Probability

Note: Probability Theory

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	References: MAT3280, STAT5005 and Probability: Theory and Examples, 4th edition, by Richard Durret	

1 Preliminary

1.1 Riemann-Stieltjes integrals

lished by Cambridge University Press.

Goal: Consider a closed interval [a,b], a < b. Let $\alpha(x)$ be a non-decreasing function on [a,b]. $f \in \text{bdd}[a,b]$. Want to define RS integral $\int_a^b f(x) d\alpha(x)$.

Let $P: a = x_0 < x_1 < \cdots < x_n = b$ be a partition of [a, b]. Let $M_k(m_k) := \sup(\inf)\{f(x) : x_{k-1} \le x \le x_k\}$, t_k is chosen arbitrarily in $[x_{k-1}, x_k]$, $k = 1, 2, \ldots, n$, we define

- upper sum $U(P, f, \alpha) := \sum_{k=1}^{n} M_k \cdot (\alpha(x_k) \alpha(\alpha_{k-1})),$
- lower sum $L(P, f, \alpha) := \sum_{k=1}^{n} m_k \cdot (\alpha(x_k) \alpha(\alpha_{k-1})),$
- Riemann–Stieltjes sum $S(P, f, \alpha) := \sum_{k=1}^{n} f(t_k) \cdot (\alpha(x_k) \alpha(\alpha_{k-1}))$.

Sandwiching inequalities: $L(P, f, \alpha) \leq S(P, f, \alpha) \leq U(P, f, \alpha), \forall P$.

Definition 1.1 (Riemann–Stieltjes integrable). If $\lim_{\|P\|\to\infty} L(P,f,\alpha) = U(P,f,\alpha)$, then we say f is Riemann-Stieltjes (RS) integrable on [a, b] w.r.t. $\alpha(x), f \in \mathcal{R}(\alpha)$. We denote the common limit as

$$\int_a^b f d\alpha$$
 or $\int_a^b f(x) d(\alpha(x))$.

An improper RS integral is defined by the double limit (provided exists and finite, i.e., converges)

$$\int_{-\infty}^{\infty} f d\alpha := \lim_{(a,b) \to (-\infty,\infty)} \int_{a}^{b} f d\alpha.$$

Note that if $\beta(x) = \alpha(x) + C$, then $\int_a^b f d\alpha = \int_a^b f d\beta$. RS integral can be defined more generally when $\alpha(x)$ is BV. RS integrals have the following properties:

- 1. (linearity) $f, g \in \mathcal{R}(\alpha), \forall c, d \in \mathbb{R}, cf + dg \in \mathcal{R}(\alpha), \int_a^b (cf + dg) d\alpha = c \int_a^b f d\alpha + d \int_a^b g d\alpha;$
- 2. (monotonicity) $f, g \in \mathcal{R}(\alpha)$, if $f \leq g, \forall x \in [a, b]$, then $\int_a^b f d\alpha \leq \int_a^b g d\alpha$; 3. (additivity of integration) Let a < c < b, if $f \in \mathcal{R}(\alpha)$ on [a, c] and on [c, b], then $f \in \mathcal{R}(\alpha)$ on [a, b], and $\int_{a}^{b} f d\alpha = \int_{a}^{c} f d\alpha + \int_{c}^{b} f d\alpha.$
- 4. (linearity of α) Suppose $\alpha_1(x)$ and $\alpha_2(x)$ are two non-decreasing functions, $f \in \mathcal{R}(\alpha_1), f \in \mathcal{R}(\alpha_2)$, then $f \in \mathcal{R}(\alpha_1 + \alpha_2)$, and $\int_a^b f d(\alpha_1 + \alpha_2) = \int_a^b f d\alpha_1 + \int_a^b f d\alpha_2$.

Theorem 1.2 (Finite discontinuities). $f \in \text{bdd}[a, b]$, and there are finitely many points $c_1, c_2, \ldots, c_n \in [a, b]$ such that f is discontinuous at each c_i . If α is continuous on $c_i, \forall i$, then $f \in \mathcal{R}(\alpha)$.

Theorem 1.3 (Continuous functions). Suppose $f \in C[a,b]$ and let $\alpha \in C^1[a,b]$ (or $\alpha'(x)$ Riemann integrable on [a,b] be nondecreasing. If $f \in \mathcal{R}(\alpha)$, then the function $f\alpha'$ is Riemann integrable on [a,b] and we have

$$\int_{a}^{b} f d\alpha = \int_{a}^{b} f(\alpha)\alpha'(x)dx.$$

There are two techniques to compute RS integral:

- 1. (integration by parts) $\int_a^b f d\alpha = f\alpha \Big|_a^b \int_a^b \alpha df$;
- 2. (change of variables) if g is strictly increasing on [c,d], $f \in \mathscr{R}(\alpha)$ on [g(c),g(d)], let $h = f \circ g$, $\beta = \alpha \circ g$, then $h \in \mathscr{R}(\beta)$ on [c,d], with $\int_{g(c)}^{g(d)} f d\alpha = \int_{c}^{d} h d\beta$.

$\mathbf{2}$ Measure Theory

2.1Probability spaces

Definition 2.1 (Concepts).

- Sample space Ω : a set.
- Field/algebra \mathscr{F}' : collection of subsets of Ω satisfying: (i) $\Omega \in \mathscr{F}'$, (ii) $A \in \mathscr{F}$ implies $A^c \in \mathscr{F}'$, (iii) $A, B \in \mathscr{F}'$ implies $A \cup B \in \mathscr{F}'$.
- σ -field/algebra \mathscr{F} : collection of subsets of Ω satisfying: (i) $\Omega \in \mathscr{F}$, (ii) $A \in \mathscr{F}$ implies $A^c \in \mathscr{F}$, (iii) $A_i \in \mathscr{F}$ implies $\bigcup_{i=1}^{\infty} A_i \in \mathscr{F}$.
- (Ω, \mathcal{F}) is a measurable space.
 - A set in \mathscr{F} is called \mathscr{F} -measurable, or simply measurable.
 - If \mathscr{G} is a subcollection of \mathscr{F} , we say that \mathscr{G} is a sub- σ -field of \mathscr{F} if \mathscr{G} is also a σ -field.
- Measure $\mu: \mathscr{F} \to \mathbb{R}$, if (i) $\mu(A) \geq 0$, (ii) $\mu(\emptyset) = 0$, (iii) $A_i \in \mathscr{F}$ disjoint, then $\mu(\biguplus_{i=1}^n A_i) = \sum_{i=1}^n \mu(A_i)$.
 - A measure μ is said to be finite if $\mu(\Omega)$ is finite.
 - If $\mu(\Omega) = 1$, then μ is a probability measure.

Let μ be a measure on (Ω, \mathcal{F}) , it has the following properties:

- (a) (monotonicity) if $A \subset B$, then $\mu(A) \leq \mu(B)$;
- (b) (addition law) $\forall A.B \in \mathscr{F}, \ \mu(A \cup B) = \mu(A) + \mu(B) \mu(A \cap B);$
- (c) (sub-additivity) $\mu(\bigcup_{i=1}^{\infty} A_i) \leq \sum_{i=1}^{\infty} \mu(A_i)$;

- (d) (lower semi-continuity/continuity from below) if $A_n \uparrow A = \bigcup_{n=1}^{\infty} A_n$, then $\mu(A_n) \uparrow \mu(A)$;
- (e) (upper semi-continuity/continuity from above) if $A_n \downarrow A = \bigcap_{n=1}^{\infty} A_n$ and $\mu(A_1) < \infty$, then $\mu(A_n) \downarrow \mu(A)$.

Definition 2.2 (liminf and limsup). Let $E_1, E_2, ...$ be an arbitrary sequence of subsets in a set Ω . The limit inferior and limit superior of $(E_i)_{i\geq 1}$ are defined respectively as

$$\liminf_{i} E_i := \bigcup_{j=1}^{\infty} \bigcap_{k \geq j} E_k = \{E_i \ i.o.\}, \qquad \limsup_{i} E_i := \bigcap_{j=1}^{\infty} \bigcup_{k \geq j} E_k = \{E_i \ e.v.\}$$

We have $\liminf_i E_i \subseteq \limsup_i E_i$. If equality holds, we say the $\liminf_i G_i = \min_i E_i$ or $\lim_i E_i$ or $\lim_i E_i$.

