

Analysis

With an emphasis on probability theory

Sudip Sinha

November 06, 2019

CONTENTS

Part 1	Real analysis	1
1.1	Convergence of sequences and series	2
1.2	Differentiation	3
1.2.1	Differentiation of functions with real powers	3
1.2.2	Example of $f \in C^\infty \setminus C^\omega$	3
1.2.3	Taylor series for multivariate functionals	3
Part 2	Probability theory	5
2.1	Elementary ideas	6
2.1.1	σ -algebras	6
2.1.2	Examples	6
2.1.3	Measurability of $\inf, \sup, \liminf, \limsup$ TODO	6
2.1.4	Method of substitution	6
2.2	Borel–Cantelli lemmas	7
2.3	Modes of convergence	8
2.4	Conditioning	9
2.4.1	In L^2 , conditional expectation is a projection	9
2.4.2	Uncorrelated does not imply independence	9
2.4.3	$\phi_{aX+bY} = \phi_{aX}\phi_{bY} \forall (a, b) \in \mathbb{R}^2$ implies independence	9
2.5	Limiting behaviour of \bar{X}_n	10
Part 3	Stochastic analysis	11
3.1	Classification of stochastic processes	12
3.2	Martingales	13
3.2.1	New martingales from old	13
3.3	Markov processes	14
3.4	Itô calculus	15
3.4.1	Step 1: $X \in \mathcal{A} \cap \mathcal{S}$ a.s.	15
3.4.2	Step 2: $X \in \mathcal{A} \cap \mathcal{B} \cap \mathcal{C}$ a.s.	15
3.4.3	Step 3: $X \in \mathcal{A} \cap \mathcal{B} \cap L^0(T \times \Omega)$	15
3.4.4	Step 4: $X \in \mathcal{A} \cap L^2(T \times \Omega)$	15
3.4.5	Step 5: $X \in \mathcal{A} \cap \left\{ X \in \mathbb{C}^{T \times \Omega} : \forall t \geq 0, \int_0^t X(s, \cdot) ds < \infty \right\}$ a.s.	15
3.4.6	Properties of the Itô integral	15
3.4.7	Itô formula for multidimensional processes	16
3.4.8	Adapted and instantly independent implies deterministic	17

3.5 Examples	18
3.5.1 Find $\mathbb{E} \left(\int_0^1 B_t^2 dt \right)^2$.	18
3.5.2 Find $\mathbb{V} \left(\int_0^1 t^2 B_t dt \right)$.	18
3.5.3 Calculate $\int_0^T e^{B_t^2} dB_t$.	18
Part 4 Infinite-dimensional analysis	19
4.1 Giuseppe Da Prato	20
4.2 Abstract Wiener spaces	24
4.3 White noise distribution theory	25
4.3.1 Characterization theorem	25
Bibliography	26

PART 1

REAL ANALYSIS

1.1 CONVERGENCE OF SEQUENCES AND SERIES

- We can only talk of *convergence of sequences* in *Hausdorff topological spaces*.
- We can only talk of *series* in *commutative groups*, because we need $+$ to be defined.
- We can only talk of *convergence of series* in *commutative Hausdorff topological groups*.
- We can only talk of *absolute convergence of series* in *normed commutative Hausdorff topological groups*.
- This is from [Wikipedia](#). Let S be the vector space of sequences. Then the partial summation $\Sigma : S \rightarrow S, (a_n) \mapsto \left(\sum_{j=1}^n a_j\right)$ is a *linear operator* on S , whose inverse is the finite difference operator, Δ . These behave as discrete analogs of integration and differentiation, only for series (functions of a natural number) instead of functions of a real variable. For example, the sequence $(1, 1, 1, \dots)$ has series $(1, 2, 3, \dots)$ as its partial summation, which is analogous to the fact that $\int_0^x 1 \, dt = x$.
- Classification of convergence of series
 1. Pointwise or uniform convergence
 2. Absolute, unconditional and conditional convergence
 - *Absolute convergence* means $\sum \|a_n\| < \infty$.
 - *Unconditional convergence* means all rearrangements of the series are convergent to the same value. That is, if $\sigma : \mathbb{N} \rightarrow \mathbb{N}$ is a permutation, then $\sum_n a_n = \sum_n a_{\sigma(n)}$.
In complete spaces, absolute convergence \implies unconditional convergence, but the converse is not true in general. In finite dimensional spaces, the converse is true by Riemann rearrangement theorem. But the Dvoretzky–Rogers theorem asserts that every infinite-dimensional Banach space admits an unconditionally convergent series that is not absolutely convergent. (see this [Wikipedia article](#))
 - *Conditional convergence* means convergent but not absolutely convergent.
 3. Depending on the space of values, for example, real number, arithmetic progression, trigonometric function, etc.

1.2 DIFFERENTIATION

1.2.1 Differentiation of functions with real powers

This idea is by Prof Sundar.

For each $n \in \mathbb{N}$, we can differentiate $f : \mathbb{R} \rightarrow \mathbb{R} : x \mapsto x^n$ using the limit definition of the derivative, by using the factorization $(x+h)^n - x^n = (x+h-x) \sum_{j=0}^{n-1} h^j x^{n-1-j}$, which gives

$$f'(x) = \lim_{h \rightarrow 0} \frac{f(x+h) - f(x)}{h} = \lim_{h \rightarrow 0} \sum_{j=0}^{n-1} (x+h)^j x^{n-1-j} = \sum_{j=0}^{n-1} x^j x^{n-1-j} = (n-1)x^{n-1}.$$

But this does not work for exponents $r \in \mathbb{R}$ in general. How can we do it?

One needs to think outside the box for this. We cannot go by definition here. We note that $x^r = e^{r \log x}$. Now use chain rule.

1.2.2 Example of $f \in C^\infty \setminus C^\omega$

How does one construct an example of a function which is smooth but not analytic? The idea is to find $f \neq 0$ such that $f^{(n)}(0) = 0 \forall n$.

Note that the graph of $x \mapsto x^n$ around $x = 0$ become flatter and flatter as $n \rightarrow \infty$.

