# Markov Chain Monte Carlo With Mixtures of Mutually Singular Distributions

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Markov chain Monte Carlo (MCMC) methods for Bayesian computation are mostly used when the dominating measure is the Lebesgue measure, the counting measure, or a product of these. Many Bayesian problems give rise to distributions that are not dominated by the Lebesgue measure or the counting measure alone. In this article we introduce a simple framework for using MCMC algorithms in Bayesian computation with mixtures of mutually singular distributions. The idea is to find a common dominating measure that allows the use of traditional Metropolis-Hastings algorithms. In particular, using our formulation, the Gibbs sampler can be used whenever the full conditionals are available. We compare our formulation with the reversible jump approach and show that the two are closely related. We give results for three examples, involving testing a normal mean, variable selection in regression, and hypothesis testing for differential gene expression under multiple conditions. This allows us to compare the three methods considered: Metropolis-Hastings with mutually singular distributions, Gibbs sampler with mutually singular distributions, and reversible jump. In our examples, we found the Gibbs sampler to be more precise and to need considerably less computer time than the other methods. In addition, the full conditionals used in the Gibbs sampler can be used to further improve the estimates of the model posterior probabilities via Rao-Blackwellization, at no extra cost.

**Key Words:** Gibbs sampler; Metropolis-Hastings algorithm; Mixture distribution; Rao-Blackwellization; Reversible jump; Singular measures.

### **1. INTRODUCTION**

Mixtures of mutually singular distributions arise quite often in statistics. For example, one could model a process that truly is a mixture of a discrete process and a continuous process. One could also be interested in model selection where the dimension of the parameter space varies, giving rise to singularities in the prior distribution. It seems that mixtures

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of mutually singular distributions are often avoided because of the measure-theoretic difficulties. Perhaps one reason is that the derivation of the density (i.e., the Radon–Nikodym derivative) is not as intuitive as when the distribution is purely discrete or continuous. However, the difficulty is not great and the goal of this article is to introduce an easily used framework that would facilitate the use of such mixtures. We are particularly interested in Markov chain Monte Carlo methods where the target distribution is of this form. We focus on Bayesian inference, where the target distribution is a posterior distribution, though the method is more general.

The Metropolis–Hastings algorithm (Metropolis et al. 1953; Hastings 1970) is a method for constructing a reversible Markov chain with a specified invariant distribution. The Metropolis–Hastings algorithm has been widely used in Bayesian inference to approximate posterior quantities of interest (Geman and Geman 1984; Gelfand and Smith 1990). In most applications, the Metropolis–Hastings algorithm is used when the dominating measure is the Lebesgue measure, the counting measure, or a product of these. However, the algorithm works for any target distribution with a  $\sigma$ -finite dominating measure (Tierney 1994, 1998)

Even though the theory of MCMC can be used with general dominating measures such as sums of mutually singular distributions, this is rarely done in practice. Simple mixtures of mutually singular distributions such as a point mass and a continuous distribution (with respect to the Lebesgue measure) have been used in Bayesian variable selection (Smith and Kohn 1996; Geweke 1996; George and McCulloch 1997). This approach was not considered in an earlier paper by George and McCulloch (1993). Such mixtures also arise with Dirichlet processes in Bayesian density estimation (Escobar and West 1995; Neal 2000) and identification of regeneration times in MCMC simulation (Brockwell and Kadane 2005). There is a need for a general formulation for MCMC computations with mixtures of mutually singular distributions. Note that the term "mixture of mutually singular distributions" was not used in these previous formulations, perhaps because they were not fully general, considering only point masses.

Here we are interested in more complicated situations where we have several distributions of different dimensions. These include Bayesian model selection where some of the parameters are allowed to lie in a hyperplane, or more generally in a submanifold of  $\mathbb{R}^n$ . There has been a great deal of work on MCMC algorithms for Bayesian model selection. Madigan and York (1995) and Raftery et al. (1997) integrated over the parameter space analytically and made the MCMC move only in the model space; the resulting method is called MCMC model composition, or MC<sup>3</sup>. This avoids the issue of mutually singular distributions, but is not applicable to all such problems. Carlin and Chib (1995) used a product space approach to keep the dimension of the parameter space fixed. Following the pioneering work of Grenander and Miller (1994) and Phillips and Smith (1995) based on jump diffusions, Green (1995) showed how to construct a reversible jump MCMC algorithm to handle cases where the dimension of the parameter space is allowed to vary. Petris and Tardella (2003) introduced a geometric approach to transdimensional MCMC and showed that it can be used to formulate the problem as a mixture of distributions with components supported by subspaces of different dimensions. In this article, we show that reversible jump can be viewed in terms of a mixture of mutually singular distributions. Our formulation is more general than Petris and Tardella (2003) as it can deal with nonnested models, and it is computationally easier as there is no need to transform the parameters. In addition, using three examples, we show how the Gibbs sampler can be used within our framework.

The article is organized as follows. In Section 2, we introduce some notation and show how one can derive densities for mixtures of mutually singular distributions. In Section 3, we briefly review the Metropolis–Hastings algorithm and show how it can be used to form an ergodic chain with a mixture of mutually singular distributions as the invariant distribution. In Section 4, we use three examples to demonstrate the methodology introduced and compare various Metropolis–Hastings algorithms including the Gibbs sampler. In Section 5, we compare our formulation with the reversible jump approach and show that the two are closely related. Finally, in Section 6 we discuss possible extensions and the limitations of our formulation.

# 2. DENSITIES FOR MIXTURES OF MUTUALLY SINGULAR DISTRIBUTIONS

In this section, we show how one can derive densities for a mixture of singular distributions. Without loss of generality, we let our sample space be the *n*-dimensional Euclidean space  $\mathbb{R}^n$ , or a subset of it. We denote the *n*-dimensional Lebesgue measure by  $\lambda_n$  and the Dirac measure concentrated at **x** by  $\delta_{\mathbf{x}}$ . We will say that a probability measure  $\Pi$  is dominated by a  $\sigma$ -finite measure  $\nu$ , if  $\Pi$  admits a density with respect to  $\nu$ . In other words, if we can write  $\Pi(\mathbf{dx}) = d\Pi/d\nu(\mathbf{x})\nu(\mathbf{dx})$ , where  $d\Pi/d\nu$  denotes the density (also known as the Radon–Nikodym derivative) of  $\Pi$  with respect to  $\nu$ . In cases where the dominating measure to be used is clear from the context, we will just use lower-case letters (e.g.,  $\pi$ , q) to denote the densities of the corresponding probability measures denoted in capital letters (e.g.  $\Pi$ , Q).

Before introducing the main result, we first recall that two measures  $v_1$  and  $v_2$  are said to be mutually singular, denoted by  $v_1 \perp v_2$ , if there exists a (measurable) set A such that  $v_1(A) = 0$  and  $v_2(A^c) = 0$  where  $A^c$  denotes the complement of A. Basically, the two measures are supported on disjoint subsets. Similarly, a countable collection of measures  $v_i, i \in I$  are said to be mutually singular if  $v_i \perp v_j$  for each  $i, j \in I, i \neq j$ . Finally, we will say that two distributions are mutually singular if the corresponding probability measures are mutually singular.

The following theorem, based on the Radon-Nikodym theorem (Billingsley 1995, p. 422), gives a way of explicitly writing down densities for mixtures of mutually singular distributions.

**Theorem 1.** Let  $\Pi$ ,  $\Pi_i$ ,  $i \in I$ , be probability measures such that  $\Pi = \sum_{i \in I} w_i \Pi_i$ ,  $w_i \in [0, 1], \sum_{i \in I} w_i = 1$  and  $\Pi_i$  is dominated by  $v_i$ , where the  $v_i$ 's are mutually singular  $\sigma$ -finite measures. Then the density (Radon-Nikodym derivative) of  $\Pi$  with respect to  $v \equiv$ 

 $\sum_{i \in I} v_i$  is given by

$$\frac{d\Pi}{d\nu}(\mathbf{x}) = \sum_{i \in I} w_i \frac{d\Pi_i}{d\nu_i}(\mathbf{x}) \mathbf{1}_{S_i}(\mathbf{x}),$$

where the  $S_i$ 's are sets such that  $v_i(S_i^c) = 0$  and  $v_j(S_i) = 0$ ,  $i \neq j$ .

The proof is given in Appendix A.1. If we let  $\pi$  and  $\pi_i$  be the densities of  $\Pi$  and  $\Pi_i$  with respect to  $\nu$  and  $\nu_i$ , respectively, then Theorem 1 can be written as

$$\pi(\mathbf{x}) = \sum_{i \in I} w_i \pi_i(\mathbf{x}) \mathbf{1}_{S_i}(\mathbf{x}).$$

Thus, the theorem says that the density  $\pi$  with respect to the global dominating measure  $\nu$ , can be expressed as a mixture of the "componentwise" densities  $\pi_i$ . In general, the componentwise densities  $\pi_i$  are known or are easy to derive, and therefore  $\pi$  can easily be derived. For example, if  $\nu_i$  is the Lebesgue measure,  $\pi_i$  is simply the usual Lebesgue density. Note, however, that the indicators  $\mathbf{1}_{S_i}$  are crucial to get a proper density. They are there to make sure that a set does not contribute to more than one component, which would not be legitimate since the  $\nu_i$ 's (and thus the  $\Pi_i$ 's) are mutually singular. If a set is counted more than once, the density might not even integrate to one! A direct consequence of this is that if  $\mathbf{x} \in S_i$ , then  $\pi(\mathbf{x}) = w_i \pi_i(\mathbf{x})$ , and even if I is large we only have to evaluate a single  $\pi_i$ . The sets  $S_i$  given in Theorem 1 need to be derived on a case-by-case basis because they depend on the supports of the measures  $\nu_i$ . However, as we will see in the example below and the ones explored in Section 4, their derivation will be clear from the context. It is also possible to get an explicit (and unique) derivation using the notion of Hausdorff measures (Gottardo and Raftery 2004).

*Example 1: Mixture of a point mass and a continuous random variable.* Let  $X_1$  be a discrete random variable equal to zero with probability one; thus its induced probability measure is the Dirac mass at zero  $\Pi_1 \equiv \delta_0$ . Let  $X_2$  be a continuous random variable with probability measure  $\Pi_2$  dominated by the Lebesgue measure  $\lambda$ ; thus  $X_2$  admits a density with respect to the Lebesgue measure. Define a third random variable Y equal to  $X_1$  with probability (1 - w) and to  $X_2$  with probability w.

