

Gaussian Processes on \mathbb{R}^n

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1 Isotropic Covariance Functions

Let $\{Z(s)\}$ be a Gaussian process on \mathbb{R}^n , *i.e.*, a collection of jointly normal random variables $Z(s)$ associated with n -dimensional locations $s \in \mathbb{R}^n$. The joint distribution of $\{Z(s)\}$ depends only on the means $\mu(s) = \mathbf{E}Z(s)$ and the covariances $C(s, t) = \mathbf{E}(Z(s) - \mu(s))(Z(t) - \mu(t))$.

The process is called *stationary* or *translation invariant* if the distribution wouldn't change under a rigid translation of the entire collection of locations, *i.e.*, if $\mu(s) = \mu(s+h)$ and $C(s+h, t+h) = C(s, t)$ for all h ; in this case $\mu(s) \equiv \mu$ is constant and $C(s, t) = C(s-t, 0)$ can only depend on the difference $h = (s-t)$ between the two locations, so must be of the form $C(s, t) = C_0(s-t)$ for some function $C_0(h) = C(h, 0)$ on \mathbb{R}^n . Not just any function $C_0(h)$ can be a covariance function; let's see what the choices are.

It's easy to see that the function c must be *even*, *i.e.*, must satisfy $C_0(h) = C_0(-h)$, since $C_0(s-t) = \mathbf{E}(Z(s) - \mu(s))(Z(t) - \mu(t)) = C_0(t-s)$. But more is true: if $\{s_j\}$ any collection of locations, then complex linear combinations $a^\top(Z - \mu) = \sum a_j(Z_j - \mu_j)$ of the centered random variables $Z_j = Z(s_j)$ (with means $\mu_j = \mu(s_j)$) must have nonnegative squared modulus $\mathbf{E} \left| \sum a_j(Z_j - \mu_j) \right|^2 = \sum a_j C_0(s_j - s_k) \bar{a}_k \geq 0$ for every set of complex numbers $\{a_j\} \subset \mathbb{C}$. A function $C_0(h)$ is called *positive semi-definite* if it always satisfies the inequality $\sum_{jk} a_j C_0(s_j - s_k) \bar{a}_k \geq 0$ for any locations s_j and complex numbers a_j ; this is equivalent to asking that $C_0(h) = C_0(-h)$ for every $h \in \mathbb{R}^n$ and that $\sum a_j C_0(s_j - s_k) a_k \geq 0$ for all *real* numbers $a_j \in \mathbb{R}$. One way to get a symmetric positive semi-definite function $C_0(h)$ is by taking the Fourier transform

$$C_0(h) = \int_{\mathbb{R}^n} e^{ih \cdot \omega} G(\omega) d^n \omega = \int_{\mathbb{R}^n} e^{ih \cdot \omega} G(d\omega) \quad (1)$$

of any positive function $G \in L_1(\mathbb{R}^n)$ or, more generally, of any finite positive Borel measure $G(d\omega)$, because then

$$\begin{aligned} \sum_{jk} a_j C_0(s_j - s_k) \bar{a}_k &= \int_{\mathbb{R}^n} \sum_{jk} (a_j e^{is_j \cdot \omega}) \overline{(a_k e^{is_k \cdot \omega})} G(d\omega) \\ &= \int_{\mathbb{R}^n} \left| \sum_j a_j e^{is_j \cdot \omega} \right|^2 G(d\omega) \geq 0. \end{aligned}$$

It turns out that this is the *only* way to get one— that every positive semi-definite function can be written in this form for some finite positive measure $G(d\omega)$, called the *spectral measure*

(if $G(d\omega) = G(\omega) d\omega$ is absolutely continuous, $G(\omega)$ is called the *spectral density*). Known as “Bochner’s Theorem,” this result is really just the Fourier inversion formula in an unfamiliar setting:

$$G(\omega) = (2\pi)^{-n} \int_{\mathbb{R}^n} e^{-ih \cdot \omega} C_0(h) d^n h.$$

Since the process $\{Z(s)\}$ is real-valued, the spectral density $G(\omega) = G(-\omega)$ must be an even function and so we can write

$$\begin{aligned} C_0(h) &= \int_{\mathbb{R}^n} \cos(h \cdot \omega) G(\omega) d^n \omega \\ G(\omega) &= (2\pi)^{-n} \int_{\mathbb{R}^n} \cos(h \cdot \omega) C_0(h) d^n h \end{aligned}$$

If the Gaussian process is also *isotropic*, or invariant under rotations, then $G(\omega) = g(|\omega|)$ must also be invariant under rotations and depend only on the length $r = |\omega|$ of the vector $\omega \in \mathbb{R}^n$. In this case we can simplify these integrals by transforming to polar coordinates.

1.1 Polar Coordinates for Probabilists

Polar coordinates are a familiar tool in two-dimensional integrals, where the change of variables from $x \in \mathbb{R}^2$ to $r = \sqrt{x_1^2 + x_2^2}$ and $\theta = \arctan x_2/x_1$ (so $x_1 = r \cos \theta$, $x_2 = r \sin \theta$) and a change from $d^2 x$ to $r dr d\theta$ lead to simple expressions for the integrals of radial functions. Equivalently, we can let σ have a uniform probability distribution (denoted by $d\sigma$) over the unit circle $S^1 = \{x : x_1^2 + x_2^2 = 1\}$, and change variables from $x = (x_1, x_2)$ to (r, σ) , with $d^2 x = dx_1 dx_2$ replaced by $2\pi r dr d\sigma$.

In three dimensions the first polar approach has its analogue in the Euler angles, while the second is simpler with uniform measure for σ on the unit sphere $S^2 \subset \mathbb{R}^3$, with $d^3 x = dx_1 dx_2 dx_3$ replaced by $4\pi r^2 dr d\sigma$. Notice that $2\pi r$ and $4\pi r^2$ are the circumference of the circle and the area of the sphere of radius r , respectively. In any number n of dimensions the sphere rS^{n-1} of radius r has area $2\pi^{n/2} r^{n-1} / \Gamma(n/2)$ (the derivative w.r.t. r of the volume $\omega_n r^n$ of the ball of radius r , where $\omega_n = \pi^{n/2} / \Gamma(1 + n/2)$ is the volume of the unit ball), and we can again evaluate integrals in polar coordinates with the uniform probability distribution $d\sigma$ for $\sigma \in S^{n-1} \subset \mathbb{R}^n$, and $d^n x = \frac{2\pi^{n/2}}{\Gamma(n/2)} r^{n-1} dr d\sigma$. This makes it easy to compute integrals of radial functions; for functions that also depend on one or more of the components x_j , it is sometimes helpful to note that the squares $\{\sigma_j^2\}$ have a Dirichlet $\text{Di}(\frac{1}{2}, \dots, \frac{1}{2})$ joint distribution, so each σ_j is distributed as the square root of a $\text{Be}(\frac{1}{2}, \frac{n-1}{2})$ random variable.