Consider $f : \mathbb{R} \rightarrow \mathbb{R} : x \mapsto e^{-\frac{1}{x}} \mathbb{1}_{x>0}(x)$. Then $f'(x) = \frac{1}{x^2} e^{-\frac{1}{x}} \mathbb{1}_{x>0}(x)$ (do the computations separately for $x > 0$ and $x = 0$). In this way, $f'(0) = 0$. Continuing, we see that $f^{(n)}(0) = 0 \forall n$. Therefore, the sum $\sum_n \frac{f^{(n)}(0)}{n!} x^n$ converges to 0, which is not the same as f .

See also [this Wikipedia article](#).

1.2.3 Taylor series for multivariate functionals

Theorem 2.1 (Taylor) *Let $\Omega \subseteq \mathbb{R}^d$ be an open set and $f : \Omega \rightarrow \mathbb{R}$ be an infinitely differentiable function at the point $p \in \Omega$. Then the Taylor series of f around p for $v \in T_p \Omega$ is given by*

$$T(p+v) = \sum_{n=0}^{\infty} \frac{1}{n!} (\langle v, D_v \rangle^n f)(p), \quad \text{where } \langle v, D_v \rangle = \sum_{j=1}^d v_j \frac{\partial}{\partial v_j}.$$

Proof. Is this a \mathbb{R} life, or is this just imaginary, caught in a landslide, no escape from \mathbb{C} . Open your eyes, look up to the sky and see. I'm trivial come, trivial go, little high, little low.

ToDo, really! □

In the differential form, this becomes

$$dT(p) = \sum_{n=1}^{\infty} \frac{1}{n!} \left(\langle dv, D_v \rangle^n f \right) (p)$$

There is another form with multi-indexes. But the above form seems way more natural to me. See more in [this](#) and [this](#) Wikipedia articles.

PART 2

PROBABILITY THEORY

2.1 ELEMENTARY IDEAS

2.1.1 σ -algebras

- $\mathcal{I} \subset \mathcal{B} \subset \overline{\mathcal{B}} = \mathcal{L} \subset 2^{\mathbb{R}}$
- $|\mathcal{I}| = |\mathcal{B}| = |\mathbb{R}| = \aleph_1, |\overline{\mathcal{B}}| = |2^{\mathbb{R}}| = \aleph_2$

2.1.2 Examples

- $\left(x \mapsto \frac{1}{\sqrt{x}} \mathbb{1}_{[0,1]}(x)\right) \in L^1 \setminus L^2$
- $\left(x \mapsto \frac{1}{x} \mathbb{1}_{[1,\infty)}(x)\right) \in L^2 \setminus L^1$

2.1.3 Measurability of $\inf, \sup, \liminf, \limsup$ TODO

Let X_n be a discrete-time stochastic process. Then $\{\inf X_t \geq c\} = \bigcap_t \{X_t \geq c\}$

2.1.4 Method of substitution

See Folland - Real Analysis Theorem (2.43).

2.2 BOREL–CANTELLI LEMMAS

Theorem 2.1 (Borel–Cantelli 1) *Let $(E_n) \subset \mathcal{F}$ such that $\sum \mathbb{P}(E_n) < \infty$. Then $\mathbb{P}(E_n \text{ i.o.}) = 0$.*

Proof. Since $\sum \mathbb{P}(E_n) < \infty$, for any fixed $n \in \mathbb{N}$, we have

$$\mathbb{P}(E_n \text{ i.o.}) = \mathbb{P}\left(\bigcap_n \bigcup_{m \geq n} E_m\right) \leq \mathbb{P}\left(\bigcup_{m \geq n} E_m\right) \leq \sum_{m \geq n} \mathbb{P}(E_m) \rightarrow 0.$$

□

Counterexample of the converse of BC1: Take $((0, 1], \lambda)$ as the probability space, and $E_n = (0, \frac{1}{n^2})$. Then $\mathbb{P}(E_n \text{ i.o.}) = 1$, but $\sum \mathbb{P}(E_n) < \infty$.

Theorem 2.2 (Borel–Cantelli 2) *Let $(E_n) \subset \mathcal{F}$ be (mutually) independent such that $\sum \mathbb{P}(E_n) = \infty$. Then $\mathbb{P}(E_n \text{ i.o.}) = 1$.*

Proof. Since $\mathbb{P}((E_n \text{ i.o.})^c) = \mathbb{P}(E_n^c \text{ ev})$, it is equivalent to prove that $\mathbb{P}(E_n^c \text{ ev}) = 0$. Using independence and the fact that $1 - x < e^{-x}$, for each fixed $n \in \mathbb{N}$, we have

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

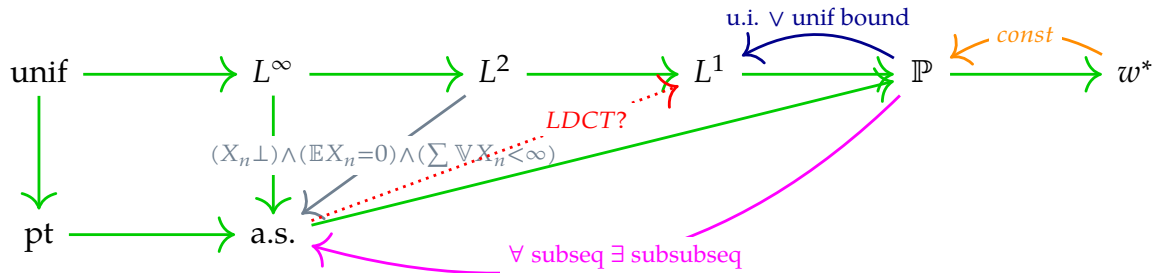
Taking $N \rightarrow \infty$, we get $\mathbb{P}(\bigcap_{m \geq n} E_m^c) \rightarrow 0$. Therefore

$$\mathbb{P}(E_n^c \text{ ev}) = \mathbb{P}\left(\bigcup_n \bigcap_{m \geq n} E_m^c\right) \leq \sum_n \mathbb{P}\left(\bigcap_{m=n}^N E_m^c\right) = 0.$$

□

2.3 MODES OF CONVERGENCE

Study this part from [Robert L Wolpert - Convergence in \$\mathbb{R}^d\$ and in metric spaces](#). In this diagram, the top row represents ‘point independent’ modes of convergence and the bottom row represents the ‘point dependent’ modes of convergence.



