## Posterior distribution

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Notation 1. When A is a square matrix, we denote by |A| its determinant. If the inverse of A exist, we denote it by  $A^{-1}$ .

# 1 Distribution of the posterior of a finite basis expansion with Gaussian coefficients

**Lemma 2.** Let  $X^T = (X_t : t \in [0,T])$  be an observation of

$$dX_t = b(X_t)dt + \sigma(X_t)dW_t,$$

where  $\sigma: \mathbb{R} \to \mathbb{R}_{>0}$  is a measurable function,  $(W_t: t \in [0,T])$  is a Brownian motion and b is equipped with the prior distribution defined by

$$b = \sum_{j=1}^{k} \theta_j \phi_j,$$

where  $\{\phi_1, \ldots, \phi_k\}$  is a linearly independent basis, and  $\theta = (\theta_1, \ldots, \theta_k)^t$  has multivariate normal distribution  $N(\mu, \Sigma)$ , with mean vector  $\mu$  and positive definite matrix  $\Sigma$ . Then the posterior distribution of  $\theta$  given  $X^T$  is  $N(\hat{\mu}, \hat{\Sigma})$ , where

$$\hat{\mu} = (S + \Sigma^{-1})^{-1} (m + \Sigma^{-1} \mu), \quad \hat{\Sigma} = (S + \Sigma^{-1})^{-1}$$

and the vector  $m = (m_1, \ldots, m_k)^t$  is defined by

$$m_l = \int_0^T \frac{\phi_l(X_t)}{\sigma(X_t)^2} dX_t, \quad l = 1, \dots, k,$$

and the symmetric  $k \times k$ -matrix S is given by

$$S_{l,l'} = \int_0^T \frac{\phi_l(X_t)\phi_{l'}(X_t)}{\sigma^2(X_t)} dt, \quad l, l' = 1, \dots, k,$$
 (1)

provided  $S + \Sigma^{-1}$  is invertible. Moreover, the marginal likelihood is given by

$$\int p(X^T \mid \theta) p(\theta) d\theta = |\Sigma^{-1} \hat{\Sigma}|^{1/2} e^{-\frac{1}{2}\mu^t \Sigma^{-1} \mu} e^{\frac{1}{2}\hat{\mu}^t \hat{\Sigma}^{-1} \hat{\mu}}.$$

*Proof.* Almost surely we have by Girsanov's theorem (e.g. Steele, 2001, chapter 13 or Chung and Williams, 1990 reprint 2014, section 9.4)

$$p(X^T \mid \theta) = \exp\left(\int_0^T \frac{b(X_t)}{\sigma(X_t)^2} dX_t - \frac{1}{2} \int_0^T \left(\frac{b(X_t)}{\sigma(X_t)}\right)^2 dt\right),\tag{2}$$

with respect to the Wiener measure. So

$$\log p(X^T \mid b) = \theta^t m - \frac{1}{2} \theta^t S \theta \tag{3}$$

and the log of the distribution of  $\theta$  with respect to the Lebesgue measure on  $\mathbb{R}^k$  is given by

$$\log p(\theta) = -\frac{k}{2}\log(2\pi) - \frac{1}{2}\log|\Sigma| - \frac{1}{2}(\theta - \mu)^t \Sigma^{-1}(\theta - \mu)$$
$$= C_1 - \frac{1}{2}\theta\Sigma^{-1}\theta + \theta^t \Sigma^{-1}\mu,$$

with

$$C_1 = -\frac{k}{2}\log(2\pi) - \frac{1}{2}\log|\Sigma| - \frac{1}{2}\mu^t \Sigma^{-1}\mu.$$

