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SUMMARY

Multivariate geostatistics is based on modelling all covariances between all possible combinations of two or more variables and their locations in a continuously indexed domain. Multivariate spatial covariance models need to be built with care, since any covariance matrix that is derived from such a model has to be nonnegative-definite. In this article, we develop a conditional approach for model construction. Starting with bivariate spatial covariance models, we demonstrate the approach’s generality, including its connection to regression and to multivariate models defined by spatial networks. We demonstrate the fitting of such models on a minimum-maximum temperature dataset.

Some key words: Asymmetry; Causal spatial models; Cross-covariance function; Kriging; Multivariate geostatistics.

1. INTRODUCTION

The conditional approach to building multivariate spatial covariance models was introduced by Royle et al. (1999) in an edited volume of case studies in Bayesian statistics, although the approach itself is relevant to all forms of inference. In that paper, pressure and wind fields are modelled as a bivariate process over a region of the globe, with the wind process conditioned on the pressure process through a physically motivated stochastic partial differential equation. This, and a univariate spatial covariance model for the pressure process, defines valid covariance and cross-covariance functions for the bivariate (wind, pressure) process. In general, such models exhibit asymmetry; that is, for \( Y_1(\cdot) \) and \( Y_2(\cdot) \) defined on \( d \)-dimensional Euclidean space \( \mathbb{R}^d \),

\[
\text{cov}(Y_1(s), Y_2(u)) \neq \text{cov}(Y_2(s), Y_1(u)); \quad s, u \in \mathbb{R}^d.
\]

Of course, it is always true that \( \text{cov}(Y_1(s), Y_2(u)) = \text{cov}(Y_2(s), Y_1(s)) \).

There are commonly used classes of multivariate spatial models that assume symmetric, stationary dependence in the cross-covariances; that is, they assume

\[
C_{12}(h) \equiv \text{cov}(Y_1(s), Y_2(s + h)) = \text{cov}(Y_2(s), Y_1(s + h)) \equiv C_{21}(h); \quad h \in \mathbb{R}^d
\]

(e.g., Gelfand et al., 2004; Cressie & Wikle, 2011, Section 4.1.5; Genton & Kleiber, 2015). The most notable of these symmetric-cross-covariance models is the linear model of coregionalization; see, for example, Journel & Huijbregts (1978, Section III.B.3), Webster et al. (1994), Wackernagel (1995), and Banerjee et al. (2004, Section 7.2). While symmetry may reduce the number of parameters or allow fast computations, it may not be supported by the underlying science or by the data. Ver Hoef & Cressie (1993) avoid making symmetry restrictions by working with (variance-based) cross-variograms. In multivariate spatial-lattice models, Sain & Cressie (2007) and Sain et al. (2011) specifically in-
clude asymmetry parameters and use them to summarize the asymmetry in the data they analyse. Martinez-Beneito (2013) gives a multivariate spatial-lattice model that can model asymmetry between different spatial processes. Other approaches used to capture asymmetry are reviewed in Genton & Kleiber (2015).

A key outcome of multivariate geostatistics is optimal spatial prediction of a hidden multivariate spatial process, \( Y(\cdot) = (Y_1(\cdot), \ldots, Y_p(\cdot))^\top \), based on multivariate noisy spatial observations, \( \{Z_q(s_{qi}) : i = 1, \ldots, m_q, q = 1, \ldots, p\} \), of the hidden processes \( \{Y_q(\cdot) : q = 1, \ldots, p\} \). Assuming additive measurement error, \( \varepsilon_q(\cdot) \), we have data \( Z_q(\cdot) = Y_q(\cdot) + \varepsilon_q(\cdot) \) at the \( m_q \) data locations, \( D_q \equiv \{s_{qi} : i = 1, \ldots, m_q\} \), for \( q = 1, \ldots, p \). Notice that we have not assumed that the data for different spatial variables are collocated.

When just one of the processes, say \( Y_1(\cdot) \), is optimally predicted using the multivariate data \( \{Z_q(s_{qi})\} \), the associated methodology is often called cokriging. Contributions to multivariate-spatial-prediction methodology include those of Myers (1982, 1992), Ver Hoef & Cressie (1993), Wackernagel (1995), Cressie & Wikle (1998), Gelfand et al. (2004), Majumdar & Gelfand (2007), Finley et al. (2008), Huang et al. (2009), and Cressie & Wikle (2011, Section 4.1.5).

Genton & Kleiber (2015) give a comprehensive review of many different ways that valid multivariate covariances can be constructed, with a brief mention of the conditional approach. Further methodological developments of the conditional approach in geostatistics can be found in Royle & Berliner (1999), and, for spatial-lattice data, Kim et al. (2001) use a conditional approach to model both the multivariate and the spatial dependence. For (regular or irregular) gridded spatial processes, Cressie & Wikle (2011, p. 234) clarify the discussion of the conditional approach given in Gelfand et al. (2004). In this article we show that a large class of multivariate spatial covariance models come naturally from conditional-probability modelling on a continuous-spatially-indexed domain.

In Section 2, we present a construction of a bivariate spatial covariance function that is based on conditional means and conditional covariances. Section 3 proves the existence of a bivariate Gaussian process with this covariance function, gives a simple example of cokriging where the continuous spatial index \( s \in \mathbb{R}^d \), and shows how to derive cross-covariance functions from marginal covariance functions. The extension of the conditional approach to more than two variables is given in Section 4. Section 5 applies the methodology to meteorological data describing minimum and maximum temperatures in the state of Colorado, U.S.A., on a given day. Finally, Section 6 contains a brief discussion.

2. Modelling Joint Dependence Through Conditioning

In this article, we introduce the conditional approach with the bivariate case, where \( \{(Y_1(s), Y_2(s)) : s \in D \subset \mathbb{R}^d\} \) are two co-varying spatial processes in a continuous-spatially-indexed domain \( D \) contained in \( d \)-dimensional Euclidean space \( \mathbb{R}^d \); the multivariate case is considered in Section 4. As was seen in Section 1, it is sometimes convenient to write the individual processes as \( Y_1(\cdot) \) and \( Y_2(\cdot) \), respectively. Then the joint probability measure of \( [Y_1(\cdot), Y_2(\cdot)] \) can be written as,

\[
[Y_1(\cdot), Y_2(\cdot)] = [Y_2(\cdot) \mid Y_1(\cdot)] [Y_1(\cdot)],
\]

where we use the convention that \([A \mid B]\) represents the conditional probability of \( A \) given \( B \), and \([B]\) represents the marginal probability of \( B \). The conditional probability in (1) is conditional on the entire process \( Y_1(\cdot) \). In this article, we are particularly interested in \( Y_2(s) \mid Y_1(\cdot) \) and \( (Y_2(s), Y_2(u)) \mid Y_1(\cdot) \), where \( s, u \in D \).
The order of the variables is a choice, but it is generally driven by the underlying science; for example, $Y_1(\cdot)$ might be a temperature field and $Y_2(\cdot)$ might be a rainfall field, where $Y_2(\cdot)$ depends to some extent on $Y_1(\cdot)$ through evapo-transpiration and the Penman-Monteith equation (e.g., Beven, 1979). In Royle & Berliner (1999), $Y_1(\cdot)$ was a pressure field and $Y_2(\cdot)$ was a wind field. For the multivariate case in Section 4, the ordering is generalized through a graphical model.

Assume that $E(Y_1(\cdot)) \equiv 0 \equiv E(Y_2(\cdot))$, although we relax this in Section 3. Consider the following model for the first two conditional moments of $[Y_2(\cdot) \mid Y_1(\cdot)]$:

\[
E(Y_2(s) \mid Y_1(\cdot)) = \int_D b(s,v)Y_1(v) \, dv; \quad s \in D, \tag{2}
\]

\[
\text{cov}(Y_2(s), Y_2(u) \mid Y_1(\cdot)) = C_{2|1}(s,u); \quad s, u \in \mathbb{R}^d,
\]

where $b(\cdot, \cdot)$ is any integrable function mapping from $\mathbb{R}^d \times \mathbb{R}^d$ into $\mathbb{R}$, and $C_{2|1}(\cdot, \cdot)$ is a univariate covariance function that does not depend functionally on $Y_1(\cdot)$. In (2), $b(\cdot, \cdot)$ may be obtained from scientific understanding of how $Y_2(\cdot)$ depends on $Y_1(\cdot)$ (e.g., how wind depends on pressure gradients; see Royle et al., 1999), and hence we call it an interaction function. It is not a kernel since it can take both positive and negative values. Further, $C_{2|1}$ in (2) is necessarily a nonnegative-definite function, and there are many classes of such functions available (e.g., Christakos, 1984; Cressie, 1993, Section 2.5; Banerjee et al., 2004, Section 2.2). Finally, suppose that $Y_1(\cdot)$ has covariance function $C_{11}(\cdot, \cdot)$, which is also necessarily nonnegative-definite (i.e., is valid). The conditional approach requires only specification of an integrable interaction function and two valid univariate spatial covariance functions $C_{2|1}$ and $C_{11}$, leading to rich classes of cross-covariance functions (e.g., Section 3.3).