1.2 Evaluating $C_0(h)$

Switching to polar coordinates $r = |\omega| \geq 0$ and $\sigma = \omega/|\omega| \in S^{n-1}$ (where $d\sigma$ denotes the uniform probability measure on the unit sphere S^{n-1} in \mathbb{R}^n), and noting that the component $\sigma_h = \sigma \cdot h/|h|$ of $\sigma \in S^{n-1}$ in the direction h again has the same distribution as the square root of a $\text{Be}(\frac{1}{2}, \frac{n-1}{2})$ random variable, writing ρ for $|h|$ and $g(|\omega|)$ for $G(\omega)$,

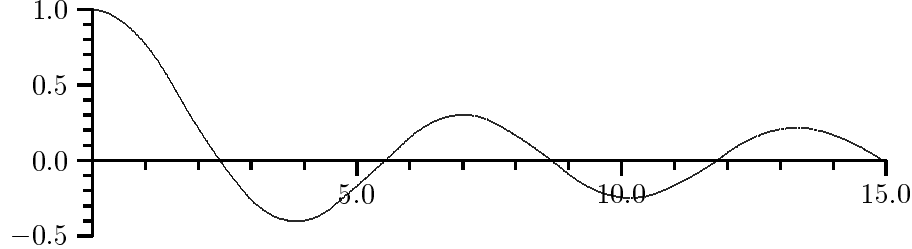
$$\begin{aligned}
C_0(h) &= \int_{\mathbb{R}^n} \cos(h \cdot \omega) g(|\omega|) d^n \omega \\
&= \iint_{\mathbb{R}_+ \times S^{n-1}} \cos(r\rho\sigma_h) g(r) \frac{2\pi^{n/2} r^{n-1}}{\Gamma(n/2)} dr d\sigma \\
&= \int_{\mathbb{R}_+} \int_0^1 \cos(r\rho\sqrt{u}) g(r) \frac{2\pi^{n/2} r^{n-1}}{\Gamma(n/2)} \frac{\Gamma(n/2)}{\Gamma(\frac{1}{2})\Gamma(\frac{n-1}{2})} u^{1/2-1} (1-u)^{(n-1)/2-1} dr du \\
&= \int_0^\infty \rho (2\pi r/\rho)^{\nu+1} J_\nu(r\rho) g(r) dr, \quad \nu \equiv \frac{n}{2} - 1 \tag{2}
\end{aligned}$$

$$\begin{aligned}
&= \int_0^\infty (r\rho/2)^{-\nu} \Gamma(\nu+1) J_\nu(r\rho) \gamma(dr) \tag{3} \\
&= \begin{cases} \int_0^\infty 2 \cos(r\rho) g(r) dr & \text{if } n = 1 \\ \int_0^\infty 2\pi r J_0(r\rho) g(r) dr & \text{if } n = 2 \\ \int_0^\infty \rho(2\pi r/\rho)^{3/2} J_{1/2}(r\rho) g(r) dr & \text{if } n = 3 \end{cases}
\end{aligned}$$

where

$$J_\nu(z) = \frac{(z/2)^\nu}{\sqrt{\pi} \Gamma(\nu + 1/2)} \int_0^\pi \cos(z \cos \theta) \sin(\theta)^{2\nu} d\theta$$

is the Bessel function of the first kind of order ν (see Watson, 1944). Bessel functions aren't as familiar as sines and cosines, but they're common in engineering and physics and are in the standard C library, the GNU Scientific library (GSL), R, Maple and Mathematica, MatLab, Python's SciPy library, *etc.*; for more details, see Abramowitz and Stegun (1964, Chapter 9). Here's a plot of $J_0(z)$:



The plot of $J_0(z)$ looks a little like a sine or cosine, but falls off like $1/\sqrt{z}$ as $z \rightarrow \infty$.

The most general isotropic covariance is given in (3), with the absolutely continuous measure $g(r) \frac{2\pi^{n/2}}{\Gamma(n/2)} r^{n-1} dr$ replaced by an arbitrary positive finite measure $\gamma(dr)$ on $[0, \infty)$. Any isotropic covariance function may be approximated by one with a discrete spectral measure $\gamma(dr) = \sum \gamma_j \delta_{r_j}(dr)$ assigning mass γ_j to finitely many points r_j :

$$\begin{aligned}
C_0(h) = C(\rho) &\approx \sum_j (2/r_j\rho)^\nu \Gamma(\nu+1) J_\nu(r_j\rho) \gamma_j \tag{4} \\
&= \begin{cases} \sum_j \gamma_j \cos(r_j\rho) & \text{if } n = 1 \\ \sum_j \gamma_j J_0(r_j\rho) & \text{if } n = 2 \\ \sum_j \gamma_j \sqrt{\pi/2r_j\rho} J_{1/2}(r_j\rho) & \text{if } n = 3 \end{cases}
\end{aligned}$$

but a more common approach is to choose small parametric families of densities $g^\theta(r)$ or measures $\gamma^\theta(dr)$.

We can recover a spectral density $g(r) = G(\omega)$ (for $r = |\omega|$) through the Fourier inversion formula, using polar coordinates with $\rho = |h| \in \mathbb{R}_+$ and $\sigma = h/|h| \in S^{n-1}$:

$$\begin{aligned} g(r) = G(\omega) &= \frac{1}{(2\pi)^n} \int \cos(-h \cdot \omega) C_0(h) d^n h \\ &= \frac{1}{(2\pi)^n} \iint_{\mathbb{R}_+ \times S^{n-1}} \cos(-r\rho\sigma_\omega) C(\rho) \frac{2\pi^{n/2} \rho^{n-1}}{\Gamma(n/2)} d\rho d\sigma \\ &= \int_0^\infty r(\rho/2\pi r)^{n/2} J_\nu(r\rho) C(\rho) d\rho, \quad \nu \equiv \frac{n}{2} - 1 \end{aligned} \tag{5}$$

$$= \begin{cases} \int_0^\infty \frac{2}{\pi} \cos(r\rho) C(\rho) d\rho & \text{if } n = 1 \\ \int_0^\infty (\rho/2\pi) J_0(r\rho) C(\rho) d\rho & \text{if } n = 2 \\ \int_0^\infty r(\rho/2\pi r)^{3/2} J_{1/2}(r\rho) C(\rho) d\rho & \text{if } n = 3 \end{cases} \tag{6}$$

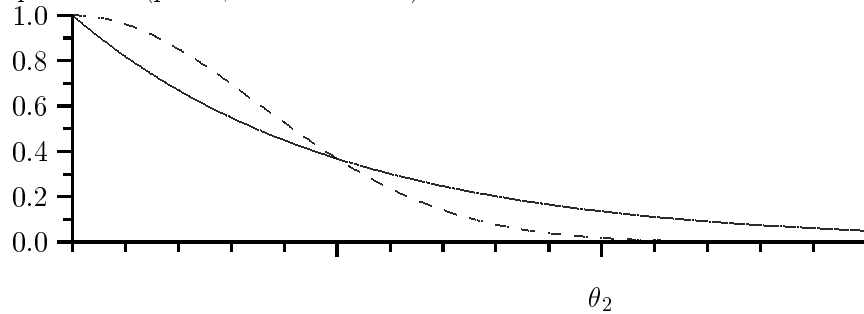
1.3 Parametric Families of Covariance Functions

