BROWNIAN MOTION AND RELATED PROCESSES

Karlin & Taylor, A First Course in Stochastic Processes, ch 7, 15

Brownian Motion: Definitions

Brownian motion can be defined and constructed in many ways. Some of these include:
1. A stochastic process $X_t$ with independent, normally distributed increments $X_t - X_s \sim N(0, t-s)$, continuous paths, and initial value $X_0 = 0$;
2. A Gaussian stochastic process with mean $E X_t = 0$, covariance $EX_sX_t = \min(s,t)$, and continuous paths;
3. A Markov process $X_t$ with initial value $X_0 = 0$, transition probability

$$P[X_t \leq r \mid X_s = x] = \int_{-\infty}^{r} e^{-(y-x)^2/2(t-s)} \frac{dy}{\sqrt{2\pi(t-s)}}$$

and continuous paths;
4.* A process $X_t$ with (a.) independent and (b.) stationary increments (i.e., the random variables $[X_{t_i} - X_{t_{i-1}}]$ are independent and have distributions depending only on $(t_i - t_{i-1})$), with continuous paths.
5. A martingale $X_t$ such that $X_t^2 - t$ is also a martingale, with initial value $X_0 = 0$ and continuous paths.

The only tricky part of constructing $X_t$ is getting continuous paths; it’s pretty easy to get a process with the right joint distribution for all times $t$. Note that definitions (4.) and (5.) don’t even mention the normal distribution; that follows from the other requirements as a consequence of the Central Limit Theorem.

Here’s one construction, for $0 \leq t \leq 1$. The idea is to construct a sequence of piecewise-linear approximations $X^{(n)}_t$ to $X_t$, with exactly the correct distribution on all dyadic rationals of degree $n$ (those of the form $\frac{i}{2^n}$). The key computation about Brownian motion is that for any $0 \leq a \leq b < c < \infty$, the conditional distribution of $X_b$ given $X_a$ and $X_c$ is normal with mean $\mu_b = \frac{(c-b)X_a + (b-a)X_c}{c-a}$ (the linear interpolate) and variance $\sigma_b^2 = \frac{(c-b)(b-a)}{c-a}$; for $a = \frac{i}{2^n}, c = \frac{i+1}{2^n}$, and $b = \frac{i+a}{2^n} = \frac{2i+a}{2^n}$, we have $\mu_b = \frac{1}{2}[X_a + X_c]$ and $\sigma_b^2 = (\frac{1}{2})^{2+n}$.

Let $z_i$ be an iid sequence of $N(0,1)$ random variables and for each $n$ define random variables $x_i^n, 1 \leq i \leq 2^n$, by:

Even $i$: $x_0 = 0 \quad x_{2i}^{n+1} = x_i^n \quad 0 \leq i \leq 2^n$
Odd $i$: $x_1 = z_1 \quad x_{2i-1}^{n+1} = \frac{1}{2}[x_{2i-2}^n + x_i^n] + (\frac{1}{2})^{1+n/2}z_{2^n+i} \quad 1 \leq i \leq 2^n$

Now define a sequence of processes $X^{(n)}_t$ by linearly interpolating the $X^{(n)}_{t/2^n} = x_i^n$'s:

$$X^{(n)}_t = (i-t2^n)x_{i-1}^n + (1-i+t2^n)x_i^n. \quad \frac{i-1}{2^n} < t \leq \frac{i}{2^n}$$

By construction $X^{(n)}_t$ is a Gaussian process with continuous paths, initial value zero, and the right probability distribution at each $n^{th}$-order dyadic rational; it remains to show that the $X^{(n)}_t$ converge uniformly a.s. and that the limit is Brownian motion. We’ll turn to that next lecture.

* With this definition the process will have mean $EX_t = t\mu$ and variance $E(X_t - t\mu)^2 = t\sigma^2$ for some constants $\mu, \sigma^2$; the rescaled process $[X_t - t\mu]/\sigma$ has the usual normalization, $\mu = 0$ and $\sigma^2 = 1$. 
Actually, constructing Brownian motion is in some sense very easy—if \( Y_t \) is any square-integrable mean-zero stochastic process starting at zero with independent increments, the function \( \sigma_s^2 = \mathbb{E}[Y_s^2] \) must be increasing since, for \( 0 \leq s \leq t \),
\[
\sigma_t^2 = \mathbb{E}
\left( X_s + (X_t - X_s) \right)^2 \]
\[
= \sigma_s^2 + \mathbb{E}(X_t - X_s)^2 \geq \sigma_s^2.
\]
If \( \sigma_s^2 \to \infty \) as \( s \to \infty \), then for each \( n \geq 1 \) and \( t \geq 0 \) we can set
\[
s_n(t) = \inf\{s : \sigma_s^2 \geq nt\}
\]
and define
\[
X_t^{(n)} = \frac{1}{\sqrt{n}} Y_{s_n(t)};
\]
for every \( n \), \( X_t^{(n)} \) has independent increments with mean zero and approximately the right covariance (exactly the right covariance if \( \sigma_s^2 \) is strictly increasing). In the limit as \( n \to \infty \), the covariance becomes exactly correct and moreover the Central Limit Theorem applies: for large \( n \) each \( s_n(t) \) becomes large, and for \( s < t \) the increment \( [Y_{s_n(t)} - Y_{s_n(s)}] \) can be thought of as the sum of very many small and independent increments. Thus the time-change (*) makes almost any independent-increment process converge to Brownian motion, and in particular we can construct Brownian Motion as a limit of random walks, Markov chains, or Poisson processes. For the simple symmetric random walk starting at zero, \( s_n = \lceil nt \rceil \) and
\[
X_t = n^{-1/2} Y_{\lfloor nt \rfloor}
\]
converges to Brownian Motion.

**Continuous Paths**

Most things we might want to compute about any random variable \( X \) defined on some probability space \((\Omega, \mathcal{F}, \mathbb{P})\) don’t really depend on \((\Omega, \mathcal{F}, \mathbb{P})\) at all, but only on the probability distribution, the induced measure \( \mu_X = \mathbb{P} \circ X^{-1} \) on the real line \((\mathbb{R}, \mathcal{B})\). The random variable \( \xi(\omega) = \omega \) on the probability space \((\mathbb{R}, \mathcal{B}, \mu_X)\) has the same probability distribution as \( X \), and so we can usually study features of \( X \) without worrying about \((\Omega, \mathcal{F}, \mathbb{P})\) by using this “canonical probability space” \((\mathbb{R}, \mathcal{B}, \mu_X)\). If we have not one but several random variables \( X_1, X_2, \ldots, X_n \), the same idea works in \( n \)-dimensional space: if \( \mu_X \) denotes the joint probability distribution, the canonical space is \((\mathbb{R}^n, \mathcal{B}^n, \mu_X)\) on which the random variables \( \xi_i(\omega) = \omega_i \) have the same joint distribution as the \( X_i \).

What about infinitely many random variables, especially the uncountable infinity of random variables \( X_t \) for Brownian Motion?

Consider the set \( \Omega \) of continuous real-valued functions on the unit interval starting at zero:
\[
\Omega = \{ \text{Continuous } \omega : [0,1] \to \mathbb{R}, \omega_0 = 0 \}.
\]
The supremum gives a natural notion of distance from one \( \omega \) to another in \( \Omega \), leading to the topological notion of open sets; let \( \mathcal{F} \) be the smallest \( \sigma \)-algebra (or Borel Field, BF) containing these open sets, i.e., containing for each \( \omega_0 \in \Omega \) and \( \epsilon > 0 \) the set
\[
[\omega \in \Omega : \sup_{0 \leq s \leq 1} |\omega(s) - \omega_0(s)| \leq \epsilon].
\]
(Don’t worry if this seems obscure). The distribution of Brownian Motion is just the probability measure \( P \) on \((\Omega, \mathcal{F})\) such that \( \xi_t(\omega) = \omega(t) \) is a Brownian Motion on \((\Omega, \mathcal{F}, \mathbb{P})\).

One way to construct \( P \), and with it Brownian Motion, is to look at the distribution \( P_n \) induced by the process \( X^{(n)} \) defined earlier; if we can show that these measures converge, we can define \( P \) to be their limit and verify that it has the right properties.
PATH CONTINUITY AND NON-DIFFERENTIABILITY

Richard Durrett, Brownian Motion and Martingales in Analysis, pp 1–7

Introduction

Last time we defined Brownian Motion in five ways, including

I. A stochastic process $X_t$ with initial value $X_0 = 0$ with
   A. Stationary independent increments $[X_t - X_s]$;
   B. Normally distributed increments $[X_t - X_s] \sim N(0, t-s)$;
   C. Continuous paths, almost surely.

Are these consistent? If a process $X_t$ has independent increments, can they also have the specified Gaussian distributions? If they do, can the process also have continuous paths? If so, is path continuity a consequence of a. and b.? As we will see, the answers are Yes, Yes, and No, respectively. To clarify the issues let’s consider other SII processes satisfying a. above; three possibilities are:

- **Brownian Motion** $X_t$: $P[X_t - X_s \in A] = \int_A e^{-x^2/2(t-s)} \frac{dx}{\sqrt{2\pi(t-s)}}$
- **Cauchy Process** $C_t$: $P[C_t - C_s \in A] = \int_A \frac{(t-s)dx}{\pi[(t-s)^2 + x^2]}$
- **Poisson Process** $N_t$: $P[N_t - N_s \in A] = \sum_{x \in A} e^{-(t-s)} \frac{(t-s)^x}{x!}$

All three distributions are possible for SII processes; to see this it is only necessary to check that for $t_1 < t_2 < t_3$, the indicated distribution for $[X_{t_3} - X_{t_1}]$ is the same as that of the sum of independent random variables with the distributions indicated for $[X_{t_2} - X_{t_1}]$ and $[X_{t_3} - X_{t_2}]$. It turns out that this is equivalent to requiring that the characteristic function $E[e^{i\alpha X_t}]$ be of the form $e^{-t\phi(a)}$; for this three distributions the characteristic functions are indeed of that form with $\phi(a) = a^2/2$, $|a|$, and $[1 - e^{-a}]$, respectively.

All three distributions are also almost-surely continuous at every point, in the sense that for every fixed $t$, $P[X_t = \lim_{s \to t} X_s] = P[C_t = \lim_{s \to t} C_s] = P[N_t = \lim_{s \to t} N_s] = 1$. Note that the Poisson process is constant except for jumps of size one, and so its paths are not continuous—they are continuous at each fixed $t$, because the jump times have continuous distributions, but any interval of length $L$ will contain at least one jump with probability $1 - e^{-L}$ and so the path will not be a.s. continuous on that interval. There is no way to construct a Poisson process with continuous paths; it turns out that there is no way to construct a Cauchy process with continuous paths, either. What about Brownian Motion?

We have constructed Brownian Motion already on the dyadic rationals $Q_2$ from an IID sequence $z_k$ of $N(0, 1)$ random variables by setting $X_0 = 0$ and $X_1 = z_1$ and, recursively, defining $X_t$ for $t = \frac{2^{i-1} + 2^i}{2^i}$ by $X_t = \frac{1}{2}[X_{(i-1)/2^i} + X_{i/2^i} + z_{i+2^i}/\sqrt{2^n}]$. Can we extend the definition to all $0 \leq t \leq 1$ by continuity, i.e., set $X_t \equiv \lim_{s \to t} X_s$?

Any continuous function $f(x)$ is uniformly continuous when restricted to a compact set like $[0, 1]$, and any uniformly continuous function $g(x)$ defined on a set $D$ can be extended to a uniformly continuous function on the closure $\overline{D}$, but in general a function that is merely continuous on a set $D$ cannot be extended to be continuous on $\overline{D}$. Pick an irrational $\omega \in (0, 1)$ (perhaps $\frac{1}{25}$) and think about the function $g(x) = 1_{[\omega,1]}(x)$ defined on the dyadic rationals $x \in Q_2$; $g(x)$ is continuous at every rational point, since $|g(x) - g(y)| < \epsilon$ whenever $|y - x| < \delta_x = |x - \omega|$, but is not uniformly continuous since no single $\delta$ will work for all $x$.