Legend

- green: automatic implication
- any other color: depends on the mentioned condition

2.4 CONDITIONING

2.4.1 In L^2 , conditional expectation is a projection

Let $(\Omega, \mathcal{F}, \mathbb{P})$ be a probability space, $\mathcal{G} \subseteq \mathcal{F}$ be a sub- σ -algebra, and $X \in L^2(\mathcal{F})$ be a random variable. Let $\pi_{L^2(\mathcal{G})}$ denote the projection operator onto $L^2(\mathcal{G})$. We know that projection operators are self-adjoint. So $\forall E \in \mathcal{G}$,

$$\begin{aligned}
 \int_E \mathbb{E}(X | \mathcal{G}) d\mathbb{P} &= \int_E X d\mathbb{P} = \int \mathbb{1}_E X d\mathbb{P} = \langle \mathbb{1}_E, X \rangle && \text{[definitions]} \\
 &= \langle \pi_{L^2(\mathcal{G})} \mathbb{1}_E, X \rangle && [E \in \mathcal{G} \implies \mathbb{1}_E \in L^2(\mathcal{G})] \\
 &= \langle \mathbb{1}_E, \pi_{L^2(\mathcal{G})}^* X \rangle = \langle \mathbb{1}_E, \pi_{L^2(\mathcal{G})} X \rangle && \text{[self-adjointness]} \\
 &= \int \mathbb{1}_E \pi_{L^2(\mathcal{G})} X d\mathbb{P} = \int_E \pi_{L^2(\mathcal{G})} X d\mathbb{P}. && \text{[definitions]}
 \end{aligned}$$

Therefore, $\mathbb{E}(X | \mathcal{G}) = \pi_{L^2(\mathcal{G})} X$ a.s.

2.4.2 Uncorrelated does not imply independence

See the wikipedia entries on [uncorrelated random variables](#) and [normally distributed and uncorrelated does not imply independent](#).

2.4.3 $\phi_{aX+bY} = \phi_{aX}\phi_{bY} \forall (a,b) \in \mathbb{R}^2$ implies independence

See [MathSx:1802289](#).

2.5 LIMITING BEHAVIOUR OF \bar{X}_n

A great resource for this chapter is the Wikipedia article on [central limit theorem](#).

Let (X_n) be a sequence of independent and identically distributed random variables with mean m and distribution μ , and let $\bar{X}_n = \frac{1}{n}S_n = \frac{1}{n} \sum_{j=1}^n X_j$.

There are three scales for which we have three theorems, namely

1. law of large numbers ($\bar{X}_n \xrightarrow{a.s.} m$),
2. law of the iterated logarithm ($\overline{\lim} \frac{\bar{X}_n}{\sqrt{\frac{2 \log \log n}{n}}} = 1$ a.s.), and
3. central limit theorem ($\sqrt{n}(\bar{X}_n - m) \xrightarrow{w^*} \mathcal{N}(0, \Sigma)$).

The idea is that we have an asymptotic expansion of \bar{X}_n given (in law) by

$$\bar{X}_n \sim m + \frac{1}{\sqrt{n}} \mathcal{N}(0, \Sigma), \quad \text{where } \Sigma \text{ is the covariance operator.}$$

The convergence to m is given by the law of large numbers, and the convergence to $\frac{1}{\sqrt{n}} \mathcal{N}(0, \Sigma)$ is given by the central limit theorem. As $n \rightarrow \infty$, the dependence on $\mathcal{N}(0, \Sigma)$ goes to zero, so this is consistent with the law of large numbers.

Central limit theorem: how to remember in 1-dim

$$\frac{\bar{X}_n - \mathbb{E}\bar{X}_n}{\sqrt{\mathbb{V}\bar{X}_n}} = \frac{\bar{X}_n - m}{\frac{\sigma}{\sqrt{n}}} \xrightarrow{w^*} \mathcal{N}(0, 1)$$

PART 3

STOCHASTIC ANALYSIS

3.1 CLASSIFICATION OF STOCHASTIC PROCESSES

This is well written in Cosma Rohilla Shalizi - Almost None of the Theory of Stochastic (2010), Chapter 1. Let X be a stochastic process given by

$$\begin{aligned} X : \mathbb{T} \times \Omega &\rightarrow \mathbb{E} \\ \mathcal{F} &\rightarrow \mathcal{X} \\ (t, \omega) &\mapsto X(t, \omega). \end{aligned}$$

The spaces are as follows.

\mathbb{T} The *index set*. Can be finite, discrete (countable) or continuous (uncountable). Can be one-sided, two-sided, spatially distributed, or sets.

$(\mathbb{E}, \mathcal{X})$ The *state space*. Requirements: measurable. Can be finite, discrete or continuous.

$(\Omega, \mathcal{F}, \mathbb{P})$ The *probability space*.