It is hard to imagine what $C_0(h)$ would look like for different choices of $g(r)$; a simple approach is to take whatever symmetric functions $G(u)$ whose Fourier transforms we can find, and see what we get. Here are some commonly used covariance families, mostly in $n = 2$ dimensions. In each case $\theta_1 = C_0(0)^{1/2}$ is an overall level parameter (in the same units as $Z(s)$), the marginal standard deviation, and θ_2 is a characteristic length scale (in the same units as s):

- Power Exponential family

$$C(\rho | \theta, p) = \theta_1^2 \exp\{-|\rho/\theta_2|^p\}, \quad 0 < p \leq 2$$

Two notable covariograms in this family are the *exponential* ($p = 1$, solid below) and the *squared exponential* ($p = 2$, dashed below):

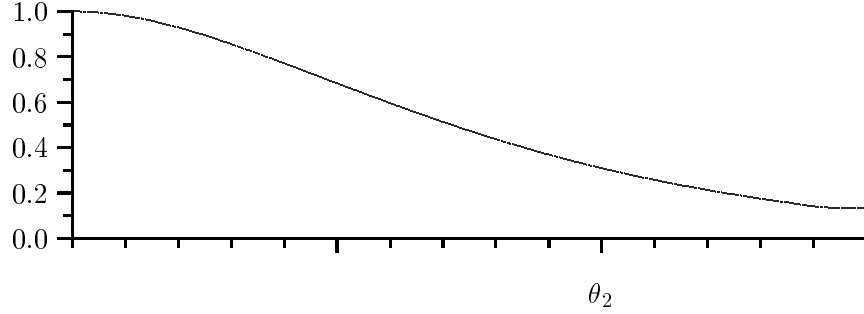


Notice that the exponential has a negative derivative at $z = 0$, so it falls off quickly at first, then slowly levels off, while the squared exponential has zero derivative near $z = 0$ then falls off very quickly. From (6) it follows that the exponential has spectral density function $g(r) = (\theta_1^2 \theta_2^2 / 2\pi) / (1 + r^2 \theta_2^2)^{3/2}$, proportional to a bivariate Cauchy density function, while the squared exponential has spectral density $g(r) = (\theta_1^2 \theta_2^2 / 4\pi) \exp(-r^2 \theta_2^2 / 4)$, proportional to a normal density.

- Matérn Family

$$C(\rho | \theta) = \theta_1^2 \frac{2^{1-\theta_3}}{\Gamma(\theta_3)} \left(\sqrt{2\theta_3} \frac{\rho}{\theta_2} \right)^{\theta_3} K_{\theta_3} \left(\sqrt{2\theta_3} \frac{\rho}{\theta_2} \right) \tag{7}$$

where $K_\nu(z)$ is the modified Bessel function of the third kind of order ν (Abramowitz and Stegun, 1964, §9.6.2):



The displayed plot has shape parameter $\theta_3 = 2$. The Matérn class is quite flexible and includes the exponential family (with $\theta_3 = \frac{1}{2}$), the squared exponential family (in the limit as $\theta_3 \rightarrow \infty$), and many others. Note authors differ in how they parameterize the Matérn class: Rasmussen and Williams (2006, p. 84) use $\nu = \theta_3$ and $\ell = \theta_2\sqrt{8\theta_3}$, for example, while Stein (1999) uses $\nu = \theta_3$ and $\alpha = 1/2\theta_2$.

In n dimensions the Matérn spectral density function is

$$g(r) = \frac{\theta_1^2}{\Gamma(\theta_3)} \left(\frac{\theta_2^2}{2\pi\theta_3} \right)^{n/2} \left(1 + \frac{\theta_2^2}{2\theta_3} r^2 \right)^{-\theta_3 - n/2},$$

proportional to the familiar n -variate Student's t density function with $2\theta_3$ degrees of freedom and scale $\sigma = \sqrt{2}/\theta_2$. This lends more insight into how the Matérn reduces to the exponential when $\theta_3 = 1/2$ and to the squared exponential when $\theta_3 \rightarrow \infty$.

Sample paths of Matérn processes have derivatives of orders less than θ_3 , but not of order θ_3 or higher. Thus all three parameters have clear interpretations: θ_1^2 is the marginal variance, θ_2 is the distance scale, and θ_3 is a smoothness parameter.

Modified Bessel functions K_ν for any $\nu > 0$ are available in most computing environments (GSL, R, Maple, Mathematica, MatLab, SciPy, *etc.*), but for the special case of half-odd-integers $\nu = n + 1/2$ they have simple polynomial expressions of degree n in $(1/z)$ (Abramowitz and Stegun, 1964, §10.2.15). For $\nu > 0$,

$$K_{\pm\nu}(z) = e^{-z} \sqrt{\pi/2z} \sum_{j=0}^n \frac{(j+n)!}{j!(n-j)!} (2z)^{-j}.$$

For example,

$$\begin{aligned} K_{1/2}(z) &= \sqrt{\pi/2z} e^{-z} \\ K_{3/2}(z) &= \sqrt{\pi/2z} e^{-z} [1 + 1/z] \\ K_{5/2}(z) &= \sqrt{\pi/2z} e^{-z} [1 + 3/z + 3/z^2], \end{aligned}$$

leading to simple expressions for the covariance of Eqn (7). These can also be computed from the simple expression for $K_{\pm 1/2} = \sqrt{\pi/2z} e^{-z}$ and the recurrence relation

$$K_{\nu+1}(z) = K_{\nu-1}(z) + \frac{2\nu}{z} K_\nu(z).$$

For this reason it is common to fix parameter $\theta_3 = 3/2$ or $5/2$ in the Matérn class imposing an intentionally selected degree of smoothness, and the estimate θ_1^2 and θ_2 from data. For example, for $\theta_3 = 3/2$ we have

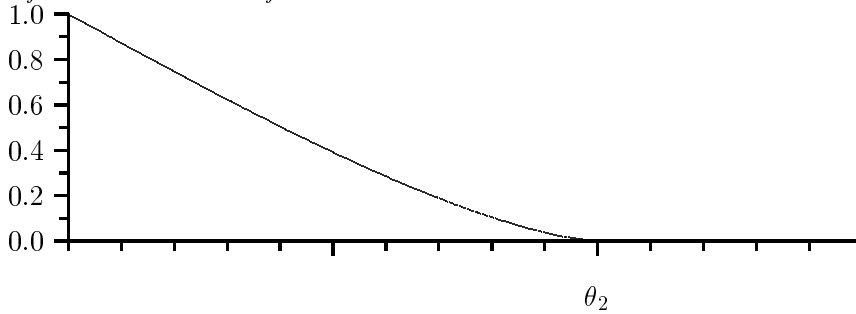
$$C(\rho | \theta) = \theta_1^2 e^{-\rho/\theta_2} [1 + \rho/\theta_2].$$

For an ardent advocate of this class see Stein (1999, §1.6,1.7,2.7,2.10).

- Spherical Family

$$C(\rho | \theta) = \begin{cases} \theta_1^2 \left[1 - \frac{2}{\pi} \left(\frac{\rho}{\theta_2} \sqrt{1 - (\frac{\rho}{\theta_2})^2} + \sin^{-1} \frac{\rho}{\theta_2} \right) \right] & \text{for } \rho < \theta_2 \\ 0 & \text{for } \rho \geq \theta_2 \end{cases}$$

The spherical covariance function is proportional to the area of intersection for two discs of diameter θ_2 with centers separated by distance ρ . In this model the Gaussian quantities Z_j and Z_k at loci s_j and s_k separated by a distance greater than θ_2 will be independent.