To extend $X_s$ continuously to $\overline{Q_2} = [0, 1]$ we must show that $X_s$ is almost surely uniformly continuous on $Q_2$, i.e., that for $a.e.$ $\omega$, there exists $\delta_\omega$ such that $0 \leq s < t \leq 1, (t - s) < \delta_\omega \Rightarrow |X_t(\omega) - X_s(\omega)| < \epsilon$. Note that this is not true for the Poisson process, despite the almost-sure continuity
at each point. The argument for Brownian Motion hinges on the Borel-Cantelli lemma and the routine calculation for normally-distributed random variables \( X \sim N(0, \sigma^2) \) and real numbers \( p > -1 \),

\[
E[|X|^p] = \int_{-\infty}^{\infty} |x|^p e^{-x^2/2\sigma^2} \frac{dx}{\sqrt{2\pi\sigma^2}}
= \frac{2}{\sqrt{\pi}} (2\sigma^2)^{p/2} \int_0^{\infty} \left( \frac{x^2}{2\sigma^2} \right)^{p/2} e^{-x^2/2\sigma^2} \frac{dx}{\sqrt{2\sigma^2}}
= \frac{2}{\sqrt{\pi}} (2\sigma^2)^{p/2} \int_0^{\infty} \left( \frac{x^2}{2\sigma^2} \right)^{(p-1)/2} e^{-x^2/2\sigma^2} x dx
= \frac{\Gamma\left(\frac{p+1}{2}\right)}{\sqrt{\pi}} (2\sigma^2)^{p/2} = c_p(\sigma^2)^{p/2}
\]

and so, for any \( \gamma > 0 \) and \( \delta > 0 \),

\[
P\left[|X_{j/2^n} - X_{i/2^n}| > \left( \frac{j-i}{2^n} \right)^{\gamma} \right. \text{ for some } 0 \leq i < j \leq 2^n, (j-i) \leq 2^{\delta n} \]
\[
\leq \sum_{i=0}^{2^n-1} \sum_{j=i+1}^{i+2^n} P\left[|X_{j/2^n} - X_{i/2^n}| > \left( \frac{j-i}{2^n} \right)^{\gamma} \right] \quad \text{(by subadditivity)}
\leq \sum_{i=0}^{2^n-1} \sum_{j=i+1}^{i+2^n} E[|X_{j/2^n} - X_{i/2^n}|^p] \quad \text{(by Chebychev)}
= \sum_{i=0}^{2^n-1} \sum_{j=i+1}^{i+2^n} c_p \left( \frac{j-i}{2^n} \right)^{p/2}
= \sum_{i=0}^{2^n-1} \sum_{j=i+1}^{i+2^n} c_p \left( \frac{j-i}{2^n} \right)^{p(1/2-\gamma)}
\leq c_p 2^{2n\gamma \delta} \left( \frac{2\delta}{2^n} \right)^{p(1/2-\gamma)} = c_p 2^{-n\epsilon}
\]

where \( \epsilon = -(1+\delta) + (1-\delta)p(1/2 - \gamma) \). For \( \gamma < 1/2 \) and \( \delta < 1 \) we can insure \( \epsilon > 0 \) by taking \( p > \frac{1+\delta}{(1/2-\gamma)(1-\delta)} \). By the Borel-Cantelli lemma, for a.e. \( \omega \exists N_\omega \forall n \geq N_\omega \forall q = i/2^n, r = j/2^n \) s.t. \(|q-r| < 2^{-n(1-\delta)}, |X_q - X_r| \leq (q-r)^\gamma \). It follows (see Durrett) that there exists a number \( c_\omega \) such that \( \forall q, r \in \mathbb{Q}_2 \cap [0,1] \),

\[
|X_q - X_r| \leq c_\omega (q-r)^\gamma
\]
i.e., that the restriction of \( X_t \) to \( \mathbb{Q}_2 \) is a.s. uniformly Hölder continuous of index \( \gamma \) for every \( \gamma < 1/2 \). In fact this is about the best we can do: \( X_t \) is a.s. not Hölder continuous of index \( 1/2 \) at any point \( t \), and in particular is not differentiable at any point \( t \). In fact, one of the Laws of the Iterated Logarithm gives

\[
1 = \limsup_{s \to t} \frac{\pm (X_t - X_s)}{\sqrt{2(t-s) \log \log 1/(t-s)}} \quad \text{so} \quad \limsup_{s \to t} \frac{|X_t - X_s|}{\sqrt{t-s}} = +\infty \text{ a.s.}
\]
Now pick any $c \in \mathbb{R}$, $c \neq 0$, and $h > 0$ and define four processes $X_k(t)$ from $X_t = X(t)$ as follows:

III. $X_1(t) = cX(t/c^2)$;
IV. $X_2(t) = tX(1/t)$ for $t > 0$, $X_2(0) = 0$;
V. $X_3(t) = X(t + h) - X(h)$;
VI. $X_4(t) = (t + 1)X(1) - X(1)$.

It is straightforward to verify that each of these is a Brownian motion satisfying $a$, $b$, $c$ above; by the way, $X_4(t)$ only depends on $X(s)$ for $0 < s \leq 1$, and yields a Brownian motion for $0 \leq t < \infty$ from our earlier construction of Brownian motion only for $0 \leq s \leq 1$.

It turns out that Property 3. above is true, not only for fixed $h > 0$, but also for random $\tau = \tau(\omega) > 0$ provided $\tau$ is a Markov time (a.k.a. stopping time); in particular, it holds for first hitting times $\tau_a = \inf\{s > 0 : X_s = a\}$. Since a Brownian Motion has probability $1/2$ of being positive at any time $t > 0$, it follows that for any time $t > 0$ and level $a \geq 0$,

$$P([X_t > a]) = P([X_t > a] \cap [\tau_a \leq t])$$
$$= P([\tau_a \leq t])P([X_t > a] | [\tau_a \leq t])$$
$$= P([\tau_a \leq t]) \left( \frac{1}{2} \right)$$

so, turning things around,

$$P([\tau_a \leq t]) = 2P([X_t > a])$$
$$= 2\Phi\left( \frac{a}{\sqrt{t}} \right)$$

Starting at $X_0 = 0$, let’s find the probability that $X_t = 0$ for any $t \in [t_0, t_1]$. One way to make this precise is to think about the Markov time $\tau = \inf\{t \geq t_0 : X_t = 0\}$ and calculate $P[\tau \leq t_1]$. If we condition on the value $a$ of $X_{t_0}$, this is just the probability that, in time $t_1 - t_0$, the Brownian motion $X_{t+t_0} - X_{t_0}$ ever reaches the value $|a|$ in time $[t_1 - t_0]$; we just calculated that this is $P[\tau \leq t_1 | X_{t_0} = a] = 2\Phi\left( \frac{|a|}{\sqrt{t_1 - t_0}} \right)$. Thus the desired probability is

$$P[\tau \leq t_1] = E[2\Phi\left( \frac{|X_{t_0}|}{\sqrt{t_1 - t_0}} \right)]$$
$$= \int 2\Phi\left( \frac{|z|}{\sqrt{t_1 - t_0}} \right) e^{-z^2/2t_0} \frac{dz}{\sqrt{2\pi t_0}}$$
$$= \frac{2}{\pi} \arccos\left( \frac{t_0}{t_1} \right) \quad \text{(see text, p.348)}.$$
use calculations for Brownian motion as approximations for other processes. For example, Lévy showed that for any $t > 0$, the Lebesgue measure of the set of times $s \leq t$ at which $X_s$ is positive exactly satisfies, for $0 \leq \theta \leq 1$, the relation

$$P\left[\frac{\lambda[s \leq t \mid X_s > 0]}{t} \leq \theta\right] = \frac{2}{\pi} \arcsin(\theta);$$

Kakutani showed that the fraction of $k \leq n$ for which $S_k > 0$ has approximately that same distribution, for any sum of i.i.d. rv’s with zero mean and finite variance.

**Processes Related to Brownian Motion**

1. **Brownian Motion with Drift.**

Let $X_t$ be a Brownian motion and let $x_0 \in \mathbb{R}$, $\mu \in \mathbb{R}$, and $\sigma^2 > 0$ be arbitrary; the process

$$X_1(t) = x_0 + \mu t + \sigma X_t$$

is called *Brownian motion with drift*. It has stationary independent increments (with the normal $N(\mu(t-s), \sigma^2(t-s))$ distribution) and continuous paths starting at $X_1(0) = x_0$.

2. **Geometric Brownian Motion.**

Let $X_1(t)$ be a Brownian motion with drift and set

$$X_2(t) = e^{X_1(t)} = e^{x_0 + \mu t + \sigma X_t}.$$

This is called *Geometric Brownian motion*, and is useful in modeling positive quantities whose fractional change is independent over different periods; it is often used in the mathematical theory of finance, and in modeling reservoir levels and related phenomena.

3. **Reflected Brownian Motion.**

Let $X(t)$ be a Brownian motion and set

$$X_3(t) = |X(t)|$$

This is called *Reflected Brownian motion*. The process is positive and Markovian: in fact, for any $x > 0$ and $y > 0$ and $t > s > 0$,

$$P[X_3(t) \leq y \mid X_3(s) = x] = P[|X(t)| \leq y \cap |X(s) = x| \mid X_3(s) = x]$$

$$+ P[|X(t)| \leq y \cap |X(s) = -x| \mid X_3(s) = x]$$

$$= \frac{1}{2}P[|X(t)| \leq y \mid X(s) = x] + \frac{1}{2}P[|X(t)| \leq y \mid X(s) = -x]$$

$$= \Phi\left(\frac{y - x}{\sqrt{t - s}}\right) - \Phi\left(\frac{-y - x}{\sqrt{t - s}}\right)$$

$$= \int_0^y p_{t-s}(z|x) \, dz$$

where the conditional pdf is given by differentiation as

$$p_u(y|x) = \frac{1}{\sqrt{2\pi u}} [e^{-(y-x)^2/2u} + e^{-(x+y)^2/2u}].$$

Let $X(s)$ be a Brownian motion, let $\lambda > 0$, and set

$$X_4(t) = e^{-\lambda t}X(e^{2\lambda t}), \quad -\infty < t < \infty.$$ 

This is called the Ornstein-Uhlenbeck Velocity Process (or the “OU process” for short). It is both Markovian and stationary, in that the joint distribution of any $k$-tuple $[X_4(t_1 + s), ..., X_4(t_k + s)]$ does not depend on $s$; for a Gaussian process like $X_4$ this follows from the fact that the mean $E[X_4(t)] \equiv 0$ is constant and the covariance

$$E[X_4(s)X_4(t)] = e^{-\lambda|t-s|}$$

depends only on the difference $|t - s|$. This process was originally conceived as the derivative $Y'(t)$ of a Gaussian process $Y(t) = Y_0 + \int_0^t X_4(s)ds$ intended to be an alternative model for the motion of pollen particles that, unlike Brownian motion, has a meaningful velocity and hence finite kinetic energy. The velocity process $X_4(t)$ has proven more useful than the position process $Y(t)$.

THE BROWNIAN BRIDGE

Introduction

Let $X$ be a Brownian Motion process and consider two processes defined as follows for $0 \leq t \leq 1$:

$$X_1(t) = X(t) - tX(1), \quad X_2(t) = (1-t)X\left(\frac{t}{1-t}\right).$$

Obviously each of these is a mean-zero Gaussian process, since $X$ is; the finite-dimensional distributions will be determined completely once we identify the covariance functions

$$E[X_1(s)X_1(t)] = E[(X(s) - sX(1))(X(t) - tX(1))]$$

$$= E[X(s)X(t) - sX(1)X(t) - X(s)tX(1) + stX(1)X(1)]$$

$$= [(s \land t) - s(1 \land t) - t(s \land 1) + st(1 \land 1)]$$

$$= s \land t - st$$

$$E[X_2(s)X_2(t)] = (1-s)(1-t)E[X\left(\frac{s}{1-s}\right)X\left(\frac{t}{1-t}\right)]$$

$$= (1-s)(1-t)\left[\left(\frac{s}{1-s}\right) \land \left(\frac{t}{1-t}\right)\right]$$

$$= (1-s)(1-t)\left(\frac{s}{1-s}\right)$$

if, say, $s \leq t$

$$= s - st$$

$$= s \land t - st$$

for any $s, t$.