- If $\mathbb{T} = \{1\}$, $\mathbb{E} = \mathbb{R}$, then X is a *random variable*.
- If $\mathbb{T} = \{1, \dots, n\}$, $\mathbb{E} = \mathbb{R}$, then X is a *random vector*.
- If $\mathbb{T} = \{1\}$, $\mathbb{E} = \mathbb{R}^d$, then X is a *random vector*.
- If $\mathbb{T} = \mathbb{N}$, $\mathbb{E} = \mathbb{R}$, then X is a *one-sided random sequence* or *one-sided discrete-time stochastic process*.
- If $\mathbb{T} = \mathbb{Z}$, $\mathbb{E} = \mathbb{R}$, then X is a *two-sided random variable* or *two-sided discrete-time stochastic process*.
- If $\mathbb{T} = \mathbb{Z}^d$, $\mathbb{E} = \mathbb{R}$, then X is a *spatial random variable*.
- If $\mathbb{T} = \mathbb{R}$, $\mathbb{E} = \mathbb{R}$, then X is a *continuous-time random variable*.
- If $\mathbb{T} = \mathcal{B}$, $\mathbb{E} = [0, \infty]$, then X is a *random set function on the reals*.
- If $\mathbb{T} = \mathcal{B} \times \mathbb{N}$, $\mathbb{E} = [0, \infty]$, then X is a *one-sided random sequence of set function on the reals*.
- *Emperical measures*. Let (Z_n) be an i.i.d. random sequence and define $\hat{\mathbb{P}}_n : \mathcal{B} \times \Omega \rightarrow \mathcal{D} : (B, \omega) \mapsto \hat{\mathbb{P}}_n(B, \omega) = \frac{1}{n} \sum_{j=1}^n \mathbb{1}_B(Z_j(\omega))$. Then $\hat{\mathbb{P}}_n$ is a *one-sided random sequence of set function on the reals*, which are in fact *probability measures*. [\mathcal{D} is the space of probability measures on \mathbb{R} .]
- If $\mathbb{T} = \mathcal{B}^d$, $\mathbb{E} = [0, \infty]$, then X is the class of set functions on \mathbb{R}^d . Let \mathcal{M} be the subclass of measures. Then a random set function with realizations in \mathcal{M} is called a *random measure*.
- If $\mathbb{T} = \mathcal{B}^d$, $|\mathbb{E}| < \infty$, then X is a *point process*.
- If $\mathbb{T} = [0, \infty)$, $\mathbb{E} = \mathbb{R}^d < \infty$. A \mathbb{E} -valued random process on \mathbb{T} with paths in $C(\mathbb{T})$ is a *continuous random process*. E.g. Wiener process.

3.2 MARTINGALES

3.2.1 New martingales from old

A stochastic process $A = (A_n)$ is called adapted if $\forall n \in \mathbb{N}, A_n \in L^0(\mathcal{F}_n)$. Let $M = (M_n)$ be a martingale. Then process $\tilde{M} = (\tilde{M}_n)$ defined by $(A \cdot M)_n = \tilde{M}_n = \sum_{j=0}^{n-1} A_j \Delta M_j$, where $\Delta M_j = M_{j+1} - M_j$, is called the *martingale transform* of M by A .

Theorem (martingale transform theorem): \tilde{M} is a martingale.

Proof.

$$\mathbb{E}(\Delta \tilde{M}_n | \mathcal{F}_n) = \mathbb{E}(A_n \Delta M_n | \mathcal{F}_n) = A_n \mathbb{E}(\Delta M_n | \mathcal{F}_n) = 0.$$

Now, let X_n be a stochastic process and τ be a stopping time. Define the stopped process $X_\tau = \sum_{j=0}^{\infty} \mathbb{1}_{\{\tau=j\}} X_j$ when $\mathbb{P}(\tau < \infty) = 1$.

Theorem (stopping time theorem): Let (M_n) be a martingale with respect to (\mathcal{F}_n) . Then $(M_{n \wedge \tau})$ is also a martingale with respect to (\mathcal{F}_n) .

Proof. Without loss of generality, assume $M_0 = 0$, otherwise we can translate by M_0 as $\tilde{M}_n = M_n - M_0$. Now, the *stake process* $A_n = \mathbb{1}_{\{\tau > n\}} = 1 - \mathbb{1}_{\{\tau \leq n\}}$ is adapted to (\mathcal{F}_n) and is bounded by n . Now,

$$\begin{aligned} (A \cdot M)_n &= \sum_{j=0}^{n-1} A_j \Delta M_j \\ &= \sum_{j=0}^{n-1} \Delta M_j - \sum_{j=0}^{n-1} \mathbb{1}_{\{\tau \leq j\}} (M_{j+1} - M_j) \\ &= M_n - M_0 - M_n \mathbb{1}_{\{\tau \leq n\}} + \sum_{j=0}^{n-1} (\mathbb{1}_{\{\tau \leq j\}} M_j - \mathbb{1}_{\{\tau \leq j-1\}} M_j) \\ &= M_n \mathbb{1}_{\{\tau > n\}} + \sum_{j=0}^{n-1} M_j \mathbb{1}_{\{\tau=j\}} \\ &= M_n \mathbb{1}_{\{\tau > n\}} + \sum_{j=0}^{n-1} M_\tau \mathbb{1}_{\{\tau=j\}} \\ &= M_n \mathbb{1}_{\{\tau > n\}} + M_\tau \sum_{j=0}^{n-1} \mathbb{1}_{\{\tau=j\}} \\ &= M_n \mathbb{1}_{\{\tau > n\}} + M_\tau \mathbb{1}_{\{\tau \leq n\}} \\ &= M_{n \wedge \tau}. \end{aligned}$$

Therefore, $(M_{n \wedge \tau})$ is a martingale transform of (M_n) . Since (A_n) is bounded and adapted, by the martingale transform theorem, $(M_{n \wedge \tau})$ is a martingale.

3.3 MARKOV PROCESSES

Equivalent definitions

Let $s \in [0, t]$. Then X is a Markov process if any of the following are true:

- $\forall E \in \mathcal{F}, \mathbb{P}(X_t \in E | \mathcal{F}_s) = \mathbb{P}(X_t \in E | X_s)$, or
- $\forall E \in \mathcal{F}, \forall f \in L^0 \cap \mathcal{B}, \mathbb{E}(f(X_t) | \mathcal{F}_s) = \mathbb{E}(f(X_t) | X_s)$.

Martingale vs Markov

See djilil.chafai.net and [MathSx:763645](https://math.stackexchange.com/questions/763645).

3.4 ITÔ CALCULUS

Notation:

In what follows, $T = [0, \infty)$, \mathcal{A} means adapted, \mathcal{B} means bounded, \mathcal{C} means continuous, and $\|\cdot\|$ denotes the L^2 -norm.