This is not quite linear. Like the exponential, it has a negative slope at $z = 0$ and falls off rapidly at first; like the squared exponential, it falls off rapidly later and in fact reaches zero. The spectral density, while available in closed form, isn't illuminating; it's best to think of the spherical process as a convolution or moving average of Gaussian white noise, integrated at each locus over the surrounding ball of diameter θ_2 .

The same idea works in any dimension d : a constant $\theta_1^2/\omega_d(\theta_2/2)^d$, times the volume of intersection of two d -balls of diameter θ_2 whose centers are separated by a distance ρ , is a valid covariance function in dimensions $n \leq d$:

$$\begin{aligned} C(\rho | \theta) &= \frac{2 \theta_1^2}{\omega_d(\theta_2/2)^d} \int_{\rho/2}^{\theta_2/2} \left\{ \omega_{d-1} [(\theta_2/2)^2 - z^2]^{(d-1)/2} \right\} dz \\ &= \theta_1^2 I_{(1-\rho^2/\theta_2^2)} \left(\frac{d+1}{2}, \frac{1}{2} \right), \quad 0 \leq \rho \leq \theta_2 \end{aligned}$$

where

$$\begin{aligned} I_z(\alpha, \beta) &:= \frac{\Gamma(\alpha + \beta)}{\Gamma(\alpha)\Gamma(\beta)} \int_0^z x^{\alpha-1} (1-x)^{\beta-1} dx \\ &= 1 - I_{1-z}(\beta, \alpha) \end{aligned}$$

is the normalized incomplete Beta function (Abramowitz and Stegun, 1964, §6.6.2), the CDF of the $\text{Be}(\alpha, \beta)$ distribution. For $d = 2$ this reduces to the equation above, for $d = 3$ to $C(\rho | \theta) = (\theta_1^2/2)[(\rho/\theta_2)^3 - 3(\rho/\theta_2) + 2]\mathbf{1}_{\{\rho \leq \theta_2\}}$, and for $d = 1$ to $\theta_1^2(1 - \rho/\theta_2)\mathbf{1}_{\{\rho \leq \theta_2\}}$. It even works for non-integer real $d \geq n$.

- Cubic (Integrated Linear)

If X_t is a stationary Gaussian process (in $n = 1$ dimension) with covariance function $\gamma(\rho)$, and $g(t)$ is an integrable function, then $Z_t := \int_{\mathbb{R}} g(t-s)X_s ds = g \star X(t)$ is also a stationary Gaussian process, with covariance $\gamma_g(\rho) = \iint g(s)\gamma(s-t+\rho)g(t) ds dt$, and typically with smoother sample paths than X_t . For example, if X_s has the Spherical covariance with $d = 1$ above and if $g(\rho) = (3/4\theta_2^2)\mathbf{1}_{\{\rho \leq \theta_2\}}$, then $Z = g \star X$ has covariance

$$\begin{aligned} \gamma(\rho) &= \frac{3\theta_1^2}{4\theta_2^2} \int_0^{\theta_2} \int_{\rho}^{\rho+\theta_2} [1 - |s-t|/\theta_2] \mathbf{1}_{\{|s-t| \leq \theta_2\}} dt ds \\ &= \begin{cases} (\theta_1/2)^2 [3(\rho/\theta_2)^3 - 6(\rho/\theta_2)^2 + 4] & 0 \leq \rho \leq \theta_2 \\ (\theta_1/2)^2 (2 - \rho/\theta_2)^3 & \theta_2 \leq \rho \leq 2\theta_2 \\ 0 & 2\theta_2 \leq \rho < \infty, \end{cases} \end{aligned}$$

a compactly-supported piecewise-cubic function of $\rho = |s-t|$ called the *nonnegative cubic covariance* when it is used in computer modeling.

1.4 Mean-Square Continuity & Smoothness

An L_2 process Z_t is called *mean square continuous* at $t \in \mathcal{T}$ if

$$\lim_{s \rightarrow t} \mathbf{E}|Z_s - Z_t|^2 = 0.$$

Evidently this property depends only on the covariance function $C(s, t)$, since

$$\mathbf{E}|Z_s - Z_t|^2 = C(s, s) - C(s, t) - C(t, s) + C(t, t)$$

or, for stationary processes,

$$= 2[C_0(0) - C_0(s-t)].$$

Thus, a stationary L_2 process Z is mean square continuous if and only if its covariance function is continuous at zero. This implies C_0 is uniformly continuous, since

$$\begin{aligned} |C_0(s) - C_0(t)| &= |\mathbf{Cov}[(Z_s - Z_t), Z_0]| \\ &\leq \{\mathbf{V}(Z_s - Z_t)\}^{1/2} \{\mathbf{V}(Z_0)\}^{1/2} \\ &= |C_0(0) - C_0(s-t)|^{1/2} |2C_0(0)|^{1/2}. \end{aligned}$$

Mean-square continuity doesn't imply path continuity— for example, the two-state continuous-time Markov process that alternates between plus and minus one after independent $\text{Ex}(1)$ waiting times has continuous covariance $C_0(h) = e^{-2|h|}$, and paths are almost-surely continuous at each specific time t , but paths do jump with probability one.

A real-valued L_2 process Z_t indexed by $t \in \mathcal{T} \subseteq \mathbb{R}$ is called *mean square differentiable* with derivative Z'_t at $t \in \mathcal{T}$ if

$$\lim_{h \rightarrow 0} \mathbf{E} \left| \frac{Z_{t+h} - Z_t}{h} - Z'_t \right|^2 = 0.$$

This too depends only on the covariance function $C_0(s - t)$ for stationary processes, since the covariance C_h of the process $Z_h(t) := (Z_{t+h} - Z_t)/h$ is

$$\begin{aligned} C_h(t) &= \mathbb{E} \left[\frac{Z_h - Z_0}{h} \frac{Z_{t+h} - Z_t}{h} \right] \\ &= h^{-2} \{2C_0(t) - C_0(t+h) - C_0(t-h)\} \\ &\rightarrow -C_0''(t) \quad \text{as } h \rightarrow 0 \end{aligned}$$

if C_0 is twice differentiable at zero. If so, then Z_t is mean square differentiable, *i.e.*, Z_h converges as $h \rightarrow 0$ to a limiting process Z_t' with covariance

$$\text{Cov}[Z_s', Z_t'] = -C_0''(s - t).$$

Similarly, Z_t is m -times mean square differentiable if and only if C_0 has $2m$ derivatives at zero, and in that case

$$\text{Cov}[Z_s^{(m)}, Z_t^{(m)}] = (-1)^m C_0^{(2m)}(s - t).$$

The same idea works for processes Z_t indexed by vector arguments $t \in \mathcal{T} \subset \mathbb{R}^n$ but now the first derivative (or gradient) Z_t' will be an n -dimensional vector, the second derivative (or Hessian) Z_t'' will be an $n \times n$ matrix, and higher-order derivatives will be higher-order tensors.