Thus both processes have continuous sample paths and mean-zero Normal finite-dimensional distributions with covariance $s \land t - st$; such a process is called a Brownian Bridge, or sometimes pinned Brownian motion. It can also be thought of as a Brownian motion conditioned on the event $X(1) = 0$. It arises (as we’ll see below) in nonparametric statistical problems, and it can be used in constructing Brownian motion and related processes. From the second definition it is clear that $X_2$ is a Markov process, but it does not have independent increments and it is not a martingale: $E[X_2(t) - X_2(s) | \mathcal{F}_s] = -\frac{t-s}{1-s}X_2(s)$. 
The Kolmogorov-Smirnov Statistic

Let \( X_i \) be independent and identically distributed from some unknown distribution \( \mu_X \) with distribution function \( F(t) = P[X_i \leq t] = \mu_X((-\infty, t]) \). If called upon to guess \( F(t) \) from observations of \( X_i \) we would no doubt consider the empirical distribution function

\[
F_n(t) = \frac{\#\{i \leq n : X_i \leq t\}}{n} = \sum_{i=1}^{n} 1_{(-\infty, t]}(X_i),
\]

a random function of \( t \) that starts at \( F_n(-\infty) = 0 \) and jumps by \( \frac{1}{n} \) at each observation \( X_i \). Kolmogorov and Smirnov studied the probability distribution of the quantity

\[
Y_n = \sup_{-\infty < s < \infty} \sqrt{n}|F_n(s) - F(s)|,
\]

the (normalized) largest deviation of the empirical distribution function from the true distribution function; it turns out that \( Y_n \) has the same probability distribution for any continuous distribution \( F(t) \), and in particular is the same as that for uniformly distributed random variables with \( F(t) = t \). What is the limiting distribution, as \( n \to \infty \)?

Regarded as a stochastic process, \( F_n \) has mean and covariance functions

\[
E[F_n(t)] = \frac{1}{n} \sum_{i=1}^{n} E1_{(-\infty, t]}(X_i)
= \frac{1}{n} \sum_{i=1}^{n} P[X_i \leq t]
= F(t)
\]

\[
E\left[\left(F_n(s) - F(s)\right)\left(F_n(t) - F(t)\right)\right] = E\left[\left(\frac{1}{n} \sum_{i=1}^{n} \left(1_{(-\infty, s]}(X_i) - F(s)\right)\right)\left(\frac{1}{n} \sum_{j=1}^{n} \left(1_{(-\infty, t]}(X_j) - F(t)\right)\right)\right]
= n^{-2} \sum_{i=1}^{n} E\left[\left(1_{(-\infty, s]}(X_i) - F(s)\right)\left(1_{(-\infty, t]}(X_i) - F(t)\right)\right]
= n^{-1} E\left[\left(1_{(-\infty, s]}(X_1) - F(s)\right)\left(1_{(-\infty, t]}(X_1) - F(t)\right)\right]
= n^{-1} E\left[1_{(-\infty, s\wedge t]}(X_1) - F(s)1_{(-\infty, t]}(X_1) - 1_{(-\infty, s]}(X_1)F(t) + F(s)F(t)\right]
= n^{-1} \left[F(s \wedge t) - F(s)F(t)\right]
= n^{-1} \left[F(s) \wedge F(t) - F(s)F(t)\right]
\]

Thus \( \sqrt{n}[F_n(t) - F(t)] \) has the same covariance function as \( X_1(F(t)) \) for a Brownian Bridge \( X_1(s) \); by the Central Limit Theorem, the finite-dimensional distributions of \( \sqrt{n}[F_n(t) - F(t)] \) converge weakly to the Normal distribution as \( n \to \infty \).

It would be nice to have something stronger— to be able to assert that any continuous functional of \( \sqrt{n}[F_n(t) - F(t)] \) converges weakly to a similar functional of the Brownian bridge, and in particular that the Kolmogorov-Smirnov statistic \( Y_n \) converges to \( Y = \sup_{0 \leq t \leq 1} |X_1(t)| \) in distribution. For this we need to develop the concept of the distribution of a stochastic process, and study weak convergence of these distributions.
Distributions

The distribution of an $\mathbb{R}^1$-valued random variable $X$ on some probability space $(\Omega, \mathcal{F}, P)$ is just the induced measure $\mu_X(B) = P[X \in B] = P \circ X^{-1}$ on the Borel sets $\mathcal{B}$ of the real line; for example, $X$ has the $N(\mu, \sigma^2)$ distribution if $\mu_X(B) = \frac{1}{\sqrt{2\pi\sigma^2}} \int_B e^{-(x-\mu)^2/2\sigma^2} \, dx$ and the Poisson distribution with mean $\lambda$ if $\mu_X(B) = \sum [e^{-\lambda x}/x! : x \in B]$.

Similarly the (joint) distribution of $n$ random variables $X_1, \ldots, X_n$ is just the occupation measure $\mu_X$ of the vector $X \in \mathbb{R}^n$, $\mu_X(B) = P[X \in B]$ on the Borel sets $\mathcal{B}^n$ in $\mathbb{R}^n$. But what about stochastic processes, where $n = \infty$? What is the distribution of Brownian motion, or of the Brownian bridge, or of the Poisson process or reflected Brownian motion?

Path spaces

If a real-valued RV takes values in $\mathbb{R}$ and a random vector in $\mathbb{R}^n$, then a real-valued stochastic process $X_t$ defined for $t \in \mathcal{T} = [0, 1]$ must take values in some set of paths $\Omega = \{\omega : \mathcal{T} \rightarrow \mathbb{R}\}$, and the distribution of $X$ must be a probability measure $\mu_X$ on some Borel Field $\mathcal{F}$ of subsets of $\Omega$. The simplest path space to consider is the set of all functions $\Omega_1 = \{\omega : \mathcal{T} \rightarrow \mathbb{R}\}$, and the cylinder sets $\mathcal{F}_1$ generated by the evaluation functionals—i.e., the smallest $\mathcal{B}$ containing sets of the form $[\omega : \omega(t) \in B]$ for each $t \in \mathcal{T}$ and Borel set $B$.

This works, after a fashion: any consistent set of finite-dimensional distributions does determine a unique measure $\mu_X$ on $\mathcal{F}_1$, and the process $X : \Omega \times \mathcal{T} \rightarrow \mathbb{R}$ defined by $X(\omega, t) = \omega_t$ does have the right probability distribution at each time $t$. Unfortunately some important sets of paths $E$ are missing from $\mathcal{F}_1$, making it impossible to calculate $\mu_X[E]$; for example, $[\omega : t \mapsto \omega_t$ is continuous] is not an event and even $[\omega : t \mapsto \omega_t$ is Lebesgue measurable] is non-measurable. We can evaluate $\omega_t$ at fixed times $t$, but the quantity $Y(\omega) = \sup_{0 \leq t \leq 1} |\omega|$ is not a random variable (it’s not $\mathcal{F}_1$-measurable) and so we can’t calculate its probability distribution. Next time we’ll look at some alternative path spaces.

Continuous Paths

For Brownian Motion and its relatives, the problem is solved by using the probability space of continuous functions $\Omega_2 = \mathcal{C} = \{\omega : \mathcal{T} \rightarrow \mathbb{R}, \ t \mapsto \omega_t$ is continuous\}]. This is a metric space in the supremum norm

$$\delta(\omega, \omega') = \sup_{0 \leq t \leq 1} |\omega_t - \omega'_t|$$

and so has a Borel $\mathcal{B} = \mathcal{F}_2$ generated by sets of the form $[\omega : \delta(\omega, \omega') < \epsilon]$ for $\epsilon > 0$ and $\omega' \in \mathcal{C}$. By the distribution of a path-continuous stochastic process $X$ we will mean the measure $\mu_X$ induced on $(\mathcal{C}, \mathcal{B})$.

PATH SPACES

If a real-valued RV takes values in $\mathbb{R}$ and a random vector in $\mathbb{R}^n$, then a real-valued stochastic process $X_t$ defined for $t \in \mathcal{T} = [0, 1]$ must take values in some set of paths $\Omega = \{\omega : \mathcal{T} \rightarrow \mathbb{R}\}$, and the distribution of $X$ must be a probability measure $\mu_X$ on some Borel Field $\mathcal{F}$ of subsets of $\Omega$. Three possible path spaces to consider are:

1. $\Omega_1 = \{\omega : \mathcal{T} \rightarrow \mathbb{R}\}$ The set of all functions $\mathcal{T} \rightarrow \mathbb{R}$
2. $\mathcal{F}_1 = \sigma[X_t^{-1}(B)]$ the cylinder sets generated by the evaluation functionals;
3. $\mathcal{F}_2 = \mathcal{B}(\mathcal{C}(\mathcal{T} : \mathbb{R}))$ the Borel sets generated by $[\omega' \in \mathcal{C} : \sup_{0 \leq s \leq 1} |\omega_s - \omega'_s| < \epsilon]$
4. $\mathcal{F}_3 = \mathcal{B}(\mathcal{D}(\mathcal{T} : \mathbb{R}))$ the Borel sets generated by Skorohod neighborhoods in $\mathcal{D}$.
The simplest one to use is \( \Omega_1 \). This works, after a fashion: any consistent set of finite dimensional distributions does determine a unique probability measure \( \mu_X \) on \( \mathcal{F}_1 \), and the process \( X : \Omega \times \mathcal{T} \rightarrow \mathbb{R} \) defined by \( X(\omega, t) = \omega_t \) does have the right probability distribution at each time \( t \). Unfortunately some important sets of paths \( E \) are missing from \( \mathcal{F}_1 \), making it impossible to calculate \( \mu_X[E] \); for example, \([ \omega : t \mapsto \omega_t \text{ is continuous}] \) is not an event and even \([ \omega : t \mapsto \omega_t \text{ is Lebesgue measurable}] \) is non-measurable. We can evaluate \( \omega_t \) at fixed times \( t \), but the quantity \( Y(\omega) = \sup_{0 \leq t \leq \varepsilon} |\omega_t| \) is not a random variable (it’s not \( \mathcal{F}_1 \)-measurable) and so we can’t calculate its probability distribution.

**Continuous Paths**

For Brownian Motion and its relatives, the problem is solved by using the probability space of continuous functions \( \Omega_2 = \mathcal{C} = [\omega : \mathcal{T} \rightarrow \mathbb{R}, t \mapsto \omega_t \text{ is continuous}] \). This is a metric space in the supremum norm

\[
\delta(\omega, \omega') = \sup_{0 \leq t \leq \varepsilon} |\omega_t - \omega'_t|
\]

and so has a Borel BF \( \mathcal{B} = \mathcal{F}_2 \) generated by sets of the form \([ \omega : \delta(\omega, \omega') < \varepsilon] \) for \( \varepsilon > 0 \) and \( \omega' \in \mathcal{C} \). By the distribution of a path-continuous stochastic process \( X \) we will mean the measure \( \mu_X \) induced on \((\mathcal{C}, \mathcal{B})\). The space \( \Omega_2 \) is suitable for processes with discontinuous paths including the Poisson process, generalized Poisson process, birth/death processes, Markov chains, the Cauchy process, and others.

**Tightness and Weak Convergence**

Any infinite sequence \( \alpha_n \subset [0, 1] \) has a limit point \( \alpha_\infty \) in \([0, 1] \), and a subsequence \( \alpha_{n_k} \rightarrow \alpha_\infty \); the proof is the so-called diagonal argument. Start with \( i = 0 \), \([a_0, b_0] = [0, 1] \), and \( n_{ij} = j \); note that \([a_i, b_i] \) contains all of the infinite sequence \( n_{ij} \). For each \( i \) let \([a_{i+1}, b_{i+1}] \) be \([a_i, \frac{a_i + b_i}{2}] \) if that contains infinitely-many of the \( n_{ij} \), and otherwise let \([a_{i+1}, b_{i+1}] \) be \([\frac{a_i + b_i}{2}, b_i] \); let \( n_{i-1,j} \) be the subsequence of \( n_{ij} \) which lie in \([a_{i+1}, b_{i+1}] \). Now the diagonal sequence \( n_{ii} \) must lie in each \([a_j, b_j] \) for \( i \geq j \), and so must be a Cauchy sequence converging to the limit \( \alpha_\infty = \cap [a_i, b_i] \) in \([0, 1] \).

