# A new reconstruction of multivariate normal orthant probabilities

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Summary. A new method is introduced for geometrically reconstructing orthant probabilities for non-singular multivariate normal distributions. Orthant probabilities are expressed in terms of those for auto-regressive sequences and an efficient method is developed for numerical approximation of the latter. The approach allows more efficient accurate evaluation of the multivariate normal cumulative distribution function than previously, for many situations where the original distribution arises from a graphical model. An implementation is available as a package for the statistical software R and an application is given to multivariate probit models.

Keywords: Cumulative distribution function; Fractional fast Fourier transform; Multivariate normal distribution; Multivariate probit model; Orthant probabilities; Orthoscheme; Polyhedral cones

#### 1. Introduction

The pervasive role of the multivariate normal distribution in probability and statistics means that its cumulative distribution function is of intrinsic interest and there is value in any additional insight, be it algebraic, geometric or numerical. Moreover, from an applied viewpoint, the cumulative distribution function is a fundamental component of the multivariate probit models that were applied by Chib and Greenberg (1998), Gueorguieva and Agresti (2001), Papathomas and O'Hagan (2005) and others, where efficient computation is important.

We seek efficient accurate computation of orthant probabilities

$$P(X \ge 0) = P(X_i \ge 0, i = 1, ..., p)$$
  $X = (X_1, ..., X_p)' \sim N_p(\mu, \Sigma)$  (1)

for non-singular multivariate normal distributions. Note that, by changing the value of  $\mu$ , we may obtain the multivariate cumulative distribution function at any point in  $R^p$ .

Many approaches have been suggested for numerical approximation of expression (1); for a review see the references in Gassmann et al. (2002) and Miwa et al. (2003). There is a line of geometric thinking running from Schläfli (1858) via Abrahamson (1964) and Miwa et al. (2003) to this paper although we exploit a somewhat different aspect of the geometry.

For many practical purposes, the efficient approach is via Monte Carlo (Genz, 1992, 1993) or random quasi-Monte Carlo (Genz and Bretz, 2002) sampling applied to a transformed version of the integral that is implied by expression (1), an implementation of which is available in R through the mytnorm package (R Development Core Team, 2006; Genz et al., 2005). However, the randomness that is inherent in such methods is sometimes more than a minor nuisance. For example, maximum likelihood estimation for multivariate probit models is much

Address for correspondence: Peter Craig, Department of Mathematical Sciences, University of Durham, South Road, Durham, DH1 3LE, UK. E-mail: P.S.Craig@durham.ac.uk more straightforward when the multivariate normal distribution function may be calculated with sufficient precision to permit use of finite differences to approximate derivatives numerically with respect to parameters. The methods that are described here and in Miwa et al. (2003) make such computations possible.

A brute force approach would require the use of a p-dimensional numerical integration scheme. Efficient general purpose schemes are difficult to derive and implement even for moderately small p and are unsuitable for higher p as the computational storage and effort that are required grow exponentially in p. We shall describe an approach which uses one-dimensional convolutions and which may scale well in many contexts.

The basis of the approach here and in Miwa et al. (2003) is that, for any factorization  $\Sigma = AA'$ , we can write  $X = \mu + AZ$  where  $Z \sim N_p(0, I)$  and so  $P(X \ge 0) = P(A^{-1}\mu + Z \in Q)$  where  $Q = \{x : Ax \ge 0\}$  is a subset of  $R^p$  of a kind that is known as a polyhedral cone.

The interesting consequence is that we can exploit the geometry of cones to represent the original orthant probability in terms of others which are easier to approximate numerically. Miwa et al. (2003) showed that the orthant probability is easily computed when the cone is 'orthoscheme' and that any cone Q may be reconstructed in terms of at most (p-1)! orthoscheme cones; the actual number required depends on the correlation structure of X. A cone is orthoscheme when there is a corresponding X which has a first-order moving average representation. We shall see in Section 3 that orthant probabilities are also fairly easily computed when the random vector may be represented as a first-order auto-regression and in Section 2 that there is a corresponding reconstruction for any orthant probability.

The intuition underlying the new reconstruction suggests that it should be particularly well suited to situations where X has a sparse graphical model. This frequently occurs when the original covariance matrix derives from hierarchical modelling or as a Bayesian network or from conditionally specified spatial modelling. Section 4 compares the efficiency of the two methods of reconstruction for some graphical models and shows that the new method is preferable in many situations.

The author has created a package called orthants (Craig, 2007) for the R statistical software, which implements the methods that are described here and in Miwa et al. (2003). The package may be used together with standard optimization functions in R to maximize the likelihood for multivariate probit models without having to turn to the hybrid Markov chain Monte Carlo expectation—maximization procedure that was introduced by Chib and Greenberg (1998). As an example, the 'six cities' data set of Chib and Greenberg (1998) is analysed in Section 5 by using this direct method.

#### 2. Cone reconstructions

In this section we consider the two geometric approaches to reconstructing an arbitrary cone in terms of other cones which correspond respectively to orthant probabilities for moving average and auto-regressive sequences.

We are concerned with polyhedral cones. A polyhedral cone is a convex subset of  $R^p$  defined by  $Ax \ge 0$  for some  $n \times p$  matrix A, i.e. it is the intersection of n half-spaces. The rows  $a'_1, \ldots, a'_n$  of A are known as the 'face vectors'. The ith face of the cone is the intersection of the cone with  $a_i^{\perp}$  (the hyperplane through the origin orthogonal to  $a_i$ ) and the boundary of the cone is the union of the faces. Faces are orthogonal if and only if their face vectors are. The face vectors may be scaled by positive constants without changing the cone.

A polyhedral cone can also be described in terms of a generating set. Vectors  $v_1, \ldots, v_n$  in  $\mathbb{R}^p$  generate a polyhedral cone by taking linear combinations of the vectors using only non-negative

coefficients; if V is the matrix that is formed by taking the vectors as columns, the cone is  $\{V\lambda:\lambda\geqslant 0\}$  and is unchanged if the vectors are scaled by positive constants.

In particular, a polyhedral cone that is defined by p linearly independent face vectors is also generated by its edge vectors  $v_1, \ldots, v_p$ . Geometrically,  $v_i$  lies in the one-dimensional intersection of all the faces other than the ith face and has positive inner product with  $a_i$ . We say that  $v_i$  is the edge opposite the ith face;  $a_i^{\perp}$  is also the span of the edge vectors other than  $v_i$ . If we choose to scale each  $v_i$  so that  $a_i'v_i = 1$ , then V is the inverse of A.

The sum of cones,  $C_1$  and  $C_2$ , is defined to be  $C_1 + C_2 = \{x_1 + x_2 : x_i \in C_i\}$ . If  $C_1$  and  $C_2$  are defined by generating sets, the union of those sets generates  $C_1 + C_2$ . If the subspaces of  $R^p$  that are spanned by  $C_1$  and  $C_2$ , or equivalently by generating sets of the two cones, are linearly independent (the intersection contains only the zero vector), the values of  $x_1$  and  $x_2$  are unique for each element of  $C_1 + C_2$ .

Returning to multivariate normal distributions, the factorization  $\Sigma = AA'$  is not unique but different choices correspond to rotations of  $R^p$  and the geometry of the cone is the same for any rotation. For a cone Q, the matrix A is not unique but the effect of different choices on  $\Sigma$  is to change only the order and standard deviations of the components of X; the correlation structure is invariant. In effect there is a one-to-one correspondence between cone geometry and correlation structure and so we can move back and forth between the two views. In the reconstructions that are described below, we first express an original orthant probability as a cone probability, then express the cone in terms of a number of other cones and finally express those cone probabilities as orthant probabilities for different multivariate normal distributions. Note that A (or equivalently  $V = A^{-1}$ ) does determine  $\Sigma$  and Q.

2.1. Reconstruction using cones corresponding to moving averages (Miwa et al., 2003) Miwa et al. (2003) expressed the original orthant probability in terms of what they called orthoscheme orthant probabilities. An orthoscheme probability is where  $\Sigma$  is tridiagonal (0 above and below the first off-diagonal band). The term orthoscheme was introduced by Schläfli (1858) and describes the geometry of the cone corresponding to the probability: faces i and j are orthogonal if |i-j| > 1. Stochastically, we can write  $X_1 = \mu_1 + \alpha_1 Z_1$  and  $X_i = \mu_i + \beta_{i-1} Z_{i-1} + \alpha_i Z_i$  for i > 1, i.e. X is a first-order moving average process. If some  $\beta_i = 0$ ,  $\Sigma$  is block diagonal and the original orthant probability factors as the product of two orthoscheme orthant probabilities.