Let $(\Omega, \mathcal{F}, \mathbb{F} = (\mathcal{F}_t), \mathbb{P})$ be a filtered probability space, $W : T \times \Omega \rightarrow \mathbb{C}$ be a \mathbb{F} -adapted Wiener martingale, and $X : T \times \Omega \rightarrow \mathbb{C}$ be a stochastic process.

3.4.1 Step 1: $X \in \mathcal{A} \cap \mathcal{S}$ a.s.

Let $X(t, \omega) = \sum_{j \geq 0} \xi_j(\omega) \mathbb{1}_{[t_j, t_{j+1})}(t)$, where $\xi_j \in L^0(\mathcal{F}_{t_j})$.

3.4.2 Step 2: $X \in \mathcal{A} \cap \mathcal{B} \cap \mathcal{C}$ a.s.

Define $X_n(t, \omega) = X\left(\frac{\lfloor nt \rfloor}{n}, \omega\right)$, $n \in \mathbb{N}$. Note that $\forall n, X_n \in \mathcal{A} \cap \mathcal{S}$, and since $X \in \mathcal{C}$, $|X_n(t, \omega) - X(t, \omega)| \rightarrow 0$ (pointwise convergence) (t, ω) -a.s. Then $\forall \varepsilon > 0$, there exists a sufficiently large $n \in \mathbb{N}$ such that $|X_n(t, \omega) - X(t, \omega)| < \varepsilon < \infty$ (bounded) (t, ω) -a.s., so $|X_n(t, \omega) - X(t, \omega)|^2 < \varepsilon^2 < \infty$ (t, ω) -a.s. Therefore, by the bounded convergence theorem, $\|X_n - X\| \rightarrow 0$.

Therefore, (X_n) is Cauchy in $L^2(T \times \Omega)$, that is, $\|X_n - X_m\| \rightarrow 0$. Now, by linearity and Itô isometry for the Itô integral for simple processes, $\|\mathcal{I}(X_n) - \mathcal{I}(X_m)\| = \|\mathcal{I}(X_n - X_m)\| = \|X_n - X_m\| \rightarrow 0$. Therefore, for $t \in T$ fixed, $(\mathcal{I}(X_n))$ is Cauchy in $L^2(\Omega)$. Since $L^2(\Omega)$ is complete, the sequence converges. Denote the limit by $\mathcal{I}(X)$, that is, $\|\mathcal{I}(X_n) - \mathcal{I}(X)\| \rightarrow 0$.

3.4.3 Step 3: $X \in \mathcal{A} \cap \mathcal{B} \cap L^0(T \times \Omega)$

3.4.4 Step 4: $X \in \mathcal{A} \cap L^2(T \times \Omega)$

3.4.5 Step 5: $X \in \mathcal{A} \cap \left\{ X \in \mathbb{C}^{T \times \Omega} : \forall t \geq 0, \int_0^t X(s, \cdot) ds < \infty \right\}$ a.s.

3.4.6 Properties of the Itô integral

In what follows, assume the following. Let $X, Y \in \mathcal{A} \cap L^2(T \times \Omega)$; $(X_n), (Y_n) \subset \mathcal{A} \cap \mathcal{S}$ such that $\|X_n - X\| \rightarrow 0$ and $\|Y_n - Y\| \rightarrow 0$. Let $z \in \mathbb{C}$.

3.4.6.1 Linearity: $\|z\mathcal{I}(X) + \mathcal{I}(Y) - \mathcal{I}(zX + Y)\| = 0$

First, note that $\|(zX_n + Y_n) - (zX + Y)\| \leq |z|\|X_n - X\| + \|Y_n - Y\| \rightarrow 0$. Now, by the linearity of the integral $\mathcal{I} : \mathcal{A} \cap \mathcal{S} \rightarrow L^2(\Omega)$, we have

$$\begin{aligned} & \|z\mathcal{I}(X) + \mathcal{I}(Y) - \mathcal{I}(zX + Y)\| \\ &= \|z\mathcal{I}(X) + \mathcal{I}(Y) - z\mathcal{I}(X_n) - \mathcal{I}(Y_n) + \mathcal{I}(zX_n + Y_n) - \mathcal{I}(zX + Y)\| \\ &\leq |z|\|\mathcal{I}(X) - \mathcal{I}(X_n)\| + \|\mathcal{I}(Y) - \mathcal{I}(Y_n)\| + \|\mathcal{I}(zX_n + Y_n) - \mathcal{I}(zX + Y)\| \rightarrow 0. \end{aligned}$$

3.4.6.2 Itô isometry: $\|\mathcal{I}(X)\| = \|X\|$

Using the isometry of the integral $\mathcal{I} : \mathcal{A} \cap \mathcal{S} \rightarrow L^2(\Omega)$, we have

$$\begin{aligned} \|\mathcal{I}(X)\| &\leq \|\mathcal{I}(X) - \mathcal{I}(X_n)\| + \|\mathcal{I}(X_n)\| \\ &= \|\mathcal{I}(X) - \mathcal{I}(X_n)\| + \|X_n\| \\ &= \|\mathcal{I}(X) - \mathcal{I}(X_n)\| + \|X_n - X\| + \|X\| \rightarrow \|X\|. \end{aligned}$$

Note that the ‘Itô isometry’ is actually a unitary transformation.