So,

$$\begin{split} \log(p(X^T \mid \theta)p(\theta)) = & C_1 + \theta^t m - \frac{1}{2}\theta^t S\theta - \frac{1}{2}\theta\Sigma^{-1}\theta + \theta^t\Sigma^{-1}\mu \\ = & C_1 + \theta^t (m + \Sigma^{-1}\mu) - \frac{1}{2}\theta^t (S + \Sigma^{-1})\theta \\ = & C_1 + \theta^t (S + \Sigma^{-1}) \Big( (S + \Sigma^{-1})^{-1} (m + \Sigma^{-1}\mu) \Big) \\ & - \frac{1}{2}\theta^t (S + \Sigma^{-1})\theta. \end{split}$$

By the Bayes formula, the posterior density of  $\theta$  is proportional to  $p(X^T \mid \theta)p(\theta)$ . It follows that  $\theta \mid X^T$  is normally distributed with mean

$$\hat{\mu} := (S + \Sigma^{-1})^{-1} (m + \Sigma^{-1} \mu).$$

and covariance matrix

$$\hat{\Sigma} := (S + \Sigma^{-1})^{-1},$$

provided  $S + \Sigma^{-1}$  is invertible. Moreover

$$\begin{split} & \int p(X^T \mid \theta) p(\theta) d\theta \\ &= \int e^{C_1} e^{\theta^t \hat{\Sigma}^{-1} \hat{\mu}} e^{-\frac{1}{2} \theta^t \hat{\Sigma}^{-1} \theta} d\theta \\ &= (2\pi)^{k/2} |\hat{\Sigma}|^{1/2} e^{\frac{1}{2} \hat{\mu}^t \hat{\Sigma}^{-1} \hat{\mu}} e^{C_1} \\ & \times \int (2\pi)^{-k/2} |\hat{\Sigma}|^{-1/2} e^{\theta^t \hat{\Sigma}^{-1} \hat{\mu}} e^{-\frac{1}{2} \theta^t \hat{\Sigma}^{-1} \theta} e^{-\frac{1}{2} \hat{\mu}^t \hat{\Sigma}^{-1} \hat{\mu}} d\theta \\ &= (2\pi)^{k/2} |\hat{\Sigma}|^{1/2} e^{\frac{1}{2} \hat{\mu}^t \hat{\Sigma}^{-1} \hat{\mu}} e^{C_1} \\ &= |\hat{\Sigma}^{-1} \hat{\Sigma}|^{1/2} e^{-\frac{1}{2} \mu^t \hat{\Sigma}^{-1} \mu} e^{\frac{1}{2} \hat{\mu}^t \hat{\Sigma}^{-1} \hat{\mu}}, \end{split}$$

using that the integrant in the third last line is the density of a multivariate normal distribution and therefore integrates to one.  $\Box$ 

Usually we refer to S as the Girsanov matrix.

## 2 The marginal maximum likelihood estimator

**Lemma 3.** Let  $\lambda > 0$ ,  $\mu \in \mathbb{R}^k$  and let  $\Sigma$  be a positive definite  $k \times k$ -matrix. Consider the prior  $\theta \sim N(\mu, \Sigma_{\lambda})$ , where  $\Sigma_{\lambda} = \lambda^2 \Sigma$  and denote its density by  $p_{\lambda}$ . Then

$$\log \int p_{\lambda}(X^{T} \mid \theta) p_{\lambda}(\theta) d\theta$$

$$= -\frac{1}{2} \log |\lambda^{2} \Sigma S + \mathbb{I}_{k}| -\frac{1}{2} \mu^{t} \Sigma^{-1} \mu + \frac{1}{2} (m + \lambda^{-2} \Sigma^{-1} \mu)^{t} (S + \lambda^{-2} \Sigma^{-1})^{-1} (m + \lambda^{-2} \Sigma^{-1} \mu).$$
(4)