Further examination of the geometry reveals that a cone is orthoscheme and each corresponding  $\beta_i$  non-zero if, and only if, the edges of the cone (suitably scaled) form a sequence of orthogonal projections, i.e., for i > 1,  $v_i$  is the orthogonal projection of  $v_{i-1}$  onto the span of  $v_i, \ldots, v_n$ , or equivalently  $v_{i-1} - v_i$  is orthogonal to  $v_j$  for  $j \ge i$ . For an orthoscheme cone, the moving average representation corresponds to factorizing  $\Sigma = AA'$  so that entries in A are 0 apart from the diagonal and first subdiagonal:  $A_{ii} = \alpha_i$  and  $A_{i+1,i} = \beta_i$ . Then, taking  $\alpha_1^* = \alpha_1^{-1}$ ,  $\gamma_i = (\alpha_i^* \alpha_i)^{-1}$  and  $\alpha_i^* = -\beta_{i-1}\alpha_{i-1}^*/\alpha_i$ , V is also lower triangular with  $V_{ij} = \alpha_i^* \gamma_j$  for  $i \ge j$  and, ignoring the scaling constants  $\gamma_j$ , column i+1 of V is obtained by zeroing the ith entry in column i. The reverse implication follows by rotating a sequence  $v_1, \ldots, v_n$  of orthogonal projections so that the resulting matrix V is lower triangular.

Miwa et al. (2003) showed how an arbitrary cone may be reconstructed in terms of orthoscheme cones. The process sequentially increases the degree to which cones in the reconstruction are orthoscheme. Their language is that a cone is 'orthoscheme of order q' if faces i and j are orthogonal for  $i \le q$  and j > i + 1. The case q = p - 2 is the original definition of orthoscheme. Their presentation is quite algebraic and the following paragraphs emphasize the geometry. The important insight is that the reconstruction is based on constructing the sequence of orthogonal projections of an edge of the original cone.

Begin by picking one edge  $(v_1 \text{ say})$  to play a special role. Define a new edge  $v^*$  by projecting  $v_1$  orthogonally onto the hyperplane  $a_1^{\perp}$  that is spanned by the other edges;  $a_1^{\perp}$  contains the face opposite  $v_1$ . Now define new polyhedral cones  $Q_2, \ldots, Q_p$ :  $Q_i$  has edges  $v_1, v^*$  and all the remaining original edges other than  $v_i$ . Each new cone is orthoscheme of order 1; it has two original face vectors  $(a_1 \text{ and } a_i)$  and p-2 new face vectors that are orthogonal to  $a_1$ . The new cones divide into two groups; disregarding boundary sets of measure 0 which have no probability content, the difference between the union of one group and the other is the original cone and within each group the cones are disjoint. Hence it is easy to compute the probability content of the original cone from the probability contents of the new cones. If  $a_1 \perp a_i$ , the edges of  $Q_i$  are linearly dependent and so  $Q_i$  has dimension less than p; it has no probability content and may be ignored subsequently.

If p > 3, we can now reconstruct each  $Q_i$  in much the same way. Project  $v^*$  onto the span of the remaining original edges other than  $v_1$  to obtain a new edge  $v_i^*$  and form p-2 new cones  $Q_{ij}$  which are orthoscheme of order 2.  $Q_{ij}$  has edges  $v_1$ ,  $v^*$ ,  $v_i^*$  and all the original edges other than  $v_i$  and  $v_j$ . It retains the original face vectors  $a_1$  and  $a_i$  and one other face vector of  $Q_i$  (orthogonal to  $a_1$ ) and has p-3 new face vectors which are orthogonal to both  $a_1$  and  $a_i$ . As before,  $Q_i$  may be expressed in terms of the new cones and its probability content straightforwardly obtained from theirs.

If p > 4, continue by reconstructing each  $Q_{ij}$ ; project  $v_i^*$  onto the span of all the remaining original edges other than  $v_1$  to obtain a new edge  $v_{ij}^*$  and construct p-3 cones  $Q_{ijk}$  with edges  $v_1, v_i^*, v_{ij}^*$  and all the original edges other than  $v_i, v_j$  and  $v_k$ . In total we now have up to (p-1)(p-2)(p-3) cones which are orthoscheme of order 3.

The original edge  $v_1$  is being projected sequentially onto the span of smaller subsets of the other original edges. For any p, the process may be continued until the final projection is onto the span of two original edges, creating a total of at most (p-1)! 'final stage' orthoscheme cones  $Q_{i_1i_2...i_{p-2}}$  having edges  $v_1, v^*, v^*_{i_1}, ..., v^*_{i_1i_2...i_{p-3}}$  and  $v_{i_{p-1}}$  where  $i_{p-1}$  is the (unique) index distinct from  $1, i_1, ..., i_{p-2}$ .

Fig. 1 illustrates the general final projection, and also the initial projection when p = 3. It shows a cross-section through the cone that is generated by the edges which are active in the final projection and how it is divided by the projection into two subcones, each having two

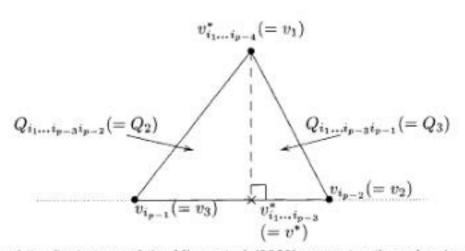


Fig. 1. Illustration of the final stage of the Miwa *et al.* (2003) reconstruction, showing a cross-section of the cone that is generated by  $v_{j_1...j_{p-2}}^*$ ,  $v_{j_{p-2}}$  and  $v_{j_{p-1}}$  in the subspace that is spanned by those vectors: cone  $Q_{j_1...j_{p-3}}$  is represented by the full line triangle and is divided into two parts by the broken line, which shows the direction of orthogonal projection of  $v_{j_1...j_{p-3}}^*$  onto the span of  $v_{i_{p-2}}$  and  $v_{i_{p-1}}$ , the resulting edge being denoted  $v_{j_1...j_{p-3}}^*$  the new cones  $Q_{j_1...j_{p-3},j_{p-2}}$  and  $Q_{j_1...j_{p-3},j_{p-1}}$  are represented respectively by the triangles falling to the left and right of the broken line; the parenthetic labelling illustrates the reconstruction for p=3

orthogonal faces. As shown, the original cone is the union of the two new cones but, if the new edge were to fall outside the original cone, the original would be expressed as the difference between the two new cones.

The description above assumes that  $v^* \neq 0$ ; if  $v^* = 0$ , then  $v_1 \perp v_j$  for j > 1, which implies that  $X_1$  is independent of  $X_2, \ldots, X_n$  and so  $P(X \geqslant 0) = P(X_1 \geqslant 0)$   $P(Y \geqslant 0)$  where  $Y' = (X_2, \ldots, X_n)$  and the procedure may be applied to  $P(Y \geqslant 0)$ . Similarly if, at some later stage,  $v^*_{i_1, \ldots, i_k} = 0$ , the corresponding orthant probability factorizes into two orthant probabilities each of which may be computed by using the reconstruction described.

## 2.2. Reconstruction using cones corresponding to auto-regressive sequences

The new reconstruction that is described here expresses the original orthant probability in terms of orthant probabilities for a number of first-order auto-regressive (Markov) sequences. The original spark for doing so was the duality between moving average and auto-regressive processes: the covariance matrix of any moving average sequence is the precision matrix of some auto-regressive sequence and vice versa. An auto-regressive sequence is one for which  $X_1 = \mu_1 + \alpha_1 Z_1$  and  $X_i = \mu_i + \beta_i (X_{i-1} - \mu_{i-1}) + \alpha_i Z_i$  for i > 1. Thus, writing  $Z = V(X - \mu)$ ,  $V = A^{-1}$  is zero except on the diagonal and first subdiagonal and so the precision matrix  $P = \sum_{i=1}^{n-1} V'V_i$  is tridiagonal. Geometrically, edges i and j of the corresponding cone are orthogonal if |i-j| > 1 and we say that the cone is edge orthoscheme.