Define $C_{qr}(s, u) \equiv \text{cov}(Y_q(s), Y_r(u))$, for $q, r = 1, 2$ and $s, u \in D$. From the two univariate spatial covariance models, $C_{2|1}$ and $C_{11}$, we have:

\[
C_{22}(s,u) \equiv \text{cov}(Y_2(s), Y_2(u)) = E(\text{cov}(Y_2(s), Y_2(u) \mid Y_1(\cdot))) + \text{cov}(E(Y_2(s) \mid Y_1(\cdot)), E(Y_2(u) \mid Y_1(\cdot))) = C_{2|1}(s,u) + \int_D \int_D b(s,v)C_{11}(v,w)b(u,w) \, dv \, dw; \quad s, u \in D. \tag{3}
\]

Importantly, the formulas for the cross-covariances are,

\[
C_{12}(s,u) = \text{cov}(Y_1(s), Y_2(u)) = \text{cov}(Y_1(s), E(Y_2(u) \mid Y_1(\cdot))) = \int_D C_{11}(s,w)b(u,w) \, dw; \quad s, u \in D, \tag{4}
\]

and

\[
C_{21}(s,u) = C_{12}(u,s); \quad s, u \in D. \tag{5}
\]

Finally, recall that

\[
C_{11}(s,u) = \text{cov}(Y_1(s), Y_1(u)); \quad s, u \in D, \tag{6}
\]

where $C_{11}(\cdot, \cdot)$ is a given nonnegative-definite function. Then (3)–(6) specifies all covariances \{C_{qr}(\cdot, \cdot)\}, and any covariance matrix obtained from them should be nonnegative-definite; see Section 3. From (4), $C_{12}(u,s) = \int_D C_{11}(u,w)b(s,w) \, dw \neq C_{2|1}(s,u)$, in general; that is, the conditional approach can capture asymmetry.
3. Bivariate stochastic processes based on conditioning

3-1. Existence of a bivariate stochastic process

Let \( \{ (Y_1^h(s), Y_2^h(s)) : s \in \mathbb{R}^d \} \) be a bivariate Gaussian process with mean 0, covariance functions \( C_{11}^h(\cdot, \cdot), C_{12}^h(\cdot, \cdot), \) and cross-covariance functions \( C_{21}^h(\cdot, \cdot), \) and \( C_{22}^h(\cdot, \cdot). \) Then for any pair of nonnegative integers \( n_1, n_2, \) such that \( n_1 + n_2 > 0; \) any locations \( \{ s_{1k} : k = 1, \ldots, n_1 \}, \{ s_{2l} : l = 1, \ldots, n_2 \}; \) and any real numbers \( \{ a_{1k} : k = 1, \ldots, n_1 \}, \{ a_{2l} : l = 1, \ldots, n_2 \}, \)

\[
\operatorname{var} \left( \sum_{k=1}^{n_1} a_{1k} Y_1^h(s_{1k}) + \sum_{l=1}^{n_2} a_{2l} Y_2^h(s_{2l}) \right) \\
= \sum_{k=1}^{n_1} \sum_{k'=1}^{n_1} a_{1k} a_{1k'} C_{11}^h(s_{1k}, s_{1k'}) + \sum_{l=1}^{n_2} \sum_{l'=1}^{n_2} a_{2l} a_{2l'} C_{22}^h(s_{2l}, s_{2l'}) \\
+ \sum_{k=1}^{n_1} \sum_{k'=1}^{n_2} a_{1k} a_{2l'} C_{12}^h(s_{1k}, s_{2l'}) + \sum_{k=1}^{n_1} \sum_{k'=1}^{n_1} a_{1k} a_{1k'} C_{21}^h(s_{2l}, s_{1k'}) \geq 0. \tag{7}
\]

Conversely, suppose that the set of functions, \( \{ C_{qr}(\cdot, \cdot) : q, r = 1, 2, \} \), has the property that \( C_{11} \) and \( C_{22} \) are nonnegative-definite; \( C_{12}(s, u) = C_{21}(u, s), \) for all \( s, u \in \mathbb{R}^d, \) and (7) holds. Then there exists a bivariate Gaussian process \( \{(Y_1(s), Y_2(s)) : s \in \mathbb{R}^d \} \) such that for \( q, r = 1, 2, \)

\[ \operatorname{cov}(Y_q(s), Y_r(u)) = C_{qr}(s, u); \quad s, u \in \mathbb{R}^d. \]

The proof of this result relies on establishing the Kolmogorov consistency conditions (e.g., Billingsley, 1995, pp. 482–484) for the finite-dimensional distributions of

\[ \{Y_1(s_{11}), \ldots, Y_1(s_{1n_1}), Y_2(s_{21}), \ldots, Y_2(s_{2n_2})\}, \]

which are specified to be Gaussian and whose second-order moments are defined by (3)–(6). The two conditions are: the finite-dimensional distributions are consistent over marginalization; and permutation of the variables’ indices does not change the probabilities of events.

Now consider \( \{ C_{qr}(\cdot, \cdot) \} \) defined by (3)–(6). The right-hand side of (3) consists of two terms: The first is \( C_{22}(\cdot, \cdot), \) which is nonnegative-definite; and the second is a quadratic form that is guaranteed to be nonnegative-definite, since \( C_{11}(\cdot, \cdot) \) in (6) is nonnegative-definite. Hence \( C_{22}(\cdot, \cdot), \) which is the sum of these two terms, is nonnegative-definite. Because the finite-dimensional distributions are Gaussian, the permutation-invariance condition is guaranteed by (5), an expression for covariances.

It only remains to establish (7): Substitute (3) and (4) into the left-hand side of (7) to obtain

\[
\sum_{k=1}^{n_1} \sum_{k'=1}^{n_1} a_{1k} a_{1k'} C_{11}(s_{1k}, s_{1k'}) + \int_D \int_D a(s) a(u) C_{11}(s, u) \, ds \, du, \tag{8}
\]

where

\[ a(s) = \sum_{k=1}^{n_1} a_{1k} \delta(s - s_{1k}) + \sum_{l=1}^{n_2} a_{2l} b(s_{2l}, s); \quad s \in \mathbb{R}^d, \]

and \( \delta(\cdot) \) is the Dirac delta function. Since both \( C_{22} \) and \( C_{11} \) are nonnegative-definite, (8) is nonnegative, which establishes (7).

Only nonnegative-definite functions for univariate processes are needed in the conditional approach. Further, the finite-dimensional distribution of \( \{ Y_1(s_{1k}), Y_2(s_{2l}) : k = 1, \ldots, n_1; l = 1, \ldots, n_2 \} \) depends critically on the finite collection of interaction functions, \( \{ b(s_{2l}, \cdot) : l = \ldots, n_2 \}. \)
In practice, any geostatistical software will discretize the continuous spatial domain \( D \) onto a fine-resolution finite grid defined by the spatial-lattice, \( D^L \equiv \{s_1, \ldots, s_n\} \), which represents the centroids of the grid cells. That is, \( Y_q(\cdot) \) is replaced with the vector \( Y_q \equiv (Y_q(s_1), \ldots, Y_q(s_n))^T \); \( q = 1, 2 \). Under this discretization, (3)–(6) become, respectively,

\[
\text{cov}(Y_2) = \Sigma_{2|1} + B\Sigma_{11}B^T, \tag{9}
\]

\[
\text{cov}(Y_1, Y_2) = \Sigma_{11}B^T, \tag{10}
\]

\[
\text{cov}(Y_2, Y_1) = B\Sigma_{11}, \tag{11}
\]

\[
\text{cov}(Y_1) = \Sigma_{11}, \tag{12}
\]

which were given by Cressie & Wikle (2011, p. 160). These same modelling equations were used by Jin et al. (2005) for bivariate spatial-lattice data. Here, \( \Sigma_{2|1} \) and \( \Sigma_{11} \) are nonnegative-definite \( n \times n \) covariance matrices obtained from \( \{C_{2|1}(s_k, s_l) : k, l = 1, \ldots, n\} \) and \( \{C_{11}(s_k, s_l) : k, l = 1, \ldots, n\} \), respectively; \( B \) is the square \( n \times n \) matrix obtained from \( \{b(s_k, s_l) : k, l = 1, \ldots, n\} \); and the \( 2n \times 2n \) joint covariance matrix,

\[
\text{cov}([Y_1^T, Y_2^T]^T) = \begin{bmatrix} \Sigma_{11} & \Sigma_{11}B^T \\ B\Sigma_{11} & \Sigma_{2|1} + B\Sigma_{11}B^T \end{bmatrix}, \tag{13}
\]

is nonnegative-definite.