Mean-square continuity and differentiability for isotropic random fields can also be taken from the spectral representation of the covariance. From

$$\begin{aligned} C_0(h) &= \int_{\mathbb{R}^n} \cos(h \cdot \omega) g(|\omega|) d\omega, \\ &= \int \int_{\mathbb{R}_+ \times S^{n-1}} \cos(r\rho\sigma_h) g(r) \frac{2\pi^{n/2} r^{n-1}}{\Gamma(n/2)} dr d\sigma \end{aligned}$$

we see that

$$-C_0''(h) = \int_{\mathbb{R}^n} \omega \omega' \cos(h \cdot \omega) g(|\omega|) d\omega$$

is well-defined and finite provided that

$$\int_{\mathbb{R}^n} |\omega|^2 g(|\omega|) d\omega = \int_0^\infty n \omega_n r^{n+1} g(r) dr < \infty$$

and, similarly, Z_t with $t \in \mathcal{T} \subset \mathbb{R}^n$ is m -times mean square differentiable if and only if

$$\int_0^\infty r^{n+2m-1} g(r) dr < \infty.$$

I suspect but don't know that there is a Sobolev-type result that Z_t will have m -times continuously differentiable sample paths if a somewhat stricter condition holds, like $\int_0^\infty r^{n+2m} g(r) dr < \infty$.

In any dimension the covariance function $C(\rho \mid \theta, p)$ of the Power Exponential family of Section (1.3) is continuous but it is not twice differentiable for any $p < 2$, so sample paths are mean-square continuous but not differentiable even once. For $p = 2$ the covariance is infinitely differentiable, so sample paths are m -times mean square differentiable for all m .

The Matérn family in Section (1.3) has spectral density proportional to $(1 + \theta_2^2 r^2 / 2\theta_3)^{-\theta_3 - n/2}$, so it will be m -times mean square differentiable if and only if $r^{n+2m-1-2\theta_3-n}$ is integrable at infinity, *i.e.*, if and only if $\theta_3 > m$, in any dimension n .

Near $\rho \approx 0$ the covariance for the Spherical Family is $C(\rho | \theta) \asymp \theta_1^2(1 - c\rho)$, for constant $c = d\Gamma(d/2)/\theta_2^2\sqrt{\pi}\Gamma((d+1)/2) > 0$, so it is continuous but not differentiable and the process is mean square continuous but not mean square differentiable.

2 Constructing GPs from their Spectral Measures

The integral representations of isotropic covariance functions can be exploited to *construct* random fields with specified covariance structure, in at least two different ways.

2.1 White-Noise Convolutions

Many processes may be constructed similarly as kernel integrals of standard Gaussian white noise,

$$Z(h) = \int_{\mathbb{R}^n} k(h-s) \zeta(ds);$$

where “standard” means that $\mathbf{E}[\zeta(ds)] = 0$ and $\mathbf{E}[\zeta(ds)^2] = ds$ (more formally, that ζ is countably additive with $\zeta(A) \sim \text{No}(0, |A|)$ for Borel $A \subset \mathbb{R}^n$ of finite Lebesgue measure $|A|$). The covariance is

$$C_0(h) = \mathbf{E}[Z(0)\overline{Z(h)}] = \int_{\mathbb{R}^n} k(h-s)\overline{k(-s)} ds$$

with spectral density

$$\begin{aligned} G(\omega) &= (2\pi)^{-n} \int e^{-i\omega \cdot h} C_0(h) dh \\ &= (2\pi)^{-n} \iint e^{-i\omega \cdot h} k(h-s)\overline{k(-s)} ds dh \\ &= (2\pi)^{-n} \left| \int e^{-i\omega \cdot x} k(x) dx \right|^2. \end{aligned}$$

Thus an isotropic kernel may be computed from the spectral density as

$$k(x) = (2\pi)^{-n/2} \int e^{i\omega \cdot x} G(\omega)^{1/2} d^n \omega$$

or, in polar coordinates,

$$\begin{aligned} k(\rho) &= \int_0^\infty r^{\nu+1} \rho^{-\nu} J_\nu(r\rho) g(r)^{1/2} dr \\ &= \begin{cases} \int_0^\infty \sqrt{\frac{2}{\pi}} \cos(r\rho) \sqrt{g(r)} dr & \text{if } n = 1 \\ \int_0^\infty J_0(r\rho) r \sqrt{g(r)} dr & \text{if } n = 2 \\ \int_0^\infty J_{1/2}(r\rho) r^{3/2} \rho^{-1/2} \sqrt{g(r)} d\rho & \text{if } n = 3 \end{cases} \end{aligned}$$

provided that the square root of the spectral density *is* the Fourier transform of a finite positive function, *i.e.*, is itself positive semidefinite. For the Matérn class, the root spectral density $\sqrt{g(r)} \propto$

$(1 + \theta_2^2 r^2 / 2\theta_3)^{-(\theta_3 + n/2)/2}$ will be another n -variate t density provided $\theta_3 > n/2$ and in this case, setting $\epsilon = (2\theta_3 - n)/4 > 0$, we find [check this]

$$k(\rho) = \frac{2\theta_1^2}{(\rho\theta_2\sqrt{2/\theta_3})^{\epsilon+n/2}} \Gamma(\epsilon + n/2) \sqrt{\Gamma(2\epsilon + n/2)} \pi^{n/4} K_\epsilon(\rho\sqrt{2\theta_3/\theta_2})$$

leads to a moving-average kernel representation for the Matérn covariance class. In any number $n \geq 1$ of dimensions the restriction $\epsilon > 0$ entails $\theta_3 > n/2 \geq 1/2$, ruling out the exponential covariance, but the squared exponential covariance (the limiting case as $\theta_3 \rightarrow \infty$) is available in any number of dimensions, with

$$k(\rho) = \theta_1^2 (\pi\theta_2^2 / 8\theta_3)^{-n/2} e^{-4\theta_3 \rho^2 / \theta_2^2}.$$

2.2 Spectral Representation

We generalize Hida and Hitsuda (1993, §III.2, pp. 42 ff). Let $\Theta \subseteq \mathbb{R}^n$ be the *spectrum*, the support of the spectral measure $G(d\omega)$ (see (1)), and let $\zeta(d\omega)$ be a Gaussian random measure on the Borel sets of Ω with control measure $G(d\omega)$ (so $\zeta(A) \sim \text{No}(0, G(A))$ for Borel sets $A \subset \Omega$ of finite G -measure). For $s \in \mathbb{R}^n$ set

$$Z_s := \int_{\Omega} e^{is \cdot \omega} \zeta(d\omega).$$

Then Z is a complex-valued Gaussian stochastic process with mean zero and autocovariance

$$\mathbb{E} Z_s \bar{Z}_t = \int_{\Omega} e^{i(s-t) \cdot \omega} G(d\omega) = C_0(s-t)$$

I think we can arrange for Z to be real-valued, but probably not ζ since if ζ were real then we would have $Z_{-t} \equiv \bar{Z}_t$.