*Proof.* It follows from lemma 2 that

$$\Sigma_{\lambda}\hat{\Sigma}_{\lambda}^{-1} = \Sigma_{\lambda}(S + \Sigma_{\lambda}^{-1}) = \Sigma_{\lambda}S + \mathbb{I}_{k} = \lambda^{2}\Sigma S + \mathbb{I}_{k}$$

and

$$\begin{split} \hat{\mu}^t \hat{\Sigma}_{\lambda}^{-1} \hat{\mu} = & (m + \Sigma_{\lambda}^{-1} \mu)^t (S + \Sigma_{\lambda}^{-1})^{-1} (S + \Sigma_{\lambda}^{-1}) (S + \Sigma_{\lambda}^{-1})^{-1} (m + \Sigma_{\lambda}^{-1} \mu) \\ = & (m + \lambda^{-2} \Sigma^{-1} \mu)^t (S + \lambda^{-2} \Sigma^{-1})^{-1} (m + \lambda^{-2} \Sigma^{-1} \mu). \end{split}$$

So it follows from the same lemma that

$$\begin{split} & \log \int p_{\lambda}(X^{T} \mid \theta) p_{\lambda}(\theta) d\theta \\ & = -\frac{1}{2} \log |\lambda^{2} \Sigma S + \mathbb{I}_{k}| - \frac{1}{2} \lambda^{-2} \mu^{t} \Sigma^{-1} \mu + \frac{1}{2} (m + \lambda^{-2} \Sigma^{-1} \mu)^{t} (S + \lambda^{-2} \Sigma^{-1})^{-1} (m + \lambda^{-2} \Sigma^{-1} \mu). \end{split}$$

So can we calculate  $(S + \lambda^{-2}\Sigma^{-1})^{-1}$  from  $(S + \Sigma^{-1})^{-1}$ ? What I found out: if A and B are symmetric matrices that commute, then there is an orthonormal matrix Q so that  $D_A = Q^T A Q$  and  $D_B = Q^T B Q$  are diagonal. In our set-up this happens when S and S and S and S are commute. They commute when S is S.

In de implementatie voor vaste  $\alpha$  kun je  $\mu^t \Sigma^{-1} \mu$  en  $\Sigma^{-1} \mu$  opslaan en hoef je maar een keer uit te rekenen.

Als  $\mu = 0$ , dan is

$$\log \int p_{\lambda}(X^{T} \mid \theta) p_{\lambda}(\theta) d\theta$$
$$= -\frac{1}{2} \log |\lambda^{2} \Sigma S + \mathbb{I}_{k}| + \frac{1}{2} m^{t} (S + \lambda^{-2} \Sigma^{-1})^{-1} m.$$

Verder hebben we

$$S + \lambda^{-2} \Sigma^{-1} = \lambda^{-2} \Sigma^{-1} \left( \lambda^2 \Sigma S + I_k \right).$$

Dus

$$\log \int p_{\lambda}(X^{T} \mid \theta) p_{\lambda}(\theta) d\theta$$
$$= -\frac{1}{2} \log |\lambda^{2} \Sigma S + \mathbb{I}_{k}| + \frac{1}{2} \lambda^{2} m^{t} \left(\lambda^{2} \Sigma S + I_{k}\right)^{-1} \Sigma m.$$

Dus de laatste formule hangt niet af van  $\Sigma^{-1}$ . De vraag is dus, zijn er slimme snelle manieren om de determinant en inverse van  $\lambda^2 \Sigma S + I_k$  uit te rekenen? Conclusie van 3 dagen aan werken is dat de determinant makkelijk uitgerekend kan worden met behulp van de eigenwaarden, maar de inverse naar het schijnt niet zo makkelijk.

**Lemma 4.** If  $\nu_1, \ldots, \nu_k$  are the eigenvalues of  $\Sigma S + I_k$ , then  $\lambda^2 \nu_1 - \lambda^2 + 1, \ldots, \lambda^2 \nu_k - \lambda^2 + 1$  are the eigenvalues of  $\lambda^2 \Sigma S + \mathbb{I}_k$ .