The book by Banerjee et al. (2015, p. 273) states that it is meaningless to talk about the joint distribution of \( Y_2(s_1) \mid Y_1(s_1) \) and \( Y_2(s_2) \mid Y_1(s_2) \), with which we agree. It also goes on to say that this “reveals the impossibility of conditioning,” with which we disagree. We have shown in this section that the conditional approach yields a well defined bivariate Gaussian process \( (Y_1(\cdot), Y_2(\cdot)) \). This implies a well defined joint distribution of the random vectors \( Y_1 \) and \( Y_2 \) (obtained from discretization) given by \( [Y_1, Y_2] = [Y_2 \mid Y_1][Y_1] \),

\[
[Y_2 \mid Y_1] \sim \text{Gau}(BY_1, \Sigma_{2|1}), \tag{14}
\]

and \( [Y_1] \sim \text{Gau}(0, \Sigma_{11}) \), where “Gau” denotes the Gaussian distribution.

Equation (14) takes the form of a linear regression on the hidden variables. However, a regression of noisy observations from \( Y_2 \) on noisy observations from \( Y_1 \) is a different, errors-in-variable model (Berkson, 1950). In (14), the conditioning is on the whole vector \( Y_1 \), but any marginal or conditional finite-dimensional distribution can be easily derived. For example, \( [Y_2(s_1) \mid Y_1(s_1)] \) can be obtained from \( [Y_1(s_1), Y_2(s_1)]/[Y_1(s_1)] \), as follows. The numerator is

\[
[Y_1(s_1), Y_2(s_1)] = \int_{\mathbb{R}} \cdots \int_{\mathbb{R}} [Y_2(s_1) \mid Y_1][Y_1]dY_1(s_2) \ldots dY_1(s_n),
\]

which from (13) is Gaussian with mean \( 0 \) and \( 2 \times 2 \) covariance matrix,

\[
\begin{bmatrix} C_{11}(s_1, s_1) + \sum_{k=1}^{n} C_{11}(s_1, s_k)b_{1k} \\ \sum_{k=1}^{n} C_{11}(s_1, s_k)b_{1k} \end{bmatrix},
\]

where \( b_{ik} \) is the \((i, k)\)th element of \( B \) in (9)–(11), and the denominator is \( \text{Gau}(0, C_{11}(s_1, s_1)) \).

We have seen above that it is not just one or a few finite-dimensional distributions that define the conditional approach, it is all of them. Further, these finite-dimensional distributions are for the hidden processes \( Y_1(\cdot) \) and \( Y_2(\cdot) \) and not for the noisy incomplete data. Banerjee et al. (2015, p. 273) state that the conditional approach is flawed and that kriging is not possible. In Section 3.2, we give a simple, one-dimensional example of the conditional approach defined by (3)–
and establish kriging and cokriging equations for predicting \( \{Y_1(s_0) : s_0 \in D^L\} \) from noisy incomplete data, \( \{Z_q(s_{qi}) : i = 1, \ldots, m_q, q = 1, 2\} \).

The incorporation of non-zero mean functions in \( \{Y_1(\cdot), Y_2(\cdot)\} \) is straightforward. Let \( \mu_1(\cdot) \) and \( \mu_2(\cdot) \) be two real-valued functions defined on \( \mathbb{R}^d \), and suppose that the finite-dimensional Gaussian distributions obtained from \( \{Y_1(s_{1k}), Y_2(s_{2l}) : k = 1, \ldots, n_1; l = 1, \ldots, n_2\} \) have means \( \{\mu_1(s_{1k}), \mu_2(s_{2l}) : k = 1, \ldots, n_1; l = 1, \ldots, n_2\} \), respectively. Then the same method of proof at the beginning of this section yields a bivariate Gaussian process \( \{Y_1(\cdot), Y_2(\cdot)\} \) with mean functions \( \{\mu_1(\cdot), \mu_2(\cdot)\} \) and covariance functions \( \{C_{qr}(\cdot, \cdot) : q, r = 1, 2\} \). Covariates \( x_1(\cdot) \) and \( x_2(\cdot) \) can then be incorporated through \( \mu_q(s) = x_q(s)^T \beta_q; s \in D, q = 1, 2 \), where \( \beta_1 \) and \( \beta_2 \) are vectors of regression coefficients of possibly different dimension.

### 3.2. Cokriging using covariances defined by the conditional approach

Section 3.1 establishes the existence of the bivariate process \( \{Y_1(\cdot), Y_2(\cdot)\} \) with \( \{C_{qr}(\cdot, \cdot)\} \) given by (3)-(6), and hence we may use cokriging for multivariate spatial prediction in the presence of incomplete, noisy data.

The aim of cokriging is to predict, say, \( Y_1(s_0) \), \( s_0 \in D \), based on \( Z_1 \) and \( Z_2 \) (Cressie, 1993, p. 138), where

\[
Z_q \equiv (Z_q(s) : s \in D_q^O)^T, \quad \text{for} \quad D_q^O \equiv \{s_{qi} : i = 1, \ldots, m_q\}; q = 1, 2.
\]

Recall that \( Z_q(s_{qi}) = Y_q(s_{qi}) + \varepsilon_q(s_{qi}), \ E(\varepsilon_q(\cdot)) = 0, \) and \( \text{var}(\varepsilon_q(\cdot)) = \sigma^2_{\varepsilon_q}, \ i = 1, \ldots, m_q, \ q = 1, 2. \) Then, assuming \( E(Y_1(\cdot)) = 0 = E(Y_2(\cdot)), \) the best predictor for \( Y_1(s_0) \) is the conditional mean, \( E(Y_1(s_0) | Z_1, Z_2) \). Under Gaussianity this is given by

\[
\hat{Y}_1(s_0) \equiv E(Y_1(s_0) | Z_1, Z_2) = \begin{bmatrix} c_{11}^T & c_{12}^T \end{bmatrix} \begin{bmatrix} C_{11} + \sigma^2_{\varepsilon_1} I_{m_1} & C_{12} \\ C_{21} & C_{22} + \sigma^2_{\varepsilon_2} I_{m_2} \end{bmatrix}^{-1} \begin{bmatrix} Z_1 \\ Z_2 \end{bmatrix},
\]

where for \( q, r = 1, 2, \)

\[
c_{1r}^T \equiv (C_{1r}(s_0, s_{ri}) : i = 1, \ldots, m_r),
\]

\[
c_{qr}^T \equiv (C_{qr}(s_{qi}, s_{rj}) : i = 1, \ldots, m_q; j = 1, \ldots, m_r),
\]

and \( I_{m_q} \) is the \( m_q \times m_q \) identity matrix. Expression (16) is called simple cokriging.

Whilst in some multivariate models, the matrices \( \{C_{qr} : q, r = 1, 2\} \) are known in closed form (Genton & Kleiber, 2015), this is not necessarily so here. Cokriging using the conditional approach may require several (analytical or numerical) integrations over \( D \) in order to compute \( \{C_{qr}\} \). This reinforces our point that the conditional approach defines a model for a bivariate stochastic process, not just a model for data defined on a collection of points within \( D \).

