the sequence $(X_n : n \in \mathbb{N})$. Fix $a < b$ and $\delta > 0$ and consider the function f which interpolates linearly the points $(-\infty, 0), (a - \delta, 0), (a, 1), (b, 1), (b + \delta, 0), (\infty, 0)$. Note that $|f(x + y) - f(x)| \leq |y|/\delta$ for all x, y . So, given $\varepsilon > 0$, for $t = (\pi/2)(\varepsilon\delta/3)^2$ and any random variable Z ,

$$|\mathbb{E}(f(Z + \sqrt{t}Y)) - \mathbb{E}(f(Z))| \leq \mathbb{E}(\sqrt{t}|Y|)/\delta = \varepsilon/3.$$

Recall from the proof of the Fourier inversion formula that

$$\mathbb{E}\left(f\left(\frac{S_n}{\sqrt{n}} + \sqrt{t}Y\right)\right) = \int_{\mathbb{R}} \left(\frac{1}{2\pi} \int_{\mathbb{R}} \phi_n(u) e^{-u^2 t/2} e^{-iuy} du\right) f(y) dy.$$

Consider a second sequence of independent random variables $(\bar{X}_n : n \in \mathbb{N})$, also independent of Y , and with $\bar{X}_n \sim N(0, 1)$ for all n . Note that $\bar{S}_n/\sqrt{n} \sim N(0, 1)$ for all n . So

$$\mathbb{E}\left(f\left(\frac{\bar{S}_n}{\sqrt{n}} + \sqrt{t}Y\right)\right) = \int_{\mathbb{R}} \left(\frac{1}{2\pi} \int_{\mathbb{R}} e^{-u^2/2} e^{-u^2 t/2} e^{-iuy} du\right) f(y) dy.$$

Now $e^{-u^2/2}f(y) \in L^1(du \otimes dy)$ and ϕ_n is bounded, with $\phi_n(u) \rightarrow e^{-u^2/2}$ for all u as $n \rightarrow \infty$, so, by dominated convergence, for n sufficiently large,

$$\left| \mathbb{E} \left(f \left(\frac{S_n}{\sqrt{n}} + \sqrt{t}Y \right) \right) - \mathbb{E} \left(f \left(\frac{\bar{S}_n}{\sqrt{n}} + \sqrt{t}Y \right) \right) \right| \leq \varepsilon/3.$$

Hence, by taking $Z = S_n/\sqrt{n}$ and then $Z = \bar{S}_n/\sqrt{n}$, we obtain

$$|\mathbb{E}(f(S_n/\sqrt{n})) - \mathbb{E}(f(\bar{S}_n/\sqrt{n}))| \leq \varepsilon.$$

But $\bar{S}_n/\sqrt{n} \sim Y$ for all n and $\varepsilon > 0$ is arbitrary, so we have shown that

$$\mathbb{E}(f(S_n/\sqrt{n})) \rightarrow \mathbb{E}(f(Y)) \quad \text{as } n \rightarrow \infty.$$

The same argument applies to the function g , defined like f , but with a, b replaced by $a + \delta, b - \delta$ respectively. Now $g \leq 1_{[a,b]} \leq f$, so

$$\mathbb{E} \left(g \left(\frac{S_n}{\sqrt{n}} \right) \right) \leq \mathbb{P} \left(\frac{S_n}{\sqrt{n}} \in [a, b] \right) \leq \mathbb{E} \left(f \left(\frac{S_n}{\sqrt{n}} \right) \right).$$

On the other hand, as $\delta \downarrow 0$,

$$\mathbb{E}(g(Y)) \uparrow \int_a^b \frac{1}{\sqrt{2\pi}} e^{-y^2/2} dy, \quad \mathbb{E}(f(Y)) \downarrow \int_a^b \frac{1}{\sqrt{2\pi}} e^{-y^2/2} dy$$

so we must have, as $n \rightarrow \infty$,

$$\mathbb{P} \left(\frac{S_n}{\sqrt{n}} \in [a, b] \right) \rightarrow \int_a^b \frac{1}{\sqrt{2\pi}} e^{-y^2/2} dy.$$

11. EXERCISES

Students should attempt Exercises 1.1–2.7 for their first supervision, then 3.1–3.13, 4.1–7.6 and 8.1–10.2 for later supervisions.