In \( \mathbb{R}^n \) any closed and bounded set \( K \) has the property that every infinite sequence \( \alpha_n \subset K \) has a limit point \( \alpha_\infty \in K \); such a set \( K \) is said to be (sequentially) compact. A set \( A \) like \([0, 1] \) whose closure is compact is sometimes called precompact or conditionally compact; every infinite sequence \( \alpha_n \subset A \) has a limit point \( \alpha_\infty \), but it is possible that \( \alpha_\infty \notin A \). In \( \mathbb{R}^n \) every bounded set is precompact, but in other metric spaces simple boundedness may not be enough; for example the functions \( f_n(x) = \sin(nx) \) are all elements of the space \( \mathcal{C} = \mathcal{C}_b(\mathcal{T}) \) of continuous bounded functions on \( \mathcal{T} = [0, 1] \) are all bounded by 1, but no subsequence converges uniformly on \( \mathcal{T} = [0, 1] \). The Arzela-Ascoli theorem asserts that a set \( A \subset \mathcal{C} \) of continuous functions is precompact if and only if the elements \( \omega \in A \) are uniformly bounded and equicontinuous, i.e., if and only if:

A. For some \( B < \infty \), \(|\omega(0)| < B \) for all \( \omega \in A \);

B. For all \( \epsilon > 0 \) there is a \( \delta > 0 \) such that \( \forall \omega \in A, \forall s, t \in \mathcal{T}, |s - t| < \delta \Rightarrow |\omega(s) - \omega(t)| < \epsilon \).

Now let \( \mu_n \) be a sequence of probability measures on the Borel sets \( \mathcal{B} \) of \([0, 1] \); the numbers \( \alpha_n = \mu_n([0, 1/2]) \) all lie in \([0, 1] \), so along some subsequence \( n_{1i} \) the numbers \( \alpha_{n_{1i}} \) converge. Along a further subsequence \( n_{2i} \) the numbers \( \mu_n([0, 1/4]) \) and \( \mu_n([1/2, 3/4]) \) also converge; along subsequence subsequences \( n_{ki} \) we can insure that \( \mu_n(A) \) converges for each interval \( A = (a, b) \). Finally, along the diagonal sequence \( \mu_{n_{1i}}(A) \) converges for every interval with dyadic-rational endpoints. Is the limit \( \mu_\infty \) a probability measure?

The surprising answer is, maybe not. Think about a sequence \( \mu_n \) of measures each giving probability one to the single point \( 2^{-n} \), the limit ought to give probability one to the limit point 0, but \( 0 \notin (0, 1] \)— and in fact the limit is \( \mu(A) = 0 \) for all \( A \subset (0, 1] \). This is the only thing that can go wrong, however:
Theorem (Prohorov). Let \( \mu_n \) be a sequence of probability measures on the Borel sets \( \mathcal{F} \) of a complete separable metric space \( \Omega \). Then some subsequence \( \mu_{n_k} \) converges weakly to a subprobability measure \( \mu_\infty \) on \( \mathcal{F} \) satisfying \( 0 \leq \mu_\infty(\Omega) \leq 1 \). If for each \( \varepsilon > 0 \) there is a compact set \( K_\varepsilon \subset \Omega \) satisfying \( \mu_n(K_\varepsilon) \geq 1 - \varepsilon \) for every \( n \), the sequence \( \mu_n \) is said to be tight and necessarily \( \mu_\infty(\Omega) = 1 \). If every convergent subsequence converges to the same limit point \( \mu_\infty \), then the entire sequence converges.

Theorem. A family \( P_n \) of probability measures on \((\mathcal{C}, \mathcal{B})\) is tight if and only if

C. For each \( \eta > 0 \) there is a \( B < \infty \) such that \( \forall n, P_n[\omega : |\omega(0)| > B] < \eta \);

D. For all \( \varepsilon > 0 \) and \( \eta > 0 \) there is a \( \delta > 0 \) such that \( \forall n, \)

\[
P_n[\omega : \sup_{|s-t| < \delta} |\omega(s) - \omega(t)| > \varepsilon] < \eta.
\]

Corollary (Kolmogorov). A family \( P_n \) of probability measures on \((\mathcal{C}, \mathcal{B})\) is tight if there exist numbers \( \alpha > 0, \beta > 0, B < \infty \), and \( C < \infty \) such that \( \forall n, \)

C. \( E_n|\omega(0)|^\beta \leq B; \)

D. \( E_n|\omega(s) - \omega(t)|^\beta \leq C|t - s|^{1+\alpha}. \)

Continuous Stochastic Processes

For each finite set \( J \subset \mathcal{T} \) let \( \mu_J \) be a probability measure on \(|J|\)-dimensional Euclidean space \( \mathbb{R}^J \) such that, for \( J \subset J' \), the measure \( \mu_J \) is the marginal for \( \mu_J' \); call such a collection of measures a “consistent finite dimensional distribution.” For example, if \( m(t) \) is any function on \( \mathcal{T} \) and \( \gamma(s, t) \) is a (positive definite) covariance function, i.e., satisfies \( \sum_{i,j \leq n} z_i \bar{z}_j \gamma(t_i, t_j) > 0 \) for every integer \( n \), complex \( z_1, \ldots, z_n \), and times \( t_i \in \mathcal{T} \), then we can construct a unique consistent finite dimensional distribution such that for each \( s, t \in \mathcal{T} \), \( \mu_{\{s, t\}} \) is bivariate Normal with mean vector and covariance matrix

\[
\begin{pmatrix}
m(s) \\
m(t)
\end{pmatrix}
\begin{pmatrix}
\gamma(s, s) & \gamma(s, t) \\
\gamma(t, s) & \gamma(t, t)
\end{pmatrix}.
\]

Obviously any measure \( \mu \) on \((\Omega_1, \mathcal{F}_1)\) or on \((\mathcal{C}, \mathcal{B})\) induces a consistent family of finite dimensional distributions. Any consistent finite dimensional distribution induces a unique measure on \((\Omega_1, \mathcal{F}_1)\), but it’s harder to induce a measure on \((\mathcal{C}, \mathcal{B})\); the Poisson distributions won’t work, for example, because Poisson sample-paths aren’t continuous. By Kolmogorov’s Corollary above,

Theorem. Let \( \{\mu_J\} \) be a consistent family of finite dimensional distributions. If there exist positive constants \( \alpha > 0, \beta > 0, \) and \( C > 0 \) such that \( \mathbb{E}|X(s) - X(t)|^\beta \leq C|t - s|^{1+\alpha} \), then \( \{\mu_J\} \) induces a unique probability measure on \((\mathcal{C}, \mathcal{B})\) and, moreover,

\[
[\omega \in \mathcal{C} : t \mapsto \omega_t \text{ is H"older continuous of index } \frac{\alpha}{\beta} - \varepsilon]
\]

has \( \mu \)-measure one for each \( \varepsilon > 0 \).

In particular, a Gaussian distribution satisfies this condition (with \( \beta = 2\alpha \)) if \( m(t) \) is H"older continuous of index \( \frac{1}{2} \) and if \( |\gamma(t, s)| \leq C'|t - s| \); this condition is satisfied by Brownian motion (with or without drift) and the Brownian bridge.
THE BROWNIAN BRIDGE REVISITED

Introduction

Last time we presented Kolmogorov’s Theorem, a corollary to a theorem of Prohorov:

**Theorem (Kolmogorov).** A family $\mu_n$ of probability measures on $(\mathcal{C}, \mathcal{B})$ is tight if there exist numbers $\alpha > 0$, $\beta > 0$, $B < \infty$, and $C < \infty$ such that $\forall n$,

1. $\int |\omega(0)|^\beta \, d\mu_n \leq B$;
2. $\int |\omega(s) - \omega(t)|^\beta \, d\mu_n \leq C |t - s|^{1+\alpha}$. In this case any limit point $\mu$ of $\{\mu_n\}$ is a probability measure giving probability one to the set of Hölder continuous functions of index $\alpha/\beta$.

Today we will use this theorem to present another construction of the Brownian Bridge; the method is quite general, and is an important tool in constructing and studying stochastic processes.

Let $z_n$ be an iid sequence of standard $N(0, 1)$ random variables and define $x_0^n = x_0 = 0$. For $n \geq 0$ and $0 \leq i < 2^n$ define

$$x_{2i}^{n+1} = x_i^n$$
$$x_{2i+1}^{n+1} = \frac{1}{2}(x_i^n + x_{i+1}^n + 2^{-n/2} z_{2^n+i})$$

$0 \leq i < 2^n$

It’s easy to verify that, with this specification, the processes $X_i^n = x_i^n + (2^n t - i)(x_i^n - x_i^0)$ for $\frac{i}{2^n} \leq t < \frac{i+1}{2^n}$ have the Brownian Bridge covariance $\Gamma_{st} = (s \wedge t) - st$ for $s, t \in 2^n \mathbb{N}$, inducing the Gaussian measures

$$\mu_n(B) = P[X^n \in B]$$
on the Borel sets $B \in \mathcal{B}$ on $\mathcal{C}$.

For $s, t$ not dyadic rationals in $2^n \mathbb{N}$, the covariance of $X^n$ (or $\mu_n$) may not quite be $\Gamma_{st}$; a tedious but straightforward calculation from the definitions shows that $\mathbb{E}X_s^n X_t^n = \Gamma_{st} = (s \wedge t) - st$ if the integer parts of $2^s$ and $2^n t$ differ ($\lfloor 2^s \rfloor \neq \lfloor 2^n t \rfloor$), for example, if $|s - t| > 2^{1-n}$, while if $\lfloor 2^s \rfloor = \lfloor 2^n t \rfloor = j$, $\mathbb{E}X_s^n X_t^n = (s \wedge t) - st - 2^{-n}(1 - (2^n t - j))(2^n s - j)$ differs from $\Gamma_{st}$ by no more than $2^{-n}$. It follows that $\mathbb{E}[(X_{t+\epsilon} - X_t)^2] = \epsilon(1 - \epsilon)(1 - 2^{-n})$; in particular, $\mathbb{E}[|X_s - X_t|^\beta] \leq c_\beta |t - s|^{3/2}$ for every $\beta > 0$ and Kolmogorov’s criteria are satisfied.

Let $\mu$ be any limit point of the family $\{\mu_n\}$; for any $0 \leq s \leq t \leq 1$ the function $\phi(\omega) = e^{i\omega_s + ib\omega_t}$ is continuous and bounded on $\mathcal{C}$, so $\int_{\mathcal{C}} \phi(\omega) \, d\mu_n$ converges to $\int_{\mathcal{C}} \phi(\omega) \, d\mu$; this is just the joint characteristic function of $X_s$ and $X_t$, which have the bivariate Normal distribution with mean $\binom{0}{0}$ and covariance $\binom{s}{s}$, so

$$\int_{\mathcal{C}} e^{i\omega_s + ib\omega_t} \, d\mu_n \to \int_{\mathcal{C}} e^{i\omega_s + ib\omega_t} \, d\mu$$
$$= e^{-\frac{1}{2}[(a^2 s(1-s) + 2abs(1-t) + b^2 t(1-t)]}$$

and under $\mu$, $X_t(\omega) = \omega_t$ is a stochastic process with

C. continuous-paths;
D. normal distribution;
E. mean zero;
F. covariance $\Gamma_{st} = (s \wedge t) - st$. 

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This uniquely determines \( \mu \) as the Brownian Bridge distribution; since \( \{ \mu_n \} \) has a unique limit point, necessarily \( \mu_n \Rightarrow \mu \). We have succeeded in constructing the Brownian Bridge. By the way, we can now construct a Brownian Motion (or Wiener process) for all time \( 0 \leq t < \infty \) by the formula

\[
W_t = (1 + t)X\left(\frac{t}{1 + t}\right).
\]

The technique used in this construction of the Brownian Bridge is quite powerful; the only features we used of the Brownian bridge were:

G. It has continuous-paths (otherwise, use Skorohod space \( D \));
H. We know how to approximate it by a sequence of processes (2.);
I. It has normally-distributed paths whose increments have mean zero (any Hölder continuous mean function would have been OK) and variance \( \mathbb{E}[(X_t - X_s)^2] = O(|t - s|) \);
J. We can identify the limit point: it is characterized by its covariance function, in this case.