To demonstrate the benefits of cokriging based on a bivariate spatial model defined by the conditional approach, we simulated data in \( D \subset \mathbb{R}^1 \), where both \( C_{11}(\cdot, \cdot) \) and \( C_{21}(\cdot, \cdot) \) are Matérn covariance functions. That is,

\[
C_{11}(s, u) \equiv \frac{\sigma^2_{11}}{2^{\nu_{11}}\Gamma(\nu_{11})}(\kappa_{11}|u - s|)^{\nu_{11}} K_{\nu_{11}}(\kappa_{11}|u - s|),
\]

\[
C_{21}(s, u) \equiv \frac{\sigma^2_{21}}{2^{\nu_{21}}\Gamma(\nu_{21})}(\kappa_{21}|u - s|)^{\nu_{21}} K_{\nu_{21}}(\kappa_{21}|u - s|),
\]

where \( \sigma^2_{11}, \sigma^2_{21} \) denote the marginal variances, \( \kappa_{11}, \kappa_{21} \) are scale parameters, \( \nu_{11}, \nu_{21} \) are smoothness parameters, and \( K_{\nu} \) is the Bessel function of the second kind of order \( \nu \). Specifically, we chose the domain \( D = [-1, 1] \subset \mathbb{R}^1 \), and we discretized \( Y_1(\cdot) \) and \( Y_2(\cdot) \) into \( n = 200 \)
grid cells, each of length 0.01. Let $D^L$ be the set of locations at the centres of the cells; then $Y_1 = (Y_1(s) : s \in D^L)^T$ and $Y_2 = (Y_2(s) : s \in D^L)^T$. We generated $Y_1$ and $Y_2$ using Gaussian distributions and (9)–(12), with $\sigma_{11}^2 = 1$, $\sigma_{21}^2 = 0.2$, $\kappa_{11} = 25$, $\kappa_{21} = 75$, $\nu_{11} = \nu_{21} = 1.5$, and interaction function,

$$b(s, v) \equiv \left\{ \begin{array}{ll} A\{1 - (|v - s - \Delta|/r)^2\}; & |v - s - \Delta| \leq r \\ 0; & \text{otherwise.} \end{array} \right. \quad (19)$$

In (19), $\Delta$ is a shift parameter that here we set equal to $-0.3$, to capture asymmetry; we also set the aperture parameter $r = 0.3$, and the scaling parameter $A = 5$.

Finally, the data $Z_1$ and $Z_2$ in (15) were generated by adding independent Gaussian measurement errors to $Y_1$ and $Y_2$ at locations $D^O_1$ and $D^O_2$, respectively. We chose $\sigma_{11}^2 = \sigma_{22}^2 = 0.25$; $D^O_2 = D^L$, so that $Z_2$ is a noisy measurement of $Y_2$ at every grid cell; and $D^O_1 \equiv D^L \cap [0, 1]$, so that $Z_1$ is a noisy measurement of only those components of $Y_1$ in the positive grid cells.

The grid cells were used to define the discretized domain over which we carried out the numerical integrations in (3) and (4). For example, $\tilde{C}_{12}(s_0, u) \equiv \sum_{k=1}^n \eta_k C_{11}(s_0, w_k) b(u, w_k)$, where $D^L \equiv \{w_k : k = 1, \ldots, n\}$ and $\{\eta_k : k = 1, \ldots, n\}$ are the grid spacings; here $\eta_1 = \eta_2 = \cdots = \eta_{200} = 0.01$. More generally, when $D^L \subset D \subset \mathbb{R}^d$, $s_0$, $u$, and $\{w_k\}$ are $d$-dimensional vectors and $\{\eta_k\}$ are $d$-dimensional volumes. The covariance matrix (13) is shown in Fig. 1, left panel, where asymmetry is clearly present. Since $\Delta < 0$, the top-left corner of $\Sigma_{22}$ reduces to that of $\Sigma_{21}$, which is due to asymmetry in the interaction function $b(s, v)$.

The benefits of cokriging become apparent when the prediction of $Y_1(s_0)$ given $Z_1$ and $Z_2$ is compared to the prediction of $Y_1(s_0)$ given only $Z_1$ (i.e., univariate kriging). In our simulation, we used the cokriging equation (16) to obtain $\hat{Y}_1 \equiv (\hat{Y}_1(s_0) : s_0 \in D^L)^T$ based on the simulated observations $Z_1$ and $Z_2$. We compared $\hat{Y}_1$ to the (univariate) kriging predictor $\tilde{Y}_1$ based only on data $Z_1$, where $\tilde{Y}_1 \equiv (\tilde{Y}_1(s_0) : s_0 \in D^L)^T$ and $\tilde{Y}_1(s_0) \equiv c_{11}^T (C_{11} + \sigma_{11}^2 I_{n_1})^{-1} Z_1$. As seen in Fig. 1, right panel, the cokriging predictor $\hat{Y}_1$ is representative of the true process $Y_1$ even on the negative grid cells where it is not observed. However, the kriging predictor $\tilde{Y}_1$ can only shrink to the mean, $E(Y_1(\cdot)) = 0$, in the spatial regions where there are no observations. While it might seem more natural to predict $Y_2(\cdot)$, since our model is based on $[Y_2(\cdot) \mid Y_1(\cdot)]$, we chose to predict $Y_1(\cdot)$ to illustrate that cokriging on either variable is possible.

### 3.3. Deriving classes of cross-covariance functions from marginal covariance functions

The conditional approach may also be used to complement the joint approach to constructing multivariate covariance functions. In particular, Genton & Kleiber (2015) posed an open problem that seems difficult when using the joint approach; “[G]iven two marginal covariances, what is the valid class of possible cross-covariances that still results in a nonnegative definite structure?”. A straightforward answer to this question is available through the conditional approach. The class of cross-covariance functions is given by (4) for any integrable function $b(s, v)$ such that the function $C_{21}(\cdot, \cdot)$ obtained from (3) is nonnegative-definite. This is potentially a very rich class of cross-covariance functions, and answering the question reduces to verifying which choice of $b(\cdot, \cdot)$ in (3) yields a nonnegative-definite $C_{21}(\cdot, \cdot)$.

For example, consider the stationary case in $D = \mathbb{R}^2$ where we have stationary covariance functions $C_{11}(h)$, $C_{21}(h)$, and interaction function $b(s, v) = b_0(v - s)$. Then from (3),

$$C_{21}(h) = C_{22}(h) - \int_{\mathbb{R}^2} \int_{\mathbb{R}^2} b_0(\tilde{v}) b_0(\tilde{w}) C_{11}(h - \tilde{v} + \tilde{w}) d\tilde{v} d\tilde{w}.$$
Fig. 1. Cokriging using spatial covariances defined by the conditional approach. Left panel: The covariance matrix (13). Right panel, top: The simulated observations $Z_1$ (open circles) and $Z_2$ (dots). Right panel, bottom: The hidden value $Y_1$ (solid line), the kriging predictor $\tilde{Y}_1$ (dashed line), and the cokriging predictor $\hat{Y}_1$ (dotted line).

Let $\omega \in \mathbb{R}^2$ denote spatial frequency, and let $\Gamma_{11}(\omega)$, $\Gamma_{22}(\omega)$, and $B_0(\omega)$ be the Fourier transforms of $C_{11}(h)$, $C_{22}(h)$, and $b_0(h)$, respectively. Then, for $C_{21}(h)$ to be a valid covariance function, it is required that $\Gamma_{22}(\omega) - B_0(\omega)B_0(-\omega)\Gamma_{11}(\omega)$ be nonnegative and integrable over $\omega \in \mathbb{R}^2$ (Cressie & Huang, 1999; Gneiting, 2002). The inequality is trivial if $\Gamma_{11}(\omega) = 0$; hence consider those $\omega \in \Omega$ for which

$$B_0(\omega)B_0(-\omega) \leq \frac{\Gamma_{22}(\omega)}{\Gamma_{11}(\omega)},$$

where $\Gamma_{11}(\omega) > 0$. Recall that $C_{11}(h)$ and $C_{22}(h)$ are covariance functions and hence, necessarily, $\Gamma_{11}(\omega) \geq 0$ and $\Gamma_{22}(\omega) \geq 0$.

Any $B_0(\cdot)$ that satisfies (20) gives the required result, since then finiteness follows from $\int \Gamma_{22}(\omega) d\omega < \infty$ being an upperbound on the integral, $\int \Gamma_{22}(\omega) - B_0(\omega)B_0(-\omega)\Gamma_{11}(\omega) d\omega$. In Appendix 1, we show how a class of valid Matérn cross-covariance functions developed by Gneiting et al. (2010) can be obtained from (20).

