open science

# Bayesian Variable Selection in Markov Mixture Models 

Roberta Paroli, Luigi Spezia

## To cite this version:

Roberta Paroli, Luigi Spezia. Bayesian Variable Selection in Markov Mixture Models. Communications in Statistics - Simulation and Computation, 2008, 37 (01), pp.25-47. 10.1080/03610910701459956 . hal-00514318

HAL Id: hal-00514318

## https://hal.science/hal-00514318

Submitted on 2 Sep 2010

HAL is a multi-disciplinary open access archive for the deposit and dissemination of scientific research documents, whether they are published or not. The documents may come from teaching and research institutions in France or abroad, or from public or private research centers.

L'archive ouverte pluridisciplinaire HAL, est destinée au dépôt et à la diffusion de documents scientifiques de niveau recherche, publiés ou non, émanant des établissements d'enseignement et de recherche français ou étrangers, des laboratoires publics ou privés.


## Bayesian Variable Selection in Markov Mixture Models

| Journal: | Communications in Statistics - Simulation and Computation |
| :---: | :---: |
| Manuscript ID: | LSSP-2007-0008.R1 |
| Manuscript Type: | Original Paper |
| Date Submitted by the Author: | 14-May-2007 |
| Complete List of Authors: | Paroli, Roberta; Università Cattolica S.C., Statistics Spezia, Luigi; Università Cà Foscari, Statistics |
| Keywords: | Gibbs variable selection, Kuo-Mallick method, Metropolized-KuoMallick method, Stochastic search variable selection |
| Abstract: | Bayesian methods for variable selection have become increasingly popular in recent years, due to advances in MCMC computational algorithms. Several methods have been proposed in literature in the case of linear and generalized linear models. In this paper we adapt some of the most popular algorithms to a class of non-linear and non-Gaussian time series models, i.e. the Markov mixture models (MMM). We also propose the "Metropolization" of the algorithm of Kuo and Mallick (1998), in order to tackle variable selection efficiently. Numerical comparisons among the competing