The same technique works for many other processes (even for infinite-dimensional ones) whenever we can verify tightness and recognize the weak limit.

**Gaussian Conditional Expectations**

Let \( X \) be a multivariate Gaussian random vector with expectation vector \( \mu \) and covariance matrix \( \Sigma \). For each subset \( I \) of indices denote by \( X_I \) the random vector with components \( X_i \), \( i \in I \), by \( \mu_I \) the expectation \( \mathbb{E}[X_I] \), and by \( \Sigma_{IJ} \) the covariance matrix \( \Sigma_{IJ} = \mathbb{E}[(X_I - \mu_I)(X_J - \mu_J)^T] \). If \( \Sigma_{JJ} \) is nonsingular, a straightforward calculation yields

\[
\mathbb{E}[X_I|X_J] = \mu_I + \Sigma_{IJ}\Sigma_{JJ}^{-1}(X_J - \mu_J)
\]

\[
\mathbb{V}[X_I|X_J] = \Sigma_{II} - \Sigma_{IJ}\Sigma_{JJ}^{-1}\Sigma_{JI}
\]

(in fact, the same formulas work even for singular \( \Sigma_{JJ} \) if we interpret \( \Sigma_{JJ}^{-1} \) as the Moore-Penrose generalized inverse). For mean-zero one-dimensional jointly normal random variables \( x \) and \( y \), we have

\[
\mathbb{E}[y|x] = (\mathbb{E}[xy]/\mathbb{E}[x^2]) x \quad \text{and} \quad \mathbb{V}[y|x] = \mathbb{E}[y^2] - \mathbb{E}[xy]^2/\mathbb{E}[x^2].
\]

From these formulas we can compute the Brownian Bridge’s conditional expectation for \( s \leq t \leq u \) as

\[
\mathbb{E}[X_t|X_s] = \frac{1-t}{1-s}X_s
\]

\[
\mathbb{V}[X_t|X_s] = \frac{(1-t)(t-s)}{1-s}
\]

\[
\mathbb{E}[X_t|X_s, X_u] = \frac{u-t}{u-s}X_s + \frac{t-s}{u-s}X_u
\]

\[
\mathbb{V}[X_t|X_s, X_u] = \frac{(u-t)(t-s)}{u-s}
\]

and, consequently,

\[
\mathbb{E}[(X_{t+\epsilon} - X_t)|\mathcal{F}_t] = -\frac{\epsilon}{1-t}X_t
\]

\[
\mathbb{V}[(X_{t+\epsilon} - X_t)|\mathcal{F}_t] = \epsilon - \frac{\epsilon^2}{1-t} = \epsilon + O(\epsilon^2)
\]

\[
\mathbb{E}[(X_{t+\epsilon} - X_t)^2|\mathcal{F}_t] = \epsilon - \frac{\epsilon^2}{1-t} + [\frac{-\epsilon}{1-t}X_t]^2 = \epsilon + O(\epsilon^2)
\]
Semimartingales
From this calculation $E[X_{t+\epsilon}|\mathcal{F}_t] = X_t - \frac{\epsilon}{t} X_t + O(\epsilon^2)$ it follows that $X_t$ is not a martingale, but that the process

$$W_t = X_t + \int_0^t \frac{X_s}{1-s} \, ds$$

is a martingale, with continuous paths and a Gaussian distribution. The covariance function turns out to be $s \wedge t$, so $W_t$ is just Brownian motion and we have the curious representation

$$X_t = X_0 + \int_0^t \frac{X_s}{1-s} \, ds + W_t$$

of $X_t$ as a semimartingale, the sum of a bounded-variation process (here $X_0 + \int_0^t \frac{X_s}{1-s} \, ds$) and a martingale. This is our first example of a diffusion and of a stochastic integral.

INTRODUCTION TO STOCHASTIC INTEGRATION
Avner Friedman, *Stochastic Differential Equations and Applications*, pp. 55–72

Introduction to Stochastic Integral Equations
Last time we constructed the distribution $\mu$ of the Brownian Bridge as a measure on the canonical space $(\mathcal{C}, \mathcal{B})$, and the B.B. itself as the canonical process $X_t(\omega) = \omega_t$ on $(\mathcal{C}, \mathcal{B}, \mu)$. Using normal distribution theory we calculated the conditional expectations $E[X_{t+\epsilon}|\mathcal{F}_t] = X_t - \frac{\epsilon}{t} X_t$ and $V[X_{t+\epsilon}|\mathcal{F}_t] = \epsilon - \frac{\epsilon^2}{t}$; since $X_t$ is Markov, these are the same as the conditional expectations given the Borel Field generated by the entire past of the process up to time $t$, $\mathcal{F}_t = \sigma[X_s : s \leq t]$. This led to the recognition that $W_t = X_t - \int_0^t \frac{X_s}{1-s} \, ds$ is a continuous-path Gaussian martingale. We calculated that the covariance function is $E W_s W_t = s \wedge t$ and so recognized $W_t$ as the Wiener process, leading to the representation

$$X_t = \int_0^t \frac{X_s}{1-s} \, ds + W_t. \quad (*)$$

This is our first example of a Stochastic Integral Equation (SIE); given a Brownian Motion $W_t$, we can try to “solve” (*) for the unknown process $X_t$. One way to proceed is to define a sequence of processes by $X^n_0 \equiv 0$ and

$$X^{n+1}_t = \int_0^t \frac{X^n_s}{1-s} \, ds + W_t;$$

upon subtracting,

$$X^{n+1}_t - X^n_t = \int_0^t \frac{X^{n-1}_s - X^n_s}{1-s} \, ds$$

so $\gamma^n_t = \sup_{s \leq t} |X^{n+1}_s - X^n_s|$ satisfies $\gamma^n_0 = \sup_{s \leq t} |W_s|$ and, for $t \leq 1 - \epsilon$,

$$\gamma^n_t \leq \int_0^t \frac{\gamma^{n-1}_s}{1-s} \, ds \leq \epsilon^{-1} \int_0^t \gamma^{n-1}_s \, ds \leq \epsilon^{-2} \int_0^t (t-s) \gamma^{n-2}_s \, ds$$

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Since \( e^{t/\epsilon} \sup_{s\leq 1} |W_s| = \sum_{n=0}^{\infty} \frac{t^n}{\epsilon^n} \gamma_1 \gamma_t^n < \infty; \gamma_t^n \to 0 \) uniformly on \( t \leq 1 - \epsilon \) and \( X_t^n \) converges uniformly on compact sets to a limit \( X_t \) satisfying (*).

Since it’s so easy to construct the Brownian Bridge by solving (*), and since all we used was the conditional mean and variance of the infinitesimal increment \( \mathbb{E}[X_{t+\epsilon} - X_t] \), maybe we can use a similar technique for other processes once we know the so-called infinitesimal mean and variance,

\[
\mathbb{E}[(X_{t+\epsilon} - X_t) | \mathcal{F}_t] = \alpha_t \epsilon + o(\epsilon) \quad \mathbb{E}[(X_{t+\epsilon} - X_t)^2 | \mathcal{F}_t] = \beta_t \epsilon + o(\epsilon).
\]

**Stochastic Integrals**

**Stieltjes Integrals**

Any finite measure \( \mu \) on \((0, 1]\) is determined uniquely by its distribution function \( G(t) = \mu((0, t]) \), since the Borel sets are generated by the half-open intervals and \( \mu((a, b]) = G(b) - G(a) \). For any bounded and continuous function \( f \),

\[
\int_0^t f(s) \mu(ds) = \lim_{n \to \infty} 2^{-n} \sum_{j=0}^{2^n t - 1} f\left(\frac{j}{2^n}\right) [G\left(\frac{j+1}{2^n}\right) - G\left(\frac{j}{2^n}\right)],
\]

justifying the \((18^{th}\)-century) Stieltjes notation \( \int_0^t f(s) dG(ds) \). If \( G(t) = \int_0^t G'(s) ds \), this is

\[
\int_0^t f(s) dG(ds) = \lim_{n \to \infty} 2^{-n} \sum_{j=0}^{2^n t - 1} f\left(\frac{j}{2^n}\right) \left[ \frac{1}{2^n} G'\left(\frac{j}{2^n}\right) + o\left(\frac{1}{2^n}\right) \right]
\]

\[= \int_0^t f(s) G'(s) ds;\]

whether or not \( G(t) \) is differentiable, the integration-by-parts formula holds for continuously differentiable functions \( f \) (note \( G(0) = 0 \)):

\[
\int_0^t f(s) dG(ds) = f(t) G(t) - \int_0^t f'(s) G(s) ds.
\]

For step functions \( f(t) \) with a constant value \( b_i \) on each of \( n \) intervals \((t_i, t_{i+1}]\), the integral is just

\[
\int_0^t f(s) dG(ds) = \sum_{i=0}^{n-1} b_i [G(t \wedge t_{i+1}) - G(t \wedge t_i)].
\]
**Wiener Integrals**

Now let \( f(t) \) be a measurable real-valued function and consider the problem of defining the Stieltjes-like integral \( M_t = \int_0^t f(s) \, dW_s \) for Brownian Motion \( W_s \); since \( W_s \) is not differentiable, we can’t use the representation \( M_t = \int_0^t f(s) \, dW_s \, ds \) as we did for differentiable functions \( G(t) \) above. The other two alternatives do work, however; for continuously differentiable \( f(t) \) we can define

\[
\int_0^t f(s) \, dW_s = f(t)W_t - \int_0^t f'(s)W_s \, ds,
\]

or for step functions we can define

\[
\int_0^t f(s) \, dW_s = \sum_{i=0}^{n-1} b_i [W_{t \wedge t_{i+1}} - W_{t \wedge t_i}].
\]

In either case \( \int_0^t f(s) \, dW_s \) is a continuous-path Gaussian martingale with mean zero and variance

\[
\mathbb{E}[\left( \int_0^t f(s) \, dW_s \right)^2] = \mathbb{E}\left[ \left( \sum_{i=0}^{n-1} b_i [W_{t_{i+1}} - W_{t_i}] \right) \left( \sum_{j=0}^{n-1} b_j [W_{t_{j+1}} - W_{t_j}] \right) \right]
\]

\[
= \sum_{i=0}^{n-1} b_i^2 \mathbb{E}[W_{t_{i+1}} - W_{t_i}]^2 \quad \text{(by independent increments)}
\]

\[
= \sum_{i=0}^{n-1} b_i^2 [t_{i+1} - t_i], \quad \text{and hence covariance}
\]

\[
\mathbb{E}[\left( \int_0^s f(u) \, dW_u \right) \left( \int_0^t g(u) \, dW_u \right)] = \int_0^{s \wedge t} f(u)g(u) \, du.
\]

**ITÔ STOCHASTIC INTEGRALS AND DIFFUSIONS**

Avner Friedman, *Stochastic Differential Equations and Applications*, pp. 55—72

**Properties of Stochastic Integrals**

Last time we constructed the so-called *Wiener stochastic integral* process \( M_t = \int_0^t f(s) \, dW_s \) for nonrandom square-integrable functions \( f(s) \). We defined the integral first for step functions \( f(t) = \sum_i b_i \mathbf{1}_{(t_i, t_{i+1}]} \) as \( M_t = \int_0^t f(s) \, dW_s = \sum_i b_i \left[ W_{t \wedge t_{i+1}} - W_{t \wedge t_i} \right] \), then extended by \( L^2 \) continuity. For continuously differentiable functions \( f(t) \) we also defined the integral by parts as \( M_t = \int_0^t f(s) \, dW_s = f(t)W_t - \int_0^t f'(s)W_s \, ds \). For continuous functions \( f(t) \) we have the \( L^2 \) convergent formula

\[
M_t = \int_0^t f(s) \, dW_s = \lim_{n \to \infty} \frac{1}{2^n} \sum_{j=0}^{2^n t - 1} f\left( \frac{j}{2^n} \right) \left[ W_{(j+1)/2^n} - W_{j/2^n} \right]
\]

(*)

which makes it easy to see why the increments \( M_t - M_s \) have Gaussian distributions with mean zero and variance \( \mathbb{E}[(M_t - M_s)^2] = \int_s^t f^2(u) \, du \).