4. Multivariate spatial models through conditioning

4.1. Definition of cross-covariance functions

In this section, we extend the conditional approach from the bivariate to the multivariate case. Initially, we work with the variables in their original ordering and subsequently show how graphical models define the general case. Now, $[Y_1(\cdot), \ldots, Y_p(\cdot)]$ can be decomposed as

$$[Y_p(\cdot) \mid Y_{p-1}(\cdot), Y_{p-2}(\cdot), \ldots, Y_1(\cdot)] [Y_{p-1}(\cdot) \mid Y_{p-2}(\cdot), \ldots, Y_1(\cdot)] \ldots [Y_1(\cdot)].$$ (21)
First, we set $\text{cov}(Y_1(s), Y_1(u)) = C_{11}(s, u)$; $s, u \in \mathbb{R}^d$. Analogous to the bivariate case $p = 2$, we define the first two conditional moments of $Y_q(\cdot)$, for $q = 1, \ldots, p$, as
\[
E(Y_q(s) \mid \{Y_r(\cdot) : r = 1, \ldots, (q - 1)\}) = \sum_{r=1}^{q-1} \int_D b_{qr}(s, v) Y_r(v) dv; \quad s \in D, \tag{22}
\]
\[
\text{cov}(Y_q(s), Y_q(u) \mid \{Y_r(\cdot) : r = 1, \ldots, (q - 1)\}) = C_{q(r < q)}(s, u); \quad s, u \in \mathbb{R}^d, \tag{23}
\]
where \{b_{qr}(\cdot, \cdot) : r = 1, \ldots, (q - 1); q = 2, \ldots, p\} are integrable functions that describe the conditional relationship of the $r$th process on the $q$th process, for $r < q$.

As a result of the decomposition in (21), we obtain from (22) and (23) the following expression for the marginal covariance functions. For $q = 1, \ldots, p$,
\[
C_{qq}(s, u) \equiv \text{cov}(Y_q(s), Y_q(u)) = C_{q(r < q)}(s, u) + \sum_{r=1}^{q-1} \sum_{r'=1}^{q-1} \int_D \int_D b_{qr}(s, v) C_{rr'}(v, w) b_{qr'}(u, w) dv dw, \tag{24}
\]
and for $r = 1, \ldots, q - 1$ the cross-covariance functions,
\[
C_{rq}(s, u) \equiv \text{cov}(Y_r(s), Y_q(u)) = \sum_{r'=1}^{q-1} \int_D \int_D b_{qr'}(u, w) C_{rr'}(s, w) dw. \tag{25}
\]
Expressions (24) and (25) depend on $C_{rr'}$, for $r, r' < q$, which are defined iteratively: Starting with $q = 2$, one obtains $C_{22}, C_{12},$ and $C_{21}$ in terms of $C_{11}$. The same idea is then repeated for $q = 3, \ldots, p$.

### 4-2. Existence of a multivariate process

Following the discussion in Section 3-1, the existence of the $p$-variate Gaussian process follows by showing that
\[
\text{var}\left(\sum_{q=1}^{p} \sum_{k=1}^{n_q} a_{qk} Y_q(s_{qk})\right) \geq 0, \tag{26}
\]
for any real numbers \{a_{qk} : k = 1, \ldots, n_q; q = 1, \ldots, p\}, any nonnegative integers \{n_q : q = 1, \ldots, p\} such that $n_1 + \cdots + n_p > 0$, and any \{s_{qk} : k = 1, \ldots, n_q; q = 1, \ldots, p\}. In Appendix 2, we demonstrate that (26) is equal to
\[
\sum_{m=1}^{n_p} \sum_{m'=1}^{n_p} a_{pm} a_{pm'} C_{p(q < p)}(s_{pm}, s_{pm'}) + \sum_{q=1}^{p-1} \sum_{r=1}^{p-1} \int_D \int_D a_{q}(s) a_{r}(u) C_{qr}(s, u) d s d u, \tag{27}
\]
where
\[
a_{q}(s) \equiv \left(\sum_{k=1}^{n_q} a_{qk} \delta(s - s_{qk}) + \sum_{m=1}^{n_p} a_{pm} b_{pq}(s_{pm}, s)\right). \tag{28}
\]
The nonnegativity of the first term in (27) follows by assumption, and the nonnegativity of the second term follows by induction; see Appendix 2.

This result implies that a multivariate spatial Gaussian model constructed using the conditional approach (22) and (23) exists, provided that the univariate covariance functions, $C_{11}(\cdot, \cdot)$ and
\{C_{q(r<q)}(\cdot,\cdot): q = 2,\ldots,p\}, are valid and the interaction functions, \{b_{qr}(\cdot,\cdot): r = 1,\ldots,q-1; q = 2,\ldots,p\}, are integrable.

4.3. A graphical-model representation

Because any joint model can be decomposed using (21), the existence result in Section 4.2 holds for systems that do not exhibit any particular structure, and recall that the univariate conditional processes themselves need not be isotropic nor stationary. Here, we give examples of implementation of the conditional approach, analogous to particular constructions of multivariate statistical models found in non-spatial settings (e.g., Cox & Wermuth, 1996).

Consider the trivariate case of \(p = 3\); clearly the decomposition (21) is not unique: The joint probability measure can be written as

\[
[Y_1(\cdot), Y_2(\cdot), Y_3(\cdot)] = [Y_3(\cdot) | Y_1(\cdot), Y_2(\cdot)][Y_2(\cdot) | Y_1(\cdot)][Y_1(\cdot)],
\]

and equally,

\[
[Y_1(\cdot), Y_2(\cdot), Y_3(\cdot)] = [Y_3(\cdot) | Y_2(\cdot), Y_3(\cdot)][Y_2(\cdot) | Y_3(\cdot)][Y_3(\cdot)].
\]

Indeed there are a total of six such possible models (in the bivariate case, there are two possible models), many of which are not physically meaningful.

When building classes of models, it is generally better to have more choice, but in the conditional approach we can be guided by a graph structure. For example, in the bivariate case considered by Royle et al. (1999), node 1 is defined by the pressure field \(Y_1(\cdot)\) and node 2 is defined by the wind field \(Y_2(\cdot)\). The direction of the edge in the directed graph is clearly from node 1 to node 2 and not the other way around, since wind fields are the result of differential pressures in the atmosphere.

As in multivariate modelling (Cox & Wermuth, 1996) and subsequently in network analysis (e.g., Lauritzen, 1996; Kolaczyk, 2009), specification of the graph structure for a multivariate model will be guided by a desire for a parsimonious model. One example of a parsimonious conditionally specified model is the spatial-moving-average model of Ver Hoef & Barry (1998), where \(p = 5\). Ver Hoef & Barry (1998) construct a bivariate model \((Y_1(\cdot), Y_2(\cdot))\) by taking moving averages of a combination of correlated processes \((Y_3(\cdot), Y_4(\cdot), Y_5(\cdot))\), where the following decomposition is implicitly assumed,

\[
[Y_1(\cdot),\ldots,Y_5(\cdot)] = [Y_1(\cdot) | Y_3(\cdot)][Y_2(\cdot) | Y_4(\cdot)][Y_4(\cdot) | Y_5(\cdot)][Y_5(\cdot) | Y_3(\cdot)]
\]

In (29), \([Y_1(\cdot) | Y_3(\cdot)]\) and \([Y_2(\cdot) | Y_4(\cdot)]\) are constructed using moving-average functions of the form \(b(s,v) = b_v(v-s)\) in (2); \([Y_3(\cdot) | Y_5(\cdot)]\) and \([Y_4(\cdot) | Y_5(\cdot)]\) are Dirac delta functions; and \(Y_5(\cdot)\) is a white-noise process. The general decomposition of \((Y_1(\cdot),\ldots,Y_p(\cdot))\) using directed graphs is given by Bishop (2006, p. 362). Should a part of the graph be undirected, it is natural to order those processes from less smooth to more smooth, motivated by the form of (3) and (24).

An important special case is commonly seen in studies of spatio-temporal phenomena. In such settings, each process is indexed by time \(t = 1,\ldots,p\), and nodes 1,\ldots,\(p\) are connected successively by directed edges. This results in the Markov factorization:

\[
[Y_1(\cdot),\ldots,Y_p(\cdot)] = [Y_1(\cdot)] \prod_{t=2}^{p}[Y_t(\cdot) | Y_{t-1}(\cdot)].
\]

4.4. The conditional approach provides spatial-model flexibility

It is easy to see that the number of covariance and interaction functions that need to be specified is \(p(p+1)/2\), which is also the number of elements in the lower triangle of a \(p \times p\) matrix.
of cross-covariance functions. Hence, the modeller using the conditional approach provides up to \( p(p + 1)/2 \) univariate nonnegative-definite and integrable functions in order to specify the \( p(p + 1)/2 \) marginal and cross-covariance functions. Critically, specification of any one function does not constrain specification of the others. This flexibility is not necessarily present in other commonly used approaches to construction of multivariate spatial models. For example, both in the convolution method of Majumdar & Gelfand (2007), and in the standard linear model of coregionalization, the user is allowed a choice of at most \( p \) covariance functions.