1.1 Show that a π -system which is also a d -system is a σ -algebra.

1.2 Show that the following sets of subsets of \mathbb{R} all generate the same σ -algebra:

(a) $\{(a, b) : a < b\}$, (b) $\{(a, b] : a < b\}$, (c) $\{(-\infty, b] : b \in \mathbb{R}\}$.

1.3 Show that a countably additive set function on a ring is both increasing and countably subadditive.

1.4 Let μ be a finite-valued additive set function on a ring \mathcal{A} . Show that μ is countably additive if and only if

$$A_n \supseteq A_{n+1} \in \mathcal{A}, n \in \mathbb{N}, \quad \bigcap_n A_n = \emptyset \quad \Rightarrow \quad \mu(A_n) \rightarrow 0.$$

1.5 Let (E, \mathcal{E}, μ) be a measure space. Show that, for any sequence of sets $(A_n : n \in \mathbb{N})$ in \mathcal{E} ,

$$\mu(\liminf A_n) \leq \liminf \mu(A_n).$$

Show that, if μ is finite, then also

$$\mu(\limsup A_n) \geq \limsup \mu(A_n)$$

and give an example to show this inequality may fail if μ is not finite.

1.6 Let $(\Omega, \mathcal{F}, \mathbb{P})$ be a probability space and $A_n, n \in \mathbb{N}$, a sequence of events. Show that $A_n, n \in \mathbb{N}$, are independent if and only if the σ -algebras they generate

$$\sigma(A_n) = \{\emptyset, A_n, A_n^c, \Omega\}$$

are independent.

1.7 Show that, for every Borel set $B \subseteq \mathbb{R}$ of finite Lebesgue measure and every $\varepsilon > 0$, there exists a finite union of disjoint intervals $A = (a_1, b_1] \cup \cdots \cup (a_n, b_n]$ such that the Lebesgue measure of $A \triangle B (= (A^c \cap B) \cup (A \cap B^c))$ is less than ε .

1.8 Let (E, \mathcal{E}, μ) be a measure space. Call a subset $N \subseteq E$ *null* if

$$N \subseteq B \text{ for some } B \in \mathcal{E} \text{ with } \mu(B) = 0.$$

Prove that the set of subsets

$$\mathcal{E}^\mu = \{A \cup N : A \in \mathcal{E}, N \text{ null}\}$$

is a σ -algebra and show that μ has a well-defined and countably additive extension to \mathcal{E}^μ given by

$$\mu(A \cup N) = \mu(A).$$

We call \mathcal{E}^μ the *completion of \mathcal{E} with respect to μ* .

2.1 Prove Proposition 2.1.1 and deduce that, for any sequence $(f_n : n \in \mathbb{N})$ of measurable functions on (E, \mathcal{E}) ,

$$\{x \in E : f_n(x) \text{ converges as } n \rightarrow \infty\} \in \mathcal{E}.$$

2.2 Let X and Y be two random variables on $(\Omega, \mathcal{F}, \mathbb{P})$ and suppose that for all $x, y \in \mathbb{R}$

$$\mathbb{P}(X \leq x, Y \leq y) = \mathbb{P}(X \leq x)\mathbb{P}(Y \leq y).$$

Show that X and Y are independent.

2.3 Let X_1, X_2, \dots be random variables with

$$X_n = \begin{cases} n^2 - 1 & \text{with probability } 1/n^2 \\ -1 & \text{with probability } 1 - 1/n^2. \end{cases}$$

Show that

$$\mathbb{E} \left(\frac{X_1 + \dots + X_n}{n} \right) = 0$$

but with probability one, as $n \rightarrow \infty$,

$$\frac{X_1 + \dots + X_n}{n} \longrightarrow -1.$$

2.4 For $s > 1$ define the zeta function by

$$\zeta(s) = \sum_{n=1}^{\infty} n^{-s}.$$

Let X and Y be independent random variables with

$$\mathbb{P}(X = n) = \mathbb{P}(Y = n) = n^{-s}/\zeta(s).$$

Show that the events

$$\{p \text{ divides } X\}, \quad p \text{ prime}$$

are independent and deduce Euler's formula

$$\frac{1}{\zeta(s)} = \prod_p \left(1 - \frac{1}{p^s} \right).$$

Prove also that

$$\mathbb{P}(X \text{ is square-free}) = 1/\zeta(2s)$$

and

$$\mathbb{P}(\text{h.c.f.}(X, Y) = n) = n^{-2s}/\zeta(2s).$$

2.5 Let X_1, X_2, \dots be independent random variables with distribution uniform on $[0, 1]$. Let A_n be the event that a record occurs at time n , that is,

$$X_n > X_m \quad \text{for all } m < n.$$

Find the probability of A_n and show that A_1, A_2, \dots are independent. Deduce that, with probability one, infinitely many records occur.