The probability measure of *Y* can be written as  $\Pi = (1 - w)\Pi_1 + w\Pi_2$ . Applying Theorem 1 we obtain the density of  $\Pi$  with respect to  $(\delta_0 + \lambda)$ , namely

$$\frac{d\Pi}{d(\delta_0+\lambda)}(x) = (1-w)\frac{d\Pi_1}{d\delta_0}(x)\mathbf{1}_{S_1}(x) + w\frac{d\Pi_2}{d\lambda}(x)\mathbf{1}_{S_2}(x) \text{ a.e. } \delta_0 + \lambda,$$

where "a.e.  $\delta_0 + \lambda$ " means almost everywhere with respect to the measure  $\delta_0 + \lambda$ , that is, everywhere except on a set that has  $(\delta_0 + \lambda)$ -measure equal to zero. If *Y* is equal to zero, we know that it comes from the first component, whereas if it is not, it must come from the continuous distribution. Thus, natural choices for  $S_1$  and  $S_2$  are {0} and  $\mathbb{R} \setminus \{0\}$ , respectively. It follows that the density of *Y* can be written as

$$\pi(x) \equiv \frac{d\Pi}{d(\delta_0 + \lambda)}(x) = (1 - w)\mathbf{1}_{\{0\}}(x) + wf(x)\mathbf{1}_{\mathbb{R}\setminus\{0\}}(x) \text{ a.e. } \delta_0 + \lambda,$$

where f(x) is the Lebesgue density of  $X_2$ . Note that it is crucial to remove {0} from  $\mathbb{R}$  in the second indicator function for the density to be a valid density with respect to  $(\delta_0 + \lambda)$ . If we did not do that, the density evaluated at 0, which corresponds to the probability that *Y* is zero, would be (1 - w) + wf(0) and not 1 - w.

# 3. MARKOV CHAIN MONTE CARLO WITH MIXTURES OF MUTUALLY SINGULAR DISTRIBUTIONS

We now consider Markov chain Monte Carlo algorithms for mixtures of mutually singular distributions. We show how it is possible to use the regular Metropolis–Hastings and Gibbs sampler algorithms, where Theorem 1 is used to derive the required densities.

Suppose that we wish to sample from a distribution for a variable  $\mathbf{x}$  with associated probability measure  $\Pi(\mathbf{dx})$  dominated by a  $\sigma$ -finite measure  $\nu$ . We use  $\pi$  to denote its density, so that  $\Pi(\mathbf{dx}) = \pi(\mathbf{x})\nu(\mathbf{dx})$ . In this article, we shall be concerned with a posterior distribution with density  $\pi(\mathbf{x}|\mathbf{y})$ , or more generally a full conditional  $\pi(\mathbf{x}|\mathbf{y}, \mathbf{z})$  when blocks of variables are updated in turns, where  $\mathbf{y}$  are the data and  $\mathbf{z}$  are the variables that are being conditioned upon.

From now on, we use  $\pi$  (**x**) to denote the target distribution, in this case  $\pi$  (**x**|**y**, **z**). One of the most widely used methods for generating such chains is the Metropolis–Hastings algorithm. MCMC methods have been widely used to generate (dependent) samples from distributions where the normalizing constant is unknown or intractable. Most applications deal with distributions where the dominating measure is either the Lebesgue measure, the counting measure, or a product of those. Not much attention has been given to measures that are not absolutely continuous with respect to the Lebesgue measure or to the counting measure.

To avoid special cases, we assume that  $\pi(\mathbf{x}) > 0$  for each  $\mathbf{x}$  in the sample space. In Markov chain Monte Carlo, one constructs a Markov chain with invariant distribution  $\pi$ . Following the notation of Tierney (1994), let Q be a Markov transition kernel of the form

$$Q(\mathbf{x}, \mathbf{dx}') = q(\mathbf{x}, \mathbf{x}')\nu(\mathbf{dx}'),$$

and define

$$\alpha(\mathbf{x}, \mathbf{x}') = \begin{cases} \min\left\{\frac{\pi(\mathbf{x}')q(\mathbf{x}', \mathbf{x})}{\pi(\mathbf{x})q(\mathbf{x}, \mathbf{x}')}, 1\right\}, & \text{if } \pi(\mathbf{x})q(\mathbf{x}, \mathbf{x}') > 0\\ 1, & \text{if } \pi(\mathbf{x})q(\mathbf{x}, \mathbf{x}') = 0. \end{cases}$$
(3.1)

Note that, in this definition, both q and  $\pi$  are densities with respect to  $\nu$ , which in our case will be the sum of several mutually singular measures. If the chain is currently at a point  $\mathbf{X}_n = \mathbf{x}$ , then a new candidate  $\mathbf{x}'$  is generated according to the distribution  $Q(\mathbf{x}, \cdot)$  and the new point is accepted with probability  $\alpha(\mathbf{x}, \mathbf{x}')$ . If the associated Metropolis–Hastings kernel, denoted by K, is  $\Pi$ -irreducible, Harris-recurrent and aperiodic, the Markov chain formed will converge to the unique stationary distribution  $\pi$  with respective probability measure  $\Pi$  (Tierney 1994). Note that the formulation given here includes the Gibbs sampler (Geman and Geman 1984; Gelfand and Smith 1990), where new observations are

drawn from the full conditional and the acceptance probability is equal to one. Most applications concern cases where the dominating measure is the Lebesgue measure, the counting measure, or a product of those. However, the general results in Tierney (1994) apply to more general distributions. This was emphasized in a more recent paper (Tierney 1998), where the author described very general conditions under which the Metropolis–Hastings algorithm is reversible.

Thus, the result applies to mixtures of mutually singular distributions. The main difficulty is that the choice of the proposal becomes limited, as mixture components put all their mass on different parts of the parameter space. For example, the usual random walk Metropolis algorithm will not be available in general as symmetric proposals do not exist. This last point will become clearer in Example 2.

Using the notation introduced in Theorem 1, we assume that the target distribution is of the form  $\Pi = \sum_{i \in I} w_i \Pi_i$ , where  $\Pi_i$  is dominated by  $v_i$ , and the  $v_i$ 's are mutually singular  $\sigma$ -finite measures. Thus  $\Pi$  is dominated by  $v = \sum_i v_i$ . In a Bayesian context, such singularities would occur when the prior itself is a mixture of mutually singular distributions, whose density with respect to v can be written as  $\pi(\mathbf{x}|\mathbf{z}) = \sum_i w_i \pi_i(\mathbf{x}|\mathbf{z}) \mathbf{1}_{S_i}$  using Theorem 1. Thus, using Bayes theorem it follows that

$$\pi(\mathbf{x}|\mathbf{y}, \mathbf{z}) \propto L(\mathbf{y}|\mathbf{x}, \mathbf{z})\pi(\mathbf{x}|\mathbf{z})$$

$$= L(\mathbf{y}|\mathbf{x}, \mathbf{z}) \sum_{i} w_{i}\pi_{i}(\mathbf{x}|\mathbf{z})\mathbf{1}_{S_{i}}(\mathbf{x})$$

$$= \sum_{i} w_{i}L(\mathbf{y}|\mathbf{x}, \mathbf{z})\pi_{i}(\mathbf{x}|\mathbf{z})\mathbf{1}_{S_{i}}(\mathbf{x})$$

$$= \sum_{i} w_{i}c_{i}\pi_{i}(\mathbf{x}|\mathbf{z}, \mathbf{y})\mathbf{1}_{S_{i}}(\mathbf{x}), \qquad (3.2)$$

where  $\pi_i$  is the "componentwise" full conditional of **x** given (**y**, **z**) and that **x** belongs to component *i* (i.e.,  $\mathbf{x} \in S_i$ ), and  $c_i$  is its normalizing constant, namely  $c_i = \int_{\mathbf{x}} L(\mathbf{y}|\mathbf{x}, \mathbf{z}) \pi_i(\mathbf{x}|\mathbf{z})v_i(\mathbf{dx})$ . Because each  $\pi_i \mathbf{1}_{S_i}$  integrates to one with respect to v it is easy to see that the global normalization constant is  $\sum_i w_i c_i$ , and it follows that

$$\pi(\mathbf{x}|\mathbf{y}, \mathbf{z}) = \sum_{i} w_{i}^{*} \pi_{i}(\mathbf{x}|\mathbf{z}, \mathbf{y}) \mathbf{1}_{S_{i}}, \qquad (3.3)$$

where  $w_i^* = c_i w_i / (\sum_i c_i w_i)$ . This implies that the posterior  $\Pi$  is a mixture of mutually singular distributions. It also implies that if each  $c_i$  is analytically available and  $\pi_i$  can be sampled from directly, one can sample from  $\Pi$  by first selecting a component *i* at random with probability  $w_i^*$  and then sampling from  $\Pi_i$ . So, if  $\Pi$  is a full conditional, and all componentwise full conditionals (the  $\Pi_i$ 's) are explicitly available, Gibbs sampling can be used. This will be illustrated with three examples in the next section.

If the full conditional is not available, the Metropolis–Hastings algorithm could be used instead. In order for the Metropolis–Hastings kernel K to be  $\Pi$ -irreducible, it is necessary to construct a  $\Pi$ -irreducible transition kernel Q. In practice, assuming that **x** belongs to component *i*, it will be convenient to write the kernel Q as

$$Q(\mathbf{x}, \mathbf{dx}') = \sum_{j \in I} p_{ij} Q_j(\mathbf{x}, \mathbf{dx}'), \qquad (3.4)$$

where the  $Q_j$ 's are mutually singular transition kernels in the sense that  $Q_j(\mathbf{x}, \mathbf{dx}') = q_j(\mathbf{x}, \mathbf{x}')v_j(\mathbf{dx}')$ , and  $p_{ij}$  is the probability of proposing a move to component j from component i. In (3.4),  $p_{ij}$  is a between-component move and  $Q_j$  is a move within component j. The  $\Pi$ -irreducibility of Q and the associated Metropolis–Hastings kernel will depend on both  $P \equiv (p_{ij})$  and the  $Q_j$ 's, and will need to be verified for each case. However, as with traditional MCMC, this is usually not hard to verify. In our framework, it will be enough to show that there is a positive probability that the chain will reach (i) any component of the mixture and (ii) any set of positive mass (with respect to  $\Pi$ ) within that component, from anywhere in the state space, in a finite number of steps. Note that the Gibbs sampler as given in (3.3) is a special case of this where  $p_{ij} \equiv w_i^*$  and  $Q_j \equiv \Pi_j$ .

To fully establish the convergence of the Metropolis–Hastings Markov chain with kernel K, one needs to show that the resulting kernel is also aperiodic and Harris recurrent. Aperiodicity is easily verified and is usually satisfied for the Metropolis–Hastings algorithm even when used as a Gibbs sampler (Tierney 1994; Roberts and Tweedie 1996). Standard MCMC results can be used to verify Harris recurrence (Tierney 1994; Chan and Geyer 1994; Roberts and Tweedie 1996; Roberts and Rosenthal 2004).

### 4. APPLICATIONS

In this section, we consider three examples that are applications of the Metropolis– Hastings algorithm for Bayesian computation with mixtures of mutually singular distributions. From now on, we denote by  $N(\mu, \sigma^2)$  the Gaussian distribution with mean  $\mu$  and variance  $\sigma^2$ . The corresponding density evaluated at x is denoted by  $N(x|\mu, \sigma^2)$ .