Definition 2.3 (generated σ -field). Let \mathscr{A} be a collection of subsets of Ω , let $\sigma(\mathscr{A})$ be the σ -field generated by \mathscr{A} , defined as

$$\sigma(\mathscr{A}) = \bigcap_{\mathscr{A} \subseteq \mathscr{F}, \mathscr{F}\sigma - \text{field}} \mathscr{F}.$$

It is well-defined because if $\{\mathscr{F}_i: i\in I\}$ are all σ -fields, I can be uncountable, then $\cap_{i\in I}\mathscr{F}_i$ is also a σ -field. The $\sigma(\mathscr{A})$ is the smallest σ -field containing \mathscr{A} . If $\Omega=\mathbb{R}$, $\mathscr{A}=\{(a,b): -\infty < a\leq b<\infty\}$, then $\sigma(\mathscr{A})$ is called the Borel field/algebra, written as $\mathscr{B}(\mathbb{R})$. Likewise, define $\mathscr{B}(\mathbb{R}^d)$ the σ -algebra generated by the open balls in \mathbb{R}^d . A set in $\mathscr{B}(\mathbb{R})$ is called a Borel set.

- $\mathscr{B}(\mathbb{R}) = \sigma(\{(a,b) : a < b\})$, which can be $\{[a,b]\}, \{(a,b]\}, \dots$
- $\mathscr{B}(\mathbb{R}^d)$ can also be generated by open d-dimensional open boxes in the form $(a_1,b_1)\times(a_2,b_2)\times\cdots\times(a_d,b_d)$.

2.1.1 Measure extension theorem

Definition 2.4 (Concepts). A pre-measure $\mu_0: \mathscr{F}_0 \to [0, \infty]$ defined on a field \mathscr{F}_0 is a set function satisfying (i) $\mu_0(\emptyset) = 0$ and (ii) if $A_i \in \mathscr{F}_0$ mutually disjoint sets in Ω and $\biguplus_i A_i \in \mathscr{F}_0$, then $\mu_0(\biguplus_{i=1}^\infty A_i) = \sum_{i=1}^\infty \mu_0(A_i)$. Let \mathscr{F} be a σ-field containing \mathscr{F}_0 . We say that $\mu: \mathscr{F} \to [0, \infty]$ is an extension of μ_0 if μ is a measure and satisfies $\mu(E) = \mu_0(E), \, \forall E \in \mathscr{F}_0$. μ_0 is said to be finite if $\mu_0(\Omega) < \infty$. It is said to be σ-finite if \exists at most countably many sets $\Omega_i \in \mathscr{F}_0$ such that $\Omega = \bigcup_{i=1}^\infty \Omega_i$ and $\mu_0(\Omega_i) < \infty$ for all i.

We note that a probability measure is automatically σ -finite. WLOG, we may assume Ω_i 's form a partition of Ω . If Ω_i 's are not mutually disjoint, set $\Omega_i = \Omega_i \setminus (\bigcup_{i=1}^{i-1} \Omega_i)$.

Theorem 2.5 (Measure extension theorem/Hahn-Kolmogorov-Carathéodory). Let \mathscr{F}_0 be a field on Ω and μ_0 be a premeasure on \mathscr{F}_0 . There is a measure μ defined on $\sigma(\mathscr{F}_0)$ that extends to μ_0 . Moreover, if μ_0 is σ -finite, then the extension is unique.

- (Counting measure) $\Omega = \mathbb{N}^+$, \mathscr{F}_0 is the collection of all finite and co-finite sets. $\mu_0(E) = \infty$ if E is infinite and |E| if E is finite. $E = \{2, 4, 6, 8, \ldots\} \notin \mathscr{F}_0$, $\mu_0(E)$ is undefined. But it is σ -finite by taking $\Omega_i = \{i\}$. μ_0 can be extended to the power set $2^{\mathbb{N}}$, the counting measure on \mathbb{N}^+ .
- If $\Omega = \mathbb{R}$, $\mathscr{F}_0 = \{\text{finite union of } (a,b], a \in \mathbb{R} \cup \{-\infty\}, b \in \mathbb{R} \cup \{\infty\}\}, \text{ then } \mathscr{F}_0 \text{ is a field. Define } \mu_0 \text{ on } \mathscr{F}_0 \text{ by } \mu_0(\biguplus_{i=1}^n (a_i,b_i]) = \sum_{i=1}^n (b_i-a_i) \text{ the total length of the set, if all } a_i \text{ and } b_i < \infty.$ Otherwise, the length is ∞ . It is σ -finite by taking $\Omega_i = (-i,i]$. μ_0 can be extended to $\mathscr{B}(\mathbb{R})$, and the extension is unique, called the Borel measure on \mathbb{R} , denoted by λ .

2.1.2 Lebesgue–Stieltjes measures

Stieltjes measure function

Definition 2.6 (distribution function). Let $(\mathbb{R}, \mathcal{B}(\mathbb{R}), \mathbb{P})$ be a probability space. The distribution function induced by P is defined as

$$F(x) := \mathbb{P}((-\infty, x]), \quad x \in \mathbb{R}.$$

df vs cdf: cdf is for rv, while df is for a probability measure.

Theorem 2.7 (Properties of df). The df F(x) of a probability measure satisfies the following properties:

- 1. F(x) is non-decreasing;
- 2. F(x) is right-continuous;
- 3. $\lim_{x \to \infty} F(x) = 1$;
- 4. $\lim_{x \to -\infty} F(x) = 0$.

Definition 2.8 (Stieltjes measure function). A function $F: \mathbb{R} \to \mathbb{R}$ is called a Stieltjes measure function if

- 1. F is non-decreasing;
- 2. F is right-continuous.

Lebesgue-Stieltjes measures Given a Stieltjes measure function, we can apply the measure extension theorem 2.5 to construct a measure on $\mathcal{B}(\mathbb{R})$. Intuitively, the Lebesgue-Stieltjes measure μ assigns a length to each Borel set on the real line, where the length of (a, b] is given by F(b) - F(a).

Theorem 2.9. Let F be a Stieltjes measure function defined on \mathbb{R} . Then \exists a unique measure μ defined on $\mathscr{B}(\mathbb{R})$ such that $\mu((a,b]) = F(b) - F(a)$. The measure μ is called the Lebesgue-Stieltjes (LS) measure induced by F.

For example,

- (Lebesgue measure) take F(x) = x, then the LS measure on $\mathscr{B}(\mathbb{R})$ is Lebesgue measure, denoted by λ , $\lambda([a,b]) = b a$. It is translation-invariant;
- (Uniform distribution) $F(x) = x \mathbf{1}_{\{0 \le x \le 1\}} + \mathbf{1}_{\{1 < x\}}$, then the LS measure μ is the uniform distribution on [0, 1], $\mu((a, b|) = b a \text{ if } 0 < a < b < 1$;
- (Continuous distribution) Let $f(x) \ge 0$ be a pdf, i.e., Riemann-integrable on R, and $\int_{\mathbb{R}} f(x) dx = 1$. Then $F(x) := \int_{-\infty}^{x} f(t) dt$ is a Stieltjes measure function. Denote the LS measure by P, $P((a,b)) = \int_{a}^{b} f(t) dt$, a < b;
- (Dirac measure, point mass) $F(x) = \mathbf{1}_{\{x \ge x_0\}}$. $P(A) = \mathbf{1}_{\{x_0 \in A\}}$;
- (Discrete distribution) $i \in \mathbb{N}^+$, $p_i \geq 0$, $\sum_{i=1}^{\infty} p_i = 1$. Define $F(x) = \sum_{i=1}^{\infty} p_i \mathbf{1}_{\{x \geq x_i\}}$. Then LS measure $P(\{x\}) = \sum_{i=1}^{\infty} p_i \mathbf{1}_{\{x = x_i\}}$.

2.1.3 Null sets and complete measures

Consider a probability space $(\Omega, \mathcal{F}, \mu)$. If $\mu(E) = 0$, we hope $\forall E' \subseteq E, \mu(E') = 0$. But E' may not be measurable.

Definition 2.10. Given a measure space $(\Omega, \mathscr{F}, \mu)$, we define a set $N \subseteq \Omega$ as a null set if $\exists E \in \mathscr{F}$ s.t. $N \subseteq E$ and $\mu(E) = 0$. The measure space $(\Omega, \mathscr{F}, \mu)$ is said to be complete if all null sets are indeed \mathscr{F} -measurable.