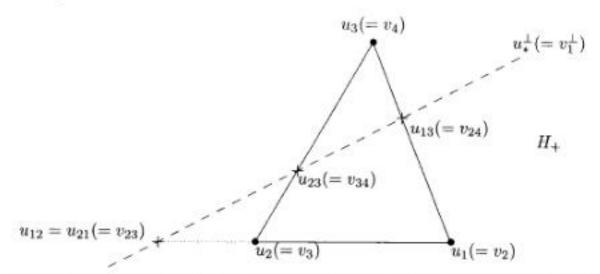
Definition 1. A cone with edges  $v_1, \ldots, v_p$  is said to be edge orthoscheme of order q if  $v_i'v_j = 0$  for  $i \le q$  and j > i + 1. If q = p - 2, the cone is said to be edge orthoscheme.

Miwa et al. (2003) successively obtained cones of higher orthoscheme order. Here, we obtain cones of higher edge orthoscheme order. Whereas their reconstruction replaces one original edge by a new edge at each stage in the process, we shall replace an original face vector by a new face vector and possibly change the sign of others. Full details are given later in the theorems and corollaries but first we give a geometric description.

Start with an original edge  $v_1$  and, for each pair of other edges, define  $v_{ij} = v_{ji}$  to be a linear combination of  $v_i$  and  $v_j$  orthogonal to  $v_1$ . Except in degenerate cases, this exists and is unique apart from norm, which is unimportant, and orientation, which is defined in theorem 1. Form new cones  $Q_2, \ldots, Q_p$  where  $Q_i$  has edges  $v_1, v_i$  and  $v_{ij}$  ( $j \neq i$ ). In  $Q_i$ , all edges except the second are orthogonal to the first, i.e. it is edge orthoscheme of order 1. The face vectors of  $Q_i$  are those of Q with two exceptions: the sign may change for some, and  $a^*$  replaces  $a_i$ ;  $a^*$  is the projection of  $a_1$  onto the span of  $a_2, \ldots, a_p$ , which is computable as a linear combination, orthogonal to  $v_1$ , of  $a_1$  and  $v_1$ . The new cones divide into two groups; the union of one together with Q equals the union of the other, both unions being effectively disjoint. If  $v_i \perp v_1$ ,  $Q_i$  is effectively empty and may subsequently be ignored.

Fig. 2 is primarily an illustration of theorem 1 but may also be seen as an illustration of this process for p = 4. It shows a cross-section through the cone that is generated by  $v_2$ ,  $v_3$  and  $v_4$  in the space that is spanned by the same vectors and also the intersection of the cross-section with the hyperplane orthogonal to  $v_1$ . The original cone Q (which is represented by the triangle  $v_4, v_2, v_3$ ) satisfies the equation  $Q \cup Q_3 = Q_2 \cup Q_4$ .

If p > 3, reconstruct each  $Q_i$  in a similar manner in terms of cones  $Q_{ij}$  where  $Q_{ij}$  has edges  $v_1$ ,  $v_i$  and  $v_{ij}$  and edges orthogonal to  $v_i$  (and  $v_1$ ) formed as linear combinations of  $v_{ij}$  and  $v_{ik}$  for k ranging from 2 to p omitting i and j. Apart from some sign changes, the face vectors of  $Q_{ij}$  are those of  $Q_i$  except that  $a_j$  has been replaced by  $a_i^*$ , which is the projection of  $a^*$  onto the span of the remaining original face vectors other than  $a_1$ . As before,  $Q_i$  may be expressed in terms of the new cones and its probability content obtained from theirs.



**Fig. 2.** Illustration of theorem 1 for n=3 and  $m_-=m_+=2$ , showing a cross-section of the cone that is generated by  $u_1$ ,  $u_2$  and  $u_3$  in the subspace that is spanned by those vectors: the subset that is orthogonal to  $u_*$  is indicated by the broken line below which lies  $H_+$ ; the original cone C is represented by the triangle with vertices  $u_1$ ,  $u_2$  and  $u_3$ ;  $C_+$  is represented by the quadrilateral with vertices  $u_1$ ,  $u_2$ ,  $u_{23}$  and  $u_{13}$  and together with the small triangle  $C_2$  forms the larger triangle  $u_{12}$ ,  $u_{13}$ ,  $u_1$  representing  $C_1$ ; the parenthetic labelling illustrates the first reconstruction stage for p=4

If p > 4, each  $Q_{ij}$  may be reconstructed in terms of cones  $Q_{ijk}$  with edges  $v_1$ ,  $v_i$ ,  $v_{ij}$  and  $v_{ijk}$  and edges orthogonal to  $v_{ij}$  (and  $v_1$  and  $v_i$ ) formed as linear combinations of  $v_{ijk}$  and  $v_{ijl}$  for l ranging from 2 to p omitting i, j and k.

The original face vector  $a_1$  is being projected sequentially onto the span of smaller subsets of the other original face vectors and the precision matrices corresponding to the cones are becoming increasingly tridiagonal. For any  $p \ge 3$ , continue until at stage p-2 the projection is onto the span of two original face vectors and each new cone is edge orthoscheme. There are up to (p-1)! such 'final stage' cones.

Theorem 1. Let  $u_*, u_1, \ldots, u_n$  be linearly independent vectors in  $R^p$  (n < p) such that  $u'_*u_i \neq 0$  for some i, and let C be the cone that is generated by  $u_1, \ldots, u_n$ . Define  $H_+ = \{u : u'_*u \geq 0\}$  and  $C_+ = C \cap H_+$ .

Order and scale  $u_1, \dots, u_n$  so that

$$u'_*u_j = \begin{cases} 1 & j \leq m_+, \\ 0 & m_+ < j \leq m_-, \\ -1 & j > m_- \end{cases}$$

and then, for each  $i \leq m_+$ , define  $u_{ij}$  orthogonal to  $u_*$  for  $j \neq i$  by

$$u_{ij} = \begin{cases} u_i - u_j & \text{for } 1 \leqslant j < i, \\ u_j - u_i & \text{for } i < j \leqslant m_+, \\ u_j & \text{for } m_+ < j \leqslant m_-, \\ u_i + u_j & \text{for } m_- < j \leqslant n. \end{cases}$$

Finally, for  $i \le m_+$ , define  $C_i$  to be the cone that is generated by  $u_i$  and  $u_{ij}$   $(j \ne i)$ . Then, if  $m_+ > 0$ ,

$$C_{+} \cup \bigcup_{\text{even } i} C_{i} = \bigcup_{\text{odd } i} C_{i}$$
 (2)

where, ignoring sets of dimension less than n, the unions on both sides are disjoint. If  $m_+ = 0$ ,  $C_+$  has dimension less than n.

Note that changing the sign of  $u_*$  provides a similar reconstruction of  $C_-$  (the intersection of C with the complement of  $H_+$ ), and so combining the two reconstructions provides a reconstruction for C.

*Proof.* Any  $u \in C$  may be written as a linear combination of  $u_1, \ldots, u_n$  with non-negative coefficients. When  $m_+ = 0$ ,  $u \in C_+$  implies that the coefficient of  $u_j$  is 0 for  $j > m_-$  so  $C_+$  has dimension  $m_-$  and the condition that  $u'_*u_i \neq 0$  for some i implies that  $m_- < n$ .

Now assume that  $m_+ > 0$ . We first show that  $C_i$  is the intersection of  $H_+$  and the set  $\{\Sigma_k b_k u_k : b_k < 0 \text{ for } k < i, b_k > 0 \text{ for } k > i\}$ . For  $u \in C_i$  if and only if  $u = c_i u_i + \sum_{j \neq i} c_j u_{ij}$  where the coefficients  $c_1, \ldots, c_n$  are non-negative, expanding  $u_{ij}$  we obtain

$$u = -\sum_{k < i} c_k u_k + \sum_{k > i} c_k u_k + \left(c_i + \sum_{k < i} c_k - \sum_{i < k \le m_+} c_k + \sum_{k > m_-} c_k\right) u_i.$$

Writing  $b_k$  for the coefficient of  $u_k$ , we see that constraining each  $c_k \ge 0$  is equivalent to constraining  $b_k \le 0$  for k < i,  $b_k \ge 0$  for k > i and  $b_i \ge -\sum_{k < i} b_k - \sum_{i < k \le m_+} b_k + \sum_{k > m_-} b_k$ ; the last constraint may be expressed as  $\sum_{k \le m_+} b_k - \sum_{k > m_-} b_k \ge 0$ , which is equivalent to  $u'_* u \ge 0$ .