Example 2: Testing a normal mean. Consider the simple Bayesian linear model,

$$y_j = \mu + \epsilon_j, \tag{4.1}$$
$$(\epsilon_j | \psi) \sim \mathcal{N}(0, \psi^{-1}).$$

We might be interested in testing if the mean  $\mu$  is equal to zero. In order to do so, we need to specify a prior distribution that allows the parameter  $\mu$  to be equal to zero. We use the prior

$$\mu \sim (1 - w)\delta_0 + w \mathcal{N}(0, \psi_{\mu}^{-1}), \tag{4.2}$$

which is a mixture of a point mass at 0 and a Gaussian distribution. Using Theorem 1 and Example 1, its density with respect to  $(\delta_0 + \lambda)$  is  $(1 - w)\mathbf{1}_{\{0\}}(\mu) + wN(\mu|0, \psi_{\mu}^{-1})\mathbf{1}_{\mathbb{R}\setminus\{0\}}(\mu)$ , where  $\delta_0$  is the Dirac mass at zero and  $\lambda$  is the one-dimensional Lebesgue measure.

We let the prior for the precision  $\psi$  be Gamma( $\xi_1, \xi_2$ ), which has mean  $\xi_1/\xi_2$  and variance  $\xi_1/\xi_2^2$ . The hyperparameters  $\psi_{\mu}, \xi_1$  and  $\xi_2$  will be kept fixed. Here the target distribution can be expressed (up to a normalizing constant) in terms of its density  $\pi(\mu, \psi | \mathbf{y})$ 

with respect to the product measure  $(\delta_0 + \lambda) \times \lambda$  as

$$\pi(\mu, \psi | \mathbf{y}) \propto L(\mu, \psi | \mathbf{y}) \pi(\mu) \pi(\psi)$$

$$\propto \psi^{J/2} \exp(-\frac{\psi}{2} \sum_{j=1}^{J} (y_j - \mu)^2)$$

$$\times \left[ (1 - w) \mathbf{1}_{\{0\}}(\mu) + w \left(\frac{\psi_{\mu}}{2\pi}\right)^{1/2} \exp(-\frac{\psi_{\mu}}{2}\mu^2) \mathbf{1}_{\mathbb{R} \setminus \{0\}}(\mu) \right]$$

$$\times \psi^{\xi_1 - 1} \exp(-\xi_2 \psi).$$

Note that here we give the distribution of  $\pi(\mu, \psi | \mathbf{y})$  explicitly (up to a normalizing constant), for illustration. However, in practice one does not need to do this, and one would simply update  $\mu$  and  $\psi$  alternatively. The update for  $\psi$  presents nothing complicated, as the dominating measure is the Lebesgue measure, and its full conditional distribution is a Gamma distribution whether or not  $\mu$  is zero. Therefore we shall be concerned only with the update of  $\mu$  conditional on  $\psi$  and  $\mathbf{y}$ .

The prior density for  $\mu$  is of the form  $w_1 \mathbf{1}_{\{0\}}(\mu) + w_2 N(\mu|0, \psi_{\mu}^{-1}) \mathbf{1}_{\mathbb{R}\setminus\{0\}}(\mu)$ . Thus, we can use (3.3) with  $w_1 = (1 - w)$ ,  $w_2 = w$ ,  $\pi_1(\boldsymbol{\mu}|\boldsymbol{\psi}) = 1$  and  $\pi_2(\boldsymbol{\mu}|\boldsymbol{\psi}) = N(\mu|0, \psi_{\mu}^{-1})$ . Now we only need to derive both componentwise full conditionals  $\pi_1(\mu|\mathbf{y}, \psi)$  and  $\pi_2(\mu|\mathbf{y}, \psi)$ . In the case of the first component, we know that  $\mu = 0$ , and it follows that

$$\pi_1(\mu | \mathbf{y}, \psi) \propto L(\mathbf{y} | 0, \psi) \cdot 1,$$

and the normalizing constant is  $c_1 = L(\mathbf{y}|0, \psi)$ . Similarly, the componentwise full conditional,  $\pi_2$ , is given by

$$\pi_2(\mu|\mathbf{y},\psi) \propto L(\mathbf{y}|\mu,\psi) \sqrt{\frac{\psi_\mu}{2\pi}} \exp\left(-\frac{\psi_\mu}{2}\mu^2\right),$$

which is just the posterior distribution of a normal mean when the prior is Gaussian. Thus, the normalizing constant,  $c_2 = \int L(\mathbf{y}|\mu, \psi) \sqrt{\psi_{\mu}/(2\pi)} \exp(-0.5\psi_{\mu}\mu^2) d\mu$ , is available analytically. After straightforward calculation of  $c_2$ , the full conditional can be obtained using (3.3) and is given by

$$(\mu|\psi,\mathbf{y}) \sim (1-w^*)\delta_0 + w^* N\left(\psi \sum_i y_i/(n\psi+\psi_\mu), (n\psi+\psi_\mu)^{-1}\right),$$
 (4.3)

where

$$w^* = 1 - \frac{1 - w}{1 - w + w\sqrt{\psi_{\mu}/(n\psi + \psi_{\mu})} \exp(0.5(\psi \sum y_i)^2/(n\psi + \psi_{\mu}))}.$$

Using this full conditional, a new value for  $\mu$  can be generated by first selecting a component at random with respective probabilities  $1 - w^*$  and  $w^*$ , and then simulating a new  $\mu$ from the componentwise full conditional for the selected component. If the first component is selected we simply set  $\mu$  to zero, and if the second is selected we generate a new value from the corresponding Gaussian distribution. In some other settings (e.g., non-Gaussian likelihood and/or prior), the componentwise full conditional  $\pi_2$  might not be explicitly available and the Metropolis–Hastings algorithm provides a good alternative. In order to use Metropolis–Hastings, we need to define an irreducible Markov transition kernel (proposal),  $Q(\mu, d\mu') = q(\mu, \mu')\nu(d\mu')$ , that is absolutely continuous with respect to  $\nu = \delta_0 + \lambda$ , where  $\delta_0$  is the Dirac measure concentrated at 0, and  $\lambda$  is the Lebesgue measure. In other words, we have to make sure that we propose a move to zero as well as to the remainder of the real line. This can be done using (3.4), which we specify with the density

$$q(\mu, \mu') = (1 - p_i)\mathbf{1}_{\{0\}}(\mu') + p_i q^*(\mu, \mu')\mathbf{1}_{\mathbb{R}\setminus\{0\}}(\mu'),$$
(4.4)

where  $p_i$  (i = 0 if  $\mu = 0$  and i = 1 if  $\mu \neq 0$ ) is the probability of proposing a move to the continuous component from component *i*. In other words, with probability  $1 - p_i$  we propose a move to zero, and with probability  $p_i$  we propose a value according to  $q^*$ , the density of a kernel absolutely continuous with respect to the Lebesgue measure. For the Metropolis-Hastings algorithm, the acceptance probability is given by

$$\alpha(\mu, \mu') = \min\left\{\frac{\pi(\mu')q(\mu', \mu)}{\pi(\mu)q(\mu, \mu')}, 1\right\}.$$

In this case it is clear that we cannot find a symmetric proposal q, as this would require that the mass going to zero be the same as the mass leaving zero. The proposal given by (4.4) is clearly not symmetric as  $q(0, \mu') \neq q(\mu', 0)$  for  $\mu' \neq 0$ . Any other proposal will have the same problem because of the singularity between the two measures  $\delta_0$  and  $\lambda$ .

We now compare the performance of the Gibbs sampler to two other kernels based on the Metropolis–Hastings algorithm. The first kernel  $K_1$  is a generalization of the random walk Metropolis kernel with proposal density given by

$$q_1(\mu, \mu') = 0.5 \cdot \mathbf{1}_{\{0\}}(\mu') + 0.5 \mathrm{N}(\mu'|\mu, \sigma_1^2) \mathbf{1}_{\mathbb{R} \setminus \{0\}}(\mu'),$$

where  $\sigma_1^2$  is a fixed number. From  $\mu$ , this proposes 0 with probability 0.5, and proposes a random walk step with probability 0.5. The second kernel,  $K_2$ , is the concatenation of two Metropolis–Hastings kernels, with proposal densities

$$q_{21}(\mu, \tilde{\mu}) = \begin{cases} \mathrm{N}(\tilde{\mu}|\hat{\mu}, \sigma_2^2) \mathbf{1}_{\mathbb{R}\setminus\{0\}}(\tilde{\mu}) & \text{if } \mu = 0\\ \mathbf{1}_{\{0\}}(\tilde{\mu}) & \text{if } \mu > 0, \end{cases}$$

where  $\hat{\mu}$  is the sample mean and  $\sigma_2^2$  is a fixed number, and

$$q_{22}(\tilde{\mu}, \mu') = \begin{cases} \mathbf{1}_{\{0\}}(\mu') & \text{if } \tilde{\mu} = 0\\ \pi_2(\mu')\mathbf{1}_{\mathbb{R}\setminus\{0\}}(\mu') & \text{if } \tilde{\mu} > 0, \end{cases}$$

where  $\pi_2$  is the componentwise full conditional of  $\mu$  for the nonzero component. From  $\mu$ , this proposes 0 if  $\mu \neq 0$  and a nonzero  $\mu'$  if  $\mu = 0$ . In the latter case, this nonzero  $\mu'$  can be simulated in two steps. The first step based on  $q_{21}$  proposes  $\tilde{\mu}$  from a N( $\hat{\mu}, \sigma_2^2$ ) density where  $\hat{\mu}$  is the sample mean and  $\sigma_2^2$  is a fixed number. If this is accepted, the second step is

a componentwise Gibbs step based on the full conditional for the component with nonzero  $\mu$ . This last step, known as a within-model move in the reversible jump literature, is not necessary for the Markov chain to be ergodic but greatly improves the sampler. It will be convenient when comparing to the reversible jump formulation (Section 5).

We randomly generated 10 observations from a Gaussian distribution with variance 1 and mean 0.5, as follows:

$$0.575, 1.808, 0.532, -0.168, 0.529, 0.888, -1.368, -0.512, 2.667, 0.874$$

We fitted the model given by Equations (4.1) and (4.2) using w = 0.5, that is, each component is equally likely a priori. We fixed  $\psi_{\mu} = 0.01$ ,  $\xi_1 = 1$  and  $\xi_2 = 0.05$ , corresponding to fairly noninformative priors. Table 1 summarizes the results. The variance proposals  $\sigma_1^2 = 0.25$  and  $\sigma_2^2 = 1.2$  were chosen to maximize the proportion of moves between components. This was done by tuning the proposal variances using several pilot runs.