*Proof.* Note that

$$\begin{aligned} 0 &= |\nu_i \, \mathbb{I}_k - (\Sigma S + \mathbb{I}_k)| \\ &\Leftrightarrow \\ 0 &= \left| \lambda^2 \nu_i \, \mathbb{I}_k - (\lambda^2 \Sigma S + \lambda^2 \, \mathbb{I}_k) \right| \\ &= \left| (\lambda^2 \nu_i - \lambda^2 + 1) \, \mathbb{I}_k - (\lambda^2 \Sigma S + \mathbb{I}_k) \right|. \end{aligned}$$

So  $\nu_i$  is an eigenvalue of  $\Sigma S + \mathbb{I}_k$  if and only if  $\lambda^2 \nu_i - \lambda^2 + 1$  is an eigenvalue of  $\lambda^2 \Sigma S + \mathbb{I}_k$ .  $\square$ 

**Lemma 5.** If  $\nu_1, \ldots, \nu_k$  are the eigenvalues of  $\Sigma S$ , then  $\lambda^2 \nu_1 + 1, \ldots, \lambda^2 \nu_k + 1$  are the eigenvalues of  $\lambda^2 \Sigma S + I_k$ .

*Proof.* Note that

$$\begin{aligned} |\nu_i \, \mathbb{I}_k - \Sigma S| &= 0 \\ \Leftrightarrow \\ 0 &= \left| \lambda^2 \nu_i \, \mathbb{I}_k - \lambda^2 \Sigma S \right| \\ &= \left| (\lambda^2 \nu_i + 1) \, \mathbb{I}_k - (\lambda^2 \Sigma S + \mathbb{I}_k) \right| \end{aligned}$$

So the eigenvalues of  $\lambda^2 \Sigma S + \mathbb{I}_k$  are easily obtained from the eigenvalues of  $\Sigma S$  or  $\Sigma S + \mathbb{I}_k$ . Note that the determinant

## 3 Random scaling

**Lemma 6.** Let  $X^T = (X_t : t \in [0,T])$  be an observation of

$$dX_t = b(X_t)dt + \sigma(X_t)dW_t,$$

where b is equipped with the prior distribution defined by

$$\lambda^{2} \sim Inverse \ Gamma(A, B) = IG(A, B)$$
  
$$\theta \mid \lambda \sim N(\mu, \lambda^{2}\Sigma)$$
  
$$b \mid \theta = \sum_{j=1}^{k} \theta_{j}\phi_{j},$$

where  $\{\phi_1, \ldots, \phi_k\}$  is a linearly independent basis. Then

$$\lambda^2 \mid \theta, X^T \sim IG\left(A + k/2, B + \frac{1}{2}(\theta - \mu)^t \Sigma^{-1}(\theta - \mu)\right).$$

*Proof.* Recall eq. (3),  $\log p(X^T \mid b) = \theta^t m - \frac{1}{2}\theta^t S\theta$ . The logarithm of the distribution of  $\theta$  given  $\lambda$  with respect to the Lebesgue measure on  $\mathbb{R}^k$  is given by (proportionality w.r.t.  $\lambda$ ),

$$\log p(\theta \mid \lambda) = C_1 - k \log \lambda - \frac{1}{2} \lambda^{-2} (\theta - \mu)^t \Sigma^{-1} (\theta - \mu).$$

for some real constant  $C_1$ , depending on  $\theta$ , but not on  $\lambda$ .

In the following,  $\propto$  means equal up to a multiplicative constant depending on  $\theta$  and  $X^T$ , but not on  $\lambda$ . By the Bayes formula,

$$p(\lambda^2 \mid \theta, X^T) \propto p(X^T \mid \lambda^2, \theta) p(\lambda^2 \mid \theta)$$

and

$$p(\lambda^2 \mid \theta) \propto p(\theta \mid \lambda^2) p(\lambda^2)$$

SO

$$p(\lambda^2 \mid \theta, X^T) \propto p(X^T \mid \lambda^2, \theta) p(\theta \mid \lambda^2) p(\lambda^2)$$