2.6 Let X_1, X_2, \dots be independent $N(0, 1)$ random variables. Prove that

$$\limsup_n (X_n / \sqrt{2 \log n}) = 1 \quad \text{a.s.}$$

2.7 Let C_n denote the n th approximation to the Cantor set C : thus $C_0 = [0, 1]$, $C_1 = [0, \frac{1}{3}] \cup [\frac{2}{3}, 1]$, $C_2 = [0, \frac{1}{9}] \cup [\frac{2}{9}, \frac{1}{3}] \cup [\frac{2}{3}, \frac{7}{9}] \cup [\frac{8}{9}, 1]$, etc. and $C_n \downarrow C$ as $n \rightarrow \infty$. Denote by F_n the distribution function of a random variable uniformly distributed on C_n . Show

- (i) $F(x) = \lim_{n \rightarrow \infty} F_n(x)$ exists for all $x \in [0, 1]$,
- (ii) F is continuous, $F(0) = 0$ and $F(1) = 1$,
- (iii) F is differentiable a.e. with $F' = 0$.

3.1 A simple function f has two representations:

$$f = \sum_{k=1}^m a_k 1_{A_k} = \sum_{j=1}^n b_j 1_{B_j}.$$

For $\varepsilon \in \{0, 1\}^m$ define $A_\varepsilon = A_1^{\varepsilon_1} \cap \cdots \cap A_m^{\varepsilon_m}$ where $A_k^0 = A_k^c$, $A_k^1 = A_k$. For $\delta \in \{0, 1\}^n$ define B_δ similarly. Then set

$$f_{\varepsilon, \delta} = \begin{cases} \sum_{k=1}^m \varepsilon_k a_k & \text{if } A_\varepsilon \cap B_\delta \neq \emptyset \\ \infty & \text{otherwise.} \end{cases}$$

Show that for any measure μ

$$\sum_{k=1}^m a_k \mu(A_k) = \sum_{\varepsilon, \delta} f_{\varepsilon, \delta} \mu(A_\varepsilon \cap B_\delta)$$

and deduce that

$$\sum_{k=1}^m a_k \mu(A_k) = \sum_{j=1}^n b_j \mu(B_j).$$

3.2 Show that any continuous function $f : \mathbb{R} \rightarrow \mathbb{R}$ is Lebesgue integrable over any finite interval.

3.3 Prove Propositions 3.3.1, 3.3.2, 3.3.3, 3.5.2.

3.4 Let X be a non-negative integer-valued random variable. Show that

$$\mathbb{E}(X) = \sum_{n=1}^{\infty} \mathbb{P}(X \geq n).$$

Deduce that, if $\mathbb{E}(X) = \infty$ and X_1, X_2, \dots is a sequence of independent random variables with the same distribution as X , then

$$\limsup(X_n/n) \geq 1 \quad \text{a.s.}$$

and indeed

$$\limsup(X_n/n) = \infty \quad \text{a.s.}$$

Now suppose that Y_1, Y_2, \dots is any sequence of independent identically distributed random variables with $\mathbb{E}|Y_1| = \infty$. Show that

$$\limsup(|Y_n|/n) = \infty \quad \text{a.s.}$$

and indeed

$$\limsup(|Y_1 + \cdots + Y_n|/n) = \infty \quad \text{a.s.}$$

3.5 For $\alpha \in (0, \infty)$ and $p \in [1, \infty)$ and for

$$f_\alpha(x) = 1/x^\alpha, \quad x > 0,$$

show carefully that

$$\begin{aligned} f_\alpha \in L^p((0, 1], dx) &\Leftrightarrow \alpha p < 1, \\ f_\alpha \in L^p([1, \infty), dx) &\Leftrightarrow \alpha p > 1. \end{aligned}$$