For  $1 \le i \le m_+ + 1$ , define  $\Delta_i$  to be the intersection of  $H_+$  and the cone that is generated by  $-u_1, -u_2, \ldots, -u_{i-1}, u_i, u_{i+1}, \ldots, u_n$ . Then  $C_+ = \Delta_1$  and  $C_i = \Delta_i \cup \Delta_{i+1}$  for  $1 \le i \le m_+$ . Observe that  $u'_+u < 0$  if  $b_{m_+} > 0$  for u in the cone that is generated by  $-u_1, \ldots, -u_{m_+}, u_{m_++1}, \ldots, u_n$  and so  $b_{m_+} = 0$  in  $\Delta_{m_++1}$ , which implies that  $\Delta_{m_++1} \subset \Delta_{m_+}$ . Therefore  $C_{m_+} = \Delta_{m_+}$ . Hence, by grouping adjacent pairs of terms in the union  $\bigcup_{i=1}^{m_+} \Delta_i$ , we find that it may be expressed as either side of equation (2). However, ignoring sets of dimension less than n, the  $\Delta_i$  are disjoint and so therefore are those unions.

Theorem 2. In addition to the conditions of theorem 1, suppose that there are  $\tilde{u}_1, \dots, \tilde{u}_q$  which, together with  $u_*, u_1, \dots, u_n$ , form a basis. Let  $C^*$  be the cone that is generated by  $u_*$  and  $\tilde{u}_1, \dots, \tilde{u}_q$  and define  $\tilde{C} = C^* + C$ ,  $\tilde{C}_i = C^* + C_i$  and  $\tilde{C}_+ = C^* + C_+$ . Then either  $m_+ = 0$ , in which case  $\tilde{C}_+$  has dimension less than p, or

$$\tilde{C}_{+} \cup \bigcup_{\text{even } i} \tilde{C}_{i} = \bigcup_{\text{odd } i} \tilde{C}_{i}$$
 (3)

where, ignoring sets of dimension less than p, the unions on both sides are disjoint. Changing the sign of  $u_*$  when applying theorem 1 leads to a reconstruction of  $\tilde{C}_- = C^* + C_-$ .

*Proof.* Since the vectors generating  $\tilde{C}$  are a basis, the subspaces that are spanned by  $C^*$  and C are linearly independent. Moreover, the subspaces that are spanned by  $C_i$  and  $C_+$  are subsets of the span of C and so are linearly independent of the span of  $C^*$ . Therefore, every vector lying in a cone appearing in equation (3) has a unique representation as the sum of a vector in  $C^*$  and a vector from the corresponding cone in equation (2). Consequently, equation (2) implies equation (3) and the disjointness property is inherited.

Corollary 1. Suppose that  $\tilde{C}$  is edge orthoscheme of order q but not of order q+1. Then  $\tilde{C}$  may be expressed as the difference of two disjoint unions of cones each of which is edge orthoscheme of order q+1.

Proof. Denote the edges of  $\tilde{C}$  by  $\tilde{u}_1, \dots, \tilde{u}_q, u_*, u_1, \dots, u_n$ . Since  $\tilde{C}$  is not orthoscheme of order q+1,  $u_*u_i \neq 0$  for some i and so we may apply theorem 2. Since  $\tilde{C}$  is edge orthoscheme of order q, each  $\tilde{u}_j$  in theorem 2 is orthogonal to each  $u_k$  and hence to all the edges of  $C_i$  in theorem 1. Moreover, all edges of  $C_i$  other than  $u_i$  are orthogonal to  $u_*$  and so each  $\tilde{C}_i$  is edge orthoscheme of order q+1. In effect  $u_*$  becomes  $\tilde{u}_{q+1}$ .

Corollary 2. Suppose that Q is edge orthoscheme of order q. Then the probability content of Q may be expressed as the difference of two sums of probabilities for cones which have edge orthoscheme order q+1.

# 3. Numerical computation

Recursive application of corollary 2 reconstructs an arbitrary cone in terms of at most (p-1)! final stage edge orthoscheme cones. However, the rescaling of  $u_i$  in the statement of theorem 1 is there only to make the proof easier to read and if implemented naively can easily lead to floating-point overflow or underflow. It is wiser to stabilize the magnitudes of the vectors that are involved although there is no uniquely appropriate scaling. Our implementation scales all edge vectors to have unit norm. Thus, to find the edges of  $C_i$ , compute  $u_{ij}$  as a linear combination of  $u_i$  and  $u_j$  orthogonal to  $u_*$ , such as  $(u'_*u_i)u_j - (u'_*u_j)u_i$ , and then scale the result to have the correct sign and norm.

Theorem 1 does not determine the ordering of  $u_1, \ldots, u_{m-}$ ; our implementation orders them to have decreasing inner product with  $u_*$ , which corresponds to increasing angle. We do not know this to be the best choice in general. However, if we exchange  $u_1$  and  $u_2$  in Fig. 2, we change the sign of  $u_{12}$ ; it is difficult to visualize exactly what happens but it appears that cones  $C_1$  and  $C_2$  become much larger in relation to  $C_+$  and the consequence is that the probability content of Q is expressed as the difference between two larger probabilities, which results in loss of precision.

The formal theory has been expressed entirely in terms of edges with no reference to faces; in effect it tells us how to compute V (and hence P) for each new cone. However, to compute the probability for an edge orthoscheme cone by using the methods in Section 3.2, we shall need to know the corresponding  $\Sigma$  and we also need to know the diagonals of  $\Sigma$  and P in Section 3.1. We could directly compute  $\Sigma = AA'$  from  $A = V^{-1}$  for each cone. However, it is not difficult to compute face vectors of each cone at the same time as the edges and, if we maintain the scaling where each has unit inner product with the opposite edge vector, we then know A and hence  $\Sigma$ . A further advantage is that, by computing A and V for each cone to some degree separately, we may detect numerical instability if the product of the resulting matrices A and V is significantly different from the identity.

Let the face vectors of  $\tilde{C}$  in theorem 2 be  $\tilde{a}_1, \dots, \tilde{a}_q, a_*$  and  $a_1, \dots, a_n$  labelled corresponding to the opposite edges and assume that edge and face vectors have been scaled as described in the previous paragraphs. Then, apart from rescaling to have unit inner product with opposite edge vectors, the face vectors of  $\tilde{C}_i$  are the same except that

- (a) for  $j \neq i$ ,  $a_j$  is now opposite  $u_{ij}$  and needs to change sign when j < i and
- (b) the face vector that is opposite u<sub>i</sub>, formerly a<sub>i</sub>, is now proportional to u<sub>\*</sub> − a<sub>\*</sub> when q = 0 and u<sub>\*</sub> − a<sub>\*</sub> − (ũ'<sub>q</sub>u<sub>\*</sub>)ã<sub>q</sub> for q > 0.

### 3.1. Testing and allowing for near orthogonality of edge vectors

In finite precision arithmetic, we cannot decide whether two edges are exactly orthogonal. Practically, we should treat edges as orthogonal if they are nearly so. Otherwise there are undesirable consequences:

- (a) we must calculate orthant probabilities for cones which have nearly parallel edges and these are difficult to approximate numerically by using the methods that are described later.
- (b) the decomposition becomes numerically unstable as the recursion progresses and
- (c) we compute orthant probabilities for a larger than necessary number of final stage cones.

Thus, we must understand when it is reasonable to treat two edges as orthogonal when applying theorem 2. For multivariate normal  $X_i$ ,  $-P_{ij}/P_{jj} = -v_i'v_j/v_j'v_j$  is the coefficient of  $X_i$  in the conditional mean of  $X_i$  given all the other components of X. We now examine the sensitivity of the orthant probability to changes in the coefficient.