Theorem 2.11 (Completion). Let $(\Omega, \mathcal{F}, \mu)$ be a measure space and \mathcal{N} be the set of all null sets. We can define a new collection of sets

$$\mathscr{F}' := \{ A \cup N : A \in \mathscr{F}, \ N \in \mathscr{N} \}.$$

We can extend the measure μ to a measure on \mathscr{F}' , denoted by μ' , by

$$\mu'(A \cup N) := \mu(A),$$

and the extended measure space $(\Omega, \mathcal{F}', \mu')$ is complete.

By enlarging the σ -field in this way, we may assume that the measure is complete without loss of generality.

2.1.4 π - λ theorem and uniqueness of measure extension

We will use the term set system to refer to a collection of subsets in Ω .

Definition 2.12. We define a π -system in Ω as a set system $\mathscr P$ that satisfies

$$A, B \in \mathscr{P} \Rightarrow A \cap B \in \mathscr{P}$$
.

A λ -system/Dynkin system in Ω is a set system \mathcal{L} that satisfies:

- (i) $\Omega \in \mathcal{L}$;
- (ii) $A \in \mathcal{L} \Rightarrow \Omega \setminus A \in \mathcal{L}$;
- (iii) If $A_i \in \mathcal{L}$, $A_i \uparrow A$, then $A \in \mathcal{L}$.
- (iii) is equivalent to: If $A_i \in \mathcal{L}$ mutually disjoint, then $\cup_i A_i \in \mathcal{L}$.
- If \mathscr{A} is both π and λ -system, then \mathscr{A} is a σ -field.

Theorem 2.13 (Dynkin's π - λ theorem). If \mathscr{P} is a π -system and \mathscr{L} is a λ -system, and $\mathscr{P} \subset \mathscr{L}$, then $\sigma(\mathscr{P}) \subset \mathscr{L}$.

We can use the π - λ theorem to prove the uniqueness of measure extension.

Theorem 2.14 (Uniqueness of measure extension). Suppose \mathscr{F}_0 is a field on a sample space Ω , and μ_1 and μ_2 are two measures on $\sigma(\mathscr{F}_0)$ such that $\mu_1(A) = \mu_2(A)$ for all $A \in \mathscr{F}_0$. Furthermore, suppose \exists a sequence of disjoint sets $\Omega_i \in \mathscr{F}_0$, for $i = 1, 2, \ldots$, such that $\biguplus_i \Omega_i = \Omega$ and $\mu_1(\Omega_i) = \mu_2(\Omega_i) < \infty$ for all i. Then $\mu_1(B) = \mu_2(B)$ for all $B \in \sigma(\mathscr{F}_0)$.

As an application of the uniqueness result, we prove that a probability measure on \mathbb{R} is uniquely determined by its Stieltjes measure function.

Theorem 2.15. Let P be a probability measure on $(\mathbb{R}, \mathscr{B}(\mathbb{R}))$, and let F(x) be its induced distribution function. If \exists another probability Q on $\mathscr{B}(\mathbb{R})$ such that

$$P((-\infty, x]) = Q((-\infty, x]), \quad \forall x \in \mathbb{R},$$

then $P(B) = Q(B), \forall B \in \mathcal{B}(\mathbb{R}).$

Multivariate version: On $(\mathbb{R}^d, \mathscr{B}(\mathbb{R}^d))$, P is uniquely determined by multivariate distribution function F(x).

2.2 Measurable functions

Given a function $f:\Omega\to\Omega'$, the inverse image of a set $A\subseteq\Omega'$ via the function f is defined as

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

- $f^{-1}(A^c) = (f^{-1}(A))^c$.
- If $(A_i)_{i\in I}$ is any collection of sets in Ω' , then $f^{-1}(\cap_{i\in I}A_i)=\bigcap_{i\in I}f^{-1}(A_i)$ and $f^{-1}(\cup_{i\in I}A_i)=\bigcup_{i\in I}f^{-1}(A_i)$.
- If $h = g \circ f$, then $h^{-1}(B) = f^{-1}(g^{-1}(B))$ for any subset B of the codomain of g.

Definition 2.16 (Measurable functions). Let (Ω, \mathscr{F}) and (Ω', \mathscr{G}) be two measurable spaces. A function $f: \Omega \to \Omega'$ is called $(\mathscr{F}, \mathscr{G})$ -measurable if $\forall B \in \mathscr{G}, f^{-1}(B) \in \mathscr{F}$. When \mathscr{G} is understood from the context, we say that f is \mathscr{F} -measurable. If both \mathscr{F} and \mathscr{G} are understood, we say that f is measurable.

When $(\Omega, \mathscr{F}) = (\mathbb{R}, \mathscr{B}(\mathbb{R}))$, we will refer to an \mathscr{F} -measurable function as a Borel measurable function.

When $(\Omega', \mathscr{G}) = (\mathbb{R}, \mathscr{B}(\mathbb{R}))$, a $(\mathscr{F}, \mathscr{B}(\mathbb{R}))$ -measurable function $X : \Omega \to \mathbb{R}$ is called a (real-valued) measurable function. If $(\Omega, \mathscr{F}, \mu)$ is a probability space, we refer to X as a random variable. Furthermore, if $(\Omega', \mathscr{G}) = (\mathbb{R}^d, \mathscr{B}(\mathbb{R}^d))$, then X is called a random vector.

Examples:

- If $\mathscr{F} = 2^{\Omega}$ or $\mathscr{G} = \{\emptyset, \Omega'\}$, then $\forall f : \Omega \to \Omega'$ is $(\mathscr{F}, \mathscr{G})$ -measurable.
- If $\mathscr{F} = \{\emptyset, \Omega\}$, and \mathscr{G} is the σ -field in which all singletons are \mathscr{G} -measurable, then an $(\mathscr{F}, \mathscr{G})$ -measurable function is a constant function.
- Let $A \subseteq \Omega$. The indicator function $\mathbf{1}_A(\omega)$ is $(\mathscr{F}, \mathscr{B}(\mathbb{R}))$ -measurable iff A is measurable.
 - Let $(\Omega, \mathscr{F}) = (\mathbb{R}, \mathscr{B}(\mathbb{R}))$, since \mathbb{Q} and Cantor set $C \in \mathscr{B}(\mathbb{R})$, $\mathbf{1}_{\mathbb{Q}}$ and $\mathbf{1}_{C}$ is Borel measurable. But Vitali set $V \notin \mathscr{B}(\mathbb{R})$, so $\mathbf{1}_{V}$ is not Borel measurable.

 σ -field generated by $f: \sigma(f) := \{f^{-1}(B) : B \in \mathcal{G}\}\$ is the smallest σ -field in Ω to make f measurable.

- If $f(\omega) = a \in \Omega'$, $\forall \omega \in \Omega$, then $\sigma(f) = {\Omega, \emptyset}$.
- If $f(\omega) = a\mathbf{1}_A + b\mathbf{1}_{A^c}$, $a \neq b$, then $\sigma(f) = \{\Omega, \emptyset, A, A^c\} = \sigma(A)$.

Induced σ -field: $\{B \subset \Omega' : f^{-1}(B) \in \mathscr{F}\}\$ is the largest σ -field in Ω' to make f measurable.

Theorem 2.17 (Composition of measurable functions). Suppose $f:(\Omega, \mathscr{F}) \to (\Omega', \mathscr{G})$ is $(\mathscr{F}, \mathscr{G})$ -measurable, and $g:(\Omega',\mathscr{G}) \to (\Omega'',\mathscr{H})$ is $(\mathscr{F},\mathscr{H})$ -measurable. Then the composed function $h=g\circ f$ is $(\mathscr{F},\mathscr{H})$ -measurable.

Theorem 2.18 (Check measurability). Let $f:(\Omega,\mathscr{F})\to (\Omega',\mathscr{G})$ be a function, $\sigma(\mathscr{C})=\mathscr{G}$. Then f is $(\mathscr{F},\mathscr{G})$ -measurable iff $f^{-1}(A)\in\mathscr{F},\,\forall A\in\mathscr{C}$.