3.6 Show that the function

$$f(x) = \frac{\sin x}{x}$$

is not Lebesgue integrable over $[1, \infty)$ but that the following limit does exist:

$$\lim_{N \rightarrow \infty} \int_1^N \frac{\sin x}{x} dx.$$

3.7 Show

$$\begin{aligned} \text{(i): } \int_0^\infty \sin(e^x)/(1 + nx^2) dx &\rightarrow 0 \quad \text{as } n \rightarrow \infty, \\ \text{(ii): } \int_0^1 (n \cos x)/(1 + n^2 x^{\frac{3}{2}}) dx &\rightarrow 0 \quad \text{as } n \rightarrow \infty. \end{aligned}$$

3.8 Let u and v be differentiable functions on $[a, b]$ with continuous derivatives u' and v' . Show that for $a < b$

$$\int_a^b u(x)v'(x)dx = \{u(b)v(b) - u(a)v(a)\} - \int_a^b u'(x)v(x)dx.$$

3.9 Prove Propositions 3.4.4, 3.4.6, 3.5.1, 3.5.2 and 3.5.3.

3.10 The moment generating function ϕ of a real-valued random variable X is defined by

$$\phi(\tau) = \mathbb{E}(e^{\tau X}), \quad \tau \in \mathbb{R}.$$

Show that the set $I = \{\tau : \phi(\tau) < \infty\}$ is an interval and find examples where I is \mathbb{R} , $\{0\}$ and $(-\infty, 1)$. Assume for simplicity that $X \geq 0$. Show that if I contains a neighbourhood of 0 then X has finite moments of all orders given by

$$\mathbb{E}(X^n) = \left(\frac{d}{d\tau} \right)^n \bigg|_{\tau=0} \phi(\tau).$$

Find a necessary and sufficient condition on the sequence of moments $m_n = \mathbb{E}(X^n)$ for I to contain a neighbourhood of 0.

3.11 Let X_1, \dots, X_n be random variables with density functions f_1, \dots, f_n respectively. Suppose that the \mathbb{R}^n -valued random variable $X = (X_1, \dots, X_n)$ also has a density function f . Show that X_1, \dots, X_n are independent if and only if

$$f(x_1, \dots, x_n) = f_1(x_1) \dots f_n(x_n) \quad \text{a.e.}$$

3.12 Let $(f_n : n \in \mathbb{N})$ be a sequence of integrable functions and suppose that $f_n \rightarrow f$ a.e. for some integrable function f . Show that, if $\|f_n\|_1 \rightarrow \|f\|_1$, then $\|f_n - f\|_1 \rightarrow 0$.

3.13 Let μ and ν be probability measures on (E, \mathcal{E}) and suppose that, for some measurable function $f : E \rightarrow [0, R]$,

$$\nu(A) = \int_A f d\mu, \quad A \in \mathcal{E}.$$

Let $(X_n : n \in \mathbb{N})$ be a sequence of independent random variables in E with law μ and let $(U_n : n \in \mathbb{N})$ be a sequence of independent $U[0, 1]$ random variables. Set

$$T = \min\{n \in \mathbb{N} : RU_n \leq f(X_n)\}, \quad Y = X_T.$$

Show that Y has law ν .

4.1 Let X be a random variable and let $1 \leq p < q < \infty$. Show that

$$\mathbb{E}(|X|^p) = \int_0^\infty p\lambda^{p-1}\mathbb{P}(|X| \geq \lambda)d\lambda$$

and deduce

$$X \in L^q(\mathbb{P}) \Rightarrow \mathbb{P}(|X| \geq \lambda) = O(\lambda^{-q}) \Rightarrow X \in L^p(\mathbb{P}).$$

4.2 Give a simple proof of Schwarz' inequality for measurable functions f and g :

$$\|fg\|_1 \leq \|f\|_2\|g\|_2.$$

4.3 Show that for independent random variables X and Y

$$\|XY\|_1 = \|X\|_1\|Y\|_1$$

and that if both X and Y are integrable then

$$\mathbb{E}(XY) = \mathbb{E}(X)\mathbb{E}(Y).$$

4.4 A *stepfunction* $f : \mathbb{R} \rightarrow \mathbb{R}$ is any finite linear combination of indicator functions of finite intervals. Show that the set of stepfunctions \mathcal{J} is dense in $L^p(\mathbb{R})$ for all $p \in [1, \infty)$: that is, for all $f \in L^p(\mathbb{R})$ and all $\varepsilon > 0$ there exists $g \in \mathcal{J}$ such that $\|f - g\|_p < \varepsilon$.