Without loss of generality (simple reordering), suppose that j = p. Write  $W = X - \mu$  and  $W_p|W_1, \ldots, W_{p-1} \sim N(\Sigma_{k < p} b_k W_k, \sigma^2)$ . Then, denoting the standard normal probability density function by  $\phi$ .

$$\left| \frac{\partial}{\partial b_i} P(X \geqslant 0) \right| \leqslant \int_{-\mu_1}^{\infty} \dots \int_{-\mu_p}^{\infty} \left| \frac{\partial}{\partial b_i} p(w_1, \dots, w_p) \right| dw_1 \dots dw_p$$

$$\leqslant \int_{R^p} \frac{|w_p - \sum_k b_k w_k|}{\sigma} \frac{|w_i|}{\sigma} \frac{1}{\sigma} \phi \left( \frac{w_p - \sum_{k < p} b_k w_k}{\sigma} \right) p(w_1, \dots, w_{p-1}) dw_1 \dots dw_p$$

$$= \frac{E[|Z_p|] E[|W_i|]}{\sigma} = \frac{2\sum_{ij}^{1/2}}{\pi \sigma}$$

and so the error that is induced by setting  $b_i = 0$  is less than  $(2\Sigma_{ii}^{1/2}/\pi\sigma)|b_i|$ . For our choice of scaling,  $\sigma^2 = P_{jj} = v_j'v_j = 1$  and  $\Sigma_{ii} = a_i'a_i$ . Therefore, we compare  $v_i'v_j(a_i'a_i)^{1/2}$  with some threshold, the correct choice of which is hardware dependent. For standard double precision, 10-12 works well and permits the computation of orthant probabilities to approximately this accuracy. Larger thresholds restrict computational accuracy whereas smaller values lead to numerical instability.

# 3.2. Approximation of orthant probabilities for auto-regressive sequences

We wish to compute  $P(X \ge 0)$  when  $X_1, \dots, X_p$  are an auto-regressive sequence. To make the calculation easier to structure, to simplify the subsequent presentation and to facilitate analysis of accuracy, we shall restrict to the case of unit variances and express the calculation in terms of  $W = X - \mu$ . Then, writing  $\rho_i = \operatorname{corr}(X_i, X_{i+1})$  and  $\sigma_i^2 = 1 - \rho_i^2$ , the sequence  $W_1, \dots, W_p$  is Markov and  $W_{i+1}|W_i \sim N(\rho_i W_i, \sigma_i^2)$ .

Taking  $\psi_1(w_1) = p(w_1) = \phi(w_1)$ , we sequentially compute approximations to the functions  $\psi_2, \dots, \psi_p$  where

$$\psi_{n+1}(w_{n+1}) = \int_{-\mu_n}^{\infty} \dots \int_{-\mu_1}^{\infty} p(w_1, \dots, w_{n+1}) \, dw_1 \dots dw_n 
= \int_{-\mu_n}^{\infty} p(w_{n+1}|w_n) \int_{-\mu_{n-1}}^{\infty} \dots \int_{-\mu_1}^{\infty} p(w_1, \dots, w_n) \, dw_1 \dots dw_{n-1} \, dw_n 
= \frac{1}{\sigma_n} \int_{-\mu_n}^{\infty} \phi\left(\frac{w_{n+1} - \rho_n w_n}{\sigma_n}\right) \psi_n(w_n) \, dw_n 
= \frac{1}{|\rho_n|\sigma_n} \int_{-|\rho_n|\mu_n}^{\infty} \phi\left\{\frac{w_{n+1} - \operatorname{sgn}(\rho_n)u}{\sigma_n}\right\} \psi_n\left(\frac{u}{|\rho_n|}\right) du$$
(5)

so that  $P(X \ge 0)$  is obtained by integrating  $\psi_p$  from  $-\mu_p$  to  $\infty$ .

The Markov property has previously been exploited, as in equation (4), by Armitage et al. (1969), Williams (1971), Hirotsu et al. (1992) and others to simplify calculation of particular multivariate normal probabilities arising from statistical tests. Hayter (2006) provided a more general framework for exploiting factorization of joint densities for efficient computation of a variety of orthant and other probabilities. Armitage et al. (1969) also considered efficient computation for their special case which has additional smoothness properties.

In principle, calculation of equation (4) appears to require the numerical evaluation of a one-dimensional integral for each value of  $w_{n+1}$ . However, equation (5) is in the form of a convolution and may be amenable to calculation by using the fast Fourier transform (FFT). The FFT approach is efficient except when  $\sigma_n$  is small; moreover there is an efficient linear filtering algorithm to handle small  $\sigma_n$ . When  $\rho_n$  is 0 (or effectively so), equation (4) reduces to integrating  $\psi_n$  and restarting the process with  $\psi_{n+1}$  proportional to  $\phi(w_{n+1})$ .

For numerical calculations, we replace the lower limit of integration in equation (4) by  $L_n = \max(-U, -\mu_n)$  and the upper limit by U. In doing so, we lose some probability from the final estimate of  $P(X \ge 0)$  but the loss is less than  $\sum_i P(|W_i| > U) = 2p \Phi(-U)$  where  $\Phi$  is the standard normal distribution function. When linear filtering, we also truncate  $\phi$  to 0 outside the interval [-U, U] and again lose less than  $2p \Phi(-U)$ . In practice, it should be easy to choose U so that the total loss is negligible; for double precision, U = 8 works well.

The strategy is to approximate  $\psi_n$  on a grid (a sequence of equally spaced values) covering  $[L_n, U)$  with a spacing  $\Delta_n$  which we try to keep close to a fixed value  $\Delta$  which is our basic control on the accuracy of approximation. The decision concerning when to use FFT or linear filtering relies on estimating the relative computational effort that is required, which depends on both hardware and the choice of FFT library. Details of this and other aspects of the algorithm may be found in a document distributed as part of the orthants (Craig, 2007) package for R.

# 3.2.1. Using the fast Fourier transform

The Fourier transform of  $\psi_{n+1}$  is

$$\begin{split} \Psi_{n+1}(t) &= \int \exp(\mathrm{i}tx) \, \psi_{n+1}(x) \, \mathrm{d}x = \frac{1}{\sigma_n} \int \exp(\mathrm{i}tv) \, \phi\left(\frac{v}{\sigma_n}\right) \, \mathrm{d}v \int_{-\mu_n}^{\infty} \exp(\mathrm{i}t\rho_n y) \, \psi_n(y) \, \mathrm{d}y \\ &= \exp\left(-\frac{\sigma_n^2 t^2}{2}\right) \Psi_n^+(\rho_n t) \end{split}$$

where  $\Psi_n^+$  is the Fourier transform of the truncation of  $\psi_n$  to 0 below  $-\mu_n$ . Thus we can compute  $\psi_{n+1}$  by finding the Fourier transform of the truncated version of  $\psi_n$ , stretching it by factor  $1/\rho_n$ , multiplying by  $\exp(-\sigma_n^2 t^2/2)$  and finally inverting the Fourier transform.

Numerically, we can approximate this operation by using the fractional FFT of Bailey and Swarztrauber (1991) to compute Fourier transforms on suitable grids of values. The fractional FFT differs from the conventional FFT in that the output fundamental frequency (the spacing of the grid of values of t) may be freely chosen rather than being determined by the spacing in the x-grid and the number of grid points. As described in detail in section 13.9 of Press et al. (1993), we may increase the accuracy of the calculation of a Fourier integral for a fixed grid spacing by using Filon's (1928) method for Fourier integrals of continuous functions. The Filonized fractional FFT is most efficient when the grid length is a power of 2.

The main numerical difficulty is that the Fourier transform of a truncated smooth function decays to zero very slowly (asymptotically at rate 1/t). Hence the decay to zero of  $\exp(-\sigma_n^2 t^2/2)$  is the limiting control on the range over which we need to approximate  $\Psi_{n+1}$  to compute  $\psi_{n+1}$ . In effect, this also determines the range over which we must evaluate  $\Psi_n^+$ . For  $\sigma_n = 1$ , approximating  $\Psi_{n+1}$  on [-U, U] should suffice but in general the range needs to be proportional to  $1/\sigma_n$ , which means that we must either use a larger grid or lose precision by increasing the spacing; we do the former as it is easy to use the Filonized fractional FFT to approximate  $\Psi_n^+$  on a grid whose length is an integer multiple  $g_n$  of the length of the grid holding the truncation of  $\psi_n$ . In principle,  $g_n$  could become arbitrarily large. However, the difficulty only arises when

 $\sigma_n$  is near to 0. In such situations, the convolution may be computed directly by linear filtering of the function  $\psi_n$  evaluated on a grid.