It was Kyoshi Itô’s observation that the same construction would also work for *random* integrands \( f(t) \), provided that always \( f(t) \) is square-integrable and independent of \( [W_{t+\varepsilon} - W_t] \); since
$W_t$ has independent increments, we can assure that independence by requiring that $f(t)$ be $\mathcal{F}_t$-measurable for each $t$. The function $f_t = W_t$ satisfies this condition, as does the Brownian Bridge $X_t = \int_0^t \frac{-X_s}{1-s} ds + W_t$ and any function $\beta_t = b_t(X_t)$ of $t$ and $X_t$.

Call a stochastic process $\beta_t$ (such as $W_t$ or $b_t(X_t)$) adapted to the family $\{\mathcal{F}_t\}$ of BF’s if $\beta_t$ is $\mathcal{F}_t$-measurable for every $t$, i.e., if $[\omega : \beta_t(\omega) \leq r] \in \mathcal{F}_t$ for each $t \geq 0$ and $r \in \mathbb{R}$, and let $L^2_t$ be the metric space of adapted processes satisfying $\mathbb{E}[\int_0^1 \beta_s^2 dt] < \infty$. Any such process $\beta_s \in L^2_t$ can be approximated by an adapted simple function, with a constant $\mathcal{F}_t$-measurable value $b_t(\omega)$ on the interval $(t_i, t_{i+1})$, for which it is easy to calculate the stochastic integral $M_t = \int_0^t \beta_s dW_s = \sum_i b_i [W_{t_{i+1}} - W_{t_i}]$; the mean and variance of the (usually non-Gaussian) Itô integral $M_t$ are:

$$\mathbb{E}[M_t] = \mathbb{E}\left[ \sum_i b_i [W_{t_{i+1}} - W_{t_i}] \right]$$

$$= \sum_i \mathbb{E}[b_i [W_{t_{i+1}} - W_{t_i}]] = 0$$

$$\mathbb{V}[M_t] = \mathbb{E}\left[ \left( \sum_i b_i [W_{t_{i+1}} - W_{t_i}] \right)^2 \right]$$

$$= \sum_i \mathbb{E}[b_i^2 [W_{t_{i+1}} - W_{t_i}]^2]$$

$$= \sum_i \mathbb{E}[b_i^2] [t_{i+1} - t_i]$$

$$= \int_0^1 \mathbb{E}[\beta_s^2] ds$$

and, more generally, $\mathbb{E}[M_t] = 0$ and $\mathbb{V}[M_t] = \mathbb{E}[M_t^2] = \int_0^t \mathbb{E}[\beta_s^2] ds$. Possibly more revealing is the conditional variance; for continuous $\beta_s$ this is

$$\mathbb{V}[M_{t+\epsilon} | \mathcal{F}_t] = \int_t^{t+\epsilon} \mathbb{E}[\beta_s^2 | \mathcal{F}_t] ds = \epsilon \mathbb{E}[\beta_t^2] + o(\epsilon),$$

so $\epsilon \mathbb{E}[\beta_t^2]$ is just the conditional variance of $[M_{t+\epsilon} - M_t]$, to first order in $\epsilon$.

If $\alpha_t \in L^1_t$ is a integrable adapted process the indefinite integral $\int_0^t \alpha_s ds$ is defined in the usual (Lebesgue or Riemann) way; for any $\mathcal{F}_0$-measurable random variable $X_0$, the sum

$$X_t = X_0 + \int_0^t \alpha_s ds + \int_0^t \beta_s dW_s$$

is a continuous-path adapted process whose increments have conditional mean and variance

$$\mathbb{E}[X_{t+\epsilon} - X_t | \mathcal{F}_t] = \int_t^{t+\epsilon} \mathbb{E}[\alpha_s | \mathcal{F}_t] ds$$

$$= c\alpha_t + o(\epsilon),$$

$$\mathbb{V}[X_{t+\epsilon} - X_t | \mathcal{F}_t] = \epsilon \mathbb{V}[\beta_s] + o(\epsilon).$$

If we restrict our attention to Markov processes $X_t$, then the conditional mean and variance $c\alpha_t$ and $\epsilon \mathbb{V}[\beta_s]$ must be not only $\mathcal{F}_t$-measurable, but $\alpha(X_t)$-measurable— so, for some functions $a_t(x)$ and $b_t(x)$, $\alpha_t = a_t(X_t)$ and $\beta_t = b_t(X_t)$. This class of functions is called diffusions:

$$X_t = X_0 + \int_0^t a_s(X_s) ds + \int_0^t b_s(X_s) dW_s.$$
We have already met several examples, including:

Brownian Bridge: \( X_0 = 0, a_t(x) = \frac{x}{t}, \) and \( b_t(x) = 1. \)

Brownian Motion with Drift: \( X_0 = x_0, a_t(x) = \mu, \) and \( b_t(x) = \sigma. \)

Geometric Brownian Motion: \( X_0 = e^{x_0}, a_t(x) = x(\mu + \frac{\sigma^2}{2}), \) and \( b_t(x) = x\sigma. \)

Reflected Brownian Motion: \( X_0 = 0, a_s(x) = \delta(x), \) and \( b_s(x) = 1. \) Here \( \delta(x) \) denotes Dirac’s delta function, the (formal) derivative of the function \( 1_{[0,\infty)}(x) \); reflected Brownian motion is properly called a diffusion with boundary, and is more complicated to study than the other processes mentioned.

The remarkable and deep fact is that all continuous-path strong Markov processes are diffusions; see Karlin & Taylor, A Second Course in Stochastic Processes, chapter 15, or Stroock & Varadhan, Multi-dimension Diffusion Processes, or Itô & McKean, Diffusion Processes and their Sample Paths (among others).

**Itô’s Formula**

For any \( \phi_t(x) \in C^{1+2}(\mathbb{R}_+ \times \mathbb{R}), \) Taylor’s formula gives

\[
\phi_{t+\epsilon}(x + \xi) = \phi_t(x) + \epsilon \frac{\partial \phi}{\partial t} + \xi \frac{\partial \phi}{\partial x} + \frac{\xi^2}{2} \frac{\partial^2 \phi}{\partial x^2} + o(\epsilon) + o(\xi^2),
\]

and in particular \( Y_t = \phi_t(X_t) \) satisfies

\[
Y_{t+\epsilon} - Y_t = \epsilon \frac{\partial \phi}{\partial t} + \left( a_t(X_t)\epsilon + b_t(X_t)[W_{t+\epsilon} - W_t]\right) \frac{\partial \phi}{\partial x} + \epsilon \frac{b_t^2(X_t)}{2} \frac{\partial^2 \phi}{\partial x^2} + o(\epsilon),
\]

so \( Y_t = \phi_t(X_t) \) is itself a diffusion with starting point \( Y_0 = \phi_0(X_0) \) and diffusion coefficients

\[
\tilde{a}_t(x) = \frac{\partial \phi}{\partial t} + a_t(x) \frac{\partial \phi}{\partial x} + \frac{b_t^2(x)}{2} \frac{\partial^2 \phi}{\partial x^2} \quad \text{and} \quad \tilde{b}_t(x) = b_t(x) \frac{\partial \phi}{\partial x}.
\]

There is a close connection between the diffusion \( X_t \) and the differential operator

\[
\mathcal{L} \phi \equiv a_t(x) \frac{\partial \phi}{\partial x} + \frac{1}{2} b_t^2(x) \frac{\partial^2 \phi}{\partial x^2},
\]

called the Generator of the process; we have just seen that \( \tilde{a}_t(x) = \frac{\partial}{\partial t} \phi_t(x) + \mathcal{L} \phi_t(x), \) for example, so for any \( \phi_t(x), \)

\[
M_t = \phi_t(X_t) - \int_0^t \left[ \frac{\partial \phi}{\partial s} + a_s(Y_s) \frac{\partial \phi}{\partial x} + \frac{1}{2} b_s^2(Y_s) \frac{\partial^2 \phi}{\partial x^2} \right] ds
\]

is a martingale. This, in fact, is the modern definition of the Diffusion Process with coefficients \( a_s(x) \) and \( b_s(x). \) Note that \( Y_t = \phi_t(X_t) \) is itself a martingale if \( \phi \) satisfies the parabolic partial differential equation \( \partial \phi / \partial t = -\mathcal{L} \phi, \) i.e.,

\[
0 \equiv \left[ \frac{\partial \phi_t(x)}{\partial t} + a_t(x) \frac{\partial \phi_t(x)}{\partial x} + \frac{1}{2} b_t^2(x) \frac{\partial^2 \phi_t(x)}{\partial x^2} \right].
\]
Examples

For example, let \( \phi_t(x) = x^2 \); then \( \mathcal{L}\phi(x) = 2a_t(x)x + b_t^2(x) \), so Itô’s formula gives

\[
X_t^2 = X_0^2 + \int_0^t (2a_s(X_s)X_s + b_s(X_s)^2) \, ds + \int_0^t 2X_s \, b_s(X_s) \, dW_s.
\]

An interesting formula for every diffusion process \( X_t \) follows from this:

\[
\int_0^t 2X_s \, dX_s = \int_0^t 2X_s \, a_s(X_s) \, ds + \int_0^t 2X_s \, b_s(X_s) \, dW_s
\]
\[
= \int_0^t 2X_s \, a_s(X_s) \, ds + X_t^2 - X_0^2 - \int_0^t (2a_s(X_s)X_s + b_s(X_s)^2) \, ds
\]
\[
= X_t^2 - X_0^2 - \int_0^t b_s(X_s)^2 \, ds; \quad (\ast)
\]

for ordinary integrals we have, of course, \( \int_0^t 2f(s) \, df(s) = f(t)^2 - f(0)^2 \), but for stochastic integrals there is an additional term.

Inference

Suppose we observe \( X_s \) for \( 0 \leq s \leq t \), and believe that \( X_t \) is a diffusion process; of course we can observe the initial value \( X_0 \), but what can we infer about the coefficients \( a_s(x) \) and \( b_s(x) \)? First let’s consider the diffusion coefficient \( b_s(x) \). There is no hope of inferring anything about \( b_s(x) \) away from the observed path (unless we make additional assumptions about the form of \( b_s(x) \)), but if we know \( a_s(x) \) and \( b_s(x) \) to be sufficiently smooth in both \( s \) and \( x \), then for small \( \epsilon > 0 \) the quadratic variation between \( s \) and \( s+\epsilon \) is

\[
Q^{s+\epsilon}_s(X) = \lim_{n \to \infty} \sum_{i=0}^{n-1} (X_{s+(i+1)\frac{\epsilon}{n}} - X_{s+i\frac{\epsilon}{n}})^2
\]
\[
= \lim_{n \to \infty} \sum_{i=0}^{n-1} \left( a_s(X_s)(\frac{\epsilon}{n}) + b_s(X_s)(W_{s+(i+1)\frac{\epsilon}{n}} - W_{s+i\frac{\epsilon}{n}}) \right)^2 + o(\epsilon)
\]
\[
= b_s(X_s)^2 \lim_{n \to \infty} \sum_{i=0}^{n-1} \left( W_{s+(i+1)\frac{\epsilon}{n}} - W_{s+i\frac{\epsilon}{n}} \right)^2 + o(\epsilon)
\]
\[
= b_s(X_s)^2 \lim_{n \to \infty} \sum_{i=0}^{n-1} \left( \frac{\epsilon}{n} \right)^2 + o(\epsilon) = \epsilon b_s(X_s)^2 \lim_{n \to \infty} \frac{\epsilon^2}{n} + o(\epsilon)
\]
\[
= \epsilon b_s(X_s)^2 + o(\epsilon),
\]

so \( b_s(X_s)^2 \) is observable as \( \lim_{\epsilon \to 0} Q^{s+\epsilon}_s(X)/\epsilon \). Since \( b_s(x) \) is observable, consider diffusions with constant diffusion coefficient:

\[
X_t = X_0 + \int_0^t a_s(X_s) \, ds + \sigma W_t.
\]