3.4.6.3 Martingale property: $\mathbb{E}(\mathcal{I}_t(X) | \mathcal{F}_s) = \mathcal{I}_s(X)$ a.s.

The martingale property of the integral $\mathcal{I} : \mathcal{A} \cap \mathcal{S} \rightarrow L^2(\Omega)$ gives $\mathbb{E}(\mathcal{I}_t(X_n) - \mathcal{I}_s(X_n) | \mathcal{F}_s) = 0$. Using this and the unitariness of the Itô isometry, we get

$$\begin{aligned} & \|\mathbb{E}(\mathcal{I}_t(X) - \mathcal{I}_s(X) | \mathcal{F}_s)\|^2 \\ &= \mathbb{E} \left| \mathbb{E}(\mathcal{I}_t(X) - \mathcal{I}_t(X_n) + \mathcal{I}_s(X_n) - \mathcal{I}_s(X) | \mathcal{F}_s) + \mathbb{E}(\mathcal{I}_t(X_n) - \mathcal{I}_s(X_n) | \mathcal{F}_s) \right|^2 \\ &= \mathbb{E} \left| \mathbb{E}(\mathcal{I}_t(X) - \mathcal{I}_t(X_n) + \mathcal{I}_s(X_n) - \mathcal{I}_s(X) | \mathcal{F}_s) + 0 \right|^2 \\ &\leq \mathbb{E} \mathbb{E} \left(\left| \mathcal{I}_t(X) - \mathcal{I}_t(X_n) + \mathcal{I}_s(X_n) - \mathcal{I}_s(X) \right|^2 | \mathcal{F}_s \right) \\ &= \mathbb{E} \left| \mathcal{I}_t(X) - \mathcal{I}_t(X_n) + \mathcal{I}_s(X_n) - \mathcal{I}_s(X) \right|^2 \\ &= \|\mathcal{I}_t(X - X_n) + \mathcal{I}_s(X_n - X)\|^2 \\ &\leq 2 \left(\|\mathcal{I}_t(X - X_n)\|^2 + \|\mathcal{I}_s(X_n - X)\|^2 \right) \quad [\|a + b\|^2 \leq 2(\|a\|^2 + \|b\|^2)] \\ &= 2 \left(\|X - X_n\|^2 + \|X_n - X\|^2 \right) = 4\|X_n - X\|^2 \rightarrow 0. \end{aligned}$$

3.4.7 Itô formula for multidimensional processes

Recall the discussion of Taylor series for multivariate functionals in section 1.2.3. This part is from (<SundarKallianpur2014>), § 5.3, 5.4 and 5.6.

A (local, continuous) semimartingale is a process X_t that can be written as $X_t = X_0 + M_t + A_t$, where

1. M_t is a mean-zero (local, continuous) martingale, and
2. A_t is an right-continuous adapted process of locally bounded variation.

This is equivalently represented in the differential form as $dX_t = dM_t + dA_t$.

Let X_t be a d -dimensional semimartingale, and let $Y_t = f(X_t)$, where $f \in C^2(\mathbb{R})$. Then

$$dY_t = df(X_t) = f'(X_t)dA_t + f'(X_t)dM_t + \frac{1}{2}f''(X_t)d\langle M_t, M_t \rangle,$$

where we use the rule $d\langle B^{(j)}, B^{(j)} \rangle_t = (dB_t^{(j)})^2 = dt$, everything else being 0.

Alternatively, as from (<Kuo2006>), we have the following.

Let $X_t = (X_t^{(1)}, \dots, X_t^{(d)})$ be a d -dimensional process and $f(t, x)$ be a functional of (t, X_t) .

Then the Itô formula becomes

$$df(t, X_t) = (\langle dt, D_t \rangle f)(t, X_t) + (\langle dx, D_x \rangle f)(t, X_t) + \frac{1}{2}(\langle dx, D_x \rangle^2 f)(t, X_t),$$

or in short, $df = (\langle dt, D_t \rangle + \langle dx, D_x \rangle + \frac{1}{2}\langle dx, D_x \rangle^2) f$.

3.4.8 Adapted and instantly independent implies deterministic

Let X_\cdot be both adapted and instantly independent. Then for any fixed t , we have $X_t = \mathbb{E}(X_t | \mathcal{F}_t) = \mathbb{E}X_t$. Therefore, X_t is constant w.r.t. ω for each t . Therefore X_\cdot must be deterministic.

3.5 EXAMPLES

3.5.1 Find $\mathbb{E} \left(\int_0^1 B_t^2 dt \right)^2$.

By Fubini's theorem,

$$\mathbb{E} \left(\int_0^1 B_t^2 dt \right)^2 = \mathbb{E} \left(\int_0^1 B_t^2 dt \int_0^1 B_s^2 ds \right) = \mathbb{E} \left(\int_0^1 \int_0^1 B_t^2 B_s^2 ds dt \right) = \int_0^1 \int_0^1 \mathbb{E}(B_t^2 B_s^2) ds dt.$$

$$\begin{aligned} \text{Now, } \forall s \in [0, t], \quad \mathbb{E}(B_t^2 B_s^2) &= \mathbb{E}(\mathbb{E}(B_t^2 B_s^2 | \mathcal{F}_s)) = \mathbb{E}(B_s^2 \mathbb{E}((B_t^2 - t) + t | \mathcal{F}_s)) \\ &= \mathbb{E}(B_s^2 ((B_s^2 - s) + t)) = \mathbb{E}(B_s^2 ((B_s^2 - s) + t)) \\ &= \mathbb{E}(B_s^4 - s B_s^2 + t B_s^2) = 3s^2 - s^2 + ts = 2s^2 + ts. \end{aligned}$$

$$\text{So} \quad \mathbb{E} \left(\int_0^1 B_t^2 dt \right)^2 = 2 \int_0^1 \int_0^t (2s^2 + ts) ds dt = \frac{7}{9}.$$

3.5.2 Find $\mathbb{V} \left(\int_0^1 t^2 B_t dt \right)$.

By Fubini's theorem, $\mathbb{E} \left(\int_0^1 t^2 B_t dt \right) = \int_0^1 t^2 \mathbb{E} B_t dt = 0$. So by Fubini's theorem (again),

$$\begin{aligned} \mathbb{V} \left(\int_0^1 t^2 B_t dt \right) &= \mathbb{E} \left(\int_0^1 t^2 B_t dt \right)^2 = \mathbb{E} \left(\int_0^1 t^2 B_t dt \int_0^1 s^2 B_s ds \right) \\ &= \mathbb{E} \left(\int_0^1 \int_0^1 t^2 s^2 B_t B_s ds dt \right) = \int_0^1 \int_0^1 t^2 s^2 \mathbb{E}(B_t B_s) ds dt \\ &= \int_0^1 \int_0^1 t^2 s^2 (t \wedge s) ds dt = 2 \int_0^1 \int_0^t t^2 s^2 ds dt = \frac{1}{14}. \end{aligned}$$

3.5.3 Calculate $\int_0^T e^{B_t^2} dB_t$.

Let $f(x) = \int_0^x e^{t^2} dt$. Then $f'(x) = e^{x^2}$ and $f''(x) = 2xe^{x^2}$.