# 3.2.2. Using linear filtering

Linear filtering is an efficient way to approximate equation (5) when  $\sigma_n$  is small since the integral will effectively only involve a small range of u for each x. Suppose that we have already calculated an approximation to  $\psi_n(x)$  so that  $\hat{\psi}_{n,k} \approx \psi_n(x_{n,k})$  for  $0 \le k < G$  where  $x_{n,k} = L_n + k\Delta_n$ . Linear filtering means that we plan (for most values of k) to compute  $\hat{\psi}_{n+1,k} = \sum_{m=0}^{M_n} c_{n,m} \hat{\psi}_{n,k_0+k+m}$  where  $k_0$  (depending on n) is an offset between the grids; for negative  $\rho_n$ , replace  $\hat{\psi}_{n,k_0+k+m}$  by  $\hat{\psi}_{n,k_0-k+m}$ . For this to be possible, the grids that are used for x and u must have the same spacing and so  $\Delta_{n+1} = |\rho_n|\Delta_n$ .

We could easily obtain coefficients  $c_{n,m}$  by using an elementary quadrature formula such as Simpson's rule to approximate equation (5). However, because  $\sigma_n$  is small, the integrand in equation (5) is not really sufficiently smooth for such a simple approach. The smoothness of the underlying multivariate normal probability density function means that  $\psi_n$  should be smooth apart from the step at  $x = -\mu_n$ , a view that is supported by some numerical experiments. This suggests the use of an approach to computing the convolution which takes advantage of the smoothness of  $\psi_n$  while allowing for the lack of smoothness in  $\phi(\cdot/\sigma_n)$ .

We replace  $\psi_n$  in equation (5) by a piecewise cubic interpolant. On each interval  $(x_{n,k}, x_{n,k+1})$ , approximate  $\psi_n(x)$  by the cubic interpolant of  $\hat{\psi}_n$  for the four nearest grid points. Hence, the contribution of the interval to  $\hat{\psi}_{n+1,j}$  may be obtained as a linear combination of the values of  $\hat{\psi}_n$  at those points. The coefficients may be calculated as fairly simple expressions in terms of the functions  $\phi$  and  $\Phi$ . To avoid excessive truncation error, those expressions need to be replaced by series approximations in certain situations.

# 4. Numerical examples

Table 1 compares accuracy and computation time for orthant probabilities computed by using the two reconstructions. It also shows the numbers  $N_{\rm m}$  and  $N_{\rm c}$  of final stage cones that are required respectively by the Miwa et al. (2003) method and the new method. The two examples considered are drawn from Miwa et al. (2003). In both cases, p=9 and  $\mu=0$  and it can be shown that the true orthant probability is 1/10. In the first example, the components of X are exchangeable (equicorrelated) with unit variance and correlation  $\rho=0.5$ . In the second,  $\Sigma^{-1}$  is tridiagonal having unit diagonal and -0.5 on the bands.

In the exchangeable case, the two methods of reconstruction require the same number of final stage cones and the new method is typically 5–10 times as long for comparable accuracy, requiring a grid up to an order of magnitude larger. The second example is in fact an auto-regressive sequence. Consequently, the new method needs just one final stage cone which is the original cone whereas  $N_{\rm m}=323$ . This disparity is enough to make the new method faster by an order of magnitude for comparable accuracy. Generally, doubling the grid size typically reduces error by an order of magnitude and roughly doubles the computation time although the computation time grows a little faster for the new method. It should be noted that, for the usual situation where the true answer is not known, additional computational effort would need to be expended to ascertain the accuracy of calculation for a particular grid size.

The time that is required for computing the reconstructions is negligible compared with the time that is spent computing final stage cone probabilities. Ratios of time per cone for the two methods are not entirely consistent owing to the variable grid sizing in our algorithm and the choice of FFT versus linear filtering but, across a range of examples including those in Table 1,

Table 1. Error and computation time depending on grid size by using the two methods of reconstruction to compute orthant probabilities for two nine-dimensional examples?

G	Results for example (i)				Results for example (ii)				
	Orthoscheme, $N_{\rm m} = 40320$		Edge orthoscheme, $N_c = 40320$		Orthoscheme, N <sub>m</sub> = 323		Edge orthoscheme, $N_{\rm C} = I$		
	Error	Time (s)	Error	Time (s)	Error	Time (s)	Error	Time (s)	
32	$-2.5 \times 10^{-5}$	8.2	5.0×10 <sup>-3</sup>	6.3	1.4×10 <sup>-5</sup>	0.07	$-4.2 \times 10^{-3}$	0.00	
64	$-1.0 \times 10^{-6}$	16.2	$1.9 \times 10^{-3}$	11.6	$6.2 \times 10^{-7}$	0.12	$-5.6 \times 10^{-4}$	0.00	
128	$-5.1 \times 10^{-8}$	32.3	$1.4 \times 10^{-4}$	24.2	$3.2 \times 10^{-8}$	0.25	$-5.4 \times 10^{-5}$	0.00	
256	$-3.2 \times 10^{-9}$	64.4	$4.7 \times 10^{-6}$	52.5	$2.0 \times 10^{-9}$	0.50	$-3.8 \times 10^{-6}$	0.00	
512	$-1.8 \times 10^{-10}$	128.7	$2.9 \times 10^{-7}$	120.3	$1.1 \times 10^{-10}$	0.99	$-2.4 \times 10^{-7}$	0.01	
1024	$-1.1 \times 10^{-11}$	258.1	$1.6 \times 10^{-8}$	279.4	$6.9 \times 10^{-12}$	1.97	$-1.7 \times 10^{-8}$	0.02	
2048	$-1.0 \times 10^{-12}$	518.0	$1.0 \times 10^{-9}$	698.3	$4.7 \times 10^{-13}$	3.96	$-2.2 \times 10^{-9}$	0.04	
4096	$-9.3 \times 10^{-14}$	1037.9	$-8.0 \times 10^{-11}$	2130.1	$7.7 \times 10^{-14}$	7.94	$-1.3 \times 10^{-10}$	0.08	

<sup>†</sup>Example (i) is a unit variance exchangeable sequence with  $\rho = 0.5$  and example (ii) is tridiagonal  $P = \Sigma^{-1}$  having  $P_{i,i} = 1$  and  $P_{i,i+1} = -0.5$ .  $N_{\rm m}$  and  $N_{\rm c}$  are the numbers of final stage cones that are required by using respectively the reconstruction of Miwa *et al.* (2003) and the new reconstruction. G is the number of grid points that were used in computing final stage cones.

it appears that for the same grid size the new method takes typically twice as long per cone. Overall, it seems that  $N_c/N_m = 0.1$  is approximately the point at which the two methods typically deliver the same accuracy for the same amount of effort.

Miwa et al. (2003) compared the performance of their method with the Monte Carlo method of Genz (1992) and found that for low accuracy the Monte Carlo method was much faster, whereas for high accuracy their method was usually faster unless the number of final stage cones became very great. We have compared both reconstruction approaches with the random quasi-Monte Carlo algorithm QRSVN of Genz and Bretz (2002) as implemented in the mytnorm package for R. Table 1 shows that both reconstruction approaches typically obtain an order of magnitude reduction in error by doubling computational time. A small simulation study using the same two examples suggested that, for QRSVN, the typical magnitude of error for a particular probability is inversely proportional to the computation time raised to power  $\frac{2}{3}$ , which is better than ordinary Monte Carlo sampling but still implies that an order of magnitude reduction in error requires a 30-fold increase in effort. A direct comparison of typical error and the errors in Table 1 showed that, for example (i), QRSVN was slower than the Miwa et al. (2003) method for absolute error less than  $10^{-8}$ ; for example (ii), QRSVN was slower than either reconstruction method for error less than  $10^{-8}$ ; for example (ii), QRSVN was slower than either reconstruction method for error less than  $10^{-4}$ .

# 4.1. Correlation structure and number of final stage cones

A key issue in comparing the efficiency of orthant probability calculations based on the two reconstructions is the number of final stage cones that are required by each method, which we lack the insight to compute theoretically. The following empirical results are subject to some numerical instability due to the threshold that was used to decide when to treat nearly orthogonal edges as orthogonal; the method of Miwa et al. (2003) requires a similar threshold.