Theorem 2.19 (Measurability on \mathbb{R}^m). A continuous function $f: \mathbb{R}^m \to \mathbb{R}^n$ is $(\mathscr{B}(\mathbb{R}^m), \mathscr{B}(\mathbb{R}^n))$ -measurable. In particular, a continuous real-valued function $f: \mathbb{R}^m \to \mathbb{R}$ is measurable.

Theorem 2.20 (Random vector). $X = (X_1, \dots, X_d)^T$ is a random vector iff X_i is a rv, $\forall i = 1, \dots, d$.

 $\sigma(X_1, \ldots, X_d)$: the smallest σ -field to make any X_i measurable, $\forall i = 1, \ldots, d$, which is $\sigma(\sigma(X_1), \ldots, \sigma(X_d))$; it is also the smallest σ -field to make $\mathbf{X} = (X_1, \ldots, X_d)^T$ measurable.

 $\sigma(X_1, X_2, \ldots) = \sigma(\sigma(X_1), \sigma(X_2), \ldots)$ or $\sigma(X: \Omega \to \mathbb{R}^{\infty}) = \sigma((-\infty, x_1] \times \cdots (-\infty, x_d], \ x \in \mathbb{R}^d, \ d = 1, 2, \ldots)$, that is, any finite collection of X_i is measurable.

2.2.1 Operations of measurable functions

If X_1, \ldots, X_n are rvs, $f: (\mathbb{R}^n, \mathscr{B}(\mathbb{R}^n)) \to (\mathbb{R}, \mathscr{B}(\mathbb{R}))$ is measurable, then $f(X_1, \ldots, X_n)$ is a rv by composition.

Theorem 2.21. If $f, g: \Omega \to \mathbb{R}$ are \mathscr{F} -measurable, then f+g, f-g, $c \cdot f$, and f/g are measurable, where c is a constant, and in f/g, we assume that $g(\omega) \neq 0$, $\forall \omega$.

Define the extended real line $\mathbb{R}^* = [-\infty, \infty]$, $\mathscr{B}(\mathbb{R}^*) = \sigma\{A, A \cup \{-\infty\}, A \cup \{\infty\}, A \cup \{\infty, -\infty\} : A \in \mathscr{B}(\mathbb{R})\} = \sigma((-\infty, x] : x \in \mathbb{R} \cup \{\infty\})$. By convention, $0 \cdot \infty = 0$ and $0 \cdot (-\infty) = 0$. If $f : \Omega \to \mathbb{R}^*$ is measurable, then it is called generalized random variable.

Theorem 2.22. Suppose $f_i: \Omega \to \mathbb{R}^*$ is measurable for $i \in \mathbb{N}^+$, then the functions $\inf_i f_i$, $\sup_i f_i$, $\lim_i \inf_i f_i$, $\lim_i \sup_i f_i$ (if exists)

are measurable, i.e., generalized rvs.

Let $\Omega_0 = \{\omega \in \Omega : \lim_n f_n(\omega) \text{ exists and finite}\} = \{\omega \in \Omega : \lim \sup_n f_n(\omega) - \lim \inf_n f_n(\omega) = 0\} \in \mathscr{F}$. If $P(\Omega_0) = 1$, then we say $\{f_n\}_{n=1}^{\infty}$ converges almost surely.

2.2.2 Random variables and distributions

Let $X: (\Omega, \mathscr{F}, \mathbb{P}) \to (\mathbb{R}, \mathscr{B}(\mathbb{R}))$ be a rv. Then its induced measure on \mathbb{R} is called the probability distribution of X, i.e., $\mathbb{P}(X \in B) := \mathbb{P}(X^{-1}(B)), \forall B \in \mathscr{B}(\mathbb{R})$. The distribution function (df) of X is defined to be $F_X(x) = \mathbb{P}(X \le x)$, $\forall x \in \mathbb{R}$. The properties in Theorem 2.7 hold.

• (Fact) If X has cdf F, F is cts, then $Y = F(X) \sim \text{Unif}(0,1)$.

Theorem 2.23. Let $\Omega=(0,1), \mathscr{F}=\mathscr{B}((0,1)),$ and \mathbb{P} is the Lebesgue measure on (0,1). Let F be an arbitrary distribution function. Define $X(\omega)=F^{-1}(\omega), \ \omega\in(0,1),$ where

$$F^{-1}(\omega) := \inf\{y \in \mathbb{R} : F(y) \ge \omega\} (= \sup\{y \in \mathbb{R} : F(y) < \omega\}).$$

Then X is regarded as a rv in Ω having df F.

It is useful to define several rvs on the same probability space, called coupling.

2.3 Statistical independence

2.4 Expectation

Lemma 2.24. Let $X \ge 0$, p > 0, we have $\mathbb{E}X^p = \int_0^\infty px^{p-1} \mathrm{Unif}(X > x) dx$.

3 Law of Large Numbers

3.1 Almost Surely Convergence

This lemma gives an equivalent relation between expectation and sum of tail probability.

Lemma 3.1. Let X_i iid and $\varepsilon > 0$, then $\sum_{n=1}^{\infty} \mathbb{P}(|X_n| > n\varepsilon) \le \varepsilon^{-1} \mathbb{E} |X_i| \le \sum_{n=0}^{\infty} \mathbb{P}(|X_n| > n\varepsilon)$.

4 Central Limit Theorem

5 Random Walks

Random walk (RW): Let X_i be iid rvs in \mathbb{R}^d . Let $S_n = \sum_{i=1}^n X_i$. Then $\{S_n : n \geq 1\}$ is called a RW. Take $S_0 = 0$. Simple random walk (SRW): If $\mathbb{P}(X_i = 1) = \mathbb{P}(X_i = -1) = 1/2$, then $\{S_n\}$ is called a SRW in \mathbb{R}^1 . If $\mathbb{P}(X_i = (1, 1)) = \mathbb{P}(X_i = (1, -1)) = \mathbb{P}(X_i = (-1, -1)) = 1/4$, then called a SRW in \mathbb{R}^2 .

5.1 Stopping Times (A.1.1)

Long-term behavior of RW

Permutable (or exchangeable): An event that does not change under finite permutation of $\{X_1, X_2, \ldots\}$.

- All events in the tail σ -field \mathcal{T} are permutable.
- $\{\omega: \mathbf{S}_n(\omega) \in B \text{ i.o.}\}\$ is permutable but not tail event.
- $\{\omega : \limsup_{n \to \infty} \mathbf{S}_n(\omega)/c_n \ge 1\}.$

Theorem 5.1 (Hewitt-Savage 0-1 law). If X_i iid and event A is permutable, then $\mathbb{P}(A) = 0$ or 1.

Theorem 5.2 (Long-term behavior of RW). For a RW in \mathbb{R} , one of the following has probability 1:

- (i) $S_n = 0$ for all n;
- (ii) $S_n \to \infty$ as $n \to \infty$;
- (iii) $S_n \to -\infty$ as $n \to \infty$;

(iv) $-\infty = \liminf_n S_n < \limsup_n S_n = \infty$.

For two levels a < b, find the probability that RW reaches b before a

Filtration: Let X_i be a sequence of rvs, $\{\mathcal{F}_n := \sigma(X_1, \dots, X_n)\}_{n=1}^{\infty}$ as an increasing sequence of σ -fields, is called a filtration. We usually take $\mathcal{F}_0 = \{\phi, \Omega\}$.

Stopping time/optional random variable/optimal time/Markov time: $\tau \in \mathbb{N}^+ \cup \{\infty\}$ is a stopping time w.r.t. $\{\mathcal{F}_n\}$ if $\{\tau = n\} \in \mathcal{F}_n$, $\forall n \in \mathbb{N}^+$. (Equivalent def: $\{\tau \leq n\} \in \mathcal{F}_n$ or $\{\tau \geq n+1\} \in \mathcal{F}_n$ for $n \in \mathbb{N}^+$)

- Constant $\tau = n$ is a stopping time.
- If τ_1, τ_2 are stopping time, then $\tau_1 \wedge \tau_2, \tau_1 \vee \tau_2, \tau_1 + \tau_2$ are stopping times.
- Hitting time of A: let A measurable, then $\tau = \inf\{n \geq 1 : S_n \in A\}$ is a stopping time.
- σ -field \mathcal{F}_N =the information known at time N. Def: \mathcal{F}_N is the collection of sets A that have $A \cup \{N = n\} \in \mathcal{F}_n$, $\forall n < \infty$. Example: $\{N \leq n\} \in \mathcal{F}_N$, i.e., N is \mathcal{F}_N -measurable.