Now, using the Itô formula, we get $d \left(\int_0^x e^{t^2} dt \right) = e^{B_t^2} dB_t + B_t e^{B_t^2} dt$, which gives us $\int_0^T e^{B_t^2} dB_t = \int_0^T e^{B_t^2} dt - \int_0^T B_t e^{B_t^2} dt$.

PART 4

INFINITE-DIMENSIONAL ANALYSIS

4.1 GIUSEPPE DA PRATO

Proposition 1.1 (Proposition 1.2 in the book) Let $a \in \mathbb{R}$, $\lambda > 0$, and $\mu = N_{a,Q}$. Then

- i. $\int_{\mathbb{R}} x N_{a,\lambda}(dx) = a$,
- ii. $\int_{\mathbb{R}} (x - a)^2 N_{a,\lambda}(dx) = \lambda$, and
- iii. $\widehat{N_{a,\lambda}}(h) := \int_{\mathbb{R}} e^{ihx} N_{a,\lambda}(dx) = e^{iah - \frac{1}{2}\lambda h^2}$, $h \in \mathbb{R}$.

Proof.

- i. $\int_{\mathbb{R}} x N_{a,\lambda}(dx) = a$,
- ii. $\int_{\mathbb{R}} (x - a)^2 N_{a,\lambda}(dx) = \lambda$, and
- iii. $\widehat{N_{a,\lambda}}(h) := \int_{\mathbb{R}} e^{ihx} N_{a,\lambda}(dx) = e^{iah - \frac{1}{2}\lambda h^2}$, $h \in \mathbb{R}$.

□

Proposition 1.2 (Proposition 1.3 in the book) Let $H \simeq \mathbb{R}^d$, $a \in H$, $Q \in L_+(H)$, and $\mu = N_{a,Q}$. Then

- i. $\int_H x N_{a,Q}(dx) = a$,
- ii. $\int_H \langle y, x - a \rangle \langle z, x - a \rangle N_{a,Q}(dx) = \langle Qy, z \rangle$, and
- iii. $\widehat{N_{a,Q}}(h) := \int_H e^{i\langle h, x \rangle} N_{a,Q}(dx) = e^{i\langle a, h \rangle - \frac{1}{2}\langle Qh, h \rangle}$, $h \in H$.

Proof. All indices vary from 1 to d .

- i. $\int_H x N_{a,Q}(dx) = a$.

$$\begin{aligned}
 \int_H x N_{a,Q}(dx) &= \int_H x \times_j N_{a_j, \lambda_j} \left(\times_i dx_i \right) \\
 &= \int_H \sum_k (x_k e_k) \prod_j N_{a_j, \lambda_j}(dx_j) && \text{[Fubini]} \\
 &= \sum_k \left(\int_{\mathbb{R}} \cdots \int_{\mathbb{R}} x_k \prod_j N_{a_j, \lambda_j}(dx_j) \right) e_k && \text{[Fubini]} \\
 &= \sum_k \int_{\mathbb{R}} x_k N_{a_k, \lambda_k}(dx_k) e_k && \left[\int_{\mathbb{R}} N_{a_j, \lambda_j}(dx_j) = 1 \ \forall j \right] \\
 &= \sum_k a_k e_k = a. && \text{[Proposition 1.1(i)]}
 \end{aligned}$$

- ii. $\int_H \langle y, x - a \rangle \langle z, x - a \rangle N_{a,Q}(dx) = \langle Qy, z \rangle$.

$$\begin{aligned}
& \int_H \langle y, x - a \rangle \langle z, x - a \rangle N_{a,Q}(\mathrm{d}x) \\
&= \int_H \sum_k y_k (x_k - a_k) \sum_l z_l (x_l - a_l) \prod_j N_{a_j, \lambda_j}(\mathrm{d}x_j) \\
&= \sum_k \sum_l y_k z_l \int_H (x_k - a_k) (x_l - a_l) \prod_j N_{a_j, \lambda_j}(\mathrm{d}x_j) \\
&= \sum_k y_k z_k \int_H (x_k - a_k)^2 \prod_j N_{a_j, \lambda_j}(\mathrm{d}x_j) + \sum_{i \neq j} y_k z_l \int_H (x_k - a_k) (x_l - a_l) \prod_j N_{a_j, \lambda_j}(\mathrm{d}x_j) \\
&= \sum_k y_k z_k \int_{\mathbb{R}} (x_k - a_k)^2 N_{a_k, \lambda_k}(\mathrm{d}x_k) + \sum_{i \neq j} y_k z_l \int_{\mathbb{R}} (x_k - a_k) N_{a_k, \lambda_k}(\mathrm{d}x_k) \int_{\mathbb{R}} (x_l - a_l) N_{a_l, \lambda_l}(\mathrm{d}x_l) \\
&= \sum_k y_k z_k \lambda_k = \langle Qy, z \rangle. \quad [\text{Proposition 1.1(ii)}]
\end{aligned}$$

$$\text{iii. } \widehat{N_{a,Q}}(h) := \int_H e^{i\langle h, x \rangle} N_{a,Q}(\mathrm{d}x) = e^{i\langle a, h \rangle - \frac{1}{2} \langle Qh, h \rangle}, h \in H.$$

$$\begin{aligned}
\widehat{N_{a,Q}}(h) &= \int_H e^{i\langle h, x \rangle} N_{a,Q}(\mathrm{d}x) \\
&= \int_H e^{i \sum_k h_k x_k} \prod_j N_{a_j, \lambda_j}(\mathrm{d}x_j) \\
&= \int_H \prod_k e^{ih_k x_k} \prod_j N_{a_j, \lambda_j}(\mathrm{d}x_j) \\
&= \prod_k \int_H e^{ih_k x_k} N_{a_k, \lambda_k}(\mathrm{d}x_k) \quad [\text{Why? This is false in general!}] \\
&= \prod_k e^{ia_k h_k - \frac{1}{2} \lambda_k h_k^2} \quad [\text{Proposition 1.1(iii)}] \\
&= e^{i \sum_k a_k h_k - \frac{1}{2} \sum_k \lambda_k h_k^2} \\
&= e^{i\langle a, h \rangle - \frac{1}{2} \langle Qh, h \rangle}.
\end{aligned}$$