The true posterior probability of the model with mean zero is 0.867, conventionally viewed as positive but not strong evidence for the null model (Kass and Raftery 1995). In each case, the null model posterior probability (model with mean zero) can be easily estimated from the MCMC output as the proportion of  $\mu$ 's equal to zero. However, note that in (4.3),  $1 - w^*$  corresponds to the probability of  $\mu = 0$  given  $\psi$  and  $\mathbf{y}$ , and we can write

$$\Pr(\mu = 0|\mathbf{y}) = \mathbb{E}_{\mu}[\mathbf{1}_{[\mu=0]}] = \mathbb{E}_{\psi}\mathbb{E}_{\mu}[\mathbf{1}_{[\mu=0]}|\psi] = \mathbb{E}_{\psi}[\Pr(\mu = 0|\mathbf{y},\psi)] = \mathbb{E}_{\psi}[1-w^*],$$

where  $\mathbb{E}_T$  is the expectation with respect to the distribution of  $T | \mathbf{y}$ . Thus, in the case of the Gibbs sampler, one can obtain a more efficient estimate by Rao–Blackwellization when averaging the  $1 - w^*$  values computed at each iteration at no extra cost. As a consequence we get four different estimates from  $K_1$ ,  $K_2$ , Gibbs and Gibbs with Rao–Blackwellization.

All four methods considered did well in estimating the posterior probability, with accurate estimates based on 10,000 iterations, and low variability of the estimates. The algorithm based on the proposal  $K_2$  performed better than the one based on  $K_1$ . It moved between components almost twice as often, not surprisingly since  $K_2$  forces moves between components. The Gibbs sampler with the simple estimate did essentially as well as Metropolis–Hastings with proposal  $K_2$  in estimating the posterior probability, and it used much less computer time. Note that the proposal  $K_2$  combines two different move types, thus this is not a fair comparison with the Gibbs sampler and this is reflected in the CPU time. Generally speaking, the Metropolis–Hastings algorithm is computationally more expensive than the Gibbs sampler because it requires more evaluations of the target distribution. Finally, the Gibbs sampler with Rao–Blackwellization performed best both in terms of variability and computing time. Overall, therefore, the Gibbs sampler performed best among these three methods for this example.

*Example 3: Robust Bayesian variable selection in regression.* Variable selection is an important problem whose purpose is to select a group of variables that best predict an outcome variable. Given a dependent variable **Y** and a set of potential regressors  $\mathbf{X}_1, \ldots, \mathbf{X}_p$ , we wish to compare models of the form  $\mathbf{Y} = \beta_0 + \mathbf{X}_{i_1}\beta_{i_1} + \cdots + \mathbf{X}_{i_q}\beta_{i_q}$ , where  $\mathbf{X}_{i_1}, \ldots, \mathbf{X}_{i_q}$ 

Table 1. Comparison of the estimated model posterior probabilities computed with each algorithm. The estimates were computed from 10,000 iterations with 1,000 burn-in iterations. The standard deviations were computed by dividing a chain of 10 million iterations into 1,000 batches of 10,000 iterations each. The "truth" was obtained from 10 million iterations on the basis of which the estimates from the three algorithms agreed to within three digits. PMBC is the percentage of time the chain moved from one component to an other. In the case of the Gibbs sampler, RB refers to estimates computed via Rao–Blackwellization.

		$MH(K_1)$		MH	$(K_2)$		Gibbs				
	Truth	Est.	sd	Est.	sd	Est.	sd	RB-Est.	RB-sd		
$\pi(\mu=0 \mathbf{y})$	0.867	0.858	0.0054	0.870	0.0031	0.862	0.0034	0.867	0.0005		
PMBC(%)		13		2.	25		24				
CPU ( $\mu s$ /iter)		4.9		:	5.0		1.3				

is a selected subset of  $X_1, \ldots, X_p$ . For this problem we take an approach similar to that of George and McCulloch (1997). However, our model explicitly tries to take account of outliers, by using *t*-distributed errors. We assume a standard linear model to describe the relationship between the response and dependent variable and the set of predictors, namely

$$y_i = \beta_0 + \sum_{j=1}^p X_{ij}\beta_j + \frac{\epsilon_i}{\sqrt{\varpi_i}}$$
$$\epsilon_i \sim \mathcal{N}(0, \psi^{-1}),$$
$$(\overline{\omega_i}|\nu) \sim \mathcal{G}a(\nu/2, \nu/2),$$

where the  $\beta_i$ 's are the unknown regression coefficients and the  $\varpi_i$ 's are independent of the  $\epsilon_i$ 's. Since the  $\varpi_i$ 's are independent of the  $\epsilon_i$ 's, we have  $\epsilon_i / \sqrt{\varpi_i} \sim t_{(\nu,0,\psi^{-1})}$ , that is, the errors have a *t* distribution with  $\nu$  degrees of freedom and scale parameter  $\psi^{-1}$ . The advantage of writing the model in this way is that, conditioning on the  $\varpi_i$ , the sampling errors are again Gaussian, but with different precisions.

In order to allow each variable to be in and out of the model, we use a mixture of a Gaussian distribution and a point mass at zero for the prior of each regression coefficient, as follows:

$$\beta_k \sim (1-w)\delta_0 + wN\left(0, \frac{S_Y^2}{S_{X_k}^2}\sigma_\beta^2\right),$$

where w is the prior probability for each variable of being in the model,  $S_{X_k}^2$  is the empirical variance of the *k*th predictor,  $S_Y^2$  is the empirical variance of the response variables and  $\sigma_\beta^2$  is a common variance parameter. The prior for the variance parameter  $\sigma_\beta^2$  is taken to be uniform on the interval [0, 1]; arguments for priors of this kind are given by Raftery et al. (1997). The prior for the scaling parameter of the *t*-distribution,  $\psi$  is taken to be improper,  $\pi(\psi) \propto \psi^{-1}$ . We also tried a spread-out proper prior for  $\psi$  and the results were almost identical. The prior for the intercept  $\beta_0$  is taken to be Gaussian with a large variance centered at the least squares estimate for the full model with all variables,  $\hat{\beta}_0$ , namely N( $\hat{\beta}_0$ , 20 se( $\hat{\beta}_0$ )<sup>2</sup>). Finally, the prior for the degrees of freedom  $\nu$  is uniform on

the set {1, 2, 4, 8, 16, 32} as suggested by Besag and Higdon (1999). For each  $\beta_k$  we use w = 0.5 which makes every model equally likely a priori.

As in Example 2, the full conditional of  $\beta_k$  can be derived explicitly. All we need is to derive both componentwise full conditionals  $\pi_1(\beta_k|\cdots)$  and  $\pi_2(\beta_k|\cdots)$  for the first and second components, respectively, where  $\beta_k|\cdots$  means  $\beta_k$  conditioning on everything else in the model. In the case of the first component  $\beta_k = 0$ , and it follows that

$$\pi_1(\beta_k|\cdots) \propto L(\mathbf{y}|\beta_0, \beta_k = 0, \boldsymbol{\beta}_{-k}, \psi) \cdot 1$$

and the normalizing constant is  $c_1 = L(\mathbf{y}|\beta_k = 0, \boldsymbol{\beta}_{-k}, \psi)$ , where  $\boldsymbol{\beta}_{-k}$  is the vector of all regression coefficients except the *k*th one. Similarly, the componentwise full conditional,  $\pi_2$ , is given by

$$\pi_2(\beta_k|\cdots) \propto L(\mathbf{y}|\beta_0, \boldsymbol{\beta}, \psi) \sqrt{\frac{\psi_{\beta_k}}{2\pi}} \exp\left(-\frac{\psi_{\beta_k}}{2}\beta_k^2\right),$$

where  $\psi_{\beta_k}^{-1} = S_Y^2/S_{\chi_k}^2 \sigma_{\beta}^2$ , which is again a Gaussian distribution. As in Example 2, the normalizing constant,  $c_2 = \int L(\mathbf{y}|\beta_0, \boldsymbol{\beta}, \psi) \sqrt{\psi_{\beta_k}/(2\pi)} \exp\left(-0.5\psi_{\beta_k}\beta_k^2\right) d\beta_k$ , is again available analytically. After straightforward calculation of  $c_2$ , the full conditional can be obtained using (3.3) and is given by

$$(\beta_k|\dots) \sim (1-w_k^*)\delta_0 + w_k^* N\left(\psi \sum_i r_{ik}/(\psi \sum_i \varpi_i X_{ik}^2 + \psi_{\beta_k}), (\psi \sum_i \varpi_i X_{ik}^2 + \psi_{\beta_k})^{-1}\right),$$

where

$$w_{k}^{*} = 1 - \frac{1 - w}{1 - w + w \sqrt{\psi_{\beta}/(\psi \sum_{i} w_{i} X_{ik}^{2} + \psi_{\beta_{k}})} \exp(0.5(\psi \sum r_{ik})^{2}/(\psi \sum_{i} \varpi_{i} X_{ik}^{2} + \psi_{\beta_{k}}))},$$

with

$$r_{ik} = \varpi_i (y_i - \beta_0 - \sum_{j \neq k} \beta_j X_{ij}) X_{ik}.$$

As in Example 2, in some nonconjugate settings (e.g., non-Gaussian likelihood and/or prior), the componentwise full conditional  $\pi_2$  might not be explicitly available and the Metropolis–Hastings algorithm provides a good alternative. We now compare the Metropolis–Hastings algorithm with the kernels  $K_1$  and  $K_2$  as defined in the previous example. For a given  $\beta_k$ , the mean of the continuous component of the proposal  $q_{21}$  was set to the least squares estimate of the corresponding coefficient based on the full model. The width of each proposal was chosen to maximize the proportion of moves between components.

To illustrate the variable selection method, we use the Stack Loss data of Brownlee (1965), previously considered by many authors including Daniel and Wood (1980), Atkinson (1985) and, in a Bayesian framework, by Hoeting et al. (1996). It consists of 21 days of operation from a plant for the oxidation of ammonia as a stage in the production of nitric acid. The response is called "stack loss," defined as the percentage of unconverted ammonia that escapes from the plant. There are three independent variables,  $X_1, X_2, X_3$ .

Table 2. Comparison of the estimated posterior inclusion probabilities of each regression coefficient for Example 3. Estimates were computed from 10,000 iterations with 1,000 burn-in iterations. The standard deviations were computed by dividing a chain of 10 million iterations into 1,000 batches of 10,000 iterations each. The "truth" was obtained from ten billion iterations on the basis of which the estimates from the three algorithms agreed to within three digits. PMBC (percent moves between components) is the percentage of time the chain moved from one component to an other. The Gibbs sampler performed better than the Metropolis–Hastings samplers. RB refers to estimates computed via Rao–Blackwellization.

	Truth	$MH(K_1)$		M	$MH(K_2)$			Gibbs				
		Est.	sd	Est		sd	-	Est.	sd	RB-Est.	RB-sd	
$\Pi(\beta_1 \neq 0   \mathbf{y})$	1.000	1.000	0.015	1.00	0	0.005		1.000	0.005	1.000	0.004	
$\Pi(\beta_2 \neq 0   \mathbf{y})$	0.839	1.000	0.315	0.91	2	0.131		0.859	0.118	0.858	0.118	
$\Pi(\beta_3 \neq 0   \mathbf{y})$	0.141	0.244	0.230	0.18	2	0.093		0.097	0.085	0.097	0.085	
PMBC % ( $\beta_1$ )		0.0			0.0			0.0				
PMBC % ( $\beta_2$ )		0.4			0.5			1.0				
PMBC % ( $\beta_3$ )		1.1			1	1.9			2.5			
CPU ( $\mu s$ /iter)		67	1		68			45				

The airflow  $X_1$  measures the rate of operation of the plant. The nitric oxides produced are absorbed in a counter-current absorption tower:  $X_2$  is the inlet temperature of cooling water circulating through coils in this tower and  $X_3$  is proportional to the concentration of acid in the tower. Small values of the response correspond to efficient absorption of the nitric oxides. The general consensus with the Stack Loss data is that the predictor  $X_3$  (acid concentration) should be dropped from the model and that observations 1, 3, 4, and 21 are outliers.