5. **Bivariate spatial models of temperature data: The conditional approach**

5.1. *The data*

We demonstrate the flexibility of the conditional approach on the bivariate minimum/maximum temperature dataset used in Genton & Kleiber (2015). The data are minimum-temperature and maximum-temperature residuals in the state of Colorado, U.S.A. (following the removal of the state-wide mean) obtained from measurements taken on September 19, 2004 at 94 weather stations in Colorado; that is, \( m_1 = m_2 = m = 94 \) and \( D_1^O = D_2^O = D^O \). Minimum temperatures are likely to have occurred in the early-morning hours of a given day, with the maximum temperatures in the afternoon of the same day. It is natural then to assume that the maximum-temperature residual later in the day is partially determined by the minimum-temperature residual in the early morning. Consequently, time induces a bivariate dependence that can be modelled naturally using a conditional approach.

5.2. *The model*

In this subsection, we propose a Bayesian hierarchical model with spatial dependence in the process model and prior distributions on the unknown parameters. Let \( Y_1(\cdot) \) and \( Y_2(\cdot) \) denote the true minimum-temperature and maximum-temperature residuals, respectively. Let \( \varepsilon(\cdot) \equiv (\varepsilon_1(\cdot), \varepsilon_2(\cdot))^T \) be the bivariate process of measurement errors or potential measurement errors, and assume that the data or potential data, \( Z(\cdot) \equiv (Z_1(\cdot), Z_2(\cdot))^T \), are related to \( Y(\cdot) \equiv (Y_1(\cdot), Y_2(\cdot))^T \) through, \( Z(\cdot) = Y(\cdot) + \varepsilon(\cdot) \). The two measurement errors are assumed to have no spatial dependence, but they could be correlated, which we model through

\[
\text{cov}(\varepsilon(s)) = \sigma^2 \begin{bmatrix} 1 & \rho \varepsilon \\ \rho \varepsilon & 1 \end{bmatrix}, \quad \text{for } \rho \varepsilon \in (-1, 1) \text{ and } s \in D.
\]

In the conditional approach, we need to specify the univariate covariance functions, \( C_{11}(s, u) \) and \( C_{21}(s, u) \), and the integrable interaction function \( b(s, v) \). We let the covariance functions be isotropic Matérn covariance functions given by (17) and (18). The smoothness parameters \( \nu_{11} \) and \( \nu_{21} \) are set equal to 1.5, to give covariance functions that are a little smoother than the exponential covariance function. We let \( b(s, v) \) be a function of displacement, \( h \equiv v - s \), and recall that \( b_0(h) \equiv b(s, v) \). We write the three different models that we fit, as:

- **Model 1 (independence):** \( b_0(h) \equiv 0, \)
- **Model 2 (pointwise dependence):** \( b_0(h) \equiv A\delta(h), \)
- **Model 3 (diffused dependence):** \( b_0(h) \equiv \begin{cases} A\{1 - ((||h - \Delta||/r)^2)^2; ||h - \Delta|| \leq r, \\ 0; \end{cases} \)

where \( \Delta = (∆_1, ∆_2)^T \) is a shift parameter that captures asymmetry. In Model 3, \( b_0(h) \) is a shifted bisquare function in \( \mathbb{R}^2 \), which is analogous to (19) given in \( \mathbb{R}^1 \).

We discretized both \( Y_1(\cdot) \) and \( Y_2(\cdot) \) using a triangulated grid with \( n_1 = n_2 = n = 968 \) vertices each. Here, the vertices of the grid define \( D^L \), an irregular spatial lattice; see Fig. 2.
Fig. 2. Left panel: State boundaries in a region of the U.S.A. (state boundaries are given by dashed lines), with the domain of interest enclosed by a bounding polygon (solid line). Right panel: The irregular triangular grid used for discretizing \( D \). The observation locations given by \( D^O \) consist of the large dots and the discretized spatial domain \( D^C \) consists of the mesh nodes.

Fig. 3. Results for Model 3. Left panel: Interpolated maps in degrees Celsius (degC) of \( E(Y_1 \mid Z_1, Z_2) \) and \( E(Y_2 \mid Z_1, Z_2) \). Right panel: Prior (light grey) and posterior (dark grey) median (solid line) and inter-quartile ranges (enclosed by dashed lines) of the interaction function \( b_o(\cdot) \) of Model 3, along a unit vector \( e \) originating at Denver (DE) in the direction of Fort Collins (FC).

Under the chosen triangulation, the integral in (2) is approximated as

\[
E(Y_2(s_l) \mid Y_1(\cdot)) \approx \sum_{k=1}^{n} \eta_k b(s_l, v_k) Y_1(v_k),
\]

where in this case \( \{\eta_k : l = 1, \ldots, n\} \) are the areas of the Voronoi tessellations constructed from the triangulated grid (e.g., Lee & Schachter, 1980).

5.3. The parameters

We employ a Bayesian hierarchical model and place prior distributions on transformations of the parameters, with each transformation chosen to account for the range of its respective parameter. The prior distributions and transformations are summarized in Table 1 (given in Appendix 3). Elicited prior distributions reflect some prior understanding of the problem: For example, we can reasonably expect that the error standard deviation \( \sigma_e \) will lie between 0.5 and 2 degrees
Celsius, and thus we place a Gamma prior on $\sigma_x^{-2}$ such that $q_{0.05}(\sigma_x) = 0.05$ and $q_{0.95}(\sigma_x) = 2$, where $q_z(\cdot)$ denotes the $z$th quantile of its argument.

5.4. Inference and model comparison

For posterior inference we employed a Gibbs sampler. Through the use of conjugate priors, we obtained conditional distributions of standard form for $\sigma_y^2, Y, A, \sigma_A^2$, and $\sigma_E^2$ (see Table 1, Appendix 3), from which we were able to sample directly. Since the conditional distributions for the remaining parameters, namely $\rho_e, \kappa_{11}, \kappa_{21}, r, \Delta_1$, and $\Delta_2$, are nonstandard, we sampled from these using a slice sampler (Neal, 2003). Slice samplers tend to be more efficient than simple Metropolis samplers, as they effectively alter the magnitude of the steps taken at each point in the chain by adapting to the local properties of the density function. However, the computational effort for generating one sample is larger than for a standard Metropolis sampler.

We generated $N = 50,000$ samples from the slice-within-Gibbs Markov chain Monte Carlo scheme, discarded the first 1,000 for burn-in, and then we thinned by only recording every 100th sample. Slice sampling was carried out using the stepping-out method (Neal, 2003, Section 4) with an interval width adapted during the burn-in period. We used marginal samplers for $\rho_e, \kappa_{11}, \kappa_{21}$, while for Model 3 we sampled $r, \Delta_1$, and $\Delta_2$ jointly since these parameters can be expected to be highly correlated. For the slice samplers, we used the R (R Core Team, 2014) package slice, available from the personal home page of Jonathan C. Rougier, University of Bristol.

A summary of our inferences on parameters is given in Table 2 (in Appendix 3), and the posterior expectations of $Y_1(s)$ and $Y_2(s)$ are shown in the left panel of Fig. 3. There the influence of the Rocky Mountains on the temperature residuals is apparent. In Fig. 3, right panel, we plot the interquartile ranges of the prior $[b_o(he)]$ and the posterior, $[b_o(he) \mid Z_1, Z_2]$, where $e$ is the unit vector originating from Denver in the direction of Fort Collins. The posterior distribution of the interaction function along this direction can be clearly distinguished from the prior distribution, suggesting that Bayesian learning has uncovered substantial interaction between maximum temperature and minimum temperature at proximal locations. Model 3, for different choices of $b_o(\cdot)$, has in fact often been used to study the dynamics in spatio-temporal processes (Kot & Schaffer, 1986; Wikle, 2002).