It follows that for some real constants  $C, \tilde{C}$  depending on  $\theta$  and  $X^T$ , but not on  $\lambda$ , we have

$$\begin{split} &\log p(\lambda^2 \mid \theta, X^T) \\ = &C + \theta^t m - \frac{1}{2} \theta^t S \theta \\ &- k \log \lambda - \frac{1}{2} \lambda^{-2} (\theta - \mu)^t \Sigma^{-1} (\theta - \mu) \\ &- (A+1) \log(\lambda^2) - \frac{B}{\lambda^2} \\ = &\tilde{C} - (A+k/2+1) \log(\lambda^2) - \frac{B + \frac{1}{2} (\theta - \mu)^t \Sigma^{-1} (\theta - \mu)}{\lambda^2}, \end{split}$$

which is up to an additive constant the logarithm of the density of the inverse gamma distribution with shape parameter A + k/2 and scale parameter  $B + \frac{1}{2}(\theta - \mu)^t \Sigma^{-1}(\theta - \mu)$ .

Lemma 7. We have

$$\log p(X^T \mid j, \lambda^2) = -\frac{1}{2} \log |\lambda^2 \Sigma S + \mathbb{I}_k| - \frac{1}{2} \mu^t \Sigma^{-1} \mu + \frac{1}{2} (m + \lambda^{-2} \Sigma^{-1} \mu)^t (S + \lambda^{-2} \Sigma^{-1})^{-1} (m + \lambda^{-2} \Sigma^{-1} \mu).$$

*Proof.* This follows from

$$p(X^T \mid j, \lambda^2) = \int p(X^T \mid j, \theta^j, \lambda^2) p(\theta^j \mid j, \lambda) d\theta^j$$

and lemma 3.

# 4 The sparsity of the Girsanov matrix with Faber-Schauder functions

The Faber-Schauder basis functions  $\psi_0, \psi_{j,k}$  are defined as follows:

$$\psi_0(x) = \begin{cases} 1 - 2x & \text{when } x \in [0, 1/2), \\ 2x - 1 & \text{when } x \in [1/2, 1], \\ 0 & \text{otherwise,} \end{cases}$$

$$\Lambda(x) = \begin{cases} 2x & \text{when } x \in [0, 1/2), \\ 2(1 - x) & \text{when } x \in [1/2, 1], \\ 0 & \text{otherwise,} \end{cases}$$

and

$$\psi_{i,k}(x) = \Lambda(2^j x - k + 1), \quad j = 0, 1, \dots, k = 1, \dots, 2^j,$$

see van der Meulen, Schauer, and van Waaij, 2018, p. 607. We say that  $\psi_0$  and  $\psi_{0,1}$  are of level zero, and the basis functions  $\psi_{j,1},\ldots,\psi_{j,2^j}$  are said to be of level j. The Girsanov matrix S defined in eq. (1) with all basis function up to and including level J is denoted by  $S^J$ . Note that  $S^J$  has  $2 + \sum_{j=1}^J 2^j = 2^{J+1}$  rows and columns, and  $2^{2J+2}$  entries.

**Definition 8.** Let  $M^n$  be an  $n \times n$ -matrix, and let  $nz(M^n)$  the number of non-zero entries of  $M^n$ . The level of sparsity of  $M^n$  is the fraction of nonzero entries,  $\frac{nz(M^n)}{n^2}$ .

The definition of a sparse matrix is vague. Usually, we mean that the number of nonzero entries grows at most linear with the number of rows. We will establish that for  $S^n$ , the number of nonzero entries grows at most like  $r \log r$  with r the number of rows.

Recall the definition of  $S_{l,l'}$  in lemma 3. Note that  $S_{l,l'} = 0$  when  $SUPP(\psi_l) \cap SUPP(\psi_{l'})$  has Lebesgue measure zero. We say that  $\psi_l$  and  $\psi_{l'}$  have non-overlapping support when their supports are either disjoint or only share a boundary point; otherwise, we say they have overlapping support.

Note that both functions of level zero,  $\psi_1$  and  $\psi_{0,1}$ , have the same support [0,1].