As in the previous example, the posterior inclusion probability for each regression coefficient, defined as the posterior probability that the coefficient is not equal to zero, can be easily estimated from the MCMC output as the proportion of zeros. Again, in the case of the Gibbs sampler one can obtain a more efficient estimate by Rao–Blackwellization when averaging the  $1 - w^*$  values computed at each iteration.

The algorithms are compared in Table 2. This time, the Gibbs sampler outperformed the other two algorithms in terms of standard error for the estimated posterior inclusion probabilities. This is not surprising. Variables move in and out of the model over the course of the MCMC run, changing the relative estimates of the  $\beta$ 's and making it hard to construct an efficient proposal. The Gibbs sampler is automatic and depends on the current value of the other coefficients currently in the model. The average number of moves between components is higher for the Gibbs sampler than for the other two algorithms, and the estimated variances for the estimate of the posterior probabilities are smaller. Finally, the computing time is significantly reduced. Note that this time the improvement from the RB estimates is not great. This is due to the fact that, in this model, there are more parameters than in the previous one and the marginalization of  $\beta$  represents only a small portion of the overall number of parameters.

In this example, we are usually interested in the posterior probabilities of all  $2^3 = 8$  possible models. These can easily be computed from the posterior samples of the  $\beta_k$ 's, and

Table 3. Estimated posterior model probabilities for the *Stack Loss* Data. The total posterior probability for the other models visited was less than 0.01.

Models	1	1,2	1,3	1,2,3
Post. prob.	0.13	0.72	0.02	0.13

are given in Table 3. The posterior probabilities are consistent with the general consensus about the data.

The posterior mode of the number of degrees of freedom,  $\nu$ , is 4, suggesting that the observations are much heavier tailed than Gaussian. The posterior weights,  $\varpi$ , from the model are summarized in Table 4. It shows that observations 1, 3, 4, and 21 are downweighted by our model.

Finally, even though we have chosen to update each regression coefficient in turn, it would be possible to update the  $\beta_k$ 's in blocks of 2 or more. In particular, the full conditionals can be derived explicitly and the Gibbs sampler could be used. For example, the full conditional of any two regression coefficients,  $\boldsymbol{\beta}_{j,k} = (\beta_j, \beta_k)$ , is now a mixture of four mutually singular distributions of the form,

$$(\boldsymbol{\beta}_{j,k}|...) \sim w_1^* \delta_0 \delta_0 + w_2^* N\left(\mu_j^*, \psi_j^{*-1}\right) \delta_0 + w_3^* \delta_0 N\left(\mu_k^*, \psi_k^{*-1}\right) + w_4^* N_2\left(\boldsymbol{\mu}_{j,k}^*, \boldsymbol{\Psi}_{j,k}^{*-1}\right),$$

where  $w_1^*$ ,  $w_2^*$ ,  $w_3^*$ ,  $w_4^*$ ,  $\mu_j^*$ ,  $\mu_k^*$ ,  $\psi_j^*$ ,  $\psi_k^*$ , and  $\mu_{j,k}^*$ ,  $\Psi_{j,k}^*$  are constants depending on the data and the parameters that are being conditioned upon. The first component sets both regression coefficients to zero, the second (resp. third) sets  $\beta_k$  (resp.  $\beta_j$ ) to zero while generating a new  $\beta_j$  (resp.  $\beta_k$ ) from the componentwise full conditional of  $[\beta_j | \beta_j > 0, \beta_k = 0]$  (resp.  $[\beta_k | \beta_j = 0, \beta_k > 0]$ ), and the last component generates new regression coefficients from the componentwise full conditional of  $[(\beta_j, \beta_k)|\beta_j > 0, \beta_k > 0]$ . See Appendix A.2 for details. Similarly, the full conditional of more than two  $\beta$ 's can be obtained but the number of components grows exponentially with the numbers of  $\beta$ 's jointly updated.

*Example 4: Three-way comparison in gene expression data.* We now consider an application that arises in the analysis of gene expression microarray data. DNA microarrays allow the monitoring of thousands of genes simultaneously under different biological or experimental conditions. One of the main tasks with microarrays is the identification of genes that are expressed differentially under the different conditions. Hedenfalk et al. (2001) con-

Table 4. Estimated posterior weights, that is, posterior means of the  $\varpi$ 's, associated with each observation of the Stack Loss Data. Observations with small weights are downweighted. Observations 1, 3, 4, and 21 have smaller weights, suggesting that they might be outliers.

Obs.	1	2	3	4	5	6	7	8	9	10	11
Weight	0.78	1.08	0.74	0.59	1.13	1.05	1.09	1.16	1.08	1.16	1.15
Obs.	12	13	14	15	16	17	18	19	20	21	
Weight	1.12	0.96	1.01	1.14	1.21	1.19	1.20	1.15	1.04	0.49	

BRCA1	BRCA2	SPORADIC
-2.74	-1.51	1.47
-2.18	0.14	-0.81
-1.74	0.10	-1.69
-1.94	0.55	-1.06
0.29	-0.45	-1.32
-1.18	-0.67	-2.00
-1.40	-0.38	-1.18
	-0.60	

Table 5. Log transformed measurements of one gene of the BRCA dataset. Each row corresponds to a different condition, and each row entry to a different tumor sample.

ducted a study to examine breast cancer tissues from patients carrying mutations in the predisposing genes, BRCA1 or BRCA2, or from patients not expected to carry a hereditary mutation. They examined 22 breast cancer tumor samples: 7 tumors with BRCA1, 8 tumors with BRCA2 and 7 SPORADIC tumors, that is, tumors with neither mutation. The goal of the experiment was to study the expression patterns of 3,226 genes under the three conditions and to detect genes whose expression changed in at least one of the conditions. For illustrative purposes, we show the results for one gene. A complete analysis would require fitting a similar model to each gene; this was done by Gottardo et al. (2006) in more general settings of gene expression experiments. The measurements for this gene in the three samples are given in Table 5.

In order to detect differential expression we consider the following model:

$$y_{ci} = \mu_c + \epsilon_{ci},$$
  
$$(\epsilon_{ci} | \psi_{\epsilon_c}) \sim N(0, \psi_{\epsilon_c}^{-1}),$$

where  $\mu_c$  represents the mean expression level of the gene under condition  $c, i = 1, ..., n_c$ and c = 1, 2, 3 (BRCA1, BRA2, SPORADIC). We wish to test the null hypothesis  $\mu_1 = \mu_2 = \mu_3$ . In this example the alternative hypothesis is more complex due to the number of possible patterns. The prior distribution needs to include all such possible patterns. We therefore consider the following prior, whose density is

$$(\boldsymbol{\mu}|\boldsymbol{\psi}_{\boldsymbol{\mu}}, \mathbf{w}) \sim w_{1} \mathbf{N}(\mu_{1}|0, \boldsymbol{\psi}_{\mu_{123}}^{-1}) \mathbf{1}_{[\mu_{1}=\mu_{2}=\mu_{3}]} + w_{2} \mathbf{N}(\mu_{1}|0, \boldsymbol{\psi}_{\mu_{1}}^{-1}) \mathbf{N}(\mu_{2}|0, \boldsymbol{\psi}_{\mu_{23}}^{-1}) \mathbf{1}_{[\mu_{1}\neq\mu_{2}=\mu_{3}]} + w_{3} \mathbf{N}(\mu_{2}|0, \boldsymbol{\psi}_{\mu_{2}}^{-1}) \mathbf{N}(\mu_{1}|0, \boldsymbol{\psi}_{\mu_{13}}^{-1}) \mathbf{1}_{[\mu_{1}=\mu_{3}\neq\mu_{2}]} + w_{4} \mathbf{N}(\mu_{3}|0, \boldsymbol{\psi}_{\mu_{3}}^{-1}) \mathbf{N}(\mu_{1}|0, \boldsymbol{\psi}_{\mu_{12}}^{-1}) \mathbf{1}_{[\mu_{1}=\mu_{2}\neq\mu_{3}]} + w_{5} \mathbf{N}(\mu_{1}|0, \boldsymbol{\psi}_{\mu_{1}}^{-1}) \mathbf{N}(\mu_{2}|0, \boldsymbol{\psi}_{\mu_{2}}^{-1}) \mathbf{N}(\mu_{3}|0, \boldsymbol{\psi}_{\mu_{3}}^{-1}) \mathbf{1}_{[\mu_{1}\neq\mu_{2}\neq\mu_{3}]}, \quad (4.5)$$

where  $\boldsymbol{\psi}_{\boldsymbol{\mu}} = (\psi_{\mu_1}, \psi_{\mu_2}, \psi_{\mu_3}, \psi_{\mu_{12}}, \psi_{\mu_{13}}, \psi_{\mu_{23}}, \psi_{\mu_{123}})$  is the vector of precisions and **w** is the vector of probabilities for the five patterns constrained to sum to one. This defines a proper distribution with respect to the following  $\sigma$ -finite measure on  $\mathbb{R}^3$ , namely

$$\nu(\cdot) = \lambda_1(\Delta \cap \cdot) + \lambda_2(P_{\mu_1} \cap \cdot) + \lambda_2(P_{\mu_2} \cap \cdot) + \lambda_2(P_{\mu_3} \cap \cdot) + \lambda_3(\cdot),$$

where  $\Delta$  is the line  $\mu_1 = \mu_2 = \mu_3$ ,  $P_{\mu_1}$  is the plane  $\mu_2 = \mu_3$ ,  $P_{\mu_2}$  is the plane  $\mu_1 = \mu_3$ ,  $P_{\mu_3}$  is the plane  $\mu_1 = \mu_2$  and  $\lambda_k$  denotes the *k*-dimensional Lebesgue measure. In this case, we defined the distribution directly in terms of the density given by (4.5). The target distribution is  $\pi(\mu, \psi | \mathbf{y})$ , but as before the update for  $\psi$  does not present anything complicated and we are thus only concerned with  $\mu$ .