The Deviance Information Criterion (DIC, Spiegelhalter et al., 2002) for the three models is given in the lower row of Table 2 in Appendix 3, from which we were able to sample directly. Since the conditional distributions for the remaining parameters, namely $\rho_e, \kappa_{11}, \kappa_{21}, r, \Delta_1$, and $\Delta_2$, are nonstandard, we sampled from these using a slice sampler (Neal, 2003). Slice samplers tend to be more efficient than simple Metropolis samplers, as they effectively alter the magnitude of the steps taken at each point in the chain by adapting to the local properties of the density function. However, the computational effort for generating one sample is larger than for a standard Metropolis sampler.

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6. Discussion

The conditional approach can be modified easily for processes indexed on different spatial domains: \{Y_1(s) : s \in D_1\} and \{Y_2(s) : s \in D_2\}, for \(D_1, D_2 \subseteq \mathbb{R}^d\). Equation (2) becomes,

\[ E(Y_2(s) \mid Y_1(\cdot)) = \int_{D_1} b(s, v)Y_1(v)dv; \quad s \in D_2. \]

For example, Cressie & Wikle (2011, p. 287) illustrate bivariate spatial dependence between Mallard breeding bird pairs in the Prairie Pothole region of North America and the El Niño phenomenon in the tropical Pacific Ocean.

In the example given in Section 5, we fitted a Bayesian hierarchical model by putting priors on the parameters in \(C_{11}(h), C_{21}(h),\) and the interaction function \(b_0(h)\). Alternatively, for an empirical hierarchical model, the parameters are considered fixed but unknown; they are then estimated and substituted into the cokriging and kriging equations given in Section 3.2. In this case, restricted maximum likelihood estimation is recommended. Numerically, this may be achieved by an expectation-maximization algorithm or a gradient search.

Even if the parameters are known or estimated offline, spatial or spatio-temporal inference with multivariate models can remain computationally challenging. When treating all variables simultaneously in joint form, sparse formulations and sparse linear-algebraic methods can greatly facilitate the computation (e.g., Zammit-Mangion et al., 2015). Sparseness is guided by graphical representations, and in this sense our multivariate spatial models are a combination of a directed acyclic graph on the indices of the individual processes and geostatistical covariance functions on the resulting conditional processes.

By constructing models through conditioning, we obtain graphical representations for which exact inference through sequential algorithms generally exist. We have already visited the ubiquitous Markov chain in (30), which can be tackled with the iterative Rauch-Tung-Striebel smoother (e.g., Rauch et al., 1965). For more general constructions, such as trees or polytrees, the sum-product or peeling algorithm may be used for exact inference. When likelihoods associated with some or all of the processes in \(\{Y_q: q = 1, \ldots, p\}\) are intractable, approximate message passing may be used to keep the computations tractable (e.g., Heskes & Zoeter, 2002), such as when the data model for \(Z_q(\cdot)\) is a spatial Poisson point process and \(Y_q(\cdot)\) is used to model the log-intensity of the process.

Reproducible code and data for the studies in Section 3.2 and Section 5 are available as part of the R package \(\text{bicon}\) from the second author’s website on github.

Acknowledgment

We would like to thank Chris Wikle for discussions on the conditional approach to modelling multivariate spatial dependence and Jonathan Rougier for discussions on slice sampling. This research was partially supported by the US National Science Foundation and the US Census Bureau through the NSF-Census Research Network program; and it was partially supported by a 2015–2017 Australian Research Council Discovery Project.

References


Multivariate spatial modelling


It can be easily shown that the inequalities, $16^\nu C_N$ are sufficient for (A1) to hold. Then, from (4), assume that they all have the same scale $\kappa$. For $C_{\nu}$, $\kappa$ and $b_0(h)$ be isotropic Matérn covariance functions on $\mathbb{R}^2$ and, for simplicity, assume that they all have the same scale $\kappa$. Then, using obvious notation, their Fourier transforms are given by

$$B_0(\omega) = \sigma_b^2 \frac{\Gamma(\nu_b + 1)\kappa^{2\nu_b}}{\pi \Gamma(\nu_b)} (\kappa^2 + \|\omega\|^2)^{-\nu_b - 1},$$

$$\Gamma_{11}(\omega) = \sigma_{11}^2 \frac{\Gamma(\nu_{11} + 1)\kappa^{2\nu_{11}}}{\pi \Gamma(\nu_{11})} (\kappa^2 + \|\omega\|^2)^{-\nu_{11} - 1},$$

$$\Gamma_{22}(\omega) = \sigma_{22}^2 \frac{\Gamma(\nu_{22} + 1)\kappa^{2\nu_{22}}}{\pi \Gamma(\nu_{22})} (\kappa^2 + \|\omega\|^2)^{-\nu_{22} - 1}.$$ 

For $C_{21}(\cdot)$ and $C_{12}(\cdot)$ to be valid cross-covariance functions, it is required that $\Gamma_{22}(\omega) - B_0(\omega)B_0(-\omega)\Gamma_{11}(\omega) \geq 0$, and hence that

$$\sigma_b^4 \leq \frac{\pi^2 \sigma_{22}^2}{\sigma_{11}^2} \frac{1}{\nu_b \kappa^{4\nu_b}} \frac{\nu_{22} \kappa^{2\nu_{22}}}{\nu_{11} \kappa^{2\nu_{11}}} (\kappa^2 + \|\omega\|^2)^{2 + 2\nu_b + \nu_{11} - \nu_{22}}.$$ (A1)

It can be easily shown that the inequalities,

$$\nu_b \geq (\nu_{22} - \nu_{11} + 2)/2,$$ (A2)

$$\sigma_b^2 \leq 2\pi \frac{\sigma_{22}}{\sigma_{11}} \frac{1}{\nu_{22} - \nu_{11} - 2} \frac{\kappa^{\nu_{22}}}{\nu_{11} \kappa^{2\nu_{11}}} \left(\frac{\nu_{22}}{\nu_{11}}\right)^{\frac{1}{2}},$$ (A3)

are sufficient for (A1) to hold. Then, from (4), $C_{12}(h)$ is also a Matérn covariance function with variance

$$\sigma_{12}^2 = \frac{1}{\pi \kappa^2} \frac{\nu_b \nu_{11}}{\nu_{22} - \nu_{11} + 1} \frac{1}{\sigma_b^2 \sigma_{11}^2},$$ (A4)

and smoothness $\nu_{12} \equiv \nu_b + \nu_{11} + 1$. Hence, from (A2),

$$\nu_{12} \geq (\nu_{11} + \nu_{22})/2.$$ (A5)

Now consider the bound on the smoothness,

$$\nu_{12} = (\nu_{11} + \nu_{22})/2,$$ (A6)
which is obtained from the bound, \( \nu_b = (\nu_{22} - \nu_{11} - 2)/2 \), in (A2). An inequality for the variance \( \sigma^2 \) is then obtained by substituting this value of \( \nu_b \) and the inequality (A3) into (A4):

\[
\sigma^2 \leq 2\sigma^2 (\frac{\nu_{11}\nu_{22}}{\nu_{11} + \nu_{22}})^{1/2}.
\] (A7)

The conditions (A6) and (A7) are those that Gneiting et al. (2010) impose in order to construct parsimonious bivariate Matérn models. Clearly, these are more restrictive than our conditions (A2) and (A3).

Generalizing these ideas to arbitrary scale parameters \( \kappa_1, \kappa_2, \kappa_b \), results in a similar problem as that encountered in the full bivariate Matérn model in \( \mathbb{R}^d \) of Gneiting et al. (2010), in the sense that one needs to find a \( \kappa_b \) and \( \nu_b \) such that the inequality,

\[
(\kappa^2_b + \| \omega \|^2)^{\nu_b+2} \geq \frac{(\kappa^2_{22} + \| \omega \|^2)^{\nu_{22}+1}}{(\kappa^2_{11} + \| \omega \|^2)^{\nu_{11}+1}},
\]

is satisfied for all \( \omega \in \mathbb{R}^d \). In this general case it is not possible to find inequality constraints for \( \kappa_b \) and \( \nu_b \) without further assumptions.