When  $j \geq 0, d \geq 0$  and  $d + j \geq 1$ , there are  $2^d$  Faber functions of level j + d that have overlapping support with  $\psi_{j,k}$ ,  $j \geq 0$ . These are

$$\psi_{j+d,(k-1)2^d+1}, \psi_{j+d,(k-1)2^d+2}, \dots, \psi_{j+d,k2^d}$$

For level 0, there are exactly two, and for level  $1, \ldots, j-1$  there is precisely one basis function with overlapping support with  $\psi_{j,k}$ .

So for  $\psi_0$  and  $\psi_{0,1}$  there are

$$2 + \sum_{d=1}^{J} 2^d = 2^{J+1}$$

basis functions  $\psi_0, \psi_{j',k'}, j' \leq J$  with overlapping support. For  $\psi_{j,k}, j \geq 1$ , there are

$$2+j-1+\sum_{d=0}^{J-j} 2^d = j+2^{J-j+1}$$

basis functions  $\psi_0, \psi_{j',k'}, j' \leq J$ , with overlapping support. When we make use of lemma 11, we see that  $S^n$  has at most

$$2 \cdot 2^{J+1} + \sum_{j=1}^{J} 2^{j} \left( j + 2^{J-j+1} \right)$$

$$= 2 \cdot 2^{J+1} + (J-1)2^{J+1} + 2 + J2^{J+1}$$

$$= (2J+1)2^{J+1} + 2$$

nonzero entries.

So the number of nonzero entries of  $S^n$  grows at most like  $r \log r$  with r the number of rows. It has level of sparsity at most

$$\frac{(2J+1)2^{J+1}+2}{2^{2J+2}} = (2J+1)2^{-J-1} + 2^{-2J-1},$$

which is of the order  $\frac{\log r}{r}$ .

### 5 Credible bands

Suppose we have a prior  $\Pi$  on  $\theta$ , where  $\theta : \mathbb{R} \to \mathbb{R}$  is a 1-periodic function. Let  $X^T = (X_t : t \in [0, T])$  be a sample path of  $dX_t = \theta(X_t)dt + dW_t$ . Consider the posterior  $\Pi(\cdot \mid X^T)$ .

**Definition 9.** A **pointwise credible band** of **credible level**  $1 - \alpha$  are two functions  $f_L : \mathbb{R} \to \mathbb{R}$  and  $f_H : \mathbb{R} \to \mathbb{R}$  so that for each  $t \in \mathbb{R}$ ,

$$\Pi(\{\theta: f_L(t) \le \theta(t) \le f_H(t)\} \mid X^T) \ge 1 - \alpha.$$

A simultaneous credible band of credible level  $1 - \alpha$  are two functions  $f_L : \mathbb{R} \to \mathbb{R}$  and  $f_H : \mathbb{R} \to \mathbb{R}$  so that

$$\Pi(\{\theta: f_L(t) \le \theta(t) \le f_H(t) \,\forall t\} \mid X^T) \ge 1 - \alpha.$$

So

simultaneous credible band  $\implies$  pointwise credible band.

The reverse does not hold necessarily.

#### 5.1 How to construct credible bands

#### 5.1.1 Exact pointwise credible bands

With Gaussian process priors you can construct exact pointwise credible bands. The posterior is of the form

$$f(t) = \sum_{k=1}^{N} \theta_k \phi_k, \quad \begin{pmatrix} \theta_1 \\ \vdots \\ \theta_N \end{pmatrix} \sim N(m, V),$$

where m is the N-dimensional mean vector and V is the  $N \times N$ -covariance matrix.