Theorem 5.3 (Wald's equation). Let X_i iid and τ be a stopping time.

- 1. (Wald's first equation) If $\mathbb{E}|X_1| < \infty$ and $\mathbb{E}\tau < \infty$, then $\mathbb{E}S_\tau = \mathbb{E}X_1\mathbb{E}\tau$.
- 2. (Wald's second equation) If $\mathbb{E}X_1 = 0$, $\mathbb{E}X_1^2 = \sigma^2 < \infty$, $\mathbb{E}\tau < \infty$, then $\mathbb{E}S_{\tau}^2 = \sigma^2 \mathbb{E}\tau$.

Example 5.4 (Results for 1-d SRW). For 1-d SRW, let $a, b \in \mathbb{Z}$, a < 0 < b. Let $N = \inf\{n : S_n \notin (a, b)\} = \inf\{n : S_n = a \text{ or } b\}$. Then

- 1. $\mathbb{E}N < \infty$,
- 2. $S_N = a \text{ or } b$,
- 3. $\mathbb{P}(S_N = a) = b/(b-a), \ \mathbb{P}(S_N = b) = -a/(b-a),$
- 4. $\mathbb{E}N = \mathbb{E}S_N^2 = (-a)b$.

5.2 Recurrence vs. Transience (A.1.2)

When RW return to 0? We consider SRW on \mathbb{R}^d and define its first, second, ..., nth returning time to the origin to be

$$\tau_1 = \inf\{m \ge 1 : \mathbf{S}_m = \mathbf{0}\},
\tau_n = \inf\{m > \tau_{n-1} : \mathbf{S}_m = \mathbf{0}\}.$$

Theorem 5.5. For any RW, the following are equivalent:

- (i) $\mathbb{P}(\tau_1 < \infty) = 1$
- (ii) $\mathbb{P}(\tau_n < \infty) = 1, \forall n = 1, 2, 3, \dots$
- (iii) $\mathbb{P}(S_m = \mathbf{0} \ i.o.) = 1$
- (iv) $\sum_{m=1}^{\infty} \mathbb{P}(S_m = \mathbf{0}) = \infty$.
- $\mathbb{P}(\tau_n < \infty) = (\mathbb{P}(\tau_1 < \infty))^n$.

Recurrent: If $\mathbb{P}(\tau_1 < \infty) = 1$, then the RW is called recurrent.

Transient: If $\mathbb{P}(\tau_1 < \infty) < 1$, then the RW is called Transient.

Theorem 5.6 (Recurrence of SRW). SRW is recurrent in $d \leq 2$ and transient in $d \geq 3$.

• We define the first time a random walk starting from \boldsymbol{a} reaches \boldsymbol{b} : $\tau_{\boldsymbol{a}\to\boldsymbol{b}}:=\inf\{m\geq 1:\boldsymbol{a}+\boldsymbol{S}_m=\boldsymbol{b}\}$. It can be proved that $\mathbb{P}(\tau_1<\infty)=1$ iff $\mathbb{P}(\tau_{\boldsymbol{a}\to\boldsymbol{b}}<\infty)=1$, $\forall \boldsymbol{a},\boldsymbol{b}$.

5.3 Reflection Principle and Arcsine Distribution (A.1.3)

What is the distribution of the time spent above 0? We consider the SRW, d = 1, and think of the sequence S_1, \ldots, S_n as being represented by a polygonal line with segments $(k-1, S_{k-1}) \to (k, S_k)$.

Theorem 5.7 (Reflection Principle).

- (Reflection principle for numbers) If x, y > 0, then the number of paths from (0, x) to (n, y) that are 0 at some time is equal to the number of paths from (0, -x) to (n, y).
- (Reflection principle for SRW) Let X_i be SRW with d=1. Then $\forall b \in \mathbb{N}^+$,

$$\mathbb{P}(\max_{1 \le k \le n} S_k \ge b) = 2\mathbb{P}(S_n > b) + \mathbb{P}(S_n = b).$$

Theorem 5.8 (Hit 0 time). $\mathbb{P}(S_1 > 0, \dots, S_{2n} > 0) = \frac{1}{2}\mathbb{P}(S_{2n} = 0)$, and $\mathbb{P}(S_1 \neq 0, \dots, S_{2n} \neq 0) = \mathbb{P}(S_{2n} = 0)$.

Arcsine distribution: a continuous distribution with density $\frac{1}{\pi\sqrt{x(1-x)}}$, $x \in (0,1)$. Define

 $L_{2n} := \sup\{m \le 2n : S_m = 0\},$ (last time at 0) $F_n := \inf\{0 \le m \le n : S_m = \max_{0 \le k \le n} S_k\},$ (first time at maximum)

 $\pi_{2n} := \text{ number of } k: 1 \leq k \leq 2n \text{ such that the line } (k-1, S_{k-1}) \to (k, S_k) \text{ is above the } x\text{-axis.}$

Theorem 5.9 (Arcsine law). $\frac{L_{2n}}{2n}, \frac{F_n}{n}, \frac{\pi_{2n}}{2n}$ all converge in distribution to the arcsine distribution.

6 Martingales

6.1 Conditional expectation A.2.1

Definition 6.1 (Conditional expectation). X is a random variable on $(\Omega, \mathcal{F}, \mathbb{P})$ with $\mathbb{E}|X| < \infty$. \mathcal{A} is a σ -field and $\mathcal{A} \subset \mathcal{F}$. We define the conditional expectation of X given \mathcal{A} , $\mathbb{E}(X|\mathcal{A})$, to be any random variable Y satisfying

- (i) Y is A-measurable, and
- (ii) $\forall A \in \mathcal{A}, \mathbb{E}[X\mathbf{1}_A] = \mathbb{E}[Y\mathbf{1}_A].$
 - For rvs X and Y, $\mathbb{E}(X|Y) := \mathbb{E}(X|\sigma(Y))$.
 - For set A, $\mathbb{P}(A|A) := \mathbb{E}(\mathbf{1}_A|A)$.
 - For set $A, B, \mathbb{P}(A|B) := \mathbb{P}(A \cap B)/\mathbb{P}(B)$ given $\mathbb{P}(B) > 0$.
- If Y satisfies (i) and (ii), then $\mathbb{E}|Y| \leq \mathbb{E}|X| < \infty$.
- $\mathbb{E}Y = \mathbb{E}X$, i.e., $\mathbb{E}[\mathbb{E}(X|\mathcal{A})] = \mathbb{E}X$.
- Uniqueness: If Y' also satisfies (i) and (ii), then Y = Y' a.s. Any such Y is said to be a version of $\mathbb{E}(X|\mathcal{A})$.
- Existence: By Radon-Nikodym theorem.

Example 6.2.

- If X is A-measurable, then $\mathbb{E}(X|A) = X$. So a constant $c = \mathbb{E}(c|A)$. (Know X, the "best guess" is X)
- If $\sigma(X) \perp \!\!\!\perp A$, then $\mathbb{E}(X|A) = \mathbb{E}X$. (Don't know anything about X, the best guess is its mean)
- Suppose $\Omega_1, \Omega_2, \ldots$ is a finite or infinite partition of Ω into disjoint sets with positive probability, and let $\mathcal{A} = \sigma(\Omega_1, \ldots)$, then $\mathbb{E}(X|\mathcal{A}) = \mathbb{E}(X\mathbf{1}_{\Omega_i})/\mathbb{P}(\Omega_i)$ on Ω_i .
 - Let $\mathcal{A} = \{\emptyset, \Omega\}$, then $\mathbb{E}(X|\mathcal{A}) = \mathbb{E}X$.
- (Bayes's formula) Let $G \in \mathcal{G}$, then $\mathbb{P}(G|A) = \int_G \mathbb{P}(A|\mathcal{G}) d\mathbb{P} / \int_{\Omega} \mathbb{P}(A|\mathcal{G}) d\mathbb{P}$. When \mathcal{G} is a σ -field generated by a partition, this reduces to the usual Bayes' formula $\mathbb{P}(G_i|A) = \mathbb{P}(A|G_i)P(G_i) / \sum_j \mathbb{P}(A|G_j)\mathbb{P}(G_j)$.