This time the prior density for  $\mu$  is a mixture of five singular Gaussian distributions, and given the Gaussian likelihood, the full conditional can once again be derived explicitly. Using (3.3), we need to derive each componentwise full conditional  $\pi_i$ , i = 1, ..., 5. Here we show only how to derive  $\pi_2$ ; other  $\pi_i$ 's can be obtained in a similar fashion. For the second component, given that  $\mu_1 \neq \mu_2 = \mu_3$ , we have

$$\pi_{2}(\boldsymbol{\mu}|\mathbf{y},\psi) \propto L(\mathbf{y}|\mu_{1},\mu_{2},\psi) \frac{\sqrt{\psi_{\mu_{1}}\psi_{\mu_{23}}}}{2\pi} \exp\left(-\frac{\psi_{\mu_{1}}}{2}\mu_{1}^{2} - \frac{\psi_{\mu_{23}}}{2}\mu_{2}^{2}\right)$$
$$\propto N(\mu_{1}|\mu_{1}^{*},\psi_{1}^{*-1})N(\mu_{2}|\mu_{23}^{*},\psi_{23}^{*-1}),$$

where the normalizing constant  $c_2$  is given by

$$c_{2} = \int_{\mu_{1}} \int_{\mu_{2}} L(\mathbf{y}|\mu_{1}, \mu_{2}, \psi) \sqrt{\psi_{\mu_{1}}\psi_{\mu_{23}}} / (2\pi) \exp(-0.5\psi_{\mu_{1}}\mu_{1}^{2} - 0.5\psi_{\mu_{23}}\mu_{2}^{2}) d\mu_{1} d\mu_{2},$$

and the parameters  $\mu_1^*$ ,  $\psi_1^*$  and  $\mu_{23}^*$ ,  $\psi_{23}^*$  are defined in Appendix A.3. Note that  $\mu_3$  does not appear in the likelihood as it is equal to  $\mu_2$ . Similarly, each componentwise full conditional can be obtained and the full conditional is again a mixture of five mutually singular Gaussian distributions with new means, precisions and proportions; see Appendix A.3.

Thus, one can simply generate a new  $\mu$  by first selecting a new component at random according to the new proportions, and then simulating  $\mu$  from the corresponding Gaussian distribution while imposing the respective constraint. For example, if the second component is selected, then new values for  $\mu_1$ ,  $\mu_2$  are generated from the Gaussian distribution above while  $\mu_3$  is set equal to  $\mu_2$ .

In nonconjugate settings the full conditional might not be available and the Metropolis– Hastings algorithm is an alternative. Here, it is harder to construct an efficient proposal for the Metropolis–Hastings algorithm than in the last example because of the greater number of mutually singular components. To try to maximize the between-component acceptance rate, we used a proposal based on local moves. The local move structure is described by the graph given in Figure 1. We use a proposal of the form given by Equations (3.4), with  $p_{ij} > 0$  if there is an edge between *i* and *j* (Figure 1). We consider only one kernel, which is similar to the kernel  $K_2$  of Examples 2 and 3, adapted to the local move structure.

From the current value of  $\mu$ , a new component is randomly chosen from among all the accessible components, based on Figure 1. Given the new component, new values are generated for  $\mu$  from a Gaussian proposal centered at the least squares estimates for the corresponding component. The width of the proposal was taken to be the same for all the Gaussian proposals, and its value was chosen by maximizing the proportion of moves between components. For example, given that the current  $\mu$  is from the first component, that is,  $\mu_1 = \mu_2 = \mu_3$ , we first select one of the three accessible components with probability 1/3 each. Then we generate new values for  $\mu$ , independently of the current values, from



Figure 1. Local move graph for the three-way comparison proposal.

a Gaussian proposal centered at the constrained least squares estimates. For example, if the second component is chosen, then the constraint is  $\mu_2 = \mu_3$ . We have chosen such an independence proposal for simplicity, but it would be possible to derive more efficient proposals. As with  $K_2$ , within a given component of the mixture each parameter is updated using a componentwise Gibbs step based on the full conditional for that component.

Similarly to the two previous examples, the model posterior probabilities can be easily estimated from the MCMC output as the proportion of time spent in the corresponding model/component. As before, in the case of the Gibbs sampler one can obtain a more efficient estimate by Rao–Blackwellization when averaging the  $w^*$  values at each iteration, as given in Appendix A.3. Table 6 summarizes the estimates of the posterior probabilities for each component using the Metropolis–Hastings algorithm and the Gibbs sampler. The estimates of the probabilities agree well but the standard deviations are much smaller for the Gibbs sampler. The average number of moves between components is much greater for the Gibbs sampler than for Metropolis–Hastings, indicating better mixing. Additionally, the Gibbs sampler required less computing time.

Note that we have chosen to update the components of  $\mu$  jointly, both for efficiency and to give an example of a block Gibbs update. However, it would be possible to update each parameter component in turn as in Example 3.

### 5. RELATIONSHIP WITH THE REVERSIBLE JUMP SAMPLER

Green (1995) introduced a Markov chain Monte Carlo method for Bayesian model determination for the situation where the dimensionality of the parameter vector is not fixed. Table 6. Comparison of the model posterior probabilities for the five models computed with each algorithm in Example 4. The estimates were computed from 10,000 iterations with 1,000 burn-in iterations. The standard deviations were computed by dividing a chain of 10 million iterations into 1,000 batches of 10,000 iterations each. The truth was estimated from 10 million iterations on the basis of which the estimates from the three algorithms agreed to within three digits. PMBC (percent of moves between components) is the percentage of time the chain moved from one component to an other. The estimates of the posterior probabilities agree well. The standard deviations are smaller for the Gibbs sampler. RB refers to estimates computed via Rao-Blackwellization. For ease of comparison and clarity, all standard deviations are multiplied by 10<sup>4</sup>.

	Truth	Mł	H			F	RJ			
		Est.	sd	Est.	sd	RB-Est.	RB-sd	Est.	sd	
$\Pi(\mu_1 = \mu_2 = \mu_3   \mathbf{y})$	0.085	0.089	90	0.087	30	0.087	20	0.089	90	
$\Pi(\mu_1 \neq \mu_2 = \mu_3   \mathbf{y})$	0.023	0.027	40	0.023	20	0.021	4	0.025	30	
$\Pi(\mu_2 \neq \mu_2 = \mu_3   \mathbf{y})$	0.850	0.826	130	0.847	40	0.851	20	0.843	120	
$\Pi(\mu_1 = \mu_2 \neq \mu_3   \mathbf{y})$	0.007	0.009	20	0.008	10	0.007	5	0.008	10	
$\Pi(\mu_1 \neq \mu_2 \neq \mu_3   \mathbf{y})$	0.035	0.050	100	0.038	20	0.035	0.3	0.034	40	
PMBC (%)	6				24				3	
CPU ( $\mu s$ /iter)	12				7				13	

Following the notation of Green (1995), we assume that we have a countable collection of models,  $\{\mathcal{M}_k : k \in \mathcal{K}\}$ . Model  $\mathcal{M}_k$  has a vector  $\boldsymbol{\theta}_k$  of unknown parameters assumed to lie in  $\mathbb{R}^{n_k}$ , where the dimension  $n_k$  may vary from model to model. Bayesian inference about k and  $\boldsymbol{\theta}_k$  is based on the joint posterior  $\pi(k, \boldsymbol{\theta}_k | \mathbf{y})$ , which can be decomposed as

$$\pi(k, \boldsymbol{\theta}_k | \mathbf{y}) \propto \pi(\mathbf{y} | \boldsymbol{\theta}_k, k) \pi(\boldsymbol{\theta}_k | k) p(k).$$

Using this formulation, the sample space can be represented by  $S = \bigcup_{k \in \mathcal{K}} \{k\} \times \mathbb{R}^{n_k}$ . Let  $\mathbf{x} = (k, \boldsymbol{\theta}_k)$ , and let  $\pi(\mathbf{x}) \equiv \pi(\mathbf{x}|\mathbf{y})$  denote the the target distribution. Even though the dimension of  $\mathbf{x}$  is allowed to change, Green (1995) showed that it is still possible to use the Metropolis–Hastings algorithm to form an irreducible and aperiodic Markov chain with stationary distribution  $\pi$ .

We now describe the reversible jump method in terms of random numbers as described in a more recent paper (Green 2003). At some current state  $\mathbf{x}$ , we generate r random numbers  $\mathbf{u}$  from a known joint density g, and then form the proposed new state as some suitable deterministic function of the current state and the random numbers:  $\mathbf{x}' = h(\mathbf{x}, \mathbf{u})$ . The reverse transformation will be made with the aid of random numbers  $\mathbf{u}'$  generated from some probability distribution g', giving  $\mathbf{x} = h'(\mathbf{x}', \mathbf{u}')$ . Assuming that the transformation from ( $\mathbf{x}, \mathbf{u}$ ) to ( $\mathbf{x}', \mathbf{u}'$ ) is a diffeomorphism, Green (1995) showed that a valid choice for the acceptance probability in the usual Metropolis–Hastings algorithm is given by

$$\alpha(\mathbf{x}, \mathbf{x}') = \min\left\{1, \frac{\pi(\mathbf{x}')g'(\mathbf{u}')}{\pi(\mathbf{x})g(\mathbf{u})} \left|\frac{\partial(\mathbf{x}', \mathbf{u}')}{\partial(\mathbf{x}, \mathbf{u})}\right|\right\},\tag{5.1}$$

where  $\pi(\mathbf{x}) \equiv \pi(k, \boldsymbol{\theta}_k | \mathbf{y}).$ 

The reversible jump formulation was introduced to handle cases where the dimension of the parameter vector  $\boldsymbol{\theta}$  can change from model to model. In Section 3, we showed how

reversible jump can also be viewed as a general Metropolis–Hastings algorithm where one fixes the number of parameters to be the same and allows some of the parameters to vanish or to lie in a hyperplane, reducing the dimension of the support. One example of this is provided by nested models with linear constraints on the parameters, and MCMC with mutually singular distributions could be applied to that case. This is also the case in Example 2, as we will show below. The Jacobian term present in (5.1) results from the change of variable induced by the diffeomorphism when a new value is proposed. If one designs a move that involves a change of variable, one should include the Jacobian term in (5.1) for the detailed balance condition to be satisfied (Green 1995).

Thinking about the problem in terms of mutually singular distributions allows us to use standard MCMC algorithms without worrying about the dimension matching. For example, we can use the usual Gibbs sampler when the full conditional is available. This is not possible with the reversible jump formulation, as pointed out by Green (1995) and Robert and Casella (1999, p. 287). We have shown that if one considers the right dominating measure, it is easy to establish that the Gibbs sampler is irreducible.