The coefficients are multivariate normally distributed, so f(t) is, as a linear combination of the coefficients, normally distributed with mean

$$\mathbb{E}[f(t)] = \sum_{k=1}^{N} \mathbb{E}[\theta_k] \phi_k(t) = \sum_{k=1}^{N} m_k \phi_k(t)$$

and variance

$$\operatorname{var}(f(t)) = \sum_{k=1}^{N} \sum_{\ell=1}^{N} \operatorname{cov}(\theta_k, \theta_\ell) \phi_k(t) \phi_\ell(t)$$
$$= \sum_{k=1}^{N} \sum_{\ell=1}^{N} V_{k\ell} \phi_k(t) \phi_\ell(t)$$

Let  $\xi_p$  be the quantile function of a standard normally distributed random variable Z, so  $\mathbb{P}(Z \leq \xi_p) = p$ . The *exact* pointwise credible band (around the posterior mean) is

$$f_L(t) = \mathbb{E}[f(t)] - \sqrt{\operatorname{var}(f(t))} \xi_{1-\alpha/2}$$

and

$$f_H(t) = \mathbb{E}[f(t)] + \sqrt{\operatorname{var}(f(t))} \xi_{1-\alpha/2}.$$

#### 5.1.2 Simulated simultaneous credible bands

Here I describe a procedure to simulate a  $1-\alpha$ -simultaneous credible band around the posterior mean.

**Algorithm 10.** Given a prior  $\Pi$  on a space of drift functions, and data  $X^T = (X_t : t \in [0, T])$ .

- 1. Calculate the posterior  $\Pi(\cdot \mid X^T)$ ,
- 2. calculate the posterior mean  $\bar{\theta} = \int \theta d\Pi(\theta \mid X^T)$  (you may use the mean function in the BayesianNonparametricStatistics.jl package),
- 3. simulate  $\theta_1, \ldots, \theta_M$  from the posterior,
- 4. for each i, calculate  $d_i = \sup\{|\theta_i(t) \bar{\theta}(t)| : t \in \mathbb{R}\}$ .
- 5. take the  $\lceil (1-\alpha) \cdot M \rceil$  functions  $\theta_{(1)}, \ldots, \theta_{(\lceil (1-\alpha)M \rceil)}$  from  $\theta_1, \ldots, \theta_M$  for which  $d_i$  is the smallest.
- 6. Define  $f_L$  and  $f_M$  as

$$f_L(t) = \min \left\{ \theta_{(1)}(t), \dots, \theta_{(\lceil (1-\alpha)M \rceil)}(t) \right\} \quad and \quad f_H(t) = \max \left\{ \theta_{(1)}(t), \dots, \theta_{(\lceil (1-\alpha)M \rceil)}(t) \right\}.$$

#### A Lemma

**Lemma 11.** For each  $J \in \mathbb{N}$ ,

$$\sum_{j=1}^{J} j2^{j} = (J-1)2^{J+1} + 2.$$

Proof. Note that

$$\sum_{j=1}^{J} j 2^{j} = \sum_{j=1}^{J} \sum_{k=j}^{J} 2^{k}$$

$$= \sum_{j=1}^{J} 2^{j} \sum_{k=0}^{J-j} 2^{k}$$

$$= \sum_{j=1}^{J} 2^{j} (2^{J-j+1} - 1)$$

$$= J2^{J+1} - (2^{J+1} - 2)$$

$$= (J-1)2^{J+1} + 2.$$

### References

Chung, K.L. and R.J. Williams (1990 reprint 2014). *Introduction to Stochastic Integration*. Modern Birkhäuser Classics. Springer New York. ISBN: 978-1-4614-9587-1. DOI: 10.1007/978-1-4614-9587-1.

- Steele, J.M. (2001). Stochastic Calculus and Financial Applications. Applications of mathematics: stochastic modelling and applied probability. Springer. ISBN: 9780387950167. DOI: 10.1007/978-1-4684-9305-4.
- van der Meulen, F.H., M. Schauer, and J. van Waaij (2018). "Adaptive nonparametric drift estimation for diffusion processes using Faber–Schauder expansions". In: *Statistical Inference for Stochastic Processes* 21.3, pp. 603–628. DOI: 10.1007/s11203-017-9163-7.