What can we learn from the path on \( 0 \leq s \leq t \) about \( a_s(x) \)? Let’s try to compute the likelihood for \( a \). For large \( n \) set \( \epsilon = t/n \) and note that

\[
X_{(i+1)\epsilon} = X_{i\epsilon} + \epsilon a_{i\epsilon} + \sigma (W_{(i+1)\epsilon} - W_{i\epsilon}) + o(\epsilon).
\]
The log likelihood for \( a \) upon observing only \( X_{i\epsilon}, i = 0, \ldots, n \), is
\[
\ell_n(a) = c_n - \frac{n}{2} \log(2\pi \epsilon \sigma^2) - \frac{1}{2\epsilon \sigma^2} \sum_{0 \leq i < n} (X_{(i+1)\epsilon} - X_{i\epsilon} - \epsilon a \epsilon (X_{i\epsilon}))^2 + o(\epsilon),
\]
for any constant \( c_n \); it’s convenient to choose \( c_n \) so that \( \ell_n(0) = 0 \), \( i.e. \), \( c_n = \frac{n}{2} \log(2\pi \epsilon \sigma^2) + \frac{1}{2\epsilon \sigma^2} \sum_{0 \leq i < n} (X_{(i+1)\epsilon} - X_{i\epsilon})^2 \), whereupon
\[
\ell_n(a) = \frac{1}{2\epsilon \sigma^2} \sum_{0 \leq i < n} \left( (X_{(i+1)\epsilon} - X_{i\epsilon})^2 - (X_{(i+1)\epsilon} - X_{i\epsilon} - \epsilon a \epsilon (X_{i\epsilon}))^2 \right) + o(\epsilon)
\]
\[
= \frac{1}{2\epsilon \sigma^2} \sum_{0 \leq i < n} \left( 2\epsilon (X_{(i+1)\epsilon} - X_{i\epsilon})a \epsilon (X_{i\epsilon}) - \epsilon^2 a \epsilon (X_{i\epsilon})^2 \right) + o(\epsilon)
\]
\[
= \sigma^{-2} \int_0^t a_s(X_s) \, dX_s - \frac{1}{2\sigma^2} \int_0^t a_s(X_s)^2 \, ds + o(\epsilon)
\]
Now we pass to the limit \( n \to \infty \) (and \( \epsilon \to 0 \)), \( \ell(a) = \sigma^{-2} \left( \int_0^t a_s(X_s) \, dX_s - \frac{1}{2} \int_0^t a_s(X_s)^2 \, ds \right) \).

**Example 1: Wiener Process**

For example, if \( X_t = X_0 + \mu t + \sigma W_t \) is Brownian motion with constant drift \( \mu \) and diffusion rate \( \sigma^2 \), then \( a_s(x) \equiv \mu \) and the log likelihood becomes
\[
\ell(\mu) = \mu(X_t - X_0)/\sigma^2 - \mu^2 t/2\sigma^2;
\]
the MLE estimate is \( \hat{\mu} = (X_t - X_0)/t \), while the Bayesian posterior distribution for a flat prior is
\[
\mu | \{X_s : 0 \leq s \leq t\} \sim N \left( \frac{X_t - X_0}{t}, \frac{\sigma^2}{t} \right).
\]

**Example 2: Ornstein-Uhlenbeck**

Now if \( X_t = X_0 - \beta \int_0^t X_s \, ds + \sigma W_t \), or \( dX_t = -\beta X_t \, dt + \sigma dW_t \), then \( a_s(x) \equiv -\beta x \) and
\[
\ell(\beta) = \frac{-\beta}{\sigma^2} \int_0^t X_s \, dX_s - \frac{\beta^2}{2\sigma^2} \int_0^t (X_s)^2 \, ds
\]
\[
= -\frac{\beta}{2\sigma^2} (X_t^2 - X_0^2 - t) - \frac{\beta^2}{2\sigma^2} \int_0^t (X_s)^2 \, ds, \quad \text{(by (*) )}
\]
so for a uniform prior we would have
\[
\beta | \{X_s : 0 \leq s \leq t\} \sim N \left( \frac{X_t^2 - X_0^2 - t}{2 \int_0^t (X_s)^2 \, ds}, \frac{\sigma^2}{\int_0^t (X_s)^2 \, ds} \right).
\]

**Testing Hypotheses**

The likelihood function provides the basis for testing hypotheses like \( H_0 : X_t \text{ is Brownian motion} \) (with no drift) against alternatives like \( H_1 : X_t \text{ is a Wiener process} \) (with constant drift) or \( H_2 : X_t \text{ is an O-U process} \) (with linear drift), by finding either \( P \)-values or posterior probabilities.
Girsanov’s Formula
To simplify expressions, rescale if necessary so that $\sigma^2 = 1$, and $X(t)$ will be the Wiener process $W(t)$ if $a \equiv 0$ vanishes. The likelihood function

$$
\zeta_t \equiv \exp[\ell(a)] = \exp\left\{ \int_0^t a_s(X_s) \, dX_s - \frac{1}{2} \int_0^t a_s(X_s)^2 \, ds \right\}
$$

plays a special role. It is a positive martingale, and so induces a new probability measure $Q$ on $(\Omega, \mathcal{F})$ by

$$
\mathbb{E}^Q[H] \equiv \mathbb{E}^P[H_G]
$$

for each $\mathcal{F}_t$-measurable $H \in L^1(\Omega, \mathcal{F}, P)$; under this new measure (i.e. on the probability space $(\Omega, \mathcal{F}, Q)$), $X_t$ is a Brownian motion (i.e. has drift 0). The technique of switching measures (from $P$ to $Q$) to eliminate drift is a powerful one, used frequently in both the theoretical study of diffusions (it leads to the strongest existence and uniqueness results) and to mathematical finance (where it goes under the name of “risk-neutral measure”, useful for computing option prices).

RANDOM MEASURES
Nonparametric Statistics and Random Measures

Let $(\mathcal{X}, \mathcal{B})$ be a measurable space and $(\Omega, \mathcal{F}, P)$ a probability space; denote by $M = M(\mathcal{X}, \mathcal{B})$ the vector space of $\sigma$-finite signed measures on $(\mathcal{X}, \mathcal{B})$. If we observe a random variable $X \in \mathcal{X}$, what can we say about its probability distribution $\mu_X \in M$? The fundamental problem of statistics is making inference about $\mu_X$ on the basis of observation of $X$. In a parametric analysis we postulate that $\mu_X$ lies in a small family of distributions $\mu_X \in \{P_\theta : \theta \in \Theta\}$ (e.g., if $\mathcal{X} = \mathbb{R}^n$, we might postulate that $\mu_X$ lies in the multivariate normal family with constant mean vector and covariance matrix $\mathbf{I} = \sigma^2 \mathbf{I}$) indexed by a parameter $\theta$ lying in a low-dimensional space $\Theta$ (e.g., $\Theta = \{(\mu, \sigma^2)\} \subset \mathbb{R}^2$). If some $\sigma$-finite measure $\nu(dx)$ dominates all the $P_\theta(dx)$, then the Radon-Nikodym derivative (or density) $L(\theta, x) = P_\theta(dx)/\nu(dx)$ is called the likelihood function and inference often proceeds either by

1. seeking the value of $\theta \in \Theta$ that maximizes $L(\theta, X)$, and studying its properties (the Frequentist approach); or by
2. specifying a “prior” probability measure $\pi(d\theta)$ on $\Theta$, calculating the conditional “posterior” distribution $\pi(d\theta|X)$, and studying its properties (the Bayesian approach).

In the nonparametric approach no finite-dimensional $\Theta$ is postulated: all probability measures $\mu \in M$ are regarded as possible distributions for $X$, and analysis proceeds either by

1. seeking the value of $\mu \in M$ that maximizes some analogue of the likelihood like $\mu(dx)/\nu(dx)$ for a reference measure $\nu$ on $\mathcal{X}$, and studying its properties (the Frequentist approach); or by
2. specifying a “prior” probability measure on the possible distributions $\mu \in M$, calculating the conditional “posterior” distribution, and studying its properties (the Bayesian approach).

We can think of $\mu$ as a “random measure,” first under a prior distribution and later under a posterior. We now turn to the study of random measures. A random measure can be thought of in at least three different ways:

1. A function $\mu : \mathcal{B} \times \Omega \to \mathbb{R}$, mapping $(B, \omega) \mapsto \mu(B, \omega) \in \mathbb{R}$;
2. A function $\mu : \Omega \to M$, mapping $\omega \mapsto \mu(\cdot, \omega) \in M$;
3. A function $\mu : B \to L^1(\Omega, \mathcal{F}, P)$, mapping $B \mapsto \mu(B, \cdot) \in L^1(\Omega, \mathcal{F}, P)$.

We omit the $\omega$ and denote the value by $\mu(B)$ in all three cases. The second perspective represents $\mu$ simply as a random variable, taking values in some abstract space $M$; sometimes that’s useful in technical arguments, but it is usually easier to think about random measures from the third perspective, as a family of ordinary random variables indexed by the Borel sets $B \in \mathcal{B}$.
Examples

Example 1: Wiener Measure

For any set $T$, any mean function $\mu : T \rightarrow \mathbb{R}$, and any real $\alpha$-positive-definite covariance function $\rho : T \times T \rightarrow \mathbb{R}$, there exists a probability space $(\Omega, \mathcal{F}, P)$ and a Gaussian process $X_t$ indexed by $t \in T$ with $E[X_t] = \mu_t$ and $E[(X_s - \mu_s)(X_t - \mu_t)] = \rho_{st}$. In particular we can take $T = \mathbb{B}$, the Borel sets in $\mathcal{X} = \mathbb{R}^+$; $\mu_B = 0$ for all $B \in \mathcal{B}$; and $\rho_{AB} = \lambda(A \cap B)$, the Lebesgue measure of the intersection. In this case the “cumulative distribution function” (or Stieltjes function) $W(t) = \mu((0, t])$ associated with the random measure $\mu$ is just the standard Wiener process, and integrals $\int f(t)\mu(dt)$ of simple or even $L^2$ functions are just the same as Wiener integrals $\int f(t)\,dW_t$. The construction is not limited to $\mathbb{R}^1$, however, and just as easily leads to $n$-dimensional Gaussian measures and Wiener integrals, and the $n$-parameter analogue of the Wiener process sometimes called the “Brownian Sheet.”

Example 2: Brownian Bridge

If we now take $T = \mathbb{B}((0, 1])$, the Borel sets in the unit interval $\mathcal{X} = (0, 1]$; $\mu_B = 0$ for all $B \in \mathcal{B}$; and $\rho_{AB} = \lambda(A \cap B) - \lambda(A)\lambda(B)$, the cumulative distribution function $B(t) = \mu((0, t])$ is just the standard Brownian Bridge process, and integrals of simple or $L^2$ functions can be written in terms of Wiener integrals as $\int f(t)\mu(dt) = \int f(t)\,dW_t - W_1 \int f(t)\,dt$ for any Wiener process $W_t = B_t + tZ$, $Z \sim N(0, 1)$ independent of $B_t$.