□

Lemma 1.3 (The operator $1 - \varepsilon Q$ in Page 14) *The operator $1 - \varepsilon Q$ is invertible with a finite positive determinant. Moreover, the inverse is bounded.*

Proof. Firstly, note that for the operator

Since $\varepsilon < \frac{1}{\lambda_1}$, we have $1 > \varepsilon \lambda_1$. Combining this with $\lambda_1 \geq \lambda_2 \geq \dots$, we get $1 > \varepsilon \lambda_1 \geq \varepsilon \lambda_2 \geq \dots$, which further implies $\infty > (1 - \varepsilon \lambda_1)^{-1} \geq (1 - \varepsilon \lambda_2)^{-1} \geq \dots$. Therefore the operator $1 - \varepsilon Q$ is invertible and $(1 - \varepsilon Q)^{-1}$ is bounded.

Now, since $(1 - \varepsilon Q)e_k = (1 - \varepsilon \lambda_k)e_k$ for every $k \in \mathbb{N}$, we have $(1 - \varepsilon Q)^{-1}e_k = (1 - \varepsilon \lambda_k)^{-1}e_k$ for every $k \in \mathbb{N}$. This gives us for every $x \in H$,

$$(1 - \varepsilon Q)^{-1}x = \sum_{k=1}^{\infty} \frac{1}{1 - \varepsilon \lambda_k} \langle x, e_k \rangle e_k.$$

Why is the infinite product finite and positive? □

Lemma 1.4 (Lemma for Proposition 1.13)

$$\int_{\mathbb{R}} e^{\frac{\varepsilon}{2}|x|^2} N_{a,\lambda}(dx) = \frac{e^{\frac{\varepsilon a^2}{2(1-\varepsilon\lambda)}}}{\sqrt{1-\varepsilon\lambda}}.$$

Proof. First, note that using completion of squares, we get

$$\begin{aligned} \frac{\varepsilon}{2}x^2 - \frac{1}{2\lambda}(x-a)^2 &= -\frac{1}{2\lambda}((1-\varepsilon\lambda)x^2 - 2ax + a^2) \\ &= -\frac{1-\varepsilon\lambda}{2\lambda} \left[x^2 - 2\frac{a}{1-\varepsilon\lambda}x + \frac{a^2}{(1-\varepsilon\lambda)^2} - \frac{a^2}{(1-\varepsilon\lambda)^2} + \frac{a^2}{1-\varepsilon\lambda} \right] \\ &= -\frac{1-\varepsilon\lambda}{2\lambda} \left[\left(x - \frac{a}{1-\varepsilon\lambda} \right)^2 - \frac{\varepsilon\lambda a^2}{(1-\varepsilon\lambda)^2} \right] \\ &= -\frac{1}{2} \frac{\left(x - \frac{a}{1-\varepsilon\lambda} \right)^2}{\frac{\lambda}{1-\varepsilon\lambda}} + \frac{\varepsilon a^2}{2(1-\varepsilon\lambda)}. \end{aligned}$$

Therefore,

$$\begin{aligned} \int_{\mathbb{R}} e^{\frac{\varepsilon}{2}x^2} N_{a,\lambda}(dx) &= \int_{\mathbb{R}} e^{\frac{\varepsilon}{2}x^2} \frac{1}{\sqrt{2\pi\lambda}} e^{-\frac{1}{2\lambda}(x-a)^2} dx \\ &= \frac{e^{\frac{\varepsilon a^2}{2(1-\varepsilon\lambda)}}}{\sqrt{1-\varepsilon\lambda}} \int_{\mathbb{R}} \frac{1}{\sqrt{2\pi\frac{\lambda}{1-\varepsilon\lambda}}} e^{-\frac{1}{2} \frac{\left(x - \frac{a}{1-\varepsilon\lambda} \right)^2}{\frac{\lambda}{1-\varepsilon\lambda}}} dx \\ &= \frac{e^{\frac{\varepsilon a^2}{2(1-\varepsilon\lambda)}}}{\sqrt{1-\varepsilon\lambda}}. \end{aligned}$$

□

Proposition 1.5 (Hint for Exercise 1.14)

$$J_m = 2^m F^{(m)}(0), \quad m \in \mathbb{N}; \quad \text{where} \quad F(\varepsilon) = \int_H e^{\frac{\varepsilon}{2}|x|^2} \mu(dx), \quad \varepsilon > 0.$$

Proof.

$$\begin{aligned}
F(\varepsilon) &= \int_H e^{\frac{\varepsilon}{2}|x|^2} \mu(dx) \\
&= \int_H \sum_{m=0}^{\infty} \frac{1}{m!} \left(\frac{\varepsilon}{2}|x|^2 \right)^m \mu(dx) \\
&= \sum_{m=0}^{\infty} \frac{\varepsilon^m}{2^m m!} \int_H |x|^{2m} \mu(dx) \quad [\text{Monotone convergence theorem}] \\
&= \sum_{m=0}^{\infty} \frac{J_m}{2^m m!} \varepsilon^m
\end{aligned}$$

□

4.2 ABSTRACT WIENER SPACES

4.3 WHITE NOISE DISTRIBUTION THEORY

4.3.1 Characterization theorem

Importance and history.

In the following, F is defined on $S_{\mathbb{C}}$, and $F(z\tilde{\xi} + \eta)$ is entire $\forall z \in \mathbb{C}$.

$$(S)_{\beta} \subset (S) \subset (L^2) \subset (S)^* \subset (S)_{\beta}^*$$

$$S \subset L^2 \subset S'$$

BIBLIOGRAPHY