*Example 2: Testing a normal mean (continued).* In this example we have two competing models,  $\mathcal{M}_0 : \mu = 0$  and  $\mathcal{M}_1 : \mu \neq 0$ . For the first model  $\theta_0 = (\psi)$ , and for the second model we have one more parameter,  $\boldsymbol{\theta}_1 = (\mu, \psi)$ . In the notation of Green (1995), the sample space is  $S = \{0\} \times \mathbb{R} \cup \{1\} \times \mathbb{R}^2$ . As  $\psi$  is common to the two models, we shall be concerned only with the update of  $\boldsymbol{\theta}$  conditional on  $\psi$ . At the current state  $\mathbf{x} = (k, \boldsymbol{\theta}_k)$ , we can generate a new value  $\mathbf{x}'$  according to

$$\mathbf{x}' = \begin{cases} 0 & \text{if } k = 1\\ (1, u) & \text{if } k = 0 \end{cases}$$

where u is a random deviate with distribution g. In this case the acceptance probability reduces to

$$\alpha(\mathbf{x}, \mathbf{x}') = \begin{cases} 1 \wedge r(\mathbf{x}, \mathbf{x}') & \text{if } k = 1\\ 1 \wedge r(\mathbf{x}', \mathbf{x})^{-1} & \text{if } k = 0, \end{cases}$$
(5.2)

where

$$r(\mathbf{x}, \mathbf{x}') = \frac{f(\mathbf{y}|\psi)(1-w)g(\mu)}{f(\mathbf{y}|\mu, \psi)N(\mu|0, \psi_{\mu}^{-1})w}.$$
(5.3)

On the other hand, following our framework, we could use the prior mixture distribution given by (4.2), whose density is  $(1 - w)\mathbf{1}_{\{0\}} + wN(\mu|0, \psi_{\mu}^{-1})(1 - \mathbf{1}_{\{0\}})$ , and the proposal

$$q(\mu, \mu') = \begin{cases} 1_{\{0\}}(\mu') & \text{if } \mu \neq 0\\ g(\mu') & \text{if } \mu = 0, \end{cases}$$

which is absolutely continuous with respect to  $(\delta_0 + \lambda)$ . It is easy to see that the acceptance probability is the same as the one given by (5.2); the two formulations are equivalent. Thinking about the problem with a common dominating measure allows us to use the usual Metropolis–Hastings algorithm and the Gibbs sampler.

Usually, in reversible jump algorithms, there are two main types of moves: betweenmodel moves and within-model moves. Even though the within-model moves are not always necessary for the algorithm to be ergodic, they can greatly improve the performance of the sampler. The Metropolis–Hastings algorithms used in Examples 2 and 3 could be seen as reversible jump algorithms, where  $q_{21}$  is the between model proposal,  $q_{22}$  is the within model proposal, and g(u) is  $N(u; \hat{\mu}, \sigma^2)$  as in the definition of  $q_{21}$ . In Example 2, where it was easy to construct an efficient proposal, the reversible jump performed slightly better than the Gibbs sampler (with mutually singular distributions) in terms of mixing, but took substantially more computer time for the same number of iterations and was not as accurate as Gibbs with Rao–Blackwellization. Overall, for the same amount of computer time, the Gibbs sampler was more efficient. In Example 3, reversible jump performed relatively poorly compared to the Gibbs sampler, with much bigger standard deviations and much more computer time. We now turn back to Example 4 where it might be possible to design better moves using the reversible jump formulation.

*Example 4: Three-way comparison (continued).* In this example we have five competing models,

$$\mathcal{M}_{1} : \mu_{1} = \mu_{2} = \mu_{3}$$
$$\mathcal{M}_{2} : \mu_{1} \neq \mu_{2} = \mu_{3}$$
$$\mathcal{M}_{3} : \mu_{1} = \mu_{3} \neq \mu_{2}$$
$$\mathcal{M}_{4} : \mu_{1} = \mu_{2} \neq \mu_{3}$$
$$\mathcal{M}_{5} : \mu_{1} \neq \mu_{2} \neq \mu_{3},$$

which correspond to the five components given by the nodes of the local move graph (Figure 1). In a reversible jump framework, each model would be viewed as having a different number of parameters: the first model has one parameter, the second has two parameters, and so on. The parameter vectors for the five models can be written as

$$\theta_{1} = (\mu_{123})$$
  

$$\theta_{2} = (\mu_{1}, \mu_{23})$$
  

$$\theta_{3} = (\mu_{2}, \mu_{13})$$
  

$$\theta_{4} = (\mu_{3}, \mu_{12})$$
  

$$\theta_{5} = (\mu_{1}, \mu_{2}, \mu_{3})$$

and the sample space is given by  $S = \bigcup_i \{i\} \times \Theta_i$ , where  $\Theta_1 = \mathbb{R}$ ,  $\Theta_2 = \Theta_3 = \Theta_4 = \mathbb{R}^2$ and  $\Theta_5 = \mathbb{R}^3$ .

Similarly to Example 3, if g as used in (5.1) is the proposal used in the Metropolis– Hastings algorithm with mutually singular distribution, that is, Gaussian centered at the least squares estimates, one can show that the acceptance probability of the reversible jump is the same as in the Metropolis-Hastings algorithm, and so the two algorithms are equivalent in this case.

One of the strengths of reversible jump is the ability to design elaborate moves, which might lead to the inclusion of a Jacobian term in the acceptance ratio. In this example, we

use the common split-merge move (Richardson and Green 1997), where some parameters are merged to form a new one. Again, we use the local move graph to jump from one model to the other.

A move between model 1 and model 2 would be made as follows,

$$\mathbf{x}' = \begin{cases} (1, \frac{\mu_1 + \mu_{23}}{2}) & \text{if } k = 2\\ (2, \mu_{123} - u, \mu_{123} + u) & \text{if } k = 1, \end{cases}$$

where u is a random deviate with distribution g. In this case the acceptance probability reduces to

$$\alpha(\mathbf{x}, \mathbf{x}') = \begin{cases} 1 \wedge r(\mathbf{x}, \mathbf{x}') & \text{if } k = 2\\\\ 1 \wedge r(\mathbf{x}, \mathbf{x}')^{-1} & \text{if } k = 1, \end{cases}$$

where

$$r(\mathbf{x}, \mathbf{x}') = \frac{f(\mathbf{y}|\psi, \mu_{123})w_1 N(\mu_{123}|0, \psi_{\mu_{123}}^{-1})}{f(\mathbf{y}|\psi, \mu_1, \mu_{23})w_2 N(\mu_1; 0, \psi_{\mu_1}^{-1}) N(\mu_{23}|0, \psi_{\mu_{23}}^{-1})} \frac{p_{12}g((\mu_1 - \mu_{23})/2)}{p_{21}} \frac{1}{2},$$

and  $p_{ij}$  is the probability of proposing a move from model *i* to model *j*. Note the Jacobian term in the acceptance ratio due to the change of variable induced by the merge move. The acceptance ratio for a move between 1 and 3 (or 4) would be the same with obvious changes in notation.

Similarly, a move between 2 and 5 would be as follows:

$$\mathbf{x}' = \begin{cases} (2, \mu_1, \frac{\mu_2 + \mu_3}{2}) & \text{if } k = 5\\ (5, \mu_1, \mu_{23} - u, \mu_{23} + u) & \text{if } k = 2, \end{cases}$$

where u is a random deviate with distribution g. In this case the acceptance probability reduces to

$$\alpha(\mathbf{x}, \mathbf{x}') = \begin{cases} 1 \wedge r(\mathbf{x}, \mathbf{x}') & \text{if } k = 5\\ \\ 1 \wedge r(\mathbf{x}, \mathbf{x}')^{-1} & \text{if } k = 2, \end{cases}$$

where

$$r(\mathbf{x}, \mathbf{x}') = \frac{f(\mathbf{y}|\psi, \mu_1, \mu_{23})w_2 N(\mu_1|0, \psi_{\mu_1}^{-1}) N(\mu_{23}|0, \psi_{\mu_{23}}^{-1})}{f(\mathbf{y}|\psi, \mu_1, \mu_2, \mu_3)w_5 N(\mu_1|0, \psi_{\mu_1}^{-1}) N(\mu_2|0, \psi_{\mu_2}^{-1}) N(\mu_3|0, \psi_{\mu_3}^{-1})} \times \frac{p_{25}g((\mu_2 - \mu_3)/2)}{p_{52}} \frac{1}{2}.$$

The acceptance ratio for a move between 5 and 3 (or 4) would be the same with obvious changes in notation. Here, we choose g to be a Gaussian distribution with mean 0 and variance  $\sigma^2$ , with  $\sigma^2$  chosen to maximize the proportion of moves between components. As with traditional reversible jump algorithms, within a given model each parameter is updated using a Gibbs step (as full conditionals are available). Table 6 shows the results obtained with the merge-split algorithm (last two columns). The performance of the algorithm is better than that of the Metropolis–Hastings algorithm with a proposal of the form (3.4). However, the improvement is not great and the Gibbs sampler with mutually singular distributions still does far better with much less computing time.

#### 6. DISCUSSION

We have introduced a framework for the use of MCMC algorithms with mixtures of mutually singular distributions. We showed how one can use the usual Metropolis– Hastings algorithm to form an ergodic chain with stationary distribution  $\Pi = \sum_{i \in I} w_i \Pi_i$ where the  $\Pi_i$ 's are mutually singular distributions. We have analyzed three examples in which the method was easy to apply.

However, because of the singularity between the different components, the choice of a good proposal is harder than in the usual setting. The same problem arises with the reversible jump formulation and there has been a great deal of work on efficient construction of reversible jump proposal distributions (Green and Mira 2001; Brooks et al. 2003). Using a simple example we have shown the relationship between our formulation and the reversible jump formulation. Indeed it was possible to derive an algorithm with the same acceptance probability. Which formulation is to be preferred will depend on the application. Our formulation is convenient in the sense that the number of parameters remains the same and we do not have to worry about dimension matching.

When full conditionals are available, the Gibbs sampler can be used and can bring great improvement. Using the Gibbs sampler, no tuning is necessary, which can be a considerable advantage in problems where the number of parameters is large. This was illustrated in the gene expression problem, Example 4. In that example we used only one gene, but in practice one would want to use the same algorithm for thousands of (conditionally independent) genes. Both computation and tuning would be a serious problem. In some cases where the full conditional is not available, it might be available for a distribution that approximates the target distribution. The full conditional of the approximate distribution could be used as a proposal in a Metropolis–Hastings algorithm for the target. This could potentially lead to a high acceptance rate; see Gottardo et al. (2006) for an application. In the past few years there has been some progress towards automatic transdimensional algorithms (Green 2003; Hastie 2004), and we believe that the Gibbs sampler with mutually singular distributions is a step in that direction.

It might be hard to write the prior as a mixture of mutually singular distributions in some complex model selection problems where there is a large difference of dimension between the models. Moreover, full conditionals might not be available and good proposals might be hard to derive. The reversible jump formulation allows for clever moves between models by introducing a Jacobian term in the acceptance ratio. For example, in the mixture problem with an unknown number of components, Richardson and Green (1997) used moment matching conditions to move from one model to the other. It would be hard to formulate this problem in terms of mutually singular distributions. Bayesian analysis of mixture models with an unknown number of components is a problem where it is hard to devise moves with high acceptance rates. There has been some effort to try to create algorithms with better properties (Stephens 2000; Cappé et al. 2003; Brooks et al. 2003).

In this article, we have estimated each model posterior probability as the proportion of time spent in that model. This is what is commonly done in transdimensional MCMC. In the case of the Gibbs sampler we have seen that a simple Rao–Blackwellization step can

improve the estimates. Other methods have also been proposed for getting better estimates of posterior model probabilities from reversible jump output; see Bartolucci et al. (2006).