4.5 Let $(X_n : n \in \mathbb{N})$ be an identically distributed sequence in $L^2(\mathbb{P})$. Show that, for $\varepsilon > 0$,

- (i): $n\mathbb{P}(|X_1| > \varepsilon\sqrt{n}) \rightarrow 0$ as $n \rightarrow \infty$,
- (ii): $n^{-\frac{1}{2}} \max_{k \leq n} |X_k| \rightarrow 0$ in probability.

5.1 Let (E, \mathcal{E}, μ) be a measure space and let $V_1 \leq V_2 \leq \dots$ be an increasing sequence of closed subspaces of $L^2 = L^2(E, \mathcal{E}, \mu)$ for $f \in L^2$, denote by f_n the orthogonal projection of f on V_n . Show that f_n converges in L^2 .

5.2 Prove Propositions 5.3.1 and 5.3.2.

6.1 Prove Proposition 6.2.2.

6.2 Find a uniformly integrable sequence of random variables $(X_n : n \in \mathbb{N})$ such that

$$X_n \rightarrow 0 \text{ a.s.} \quad \text{and} \quad \mathbb{E}(\sup_n |X_n|) = \infty.$$

6.3 Let $(X_n : n \in \mathbb{N})$ be an identically distributed sequence in $L^2(\mathbb{P})$. Show that

$$\mathbb{E}(\max_{k \leq n} |X_k|)/\sqrt{n} \rightarrow 0 \quad \text{as} \quad n \rightarrow \infty.$$

7.1 Show that the Fourier transform of a finite Borel measure is a bounded continuous function.

7.2 Let μ be a Borel measure on \mathbb{R} of finite total mass. Suppose the Fourier transform $\hat{\mu}$ is Lebesgue integrable. Show that μ has a continuous density function f with respect to Lebesgue measure:

$$\mu(A) = \int_A f(x) dx.$$

7.3 Show that there do not exist independent identically distributed random variables X, Y such that

$$X - Y \sim U[-1, 1].$$

7.4 The Cauchy distribution has density function

$$f(x) = \frac{1}{\pi(1+x^2)}, \quad x \in \mathbb{R}.$$

Show that the corresponding characteristic function is given by

$$\varphi(u) = e^{-|u|}.$$

Show also that, if X_1, \dots, X_n are independent Cauchy random variables, then $(X_1 + \dots + X_n)/n$ is also Cauchy. Comment on this in the light of the strong law of large numbers and central limit theorem.

7.5 For a finite Borel measure μ on the line show that, if $\int |x|^k d\mu(x) < \infty$, then the Fourier transform $\hat{\mu}$ of μ has a k th continuous derivative, which at 0 is given by

$$\hat{\mu}^{(k)}(0) = i^k \int x^k d\mu(x).$$

7.6 (i) Show that for any real numbers a, b one has $\int_a^b e^{itx} dx \rightarrow 0$ as $|t| \rightarrow \infty$.

(ii) Suppose that μ is a finite Borel measure on \mathbb{R} which has a density f with respect to Lebesgue measure. Show that its Fourier transform

$$\hat{\mu}(t) = \int_{-\infty}^{\infty} e^{itx} f(x) dx$$

tends to 0 as $|t| \rightarrow \infty$. This is the *Riemann–Lebesgue Lemma*.

(iii) Suppose that the density f of μ has an integrable and continuous derivative f' . Show that

$$\hat{\mu}(t) = o(t^{-1}), \quad \text{i.e., } t\hat{\mu}(t) \rightarrow 0 \text{ as } |t| \rightarrow \infty.$$

Extend to higher derivatives.

8.1 Prove Proposition 8.1.1.

8.2 Suppose that X_1, \dots, X_n are jointly Gaussian random variables with

$$\mathbb{E}(X_i) = \mu_i, \quad \text{cov}(X_i, X_j) = \Sigma_{ij}$$

and that the matrix $\Sigma = (\Sigma_{ij})$ is invertible. Set $Y = \Sigma^{-\frac{1}{2}}(X - \mu)$. Show that Y_1, \dots, Y_n are independent $N(0, 1)$ random variables.

Show that we can write X_2 in the form $X_2 = aX_1 + Z$ where Z is independent of X_1 and determine the distribution of Z .