For the special case of $n = 1$, we can start with independent real-valued Gaussian random measures $\zeta_0(d\omega)$ and $\zeta_1(d\omega)$ on $\Omega_+ = \Omega \cap \mathbb{R}_+$, each with control measure $\frac{1}{2}G(d\omega)$ restricted to Ω_+ , and set

$$\zeta(A) = [\zeta_0(A_+) + \zeta_0(A_-)] + i [\zeta_1(A_+) - \zeta_1(A_-)]$$

(where of course $A_+ = A \cap \mathbb{R}_+$ and $A_- = A \cap \mathbb{R}_-$). Then ζ is a complex Gaussian measure on \mathbb{R} with mean zero and covariance

$$\mathbb{E} \zeta(ds) \bar{\zeta}(dt) = \delta_s(dt) G(ds)$$

Set

$$Z_s = \int_{\Omega} e^{i\omega s} \zeta(d\omega) = \int_{\Omega_+} \cos(\omega s) \zeta_0(d\omega) - \int_{\Omega_+} \sin(\omega s) \zeta_1(d\omega) \quad (8)$$

I'm not sure how to think of this in Polar coordinates, or how to extend it to ≥ 2 dimensions, but I expect that's possible useful interesting and, probably, old. The Free Euclidean Field should be an example.

3 Green's Functions

In d dimensions the negative Laplacian is given in polar coordinates by

$$-\Delta f(r, \theta) \equiv - \left\{ \frac{d-1}{r} \frac{\partial f}{\partial r} + \frac{\partial^2 f}{\partial r^2} + \frac{1}{r^2} \Delta_{S^{d-1}} f \right\}$$

where $\Delta_{S^{d-1}}$ denotes the Laplace-Beltrami operator on the $(d-1)$ sphere. The radial solutions to $-\Delta f(r) = 0$ are of the form $f(r) = a + br$ in $d = 1$ dimension, $f(r) = a + b \log r$ in $d = 2$ dimensions, and $f(r) = a + br^{2-d}$ in $d > 2$ dimensions. For $\beta > 0$ the formal solutions to the Green's Function equation $(-\Delta + \beta^2)\gamma(s-t) = \delta_t(s)$ are:

$$\begin{aligned} d = 1 & \quad g(s, t) = \frac{1}{2\beta} e^{-|s-t|} \\ d = 2 & \quad g(s, t) = \frac{1}{2\pi} \log |s-t| \\ d = 3 & \quad g(s, t) = \frac{1}{4\pi} |s-t|^{-1} \\ d > 3 & \quad g(s, t) = \frac{\Gamma(d/2)}{2\pi^{d/2}} |s-t|^{2-d} \end{aligned}$$

4 Compact Ranges

Fourier transforms and harmonic analysis are most commonly studied on Euclidean space \mathbb{R}^d , but many of the ideas extend to bounded spaces like the closed unit ball $B^d \subset \mathbb{R}^d$ or its boundary the unit sphere S^{d-1} (the disc and circle, for $d = 2$), or even to locally compact groups.

In this section we'll consider functions on the unit ball B^d . Because it is compact, functions in $L_2(B^d)$ will have Fourier *series* with a discrete index set, rather than the Fourier *transforms* of Section (1).

4.1 Fourier-Bessel Expansions in $d = 2$ Dimensions

The key for $d = 2$ is the orthogonality relationship for Bessel functions of the first kind (Abramowitz and Stegun, 1964, §9.1):

$$\int_0^1 J_\alpha(rj_{\alpha,n}) J_\alpha(rj_{\alpha,n'}) r dr = \frac{1}{2} \delta_{nn'} [J_{\alpha+1}(j_{\alpha,n})]^2,$$

where $j_{\alpha,n}$ is the n^{th} positive zero of $J_\alpha(x)$, for any $n, n' \in \mathbb{Z}$ and any $\alpha \in \mathbb{R}$. For integers α , $J_\alpha(-x) = (-1)^\alpha J_\alpha(x)$, so $[J_{\alpha+1}(j_{\alpha,n})]^2 = [J_{|\alpha|+1}(j_{\alpha,n})]^2$ and $j_{\alpha,n} = j_{-\alpha,n}$. It follows that for $m \in \mathbb{Z}$ and $n \in \mathbb{N}$ the functions

$$\phi_{mn}(r, \theta) \equiv \frac{1}{\sqrt{\pi} J_{|m|+1}(j_{m,n})} J_m(rj_{m,n}) e^{im\theta} \quad (9)$$

are orthonormal in $L_2(B^2)$. In fact they form a CONS, so any function $f \in L_2(B^2)$ has a convergent expansion

$$f(r, \theta) = \sum_{\substack{m \in \mathbb{Z} \\ n \in \mathbb{N}}} a_{mn} \phi_{mn}(r, \theta) = \sum_{\substack{m \in \mathbb{Z} \\ n \in \mathbb{N}}} \frac{a_{mn}}{c_{mn}} J_m(rj_{m,n}) e^{im\theta}$$

where

$$c_{mn} = \sqrt{\pi} J_{|m|+1}(j_{m,n}), \quad a_{mn} = \int_{B^2} f(r, \theta) \overline{\phi_{mn}(r, \theta)} r dr d\theta. \quad (10)$$

Note a_{mn} from (10) can also be evaluated in *Cartesian* coordinates if, for example, f is known on a rectangular grid containing B^2 — just multiply the integrand by $\mathbf{1}_{B^2}$ and replace $r dr d\theta$ with $dx dy$. Both the Bessel functions J_m and the zeros $j_{m,n}$ are available in \mathbb{R} . From Bessel's differential equation

$$z^2 \frac{\partial^2}{\partial z^2} J_m(z) + z \frac{\partial}{\partial z} J_m(z) + (z^2 - m^2) J_m(z) = 0$$

(Abramowitz and Stegun, 1964, §9.1.1), ϕ_{mn} is an eigenfunction of the Laplacian with eigenvalue $-j_{m,n}^2$ satisfying Dirichlet boundary conditions $\phi_{mn}(1, \theta) \equiv 0$ on $S^1 = \partial B^2$, so the negative Laplacian of f is given in polar coordinates by

$$\begin{aligned} -\Delta f(r, \theta) &\equiv - \left\{ \frac{1}{r} \frac{\partial f}{\partial r} + \frac{\partial^2 f}{\partial r^2} + \frac{1}{r^2} \frac{\partial^2 f}{\partial \theta^2} \right\} \\ &= \sum (j_{m,n})^2 (a_{mn}/c_{mn}) J_m(r j_{m,n}) e^{im\theta}. \end{aligned}$$

The squared L_2 and Sobolev H_1 norms of f are

$$\langle f, f \rangle = \sum |a_{mn}|^2, \quad \|f\|_1^2 = \langle (-\Delta + I)f, f \rangle = \sum (1 + j_{m,n}^2) |a_{mn}|^2,$$

so one reasonable scale-invariant measure of roughness would be

$$\frac{\langle \nabla f, \nabla f \rangle}{\langle f, f \rangle} = \frac{\langle -\Delta f, f \rangle}{\langle f, f \rangle} = \frac{\sum j_{m,n}^2 |a_{mn}|^2}{\sum |a_{mn}|^2}$$

or, similarly but more generally, the ratio of squared Sobolev norms

$$\frac{\|f\|_s^2}{\|f\|_0^2} = \frac{\langle (-\Delta + I)^s f, f \rangle}{\langle f, f \rangle} = \frac{\sum (1 + j_{m,n}^2)^s |a_{mn}|^2}{\sum |a_{mn}|^2}$$

for any $s > 0$. Another measure depending only on angular roughness would be

$$\frac{\sum |m|^{2s} |a_{mn}|^2}{\sum |a_{mn}|^2}$$

for any $s > 0$. Similar results hold for $d = 3$ (using the spherical Bessel function j_α) or for other boundary conditions. Each of these measures is the moment of some “roughness” aspect, for the discrete probability distribution with masses proportional to the squared Fourier coefficients $|a_{mn}|^2$.