Theorem 6.3 (Properties). If $\mathbb{E}|X|$, $\mathbb{E}|X_n|$, $\mathbb{E}|Y| < \infty$, then

- (a) (linearity) $\mathbb{E}(aX + Y|\mathcal{A}) = a\mathbb{E}(X|\mathcal{A}) + \mathbb{E}(Y|\mathcal{A}).$
- (b) (monotonicity) If $X \leq Y$, then $\mathbb{E}(X|\mathcal{A}) \leq \mathbb{E}(Y|\mathcal{A})$.
- (c) (MCT) If $X_n \geq 0$, $X_n \uparrow X$, and $\mathbb{E}X < \infty$, then $E(X_n | A) \uparrow E(X | A)$.
- (d) (Fatou) If $X_n \geq 0$, $\mathbb{E}X_n < \infty$, and $\mathbb{E}[\liminf_n X_n] < \infty$, then $\liminf_n \mathbb{E}(X_n | \mathcal{A}) \leq \mathbb{E}[\liminf_n X_n | \mathcal{A}]$.
- (e) (DCT) If $X_n \to X$ a.s., $|X_n| \le Y$, $\mathbb{E}|Y| < \infty$, then $\mathbb{E}(X_n|\mathcal{A}) \to \mathbb{E}(X|\mathcal{A})$.
- (f) If X is A-measurable, $\mathbb{E}|XY| < \infty$, then $\mathbb{E}(XY|A) = X\mathbb{E}(Y|A)$.
- (g) (Tower property) If $A_1 \subset A_2$, then $\mathbb{E}[\mathbb{E}(X|A_1)|A_2] = \mathbb{E}(X|A_1)$, and $\mathbb{E}[\mathbb{E}(X|A_2)|A_1] = \mathbb{E}(X|A_1)$.

Theorem 6.4 (Inequalities). If $\mathbb{E}|X|$, $\mathbb{E}|Y| < \infty$, then

- (i) (Jensen's inequality) If φ is convex, $\mathbb{E}|\varphi(X)| < \infty$, then $\varphi(\mathbb{E}[X|\mathcal{A}]) \leq \mathbb{E}[\varphi(X)|\mathcal{A}]$.
- (ii) (Markov's inequality) If $X \ge 0$, a > 0, then $\mathbb{P}(X \ge a | \mathcal{A}) \le a^{-1} \mathbb{E}(X | \mathcal{A})$.
- (iii) (Chebyshev's inequality) If a > 0, then $\mathbb{P}(|X| \ge a|\mathcal{A}) \le a^{-2}\mathbb{E}(X^2|\mathcal{A})$.
- (iv) (Hölder's Inequality) If $p \ge 1$, $p^{-1} + q^{-1} = 1$, and $\mathbb{E}|X|^p$, $\mathbb{E}|Y|^p < \infty$, then

$$|\mathbb{E}(XY|\mathcal{A})| \leq {\mathbb{E}(|X|^p |\mathcal{A})}^{1/p} {\mathbb{E}(|Y|^q |\mathcal{A})}^{1/q}.$$

(v) (Minkowski inequality) If $p \geq 1$, $\mathbb{E}|X|^p$, $\mathbb{E}|Y|^p < \infty$, then

$$\{\mathbb{E}(|X+Y|^p|A)\}^{1/p} \le \{\mathbb{E}(|X|^p|A)\}^{1/p} + \{\mathbb{E}(|Y|^p|A)\}^{1/p}.$$

(vi) (Triangular inequality) If $\mathbb{E}X^2 < \infty$. Then for any A-measurable Y with $\mathbb{E}Y^2 < \infty$, we have

$$\left\|X - \mathbb{E}(X|\mathcal{A})\right\|^{2} \leq \left\|X - Y\right\|^{2}.$$

By (vi), $\mathbb{E}(X|\mathcal{A})$ is the projection of X onto $\mathcal{L}^2(\mathcal{A})$, that is, $\mathbb{E}(X|\mathcal{A}) = \arg\min_Y \|X - Y\|^2$, for \mathcal{A} -measurable Y.

6.2 Martingales

Definition 6.5. Let $\{\mathcal{F}_n\}$ be a filtration, an increasing sequence of σ -fields. A sequence $\{S_n\}$ is said to be adapted to $\{\mathcal{F}_n\}$ if S_n is \mathcal{F}_n -measurable. $\{S_n\}$ is called a martingale w.r.t. $\{\mathcal{F}_n\}$ if

- (i) $\mathbb{E}|S_n| < \infty$.
- (ii) $\{S_n\}$ is adapted to $\{\mathcal{F}_n\}$.
- (iii) $\mathbb{E}(S_n|\mathcal{F}_{n-1}) = S_{n-1}$.

If in (iii), $\mathbb{E}(S_n|\mathcal{F}_{n-1}) \leq S_{n-1}$ (or \geq), then $\{S_n\}$ is said to be a supermartingale (or submartingale).

Simple facts:

- If $\{S_n\}$ is a martingale, then $\mathbb{E}S_1 = \cdots = \mathbb{E}S_n = \cdots$, and $\mathbb{E}|S_1| \leq \mathbb{E}|S_2| \leq \cdots$.
- If $\{S_n\}$ is a supermartingale, then $\mathbb{E}S_1 \geq \mathbb{E}S_2 \geq \cdots$; if submartingale, $\mathbb{E}S_1 \leq \mathbb{E}S_2 \leq \cdots$.
- If $\{S_n\}$ is a supermartingale, then $\{-S_n\}$ is a submartingale, and vice versa.
- Let n > m, if $\{S_n\}$ is a
 - martingale $\Longrightarrow \mathbb{E}(S_n|\mathcal{F}_m) = S_m$
 - supermartingale $\implies \mathbb{E}(S_n|\mathcal{F}_m) \leq S_m$
 - submartingale $\implies \mathbb{E}(S_n|\mathcal{F}_m) \geq S_m$

Theorem 6.6 (Martingale transforms).

- (1) If $\{S_n\}$ is a martingale and φ is a convex (concave) function such that $\mathbb{E}|\varphi(S_n)| < \infty$, then $\varphi(S_n)$ is a submartingale (supermartingale).
- (2) If S_n is a submartingale and φ is an increasing convex (concave) function such that $\mathbb{E}|\varphi(S_n)| < \infty$, then $\varphi(S_n)$ is a submartingale (supermartingale).
- From (1), if $p \ge 1$ and $\mathbb{E}|S_n|^p < \infty$, then $|S_n|^p$ is a submartingale.
- From (2), if S_n is a submartingale, then $(X_n a)^+$ is a submartingale; if X_n is a supermartingale, then $X_n \wedge a$ is a supermartingale.

6.3 Martingale convergence

Predictable sequence: H_n , $n \ge 2$, which is \mathcal{F}_{n-1} -measurable. $(H \cdot S)_n = \sum_{m=1}^n H_m(S_m - S_{m-1})$.

- If $\{S_n\}$ is a supermartingale (submartingale), H_n $(n \ge 2)$ is predictable, $H_n \ge 0$, and H_n is bounded. Then $(H \cdot S)_n$ is a supermartingale (submartingale). For martingale, it is true without assuming $H_n \ge 0$.
- Let N be a stopping time, and $H_n = \mathbf{1}_{\{n \leq N\}}$, then $S_{n \wedge N} = (H \cdot S)_n + S_0$ is a supermartingale/submartingale/martingale as S_n is.

Let a < b, $N_0 = 1$, for $k \ge 1$, let $N_{2k-1} = \inf\{m > N_{2k-2} : S_m \le a\}$ and $N_{2k} = \inf\{m > N_{2k-1} : S_m \ge b\}$, all are stopping times. Let $H_m = \mathbf{1}\{N_{2k-1} < m < N_{2k} \text{ for some } k\}$ be the indicator of climbing. Let $U_n = \sup\{k : N_{2k} \le n\}$ be the number of upcrossings by time n. From the picture we have

- (1) H_m is predictable;
- (2) $(b-a)U_n \le \sum_{m=2}^n H_m(S_m S_{m-1}) \Rightarrow (b-a)\mathbb{E}U_n \le \mathbb{E}(S_n S_1).$

Theorem 6.7 (Upcrossing inequality). $\{S_n\}$ is a submartingale. a < b are two constants. For U_n defined above,

$$\mathbb{E}U_n \le \frac{1}{b-a} \left[\mathbb{E}(S_n - a)^+ - \mathbb{E}(S_1 - a)^+ \right]$$

Theorem 6.8 (Martingale convergence theorem). Suppose S_n is a submartingale and $\liminf \mathbb{E} S_n^+ < \infty$ (or $\sup \mathbb{E} S_n^+ < \infty$), then $S_n \to S$ a.s. with $\mathbb{E} |S| < \infty$.