Our formulation is different from the product space approach (Carlin and Chib 1995; Besag 1997; Godsill 2001; Dellaportas et al. 2002). In our approach we also keep the number of parameters fixed but we allow the dimension of the support of the distribution to vary. We only need to store as many parameters as in the reversible jump formulation. With our formulation, even though the number of parameters across models is the same, there is some redundancy in the parameters. Instead of varying the number of parameters, we vary the dimension of the support of the distribution. This is different from the product space approach where one needs to introduce pseudo-priors in order to keep the parameters that are not in the model currently visited. Our formulation is also more general than that of Petris and Tardella (2003) as (1) the Gibbs sampler can be used, and (2) it can deal with nonnested models. In the gene expression example, if the prior was restricted to the second  $(\mu_1 \neq \mu_2 = \mu_3)$  and third  $(\mu_2 \neq \mu_3 = \mu_1)$  components, their formulation could not be used because the models are not nested, whereas ours could.

#### A. APPENDIX

#### A.1 PROOF OF THEOREM 1

**Proof:** Since  $v_i \perp v_j$  for  $i \neq j$ , we know that there exist sets  $S_i$  such that  $v_i(S_i^c) = 0$ and  $v_j(S_i) = 0$ . Using the assumption that  $\Pi_i$  is dominated by  $v_i$ , we have by the Radon– Nikodym Theorem,

$$\Pi_{i}(A) = \int_{A} \frac{d\Pi_{i}}{dv_{i}}(x)v_{i}(dx)$$
  
= 
$$\int_{A} \frac{d\Pi_{i}}{dv_{i}}(x)\mathbf{1}_{S_{i}}(x)v_{i}(dx)$$
  
= 
$$\int_{A} \frac{d\Pi_{i}}{dv_{i}}(x)\mathbf{1}_{S_{i}}(x)(\sum_{k \in I} v_{k})(dx)$$

The result follows from the fact that  $\Pi = \sum_{i \in I} w_i \Pi_i$  and the linearity of the integral operator.

#### A.2 FULL CONDITIONAL FOR TWO REGRESSION COEFFICIENTS

When two regression coefficients  $\beta_j$ ,  $\beta_k$ , are jointly updated, the prior used is of the form,

$$\Pi(\beta_{j}, \beta_{k} | \sigma_{\beta}^{2}) = [(1 - w_{j})\delta_{0} + w_{j}N(0, \psi_{\beta_{j}}^{-1})][(1 - w_{k})\delta_{0} + w_{k}N(0, \psi_{\beta_{k}}^{-1})]$$
  
=  $(1 - w_{j})(1 - w_{k})\delta_{0}\delta_{0}$   
+ $(1 - w_{j})w_{k}\delta_{0}N(0, \psi_{\beta_{k}}^{-1})$   
+ $w_{j}(1 - w_{k})N(0, \psi_{\beta_{j}}^{-1})\delta_{0}$   
+ $w_{j}w_{k}N(0, \psi_{\beta_{j}}^{-1})N(0, \psi_{\beta_{k}}^{-1}),$ 

where  $\psi_{\beta_j}^{-1} = S_Y^2/S_{\chi_j}^2\sigma_{\beta}^2$ . Thus we can write the prior as a mixture of four mutually singular distributions with respective weights  $u_1 = (1 - w_j)(1 - w_k)$ ,  $u_2 = (1 - w_j)w_k$ ,  $u_3 = w_j(1 - w_k)$  and  $u_4 = w_jw_k$ . Again, given the full conjugacy of the prior to the likelihood, the full conditional can be derived explicitly and is of the form,

$$(\boldsymbol{\beta}_{j,k}|...) \propto u_1 c_1 \delta_0 \delta_0 + u_2 c_2 N\left(\mu_j^*, \psi_j^{*-1}\right) \delta_0 + u_3 c_3 \delta_0 N\left(\mu_k^*, \psi_k^{*-1}\right) + u_4 c_4 N_2\left(\boldsymbol{\mu}_{j,k}^*, \boldsymbol{\Psi}_{j,k}^{*-1}\right), \qquad (A.1)$$

where  $\psi_j^* = (\psi \sum_i \varpi_i X_{ij}^2 + \psi_{\beta_j})^{-1}$  and  $\mu_j^* = \psi_j^{*-1} (\psi \sum_i r_{ij})$ , similarly for  $\psi_k^*$  and  $\mu_k^*$  replacing the index *j* with a *k*. Also

$$r_{ij} = \varpi_i (y_i - \beta_0 - \sum_{l \neq j,k} \beta_l X_{il}) X_{ij},$$

$$c_1 = \left(\frac{\psi}{2\pi}\right)^{n/2} \exp\left(-0.5\psi \sum_i \varpi_i (y_i - \beta_0 - \sum_{l \neq j,k} \beta_l X_{il})^2\right),$$

$$c_2 = c_1 \sqrt{\frac{\psi \beta_j}{\psi_j^*}} \exp\left(0.5\psi_j^{*-1}(\psi \sum_i r_{ij})^2\right),$$

$$c_3 = c_1 \sqrt{\frac{\psi \beta_k}{\psi_k^*}} \exp\left(0.5\psi_k^{*-1}(\psi \sum_i r_{ik})^2\right),$$

$$\Psi_{j,k}^* = \mathbf{D}_{j,k} + \psi \tilde{\mathbf{X}}_{j,k}^t \tilde{\mathbf{X}}_{j,k},$$

$$\mu_{j,k}^* = \Psi_{j,k}^{*-1} \tilde{\mathbf{X}}_{j,k}^t \mathbf{R}_{j,k},$$

$$c_4 = c_1 \sqrt{\frac{\psi \beta_j \psi \beta_k}{|\Psi_{j,k}^*|}} \exp\left(-0.5\psi \tilde{\mathbf{R}}_{j,k}^t \tilde{\mathbf{X}}_{j,k} \Psi_{j,k}^{*-1} \tilde{\mathbf{X}}_{j,k}^t \tilde{\mathbf{R}}_{j,k}\right),$$

where  $R_{j,k} = \sqrt{\psi} (\tilde{\mathbf{Y}} - \tilde{\mathbf{X}}_{-\{j,k\}}^t \boldsymbol{\beta}_{-\{j,k\}})$ ,  $\mathbf{D}_{j,k} = \text{diag}(\psi_{\beta_j}, \psi_{\beta_k})$ ,  $\tilde{\mathbf{X}}_{j,k}$  is the matrix formed with the *j*th and *k*th columns of  $\tilde{\mathbf{X}}$ ,  $\tilde{\mathbf{X}}_{-\{j,k\}}$  is the matrix formed by removing the *j*th and *k*th columns of  $\tilde{\mathbf{X}}$ , and finally  $\tilde{\mathbf{Y}}$  (resp.  $\tilde{\mathbf{X}}$ ) is obtained by multiplying each element (resp. each row) by the corresponding  $\sqrt{\varpi_i}$ . Then, the full conditional is obtained by normalizing (A.1) so that the normalized weights, the *w*\*'s, sum up to one. In this case, the global dominating measure can be written as  $\nu = \delta_0 \times \delta_0 + \delta_0 \times \lambda + \lambda \times \delta_0 + \lambda \times \lambda$ , where  $\nu_1 \times \nu_2$ represents the product measure of  $\nu_1$  with  $\nu_2$ .

#### A.3 FULL CONDITIONAL IN THE THREE-WAY COMPARISON

In Example 4, the full conditional is again available, and is given by

$$(\boldsymbol{\mu}|\boldsymbol{\psi}_{\boldsymbol{\mu}}, \mathbf{w}) \propto w_{1}k_{123}N(\mu_{1}|\mu_{123}^{*}, \psi_{123}^{*}^{-1})\mathbf{1}_{[\mu_{1}=\mu_{2}=\mu_{3}]} + w_{2}k_{1}k_{23}N(\mu_{1}|\mu_{1}^{*}, \psi_{1}^{*-1})N(\mu_{2}|\mu_{23}^{*}, \psi_{23}^{*-1})\mathbf{1}_{[\mu_{1}\neq\mu_{2}=\mu_{3}]} + w_{3}k_{2}k_{13}N(\mu_{2}|\mu_{2}^{*}, \psi_{2}^{*-1})N(\mu_{1}|\mu_{13}^{*}, \psi_{13}^{*-1})\mathbf{1}_{[\mu_{1}=\mu_{3}\neq\mu_{2}]} + w_{4}k_{3}k_{12}N(\mu_{3}|\mu_{3}^{*}, \psi_{3}^{-1})N(\mu_{1}|\mu_{12}^{*}, \psi_{12}^{*-1})\mathbf{1}_{[\mu_{1}=\mu_{2}\neq\mu_{3}]}$$
(A.2)  
$$+ w_{5}k_{1}k_{2}k_{3}N(\mu_{1}|\mu_{1}^{*}, \psi_{1}^{*-1})N(\mu_{2}|\mu_{2}^{*}, \psi_{2}^{*-1})N(\mu_{3}|\mu_{3}^{*}, \psi_{3}^{*-1})\mathbf{1}_{[\mu_{1}\neq\mu_{2}\neq\mu_{3}]}$$

where

$$\psi_c^* = n_c \psi + \psi_{\mu_c}, \quad \mu_c^* = \psi \psi_c^{*-1} \sum_{i=1}^{n_c} y_{ci},$$
$$\psi_{sr}^* = (n_s + n_r) \psi + \psi_{\mu_{sr}}, \quad \mu_{sr}^* = \psi \psi_{rs}^{*-1} (\sum_{i=1}^{n_s} y_{si} + \sum_{i=1}^{n_r} y_{ri}),$$

and

$$\psi_{123}^* = (n_1 + n_2 + n_3)\psi + \psi_{\mu_{123}}, \quad \mu_{123}^* = \psi \psi_{123}^{*-1} \sum_{i,j} y_{ij}.$$

The constants  $k_i$ ,  $k_{sr}$ ,  $k_{123}$  are given by

$$k_{i} = \sqrt{\frac{\psi_{\mu_{i}}}{\psi_{i}^{*}}} \exp\left\{-0.5\psi \sum_{j=1}^{n_{i}} y_{ij}^{2} + 0.5\psi_{i}^{*-1} \left(\psi \sum_{j=1}^{n_{i}} y_{ij}\right)^{2}\right\},$$

$$k_{sr} = \sqrt{\frac{\psi_{\mu_{sr}}}{\psi_{sr}^{*}}} \exp\left\{-0.5\psi \left(\sum_{j=1}^{n_{s}} y_{sj}^{2} + \sum_{j=1}^{n_{r}} y_{rj}^{2}\right) + 0.5\psi_{sr}^{*-1}\psi^{2} \left(\sum_{j=1}^{n_{s}} y_{sj} + \sum_{j=1}^{n_{r}} y_{rj}\right)^{2}\right\}$$

and

$$k_{123} = \sqrt{\frac{\psi_{\mu_{123}}}{\psi_{123}^*}} \exp\left\{-0.5\psi \sum_{i,j} y_{ij}^2 + 0.5\psi_{123}^{*-1}\psi^2 \left(\sum_{i,j} y_{ij}\right)^2\right\}.$$

It follows that the full conditional is again a mixture of five mutually singular distributions, whose weights,  $w^*$ 's, can be obtained by normalizing (A.2).

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