8.3 Let X_1, \dots, X_n be independent $N(0, 1)$ random variables. Show that

$$\left(\bar{X}, \sum_{m=1}^n (X_m - \bar{X})^2 \right) \quad \text{and} \quad \left(X_n/\sqrt{n}, \sum_{m=1}^{n-1} X_m^2 \right)$$

have the same distribution, where $\bar{X} = (X_1 + \dots + X_n)/n$.

9.1 Let (E, \mathcal{E}, μ) be a measure space and $\tau : E \rightarrow E$ a measure-preserving transformation. Show that

$$\mathcal{E}_\tau := \{A \in \mathcal{E} : \tau^{-1}(A) = A\}$$

is a σ -algebra, and that a measurable function f is \mathcal{E}_τ -measurable if and only if it is *invariant*, that is $f \circ \tau = f$.

9.2 Prove Propositions 9.1.1 and 9.1.2.

9.3 For $E = [0, 1)$, $a \in E$ and $\mu(dx) = dx$, show that

$$\tau(x) = x + a \pmod{1}$$

is measure-preserving. Determine for which values of a the transformation τ is ergodic.

Let f be an integrable function on $[0, 1)$. Determine for each value of a the limit

$$\bar{f} = \lim_{n \rightarrow \infty} \frac{1}{n} (f + f \circ \tau + \dots + f \circ \tau^{n-1}).$$

9.4 Show that

$$\tau(x) = 2x \pmod{1}$$

is another measure-preserving transformation of Lebesgue measure on $[0, 1)$, and that τ is ergodic. Find \bar{f} for each integrable function f .

9.5 Call a sequence of random variables $(X_n : n \in \mathbb{N})$ on a probability space $(\Omega, \mathcal{F}, \mathbb{P})$ *stationary* if for each $n, k \in \mathbb{N}$ the random vectors (X_1, \dots, X_n) and $(X_{k+1}, \dots, X_{k+n})$ have the same distribution: for $A_1, \dots, A_n \in \mathcal{B}$,

$$\mathbb{P}(X_1 \in A_1, \dots, X_n \in A_n) = \mathbb{P}(X_{k+1} \in A_1, \dots, X_{k+n} \in A_n).$$

Show that, if $(X_n : n \in \mathbb{N})$ is a stationary sequence and $X_1 \in L^p$, for some $p \in [1, \infty)$, then

$$\frac{1}{n} \sum_{i=1}^n X_i \rightarrow X \quad \text{a.s. and in } L^p,$$

for some random variable $X \in L^p$ and find $\mathbb{E}[X]$.

10.1 Let f be a bounded continuous function on $(0, \infty)$, having Laplace transform

$$\hat{f}(\lambda) = \int_0^\infty e^{-\lambda x} f(x) dx, \quad \lambda \in (0, \infty).$$

Let $(X_n : n \in \mathbb{N})$ be a sequence of independent exponential random variables, of parameter λ . Show that \hat{f} has derivatives of all orders on $(0, \infty)$ and that, for all $n \in \mathbb{N}$, for some $C(\lambda, n) \neq 0$ independent of f , we have

$$(d/d\lambda)^{n-1} \hat{f}(\lambda) = C(\lambda, n) \mathbb{E}(f(S_n))$$

where $S_n = X_1 + \cdots + X_n$. Deduce that if $\hat{f} \equiv 0$ then also $f \equiv 0$.

10.2 For each $n \in \mathbb{N}$, there is a unique probability measure μ_n on the unit sphere $S^{n-1} = \{x \in \mathbb{R}^n : |x| = 1\}$ such that $\mu_n(A) = \mu_n(UA)$ for all Borel sets A and all orthogonal $n \times n$ matrices U . Fix $k \in \mathbb{N}$ and, for $n \geq k$, let γ_n denote the probability measure on \mathbb{R}^k which is the law of $\sqrt{n}(x^1, \dots, x^k)$ under μ_n . Show

- (i) if $X \sim N(0, I_n)$ then $X/|X| \sim \mu_n$,
- (ii) if $(X_n : n \in \mathbb{N})$ is a sequence of independent $N(0, 1)$ random variables and if $R_n = (X_1^2 + \cdots + X_n^2)^{\frac{1}{2}}$ then $R_n/\sqrt{n} \rightarrow 1$ a.s.,
- (iii) for all bounded continuous functions f on \mathbb{R}^k , $\gamma_n(f) \rightarrow \gamma(f)$, where γ is the standard Gaussian distribution on \mathbb{R}^k .