Theorem 6.9. Suppose $\{S_n\}$ is a supermartingale. If $S_n \geq 0$, then $S_n \to S$ a.s. and $\mathbb{E}S \leq \mathbb{E}S_1 < \infty$.

Corollary 6.10 (WLLN for martingales). Suppose X_i are identically distributed and $\mathbb{E}|X_1| < \infty$. Let $\mathcal{F}_n = \sigma(X_1, \ldots, X_n)$. Let $S_1 = X_1$, and $S_n = S_{n-1} + X_n - \mathbb{E}(X_n | \mathcal{F}_{n-1})$, $n \geq 2$. Then $\{S_n\}$ is a martingale and $S_n/n \stackrel{\text{p}}{\to} 0$.

Example 6.11 (Branching processes). Let $\xi_i^n \in \mathbb{N}$, $i, n \geq 1$ be iid (n: time, i: the ith parent). Define Z_n by $Z_0 = 1$, and

$$Z_{n+1} = \begin{cases} \xi_1^{n+1} + \dots + \xi_{Z_n}^{n+1} & \text{if } Z_n > 0\\ 0 & \text{if } Z_n = 0. \end{cases}$$

 Z_n is called a Galton-Watson process. Let $\mathcal{F}_n = \sigma(\xi_i^m : i \ge 1, 1 \le m \le n)$ and $\mu = \mathbb{E}\xi_i^m \in (0, \infty)$, then $\{Z_n/\mu^n\}$ is a martingale w.r.t. \mathcal{F}_n .

- If $\mu < 1$, then $Z_n = 0$ for all n sufficiently large, so $Z_n/\mu^n \to 0$.
- If $\mu = 1$ and $\mathbb{P}(\xi_i^m = 1) < 1$, then $Z_n = 0$ for all n sufficiently large.

6.4 Doob's inequality; \mathcal{L}^p convergence; CLT

Theorem 6.12 (Doob's inequality). $\{S_n\}$ is a submartingale w.r.t. $\{\mathcal{F}_n\}$. Then $\forall x > 0$,

$$\mathbb{P}\left(\max_{1\leq k\leq n}S_k\geq x\right)\leq \frac{1}{x}\mathbb{E}\left[S_n\mathbf{1}\{\max_{1\leq k\leq n}S_k\geq x\}\right]\leq \frac{\mathbb{E}S_n^+}{x}.$$

7 Techniques

7.1 Convergence

7.1.1 Convergence of random series

Let X_i be a sequence of rvs. $S_n = \sum_{i=1}^n X_i$. By 0-1 law, $\mathbb{P}(\lim_n S_n \text{ exists}) = 0$ or 1.

To show the convergence of random series:

1.

To show the divergence of random series:

1. If SLLN holds, $S_n/n \to \mu$ a.s., if $\mu > 0$, then $S_n \to \infty$ a.s.

2. Use S'_n , $S'_n = S_n$ e.v. a.s., that is, $\mathbb{P}(S_n \neq S'_n \ i.o.) = 0$. And show $S'_n \to \infty$ a.s. Next, we consider $S_n/f(n)$.

A Proofs

A.1 Proofs - 5

A.1.1 Proofs - 5.1

<u>Proof of Theorem 5.2.</u> By the 0-1 law 5.1, $\{\limsup_n S_n \geq c\}$ has probability 0 or 1, meaning that $\limsup_n S_n = c \in [-\infty, \infty]$ w.p.1. Since $S_n \stackrel{\text{d}}{=} S_{n+1} - X_1$, we have $c = c - X_1$.

(i) If $c \in \mathbb{R}$, then $X_1 \equiv 0$ a.s., so $S_n = 0$ for all n a.s.

If $X_1 \neq 0$ a.s., then $c = -\infty$ or ∞ ,

- (ii) If $c = \infty$, and $\liminf_n S_n = \infty$, then case (ii);
- (iii) If $c = -\infty$, and $\liminf_n S_n = -\infty$, then case (iii);
- (iv) If $c = \infty$, and $\liminf_n S_n = -\infty$, then case (iv).

Proof of Theorem 5.3. Prove 1: First suppose $X_i \geq 0$. We have

$$\mathbb{E} S_{\tau} = \mathbb{E} \sum_{i=1}^{\tau} X_i = \mathbb{E} \sum_{i=1}^{\infty} X_i \mathbf{1}_{\{\tau \geq i\}} = \sum_{i=1}^{\infty} \mathbb{E} X_i \mathbf{1}_{\{\tau \geq i\}} = \sum_{i=1}^{\infty} \mathbb{E} X_i \mathbb{E} \mathbf{1}_{\{\tau \geq i\}} = \mathbb{E} X_1 \mathbb{E} \tau,$$

where the 3rd equality uses Fubini by $X_i \ge 0$, and the 4th uses $\{\tau \ge i\} \in \mathcal{F}_{i-1}$. For general case, since $\sum_{i=1}^{\infty} \mathbb{E} |X_i| \mathbf{1}_{\{\tau \ge i\}} = \sum_{i=1}^{\infty} \mathbb{E} |X_i| \mathbb{E} \mathbf{1}_{\{\tau \ge i\}} < \infty$, we can still use the Fubini.

Prove 2: If $\tau < n$, then $\tau \wedge n = \tau \wedge (n-1)$, so $S_{\tau \wedge n}^2 = S_{\tau \wedge (n-1)}^2$; if $\tau \ge n$, we have $\tau \wedge n = n$ and $\tau \wedge (n-1) = n-1$, so $S_{\tau \wedge n}^2 = S_n^2 = (S_{n-1} + X_n)^2 = (S_{\tau \wedge (n-1)} + X_n)^2$. Hence write

$$S_{\tau \wedge n}^2 = S_{\tau \wedge (n-1)}^2 + (2X_n S_{n-1} + X_n^2) \mathbf{1}_{\{\tau \ge n\}}.$$

Note that all the following expectations exist,

$$\mathbb{E}S_{\tau \wedge n}^{2} = \mathbb{E}S_{\tau \wedge (n-1)}^{2} + \mathbb{E}\left(2X_{n}S_{n-1}1_{\{\tau \geq n\}}\right) + \mathbb{E}[X_{n}^{2}1_{\{\tau \geq n\}}]$$

$$= \mathbb{E}S_{\tau \wedge (n-1)}^{2} + \sigma^{2}\mathbb{P}(\tau \geq n) \qquad \text{(stopping time, independence, and } \mathbb{E}X_{i} = 0)$$

$$= \dots \qquad \text{(reduce to } n-2, n-3, \dots)$$

$$= \sigma^{2}\sum_{i=1}^{n}\mathbb{P}(\tau \geq i).$$

In the 1st line, the expectation $\mathbb{E}X_nS_{n-1}$ exists since both rvs are in \mathcal{L}^2 . By the last line, $\|S_{\tau\wedge n} - S_{\tau\wedge m}\|^2 = \sigma^2 \sum_{i=m+1}^n \mathbb{P}(\tau \geq i) \to 0$ as $n, m \to \infty$, $\{S_{\tau\wedge n}\}_n$ is a Cauchy sequence in \mathcal{L}^2 , so letting $n \to \infty$ gives the result.

Proof of Example 5.4. 1. For any positive integer k, by dividing the interval (0, k(b-a)) into k subintervals of equal length and considering an extreme case behavior (keep going upwards) of the random walk within each subinterval, we obtain

$$\mathbb{E}N = \sum_{i=0}^{\infty} \mathbb{P}(N > i) \le (b - a) \sum_{k=0}^{\infty} \mathbb{P}(N > k(b - a))$$

$$\le (b - a) \sum_{k=0}^{\infty} \mathbb{P}((X_{(j-1)(b-a)+1}, \dots, X_{j(b-a)}) \ne (1, \dots, 1), j = 1, \dots, k)$$

$$\le (b - a) \sum_{k=0}^{\infty} \left(1 - \frac{1}{2^{b-a}}\right)^k < \infty.$$

- 2. It is obvious.
- 3. By Wald's first equation 5.3, $0 = \mathbb{E}S_N = a\mathbb{P}(S_N = a) + b\mathbb{P}(S_N = b)$, we also have $1 = \mathbb{P}(S_N = a) + \mathbb{P}(S_N = b)$, so solve for the result.
- 4. By Wald's second equation 5.3 and $\sigma=1$, we have $\mathbb{E}N=\mathbb{E}S_N^2=a^2\mathbb{P}(S_N=a)+b^2\mathbb{P}(S_N=b)$, and use 3.