4.2 Fourier Expansions in $d = 1$ Dimension

This section has a brief review of one-dimensional Sobolev spaces on an interval, for context and intuition.

Let $\mathcal{T} = [0, T]$ be a closed interval in \mathbb{R}^1 . The functions

$$\phi_n(x) = \sin(n\pi x/T) \sqrt{2/T}$$

are a CONS in $L_2(\mathcal{T}, dx)$ of eigenfunctions of $-\Delta$ (with eigenvalues $n^2\pi^2/T^2$), with Dirichlet boundary conditions at $\partial\mathcal{T}$, so every function $f \in L_2(\mathcal{T}, dx)$ has a convergent expansion of the form

$$f(x) = \sum_{n \in \mathbb{N}} a_n \phi_n(x), \tag{11}$$

with coefficients

$$a_n := \int_{\mathcal{T}} f(x) \phi_n(x) dx.$$

For $s \in \mathbb{R}$ the closed linear span H_s of $\{\phi_n\}$ in the norm

$$\begin{aligned} \|f\|_s &:= \left\{ \sum |a_n|^2 (1 + n^2 \pi^2 / T^2)^s \right\}^{\frac{1}{2}} \\ &= \langle (-\Delta + I)^s f, f \rangle^{\frac{1}{2}} \end{aligned} \quad (12)$$

is a Banach space (complete separable normed metric space) whose dual space is H_{-s} . The self-dual Hilbert space H_0 is just $L_2(\mathcal{T}, dx)$.

Differentiating (11) gives

$$\begin{aligned} \|f'\|_0^2 &= \left\| \sum_{n \in \mathbb{N}} a_n \frac{n\pi}{T} \cos(n\pi \cdot / T) \sqrt{2/T} \right\|_0^2 \\ &= \sum_{n, m \in \mathbb{N}} a_n a_m \frac{nm\pi^2}{T^2} \frac{2}{T} \int_0^T \cos(n\pi x/T) \cos(m\pi x/T) dx \\ &= \sum_{n \in \mathbb{N}} a_n^2 \frac{n^2 \pi^2}{T^2} = \sum_{n \in \mathbb{N}} |a_n|^2 n^2 \pi^2 / T^2 \\ &= \|f\|_1^2 - \|f\|_0^2, \end{aligned}$$

so $f \in H_1$ if and only if $f \in H_0$ and $f' \in H_0$. Similarly one can show that $f \in H_s \Rightarrow f' \in H_{s-1}$ for any $s \in \mathbb{R}$. For $s > 0$, H_s consists of those functions with s derivatives in L_2 , while H_{-s} consists of s^{th} derivatives of L_2 functions.

In particular, for $t \in \mathcal{T}$ the Dirac delta function $\delta_t \in H_{-1}$ but $\delta_t \notin H_0$. For example, for $T = \pi$ and $t = \pi/2$, $f \equiv \delta_t$ has coefficients

$$a_n \equiv \int_0^\pi f(x) \sin(nx) \sqrt{2/\pi} dx = \begin{cases} \sqrt{2/\pi} & n \equiv 1 \pmod{4} \\ -\sqrt{2/\pi} & n \equiv 3 \pmod{4} \end{cases},$$

otherwise zero, with eigenvalues n^2 , so

$$\|\delta_t\|_s^2 = \frac{2}{\pi} \sum_{k=1}^{\infty} (1 + (2k+1)^2)^s.$$

This is only finite for $s < -\frac{1}{2}$, so $\delta_t \in H_s$ only for $s < -\frac{1}{2}$.

The boundary $\partial\mathcal{T}$ and boundary conditions play a subtle rôle in $H_s(\mathcal{T})$. For example, in the interior of $\mathcal{T} = [0, T]$ for $T > 0$ the function

$$f(x) = x(T - x)$$

satisfies $f''(x) \equiv -2$ and hence

$$(-\Delta + I)f(x) = x(T - x) + 2$$

and so, at least for non-negative integers s ,

$$(-\Delta + I)^s f(x) = x(T - x) + 2s.$$

One might expect $f \in H_s$ for all $s > 0$, with

$$\begin{aligned} \|f\|_s^2 &\equiv \langle (-\Delta + I)^s f, f \rangle \\ &\stackrel{?}{=} \int_0^T [x(T - x) + 2s][x(T - x)] dx \\ &= T^5/30 + sT^3/3, \end{aligned}$$

suggesting (wrongly) that $f \in H_s$ for all $s > 0$. In fact the coefficients are

$$a_n := \int_{\mathcal{T}} f(x) \phi_n(x) dx = 4\sqrt{2}T^{5/2}/(n\pi)^3$$

for odd n , and zero for even n , so

$$\begin{aligned} \|f\|_s^2 &\equiv \sum_{n=1}^{\infty} |a_n|^2 (n^2\pi^2/T^2 + 1)^s \\ &= \frac{32T^5}{\pi^6} \sum_{k=1}^{\infty} \frac{((2k+1)^2\pi^2/T^2 + 1)^s}{(2k+1)^6} \end{aligned}$$

which is only finite for $s < 5/2$. The problem is that even the first derivative f' fails to satisfy Dirichlet boundary conditions $f(\partial\mathcal{T}) = 0$. Integrating by parts,

$$\begin{aligned} \int_{\mathcal{T}} f'(x) f'(x) dx &= f'(x)f(x)|_0^T - \int_{\mathcal{T}} f''(x) f(x) dx \\ &= - \int_{\mathcal{T}} f''(x) f(x) dx = \langle (-\Delta)f, f \rangle \\ \int_{\mathcal{T}} f''(x) f''(x) dx &= f''(x)f'(x)|_0^T - \int_{\mathcal{T}} f'''(x) f'(x) dx \\ &= 4T - f'''(x)f(x)|_0^T + \int_{\mathcal{T}} f''''(x) f(x) dx \\ &= \langle (-\Delta)^2 f, f \rangle + 4T \end{aligned}$$

so for $s \geq 2$ a non-zero boundary term enters the integration by parts and the naive calculation fails. Intuitively, what happens is that $f(x)$ takes the constant value of zero outside \mathcal{T} and so derivatives beyond the first pick up something like a Dirac delta function *at the boundary*.