**APPENDIX 2**

*Proof for multivariate-process existence*

Here, we prove by induction that (26) holds for for any real numbers \( \{a_{qk} : k = 1, \ldots, n_q; q = 1, \ldots, p\} \), any nonnegative integers \( \{n_q : q = 1, \ldots, p\} \) such that \( n_1 + \cdots + n_p > 0 \), and any \( \{s_{pq} : k = 1, \ldots, n_q; q = 1, \ldots, p\} \). We have already shown, through (8), that there exists a bivariate stochastic process and hence the variance of any linear combination of the two processes is nonnegative. Now, assume that \( (Y_1(\cdot), \ldots, Y_{p-1}(\cdot))' \) is a well defined \( (p-1) \)-variate stochastic process. We re-write (26) as:

\[
\text{var} \left( \sum_{q=1}^{p-1} \sum_{k=1}^{n_q} a_{qk} Y_q(s_{pq}) + \sum_{m=1}^{n_p} a_{pm} Y_p(s_{pm}) \right).
\]

Then, following the definitions for the marginal and cross-covariances in (24) and (25) and using standard identities, we obtain the following expression for (26):

\[
\sum_{m=1}^{n_p} \sum_{m'=1}^{n_{p'}} a_{pm} a_{pm'} C_p(q<p)(s_{pm}, s_{pm'})
\]

\[
\sum_{q=1}^{p-1} \sum_{k=1}^{n_q} \sum_{m=1}^{n_{m'}} a_{pq} a_{pm'} \int_D \int_D b_{pq}(s_{pm}, v) C_q(v, w) b_{pr}(s_{pm'}, w) \, dv \, dw
\]

\[
\sum_{q=1}^{p-1} \sum_{k=1}^{n_q} \sum_{m=1}^{n_{m'}} a_{pq} a_{pm'} \int_D b_{pq}(s_{pm}, w) C_q(s_{pq}, w) \, dw
\]

\[
\sum_{q=1}^{p-1} \sum_{k=1}^{n_q} \sum_{m=1}^{n_{m'}} a_{pq} a_{pm} \int_D b_{pq}(s_{pm}, v) C_q(v, s_{pk'}) \, dv
\]

\[
\sum_{q=1}^{p-1} \sum_{k=1}^{n_q} \sum_{m=1}^{n_{m'}} a_{pq} a_{pm} C_q(s_{pq}, s_{pk'})
\]
which can be reduced to
\[
\begin{align*}
&\sum_{m=1}^{n_p} \sum_{m'=1}^{n_p} a_{pm}a_{pm'}C_{p|(q<p)}(s_{pm}, s_{pm'}) \\
&\quad + \sum_{q=1}^{p-1} \sum_{r=1}^{p-1} \int_D \int_D \left( \sum_{k=1}^{n_q} a_{sk} \delta(s - s_{sk}) + \sum_{m=1}^{n_p} b_{pq}(s_{pm}, s) \right) \\
&\quad \times \left( \sum_{k'=1}^{n_q} a_{sk'} \delta(u - s_{sk'}) + \sum_{m'=1}^{n_p} b_{pq}(s_{pm'}, u) \right) C_{qr}(s, u) \, ds \, du \\
&= \sum_{m=1}^{n_p} \sum_{m'=1}^{n_p} a_{pm}a_{pm'}C_{p|(q<p)}(s_{pm}, s_{pm'}) + \sum_{q=1}^{p-1} \sum_{r=1}^{p-1} \int_D \int_D a_q(s)a_r(u)C_{qr}(s, u) \, ds \, du. \quad (B1)
\end{align*}
\]

The first term in (B1) is nonnegative by assumption while the second term is nonnegative since \((Y_1(\cdot), \ldots, Y_{p-1})'\) is well defined.

**Appendix 3**

Prior specifications and posterior distributional summaries

Table 1. Parameters used in Section 5 and the models in which they appear. Each parameter, \(\theta\), is (possibly) transformed to \(g(\theta)\), whose prior is given in the last column. Transformations are used to achieve the desired range. The notation \(\text{Gau}(\mu, \sigma^2)\) is used to denote a Gaussian distribution with mean \(\mu\) and variance \(\sigma^2\), whose cumulative distribution function is \(\Phi((\cdot - \mu)/\sigma)\), and \(\Gamma(\alpha, \beta)\) is used to denote a Gamma distribution with shape \(\alpha\) and rate (inverse scale) \(\beta\).

<table>
<thead>
<tr>
<th>Model</th>
<th>Parameter, (\theta)</th>
<th>Range</th>
<th>Transformation, (g(\theta))</th>
<th>Prior over (g(\theta))</th>
</tr>
</thead>
<tbody>
<tr>
<td>1,2,3</td>
<td>(\sigma^2_x)</td>
<td>(\mathbb{R}^+)</td>
<td>((\theta)^{-1})</td>
<td>(\Gamma(1.79, 1.10))</td>
</tr>
<tr>
<td>1,2,3</td>
<td>(\rho_e)</td>
<td>((-1, 1))</td>
<td>(\Phi^{-1}((\theta + 1)/2))</td>
<td>(\text{Gau}(0.00, 1.00))</td>
</tr>
<tr>
<td>1,2,3</td>
<td>(\sigma^2_{11})</td>
<td>(\mathbb{R}^+)</td>
<td>((\theta)^{-1})</td>
<td>(\Gamma(0.84, 2.68))</td>
</tr>
<tr>
<td>1,2,3</td>
<td>(\sigma^2_{21})</td>
<td>(\mathbb{R}^+)</td>
<td>((\theta)^{-1})</td>
<td>(\Gamma(0.84, 2.68))</td>
</tr>
<tr>
<td>1,2,3</td>
<td>(\kappa_{11})</td>
<td>(\mathbb{R}^+)</td>
<td>(\log(\theta))</td>
<td>(\text{Gau}(0.69, 1.00))</td>
</tr>
<tr>
<td>1,2,3</td>
<td>(\kappa_{21})</td>
<td>(\mathbb{R}^+)</td>
<td>(\log(\theta))</td>
<td>(\text{Gau}(0.69, 1.00))</td>
</tr>
<tr>
<td>2, 3</td>
<td>(A)</td>
<td>(\mathbb{R})</td>
<td>(\theta)</td>
<td>(\text{Gau}(0.00, 0.04))</td>
</tr>
<tr>
<td>3</td>
<td>(r)</td>
<td>(\mathbb{R}^+)</td>
<td>(\log(\theta))</td>
<td>(\text{Gau}(1.00, 0.25))</td>
</tr>
<tr>
<td>3</td>
<td>(\Delta_1)</td>
<td>(\mathbb{R})</td>
<td>(\theta)</td>
<td>(\text{Gau}(0.00, 0.25))</td>
</tr>
<tr>
<td>3</td>
<td>(\Delta_2)</td>
<td>(\mathbb{R})</td>
<td>(\theta)</td>
<td>(\text{Gau}(0.00, 0.25))</td>
</tr>
</tbody>
</table>
Table 2. Posterior distributional summaries for all unknown parameters for each model. Entries for all except the last row are in the format median [lower quartile, upper quartile].

<table>
<thead>
<tr>
<th>Parameter</th>
<th>Model 1</th>
<th>Model 2</th>
<th>Model 3</th>
</tr>
</thead>
<tbody>
<tr>
<td>$\sigma_\varepsilon^2$</td>
<td>9.21 [8.54, 10.04]</td>
<td>9.3 [8.6, 10.21]</td>
<td>8.9 [8, 9.63]</td>
</tr>
<tr>
<td>$\rho$</td>
<td>$-0.14 [-0.23, -0.06]$</td>
<td>$-0.18 [-0.26, -0.11]$</td>
<td>$-0.16 [-0.24, -0.08]$</td>
</tr>
<tr>
<td>$\sigma_{21}^2$</td>
<td>26.2 [17.43, 41.09]</td>
<td>14.2 [8.46, 26.58]</td>
<td>4.02 [2.61, 6.14]</td>
</tr>
<tr>
<td>$\kappa_{11}$</td>
<td>0.96 [0.74, 1.26]</td>
<td>1 [0.8, 1.26]</td>
<td>1.03 [0.83, 1.25]</td>
</tr>
<tr>
<td>$\kappa_{21}$</td>
<td>0.81 [0.63, 1]</td>
<td>0.62 [0.46, 0.81]</td>
<td>3.65 [1.16, 6.72]</td>
</tr>
<tr>
<td>$A$</td>
<td>0.44 [0.3, 0.55]</td>
<td>0.32 [0.24, 0.4]</td>
<td></td>
</tr>
<tr>
<td>$r$</td>
<td></td>
<td>2.15 [1.82, 2.73]</td>
<td></td>
</tr>
<tr>
<td>$\Delta_1$</td>
<td></td>
<td>0.04 [-0.09, 0.14]</td>
<td></td>
</tr>
<tr>
<td>$\Delta_2$</td>
<td></td>
<td>$-0.07 [-0.726, 0.12]$</td>
<td></td>
</tr>
</tbody>
</table>

$DIC$ 990.84 985.17 982.45