A.1.2 Proofs - 5.2

Proof of Theorem 5.5. We have

$$\mathbb{P}(\tau_{2} < \infty) = \mathbb{P}(\tau_{1} < \infty, \tau_{2} - \tau_{1} < \infty)$$

$$= \sum_{m,n=1}^{\infty} \mathbb{P}(\tau_{1} = m, \tau_{2} - \tau_{1} = n)$$

$$= \sum_{m,n=1}^{\infty} \mathbb{P}(\boldsymbol{X}_{1} + \dots + \boldsymbol{X}_{m} = \boldsymbol{0}, \boldsymbol{X}_{1} + \dots + \boldsymbol{X}_{u} \neq \boldsymbol{0}, \forall 1 \leq u < m;$$

$$\boldsymbol{X}_{m+1} + \dots + \boldsymbol{X}_{m+n} = \boldsymbol{0}, \boldsymbol{X}_{m+1} + \dots + \boldsymbol{X}_{m+v} \neq \boldsymbol{0}, \forall 1 \leq v < n)$$

$$= \sum_{m,n=1}^{\infty} \mathbb{P}(\tau_{1} = m)\mathbb{P}(\tau_{1} = n)$$

$$= (\mathbb{P}(\tau_{1} < \infty))^{2}.$$
(iid)

Similarly, we can prove $\mathbb{P}(\tau_n < \infty) = (\mathbb{P}(\tau_1 < \infty))^n$. So (i) and (ii) are equivalent. They are equivalent to (iii) by examining their meanings. Finally,

$$\sum_{m=0}^{\infty} \mathbb{P}(\boldsymbol{S}_m = \boldsymbol{0}) = \sum_{m=0}^{\infty} \mathbb{E} \boldsymbol{1}_{\{\boldsymbol{S}_m = \boldsymbol{0}\}} = \mathbb{E} \sum_{m=0}^{\infty} \boldsymbol{1}_{\{\boldsymbol{S}_m = \boldsymbol{0}\}} = \mathbb{E} \sum_{n=0}^{\infty} \boldsymbol{1}_{\{\tau_n < \infty\}}$$
$$= \sum_{n=0}^{\infty} \mathbb{P}(\tau_n < \infty) = \sum_{n=0}^{\infty} \mathbb{P}(\tau_1 < \infty)^n = \frac{1}{1 - \mathbb{P}(\tau_1 < \infty)}.$$

So (i) and (iv) are equivalent.

Proof of Theorem 5.6. In d = 1, use (iv) in Theorem 5.5 to show.

$$\sum_{m=1}^{\infty} P(S_m = 0) = \sum_{n=1}^{\infty} P(S_{2n} = 0)$$
 (can only return to 0 at even steps)
$$= \sum_{n=1}^{\infty} \binom{2n}{n} (\frac{1}{2})^{2n}$$
 (combinatorics)
$$\sim \sum_{n=1}^{\infty} \frac{\sqrt{2\pi 2n} (\frac{2n}{e})^{2n}}{(\sqrt{2\pi n} (\frac{n}{e})^n)^2} \frac{1}{2^{2n}}$$
 (Stirling's formula)
$$= \sum_{n=1}^{\infty} \frac{1}{\sqrt{\pi n}} = \infty.$$

In d=2, note that in order for $S_{2n}=0$, we must for some $0 \le m \le n$ have m up steps, m down steps, n-m to the left, and n-m to the right, so

$$\mathbb{P}(S_{2n} = \mathbf{0}) = \frac{1}{4^{2n}} \sum_{m=0}^{n} \binom{2n}{m} \binom{2n-m}{m} \binom{2n-2m}{n-m} = \frac{1}{4^{2n}} \sum_{m=0}^{n} \frac{(2n)!}{m!m!(n-m)!(n-m)!}$$
$$= \frac{1}{4^{2n}} \binom{2n}{n} \sum_{m=0}^{n} \binom{n}{m} \binom{n}{n-m} = 4^{-2n} \binom{2n}{n}^2 \times n^{-1}.$$

by Stirling's formula. So its sum is ∞ , still recurrent.

For d=3, more complicated combinatorics give $\mathbb{P}(S_{2n}=0) \approx \frac{1}{n^{3/2}}$, summing up to a finite number; hence transient. In even higher dimensions, the probabilities become even smaller; hence all transient.

A.1.3 Proofs - 5.3

Proof of Theorem 5.7. To show the first result, suppose $(0, s_0), (1, s_1), \ldots, (n, s_n)$ is a path from (0, x) to (0, y). Let $K = \inf\{k : s_k = 0\}$. Let $s'_k = -s_k$ for $k \le K$ and $s'_k = s_k$ for $K \le k \le n$. Then $(k, s'_k), 0 \le k \le n$, is a path from (0, -x) to (n, y). Conversely, given a path from (0, -x) to (n, y), we can also construct a reflected path. We have a one-toone correspondence between the two classes of paths, so their numbers must be equal.

Then, we have

$$\mathbb{P}(\max_{1 \le k \le n} S_k \ge b) = P(\max_{1 \le k \le n} S_k \ge b, S_n > b) + P(\max_{1 \le k \le n} S_k \ge b, S_n < b) + P(\max_{1 \le k \le n} S_k \ge b, S_n = b)$$

$$= \mathbb{P}(S_n > b) + \mathbb{P}(\max_{1 \le k \le n} S_k \ge b, S_n > b) + \mathbb{P}(S_n = b)$$

$$= 2\mathbb{P}(S_n > b) + \mathbb{P}(S_n = b),$$

which prove the second result.

<u>Proof of Theorem 5.8.</u> To count the number of paths from (0,0) to (n,x), denote $a,b \in \mathbb{N}$ be the number of positive steps and b negative steps, respectively. n = a + b, and x = a - b, where $x \in [-n,n]$, and n - x is even. The number of paths from (0,0) to (n,x) is $N_{n,x} = \binom{n}{a}$.

Since

$$\mathbb{P}(S_1 > 0, \dots, S_{2n} > 0) = \sum_{r=1}^{\infty} \mathbb{P}(S_1 > 0, \dots, S_{2n-1} > 0, S_{2n} = 2r).$$

Now we count the number of paths of $(1,1) \to (2n,2r)$, that are never 0. Since the total number of paths of $(1,1) \to (2n,2r)$ is $N_{2n-1,2r-1}$, the number of these paths touching 0 is the number of paths of $(1,-1) \to (2n,2r)$, i.e., $N_{2n-1,2r+1}$, by reflection principle, we have the number of paths of $(1,1) \to (2n,2r)$ never touching 0 is $N_{2n-1,2r-1} - N_{2n-1,2r+1}$. Hence

$$\sum_{r=1}^{\infty} \mathbb{P}(S_1 > 0, \dots, S_{2n-1} > 0, S_{2n} = 2r) = \sum_{r=1}^{\infty} \frac{1}{2} \frac{1}{2^{2n-1}} (N_{2n-1, 2r-1} - N_{2n-1, 2r+1}) = \frac{1}{2^{2n}} N_{2n-1, 1},$$

where the 1/2 in the 2nd term guarantees $S_1 > 0$. Since $\mathbb{P}(S_{2n} = 0) = \frac{1}{2^{2n}} N_{2n-1,-1} + \frac{1}{2^{2n}} N_{2n-1,1} = 2 \cdot \frac{1}{2^{2n}} N_{2n-1,1}$,

$$\mathbb{P}(S_1 > 0, \dots, S_{2n} > 0) = \frac{1}{2^{2n}} N_{2n-1,1} = \frac{1}{2} \mathbb{P}(S_{2n} = 0).$$

Symmetry implies $\mathbb{P}(S_1 < 0, \dots, S_{2n} < 0) = (1/2)\mathbb{P}(S_{2n} = 0)$. Then the proof is completed.

A.2 Proofs - 6

A.2.1 Proofs - 6.1

References