In general, the number of cones that is required depends on which edge or face is chosen to be the special one to start the algorithm. For the new method, the ratio of worst to best number of cones varied by two orders of magnitude between examples but was generally greater than the typical value of 7 for the method of Miwa et al. (2003). The empirical evidence suggests that the optimal choice is a parentless or childless node for directed graphs and is a corner for the spatial lattice. The following results relate to the optimal choice.

Exchangeability leads to worst case behaviour for both methods, the number of cones being  $N_e = (p-1)!$ . For moving average and auto-regressive sequences respectively, the Miwa et al. (2003) and new methods require just a single cone. Duality of the two approaches to reconstruction indicates that the number of cones that was required by Miwa et al. (2003) for an auto-regressive sequence should be the same as for the new method applied to a moving average sequence. Empirically, some differences were observed and are believed to be due to different thresholding for near orthogonality in the two implementations. The original intuition underlying the development of the new reconstruction was that  $N_c$  would be substantially less than  $N_m$  for X having a sparse directed or undirected graphical model. This turns out to be true in many but not all cases.

Fig. 3 shows results for auto-regressive sequences, the directed net and binary tree graphs (Figs 4(a) and 4(b)) and lattice graphs (Fig. 4(c)) with undirected and directed edges. Fig. 3(a) compares the number of cones  $N_{\rm m}$  that were required by Miwa et al. (2003) with  $N_{\rm e}$  whereas Fig. 3(b) compares the number  $N_{\rm c}$  that is required by the new method with  $N_{\rm m}$ . We see that  $N_{\rm m}$  is often very much lower than the worst case and  $N_{\rm c}$  often much lower than  $N_{\rm m}$ , both ratios generally decreasing as the dimension p increases for a particular graph. Overall, the ratios are lowest for the auto-regressive case, lower for the binary tree than for the directed net and generally lower for the directed net than for the spatial lattices where the ratios are nearly always slightly lower for the undirected than the directed version of the same graph. The shape of the lattice affects the outcome for lattices of the same dimension.

For 'loop' graphs (Fig. 4(d)), p = m + n. Investigating all loops with  $p \le 15$ , we found that  $N_c = 2^p/4$  whereas  $N_m \approx 3.24^p/25$ . Thus, for higher p, loop graphs would lie between autoregressions and binary trees in Fig. 3(a) and near binary trees in Fig. 3(b).

All star graphs (Fig. 4(e)) having up to 13 nodes were investigated. Empirically, we found that  $N_c$  depends on the number of edges on all except the longest branch. Labelling the branches so that  $m_1$  is the largest,  $N_c$  is the multinomial coefficient

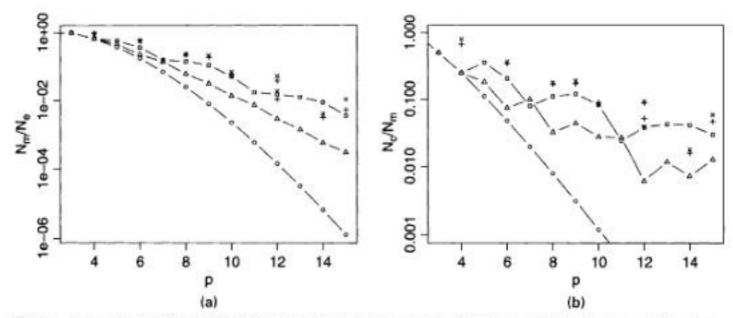


Fig. 3. Dependence of the optimal number of final stage cones on dimension p and various correlation structures ( $\bigcirc$ , auto-regression;  $\square$ , directed net;  $\triangle$ , binary tree; +, lattice; ×, directed lattice): (a) number of cones,  $N_m$ , that are required by the method of Miwa *et al.* (2003) relative to the worst case number,  $N_e = (p-1)!$ ; (b) number required by the new method,  $N_c$ , relative to  $N_m$ 

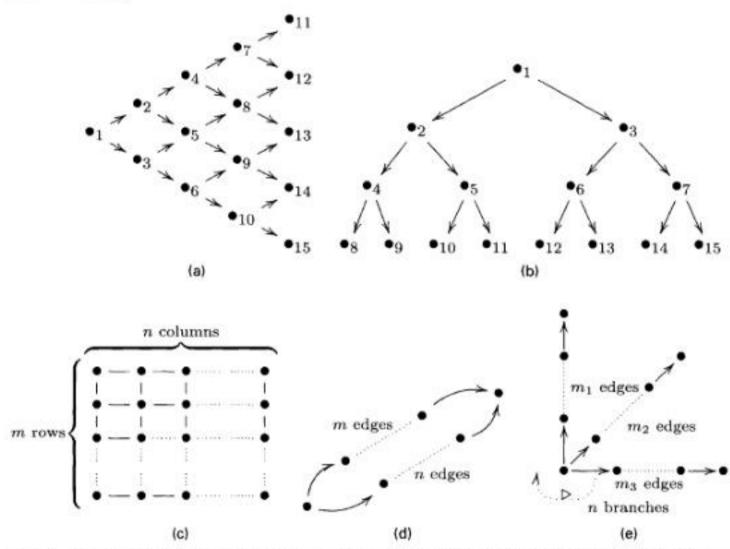


Fig. 4. Graphical models for which numbers of final stage cones are considered (for models (a) and (b), the graph that is used for p < 15 is obtained by deleting higher numbered nodes and the edges that are connected to them): (a) directed net (p = 15); (b) binary tree (p = 15); (c) lattice (p = mn); (d) loop (p = m + n); (e) star  $(p = m_1 + ... + m_n + 1)$ 

$$(m_2+\ldots+m_n)!\Big/\prod_{i=2}^n m_i!.$$

For fixed p, the maximum of  $N_c$  is (p-2)!, which is attained by the graph having p-1 branches each of length 1; this is an example of a sparse graph which requires many cones. In contrast, an auto-regression is a star graph (n=1) for which  $N_c=1$ . For any star, merging branches i and j reduces  $N_c$  by a factor of at least  $\binom{m_i+m_j}{m_j}$ . In all cases with p>2,  $N_c< N_m$ , and frequently much less. A purely empirical approximation to the number of Miwa cones is  $N_m \approx 0.03 \sqrt{(8^p N_c)}$ , which was almost always accurate to within a factor of 2 although this is unlikely to extrapolate to substantially higher p.

The overall picture which emerges is that the new method is generally competitive for these graphs for  $p \ge 7$  and nearly always preferable for high p although  $N_c$  may then be so large as to be impractical.

# A statistical application

The frequently cited paper by Chib and Greenberg (1998) on multivariate probit modelling includes an example which we shall reconsider here. In probit modelling, binary outcomes are obtained by thresholding latent normally distributed continuous variables.

Chib and Greenberg (1998) considered a subset of the data from the 'six cities' longitudinal study on health effects of air pollution. Their data recorded whether or not each of 537 children suffered from wheeze at ages 7, 8, 9 and 10 years and whether or not the mother smoked during the first year of the study. Denoting the outcome (wheezing status) for child i at age j by  $y_{ij}$  (j = 7, 8, 9, 10) and the mother's (binary) smoking status by  $s_i$ , the model is that  $y_{ij} = 1$  when  $w_{ij} > 0$  where the vector  $w_i$  follows a multivariate normal distribution with variance matrix  $\Sigma$  and the covariates are incorporated through  $E[w_{ij}] = \beta_1 + \beta_2(j-9) + \beta_3 s_i + \beta_4(j-9) s_i$ . For identifiability,  $\Sigma$  is taken to be a correlation matrix and three different correlation structures were considered by Chib and Greenberg (1998): independence (no correlation parameters), exchangeability (one parameter) and the general correlation matrix (six parameters). Here, in addition, we consider auto-regressive structure (three parameters).

Chib and Greenberg (1998) considered maximum likelihood estimation, using a Monte Carlo expectation–maximization algorithm where the expectation step is with respect to the latent variables  $w_i$  and is approximated by simulation from the conditional distribution of  $w_i$  given  $y_i$  and the parameters. That conditional distribution is a truncated multivariate normal distribution; there is no direct sampling method and Gibbs sampling is used.