Something similar happens in $d = 2$ dimensions for functions f that satisfy Dirichlet boundary conditions at ∂B^2 but whose derivatives do not, such as

$$f(x, y) = (1 - x^2 - y^2) \vee 0$$

for which $-\Delta f \equiv 4$ and hence $(-\Delta + I)^s f = f + 4s$ in the interior. Still, $f \in H_s$ only for $s < 5/2$ because by circular symmetry its Fourier-Bessel coefficients are $a_{mn} = 0$ for $m \neq 0$, while

$$a_{0n} = \frac{2\pi}{\sqrt{\pi} J_1(j_{0,n})} \int_0^1 (r - r^3) J_0(r j_{0,n}) dr = \frac{8\sqrt{\pi}}{j_{0,n}^3}, \quad \text{so}$$

$$\|f\|_s^2 = 64\pi \sum_{n=1}^{\infty} (1 + \lambda_{0n}^2)^s j_{0,n}^{-6},$$

which is only finite for $s < 5/2$ since $j_{0,n} \asymp n\pi$.

5 Stochastic Expansions

5.1 Brownian Bridge on One Dimensional Interval

As in Sec. 4.2, let

$$\phi_n(x) = \sin(n\pi x/T) \sqrt{2/T}$$

be a CONS in $L_2(\mathcal{T}, dx)$ for $\mathcal{T} = [0, T] \subset \mathbb{R}^1$, and consider the *Brownian Bridge* Gaussian stochastic process X_t with mean and covariance

$$\mathbb{E}X_t \equiv 0 \quad \text{Cov}(X_s, X_t) = \mathbb{E}X_s X_t = (s \wedge t)(1 - (s \vee t)/T),$$

identical to unit-rate Brownian motion conditioned to reach $X_T = 0$. The Fourier coefficients for a random path $\{X_t\}$ are Gaussian random variables

$$A_n := \int_{\mathcal{T}} Z_t \phi_n(s) dt$$

with mean zero and covariance

$$\begin{aligned} \mathbb{E}A_m A_n &= \frac{2}{T^2} \int_{\mathcal{T}^2} (s \wedge t)(T - (s \vee t)) \sin(m\pi s/T) \sin(n\pi t/T) ds dt \\ &= \begin{cases} 0 & m \neq n \\ \frac{T^2}{n^2\pi^2} & m = n \end{cases} \end{aligned}$$

so the $\{A_n\} \stackrel{\text{ind}}{\sim} \text{No}(0, (\frac{T}{n\pi})^2)$ are independent and the expected squared H_s norm of the sample path $s \rightsquigarrow Z_s$ is

$$\begin{aligned} \mathbb{E}\|Z_{\bullet}\|_s^2 &= \sum_{n \in \mathbb{N}} \frac{T^2(1 + n^2\pi^2/T^2)^s}{n^2\pi^2} \\ &< \infty \text{ if and only if } s < \frac{1}{2}, \end{aligned}$$

with $\mathbb{E}\|Z_{\bullet}\|_0^2 = T^2/6$. Roughly speaking, Brownian motion paths have just under one-half of a derivative.

Beginning with iid $\zeta_n \stackrel{\text{iid}}{\sim} \text{No}(0, 1)$ one can *construct* a Brownian Bridge process as the sum

$$Z_t := \sum_{n=1}^{\infty} \frac{T}{n\pi} \zeta_n \phi_n(t)$$

or fractional derivatives of it by

$$Z_t^\sigma := \sum_{n=1}^{\infty} \left[\frac{T}{n\pi} \right]^{1-\sigma} \zeta_n \phi_n(t),$$

which will be centered Gaussian processes with about $\frac{1}{2}-\sigma$ derivatives— smooth paths, for $\sigma \ll 0$ or, for $\sigma > 0$, generalized stochastic processes that only make sense as convolutions $Z[\psi] = \int Z_t \psi(t) dt$ for $\psi \in H_s$ with $s > \sigma - \frac{1}{2}$.

5.2 Stochastic Expansions in Two Dimensions

As in Section (4.1) let

$$\phi_{mn}(r, \theta) \equiv \frac{1}{\sqrt{\pi} J_{|m|+1}(j_{m,n})} J_m(r j_{m,n}) e^{im\theta} \quad (9)$$

be a CONS in $L_2(B^2)$, and let $\zeta_{mn} \stackrel{\text{iid}}{\sim} \text{No}(0, 1)$ for $m \in \mathbb{Z}$, $n \in \mathbb{N}$. Fix $\sigma \in \mathbb{R}$ and set

$$Z^\sigma(r, \theta) := \sum [j_{m,n}]^\sigma \zeta_{mn} \phi_{mn}(r, \theta)$$

for $(r, \theta) \in B^2$. The expected squared H_s norm of Z^σ is

$$\begin{aligned} \mathbf{E}\langle Z^\sigma, Z^\sigma \rangle_s &= \mathbf{E}\langle (-\Delta + I)^s Z^\sigma, Z^\sigma \rangle \\ &= \sum (j_{m,n}^2 + 1)^s j_{m,n}^\sigma \end{aligned}$$

Since $j_{m,n} \asymp \pi(n + |m|/2 - 1/4)$ (Abramowitz and Stegun, 1964, §9.5.12), this is finite by the integral test if and only if

$$\infty > \int_1^\infty (r^2 + 1)^s r^\sigma r dr,$$

so $Z^\sigma \in H_s$ almost-surely for $\sigma < -2(s + 1)$ — or $\sigma < -2$ for $s = 0$.

5.3 Reproducing Kernel Hilbert Space (RKHS)

Now let X_s be an arbitrary mean-zero Gaussian process with (positive definite) covariance kernel

$$k(s, t) = \mathbf{E}X_s X_t$$

for $s, t \in \mathcal{T}$, and consider the Hilbert space \mathcal{H} consisting of limits of linear combinations of the $\{X_{s_j}\}$ for $\{s_j\} \subset \mathcal{T}$. For tame enough functions $f, g : \mathcal{T} \rightarrow \mathbb{R}$ this space will include elements $X_f \equiv \int_{\mathcal{T}} f(s) X_s ds$ and $X_g \equiv \int_{\mathcal{T}} g(t) X_t dt$, each Gaussian with mean zero and covariance

$$\langle X_f, X_g \rangle \equiv \int_{\mathcal{T}^2} f(s) k(s, t) g(t) ds dt,$$

the inner-product in $L_2(\mathcal{T})$ of f and Kg for the operator $K : L_2(\mathcal{T}) \rightarrow L_2(\mathcal{T})$ with kernel k ,

$$Kg(s) \equiv \int_{\mathcal{T}} k(s, t) g(t) dt.$$

If $k \in L_2(\mathcal{T}^2)$ that operator is Hilbert-Schmidt, with a complete orthonormal set $\{\phi_n\}$ of eigenfunctions

$$K \phi_n(s) \equiv \int_{\mathcal{T}} k(s, t) \phi_n(t) dt = \lambda_n \phi_n(s)$$

whose eigenvalues are square-summable, with

$$\int_{\mathcal{T}^2} k(s, t)^2 ds dt = \sum \lambda_n^2 < \infty.$$

In this case for iid $\{\zeta_n\} \stackrel{\text{iid}}{\sim} \text{No}(0, 1)$ the series

$$Z_s \equiv \sum \lambda_n \zeta_n \phi_n(s)$$

converges almost-surely in L_2 to a Gaussian process with mean zero and covariance k , the so-called *Karhunen-Loève* expansion. The example of Sec. 5.1 was an example of this. As in that example, often k is the Greens function for a differential operator (here $-\partial^2/\partial x^2$ with Dirichlet boundary conditions), and the eigenfunctions may be found as solutions to differential equations.

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