Example 3: The Gamma Process

Preliminaries: Gamma, Beta, and Dirichlet Distributions

If $X \sim \text{Ga}(\alpha, 1)$ and $Y \sim \text{Ga}(\beta, 1)$ are independent Gamma random variables, then $X$ and $Y$ have joint density function

$$f(x, y)\,dx\,dy = \frac{x^{\alpha-1}e^{-x}}{\Gamma(\alpha)} \frac{y^{\beta-1}e^{-y}}{\Gamma(\beta)} 1_{\mathbb{R}_+}(x) 1_{\mathbb{R}_+}(y)\,dx\,dy$$

so $W = X + Y$ and $Z = \frac{X}{X + Y}$ have joint distribution

$$f(w, z)\,dw\,dz = \frac{(zw)^{\alpha-1}e^{-zw}}{\Gamma(\alpha)} \frac{((1-z)w)^{\beta-1}e^{-(1-z)w}}{\Gamma(\beta)} 1_{\mathbb{R}_+}(w)\,dw 1_{[0,1]}(z)\,dz$$

$$= \frac{\Gamma(\alpha + \beta)}{\Gamma(\alpha)\Gamma(\beta)} (z)^{\alpha-1}(1-z)^{\beta-1} 1_{[0,1]}(z)\,dz \frac{w^{\alpha+\beta-2}e^{-w}}{\Gamma(\alpha + \beta)} 1_{\mathbb{R}_+}(w)\,dw$$

It follows that $W$ and $Z$ are independent with $\text{Ga}(\alpha + \beta, 1)$ and $\text{Be}(\alpha, \beta)$ distributions, respectively; thus the conditional distribution of $X$, given $X + Y = W$, is that of $W$ times an independent $\text{Be}(\alpha, \beta)$ variable. We will need this for $\alpha = \beta = 1/2^n$.

The Construction

Let $\alpha$ be a $\sigma$-finite nonnegative measure on the space $(\mathcal{X}, \mathcal{B})$, and let $(\Omega, \mathcal{F}, P)$ be a probability space; the Gamma Process with mean $\alpha$ is a random measure $\nu : \mathcal{B} \times \Omega \rightarrow \mathbb{R}$ which assigns independent Gamma random variables $\nu(\Lambda_i) \sim \text{Ga}(\alpha_i, 1)$ to disjoint sets $\Lambda_i \in \mathcal{B}$ with finite measures $\alpha(\Lambda_i) = \alpha_i < \infty$. Here is an explicit construction of $\nu$ for $\mathcal{X} = \mathbb{R}_+$:

Let $z_i^n$ be a doubly-indexed independent family of random variables with the Beta distribution $\text{Be}(\frac{1}{2^n}, \frac{1}{2^n})$; the $z_i^n$ are independent with uniform distributions, the $z_i^1$ have the $\text{Be}(1/2, 1/2)$, etc. Define a stochastic process $X_t$ at integer times $t \in \mathbb{N}$ by $X_t = x^n_t$ where

$$x^n_t = \sum_{i=1}^t -\log(z_i^0),$$

so $X_t$ has independent increments $[X_t - X_s]$ with the $\Gamma(t-s, 1)$ distribution for integers $s, t$. For successive $n$ define $X_t$ at dyadic rational times recursively by $X_t = x^n_t$, $t = i/2^n$, where
\[
x_{2t}^{n+1} = x_t^n \\
x_{2t+1}^{n+1} = x_t^n + (x_{i+1}^n - x_t^n)z_{i+1}^n
\]

This defines \( X_t \) for all dyadic rational \( t \); by our preliminary observation above, the increments \([X_t - X_s]\) are independent with the \( \Gamma(t-s, 1) \) distribution. The process \( X_t \) is nonnegative and nondecreasing, so we can extend the definition to all of \( \mathbb{R}_+ \) by requiring right-continuity: \( X_t = \inf[x^n_i : t \leq i/2^n] \). We will see below that right-continuity is the best we can hope for, i.e., that the process \( X_t \) does not have continuous sample paths (in fact, it has infinitely many jumps in every open interval \((t, t + \epsilon)\) almost surely!) For both rational and irrational \( s < t \), the increments \([X_t - X_s]\) are independent with the Gamma \( \Gamma((t-s), 1) \) distributions, and hence with finite means \( \mathbb{E}[X_t - X_s] = (t-s) \) and variances \( \mathbb{V}[X_t - X_s] = (t-s) \).

Given any \( \sigma \)-finite measure \( \alpha \) on \( \mathbb{R}_+ \), define a right-continuous function \( A(x) = \alpha(0, x] \) and a random measure \( \nu \) by:

\[
\nu((s, t]) = X_{A(t)} - X_{A(s)}
\]

for the standard Gamma process \( X_t \) defined above; we extend by additivity to the field generated by the half-open intervals \((s, t]\), and by continuity to the Borel sets with finite \( \alpha \)-measure, upon noting that

\[
\mathbb{E}\nu((s, t]) = [A(t) - A(s)] = \alpha((s, t]) \\
\mathbb{V}\nu((s, t]) = [A(t) - A(s)] = \alpha((s, t]),
\]

so (by \( L^2 \) continuity) \( \mathbb{E}[\nu(B)] = \mathbb{V}[\nu(B)] = \alpha(B) \) for all \( B \in \mathcal{B} \). We will see below that, almost surely, \( \nu \) is a discrete measure concentrated on a (random) countable set of points \( \tau_i(\omega) \).

**Example 4: The Dirichlet Process**

Now let \( \alpha \) be a finite nonnegative measure on \((\mathcal{X}, \mathcal{B})\) and let \( \nu(dx) \) be a Gamma process random measure with mean \( \mathbb{E}[\nu(dx)] = \alpha(dx) \); since \( \alpha(\mathbb{R}_+) < \infty \), \( \nu(\mathbb{R}_+) \) is a well-defined random variable and we can construct

\[
\mu(A) = \frac{\nu(A)}{\nu(\mathbb{R}_+)}
\]

for all \( A \in \mathcal{B} \). Each random variable \( \mu(A) \) has a Beta \( \text{Be}(\alpha(A), \alpha(A^c)) \) distribution, and for any partition \( \Lambda_i \) of \( \mathcal{X} \) into \( n \) disjoint sets, the \( n \)-variate random variables \( X_i = \mu(\Lambda_i) \) have the Dirichlet \( D(\alpha_1, \ldots, \alpha_n) \) distribution with parameters \( \alpha_i = \alpha(\Lambda_i) \). Just as the Gamma process \( \nu \) was almost-surely concentrated on a countable set of points \( \tau_i \), so too is the Dirichlet process \( \mu \)... in fact, it is the same set \( \{\tau_i\}! \) The Dirichlet process is, almost surely, a discrete distribution.

The Dirichlet Process is an important example, because of its use in nonparametric Bayesian statistics. The principal result is this:

**Theorem.** Let \( \mu \sim \text{Dir}(\alpha_o) \) for some finite measure \( \alpha_o \) and let \( X_1, X_2, \ldots, X_n \) be independent observations all with distribution \( \mu_o \). Then, conditional on \( X_1, \ldots, X_n \), \( \mu \sim \text{Dir}(\alpha_n) \) for the measure \( \alpha_n(dx) = \alpha_o(dx) + \sum_{i=1}^n \delta(x - X_i) \) equal to \( \alpha_o(dx) \) plus a unit point mass at each observed \( X_i \).

**Corollary.** Under the same conditions, the predictive distribution for \( X_{n+1} \) assigns mass \( \frac{1}{n+\alpha(\mathbb{R})} \) to \( x = X_i \) for each \( 1 \leq i \leq n \) and the rest of the mass \( \frac{\alpha(\mathbb{R})}{n+\alpha(\mathbb{R})} \) to the prior mean, \( \frac{\alpha(dx)}{\alpha(\mathbb{R})} \).

Note that, from the corollary, the probability of a tie among the first \( n \) variables is

\[
1 - \prod_{i=1}^n \left( \frac{\alpha(\mathbb{R})}{\alpha(\mathbb{R}) + i - 1} \right) = 1 - \frac{\alpha(\mathbb{R})^n \Gamma(\alpha(\mathbb{R}))}{\Gamma(\alpha(\mathbb{R}) + n)}.
\]
arbitrarily close to 1 for large enough \( n \); if the \( X_i \) “really” come from any continuous distribution, no ties will be observed no matter how large \( n \) might be. If \( f(x) \) is any density at all and \( \epsilon > 0 \), the posterior for a Bayesian prior giving probability \( \epsilon \) to \( \mu_X(dx) = f(x) \, dx \) and \( 1-\epsilon \) to \( \mu \sim \text{Dir}(\alpha) \) will eventually be concentrated on \( f(x) \). This proves that Bayesian analysis can be inconsistent. See me for more references if you’re interested in this point.

**Path Discontinuity and SII Processes**

A celebrated theorem of Lévy and Khinchine asserts that every stationary, independent increment (SII) process \( X_t \) has a characteristic function of the form:

\[
E[e^{i\lambda X_t}] = e^{i\lambda x_0 + it\lambda m - t\frac{\lambda^2}{2} + t\int \left( e^{i\lambda u} - 1 \right) \nu(du)}
\]

for some initial value \( x_0 \), drift \( m \), diffusion constant \( \sigma^2 \), and jump rate (“Lévy”) measure \( \nu \). If \( \nu(du) \equiv 0 \) then \( X_t \) is simply Brownian motion with drift, \( X_t = x_0 + mt + \sigma W_t \); \( \nu(E) \) is the rate at which the process jumps by amounts \( u \in E \). The total jump rate \( \nu(\mathbb{R}) \) need not be finite, but \( \nu \) must satisfy \( \int [1 \wedge |u|] \nu(du) < \infty \); it’s OK to have infinitely many tiny jumps if they’re small enough to have a finite sum, almost surely. If \( \mu(\mathbb{R}) < \infty \) then we can interpret the process as one with exponentially distributed waiting times (with means \( 1/\mu(\mathbb{R}) \)) between successive jumps, which are randomly distributed with distribution \( \mu(du)/\mu(\mathbb{R}) \). With a little more work it’s possible to make sense of processes with jump measures satisfying only the weaker condition \( \int [1 \wedge u^2] \nu(du) < \infty \), but the argument gets more subtle. Ask if you need references.

A standard Poisson process \( N(t) \), for example, has characteristic function

\[
E[e^{i\lambda X_t}] = \sum_{k=0}^{\infty} e^{-t} \frac{t^k}{k!} e^{i\lambda k}
\]

for \( \lambda x_0 = 0 \), drift \( m = 0 \), \( \sigma^2 = 0 \), and \( \nu(du) = \delta(u-1) \, du \). A generalized Poisson process is a sum \( X_t = \sum_i u_i N_i(r_i t) \) of re-scaled independent Poissons, which takes jumps of size \( u_i \) at rate \( r_i \); it has measure \( \nu(du) = \sum_i r_i \delta(u_i - u) \, du \). Any SII process can be approximated by the sum of Brownian motion with drift and a generalized Poisson Process. In particular, its paths will be continuous if and only if it is Brownian, and if not we can find the rate of jumps by identifying the Lévy measure \( \nu \).

For example, the characteristic function of the standard Gamma Process \( X_t \sim \text{Ga}(t, 1) \) is \( E[e^{i\lambda X_t}] = (1 - i\lambda)^{-t} \); it has no drift or diffusion part, and has Lévy measure \( \nu(du) \) satisfying

\[
-t \log(1 - i\lambda) = t \int \left( e^{iu\lambda} - 1 \right) \nu(du)
\]

or, after differentiating with respect to \( \lambda \),

\[
\frac{it}{1 - i\lambda} = t \int iue^{iu\lambda} \nu(du)
\]

But \( it \int_0^\infty e^{-u(1-i\lambda)} \, du = it(1 - i\lambda)^{-1} \), so

\[
ite^{-u}_{(0,\infty)}(u) \, du = tiu \, \nu(du)
\]

\[
\nu(du) = u^{-1}e^{-u} \, du \quad (u > 0).
\]

This is not a finite measure, so the Gamma process jumps infinitely often in every time interval; the rate of jumps bigger than \( \epsilon \) is \( \int_\epsilon^\infty e^{-u}/u \, du \), finite for every \( \epsilon > 0 \), and the mean sum of all jumps in time \( t \) is \( t \int_0^\infty u e^{-u}/u \, du = t \).
As interesting exercises, try to find:

1. The Lévy measure $\nu(du)$ for the Cauchy process with characteristic function $E[e^{i\lambda X_t}] = e^{-t|\lambda|}$;

2. The joint distribution for the largest jump of the Gamma process $X_t \sim \text{Ga}(\alpha(dt))$ in the time interval $(0, t]$ and the time $\tau$ at which it occurs (hint: do $\alpha(dt) = dt$ first);

3. The joint distribution for the largest jump of the Dirichlet process $X_t \sim \text{Dir}(\alpha(dt))$ in the time interval $(0, t]$, and the time $\tau$ at which it occurs.