It is not trivial to implement the Chib and Greenberg (1998) method in a statistical package such as R or S-PLUS. However, the likelihood for a single observation may be calculated directly as an orthant probability and so computation of the likelihood function takes only a few lines of code in R provided that a suitable function is available for computing orthant probabilities. There is already the mytnorm package, which is based on Monte Carlo integration (see Genz et al. (2005)). However, it is not suitable for use of standard optimization routines to maximize the likelihood, which require the objective function to be deterministic and which, for efficiency, rely on the use of finite differences to approximate partial derivatives of the likelihood function. The methods that are presented here and in Miwa et al. (2003) are suitable for that approach and the various models were easily estimated by using R.

Table 2 shows the results that were obtained for the four models fitted to the six cities data. The labelling  $M_i$  coincides with that used by Chib and Greenberg (1998) and the estimates here

β	Results for the following models:										
	$\mathcal{M}_1$		Mar		$\mathcal{M}_2$		$\mathcal{M}_3$				
	-1.122 -0.078 0.159	0.062 0.031 0.101	-1.130 -0.079 0.155	0.062 0.036 0.100	-1.119 -0.078 0.161	0.062 0.030 0.100	-1.126 -0.077 0.171	0.047 0.038 0.076			
ρ	0.037 0.585 0.524	0.051 0.066 0.072	0.039 0.671 0.728	0.058 0.060 0.049	0.039 0.599	0.049 0.041	0.037	0.061			
	0.579 0.687 0.559 0.631	0.074 0.056 0.074 0.067	0.623	0.060							
$L(\hat{\theta})$	-794.74		-802.70		-797.67		-909.72				

Table 2. Maximum of likelihood, maximum likelihood estimates and standard errors for the four models that were fitted to the six cities wheeze data?

 $\dagger M_1$  is the unrestricted model,  $M_{ar}$  is the auto-regressive model,  $M_2$  is the exchangeability model and  $M_3$  is the independence model. differ only by at most 0.001. The standard errors also agree except for  $\mathcal{M}_2$  where it appears that their values are erroneous as they differ substantially from the Bayesian posterior standard deviations that they quoted for that model. Their method does not provide a value for the loglikelihood at the maximum although they could, for example, have used the methods that are available in the mytnorm package to compute them with very good accuracy.

We note also that the availability of derivatives of the log-likelihood function makes it straightforward to implement for this problem a Bayesian analysis using a Laplace approximation to the posterior distribution as the basis for a Metropolis-Hasting independence sampler written in R; hence it was easy to verify the Bayes estimates that were provided in Chib and Greenberg (1998), which they obtained by random-walk Metropolis-Hastings steps with latent variables. The availability of the likelihood function also makes it possible to use the method of Chib and Jeliazkov (2001) to compute marginal likelihoods without having to use kernel density estimation of the posterior. The values that were obtained agreed with those of Chib and Greenberg (1998) up to the first decimal place for  $M_1$  and  $M_3$  but differed for  $M_4$  where the value that was obtained here was -926.8 compared with their value of -931.2.

Finally, we note that Chib and Jeliazkov (2006) used auto-regressive normally distributed latent variables as the basis for a semiparametric analysis of longitudinal binary data and that such calculations should also benefit from the availability of efficient methods, such as proposed here, for computing orthant probabilities when the precision matrix is sparse. A detailed account is beyond the scope of this paper.

# 6. Possible developments

Geometrically, there are some interesting possible areas of further development. First, exchangeable sequences correspond to cones where the angle between each pair of edges is the same as is the angle between each pair of faces. It is conceivable that there might be another method of reconstruction which uses such cones. More generally, first-order moving averages, first-order auto-regressions and exchangeable sequences have in common the fact that the sequence of conditional means  $E_k = E[X_k | X_1, \dots, X_{k-1}]$  is Markov. In the case of first-order moving averages,  $E_2, \dots, E_p$  are independent, which is why orthant probability computation is fundamentally more efficient for moving average sequences than for auto-regressive sequences. The general situation in which the sequence  $E_2, \dots, E_p$  is Markov is for a (non-stationary) auto-regressive moving average ARMA(1,1) process. In principle, it should be possible to find a more efficient reconstruction of an arbitrary cone in terms of such cones. However, the geometry is more complex and it is not obvious how to exploit it.

Alternatively, we could seek ways to augment the quantities considered to increase computational efficiency. For example, the orthant probability for an exchangeable sequence may be computed efficiently by using Hayter's (2006) method if we augment the sequence by the 'population mean' so that the directed acyclic graph becomes a tree. The idea may be extended to provide efficient computation for the structures in Papathomas and O'Hagan (2005).

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#### References

- Abrahamson, I. G. (1964) Orthant probabilities for the quadrivariate normal distribution. Ann. Math. Statist., 35, 1685-1703.
- Armitage, P., McPherson, C. K. and Rowe, B. C. (1969) Repeated significance tests on accumulating data. J. R. Statist. Soc. A, 132, 235-244.
- Bailey, D. H. and Swarztrauber, P. N. (1991) The fractional Fourier transform and applications. SIAM Rev., 33,
- Chib, S. and Greenberg, E. (1998) Analysis of multivariate probit models. Biometrika, 85, 347–361.
- Chib, S. and Jeliazkov, I. (2001) Marginal likelihood from the Metropolis-Hastings output. J. Am. Statist. Ass., **96**, 270–281.
- Chib, S. and Jeliazkov, I. (2006) Inference in semiparametric dynamic models for binary longitudinal data. J. Am. Statist. Ass., 101, 685-700.
- Craig, P. (2007) Orthants: multivariate normal orthant probabilities. University of Durham, Durham. (Available from http://www.maths.dur.ac.uk/~dma0psc/orthants/.)
- Filon, L. (1928) On a quadrature formula for trigonometric integrals. Proc. R. Soc. Edinb., 49, 38-47.
- Gassmann, H. I., Deák, I. and Szántai, T. (2002) Computing multivariate normal probabilities: a new look. J. Computnl Graph. Statist., 11, 920–949.
- Genz, A. (1992) Numerical computation of multivariate normal probabilities. J. Computal Graph. Statist., 1, 141-150.
- Genz, A. (1993) Comparison of methods for the computation of multivariate normal probabilities. Comput. Sci. Statist., 25, 400-405.
- Genz, A. and Bretz, F. (2002) Comparison of methods for the computation of multivariate t probabilities. J. Computnl Graph. Statist., 11, 950–971.
- Genz, A., Bretz, F. and Torsten Hothorn, R. (2005) mytnorm: multivariate normal and t distribution. Friedrich-Alexander-Universität, Erlangen-Nürnberg. (Available from http://cran.r-project.org/src/contrib/ Descriptions/mvtnorm.html.)
- Gueorguieva, R. V. and Agresti, A. (2001) A correlated probit model for joint modeling of clustered binary and continuous responses. J. Am. Statist. Ass., 96, 1102-1112.
- Hayter, A. J. (2006) Recursive integration methodologies with statistical applications. J. Statist. Planng Inf., 136, 2284-2296.
- Hirotsu, C., Kuriki, S. and Hayter, A. J. (1992) Multiple comparison procedures based on the maximal component of the cumulative chi-squared statistic. Biometrika, 79, 381-392.
- Miwa, T., Hayter, A. J. and Kuriki, S. (2003) The evaluation of general non-centred orthant probabilities. J. R. Statist. Soc. B, 65, 223-234.
- Papathomas, M. and O'Hagan, A. (2005) Updating beliefs for binary variables. J. Statist. Planng Inf., 135, 324-
- Press, W. H., Flannery, B. P., Teukolsky, S. A. and Vetterling, W. T. (1993) Numerical Recipes in C: the Art of Scientific Computing, 2nd edn. Cambridge: Cambridge University Press.
- R Development Core Team (2006) R: a Language and Environment for Statistical Computing. Vienna: R Foundation for Statistical Computing.
- Schläfli, L. (1858) On the multiple integral  $\int_{-\infty}^{n} dx dy \dots dz$ , whose limits are  $p_1 = a_1x + b_1y + \dots + b_1z > 0$ ,  $p_2 > 1$  $0, \dots, p_n > 0$  and  $x^2 + y^2 + \dots + z^2 < 1$ . Q. J. Pure Appl. Math., 2, 269-301.

  Williams, D. (1971) A test for differences between treatment means when several dose levels are compared with
- a zero dose control. Biometrics, 27, 103-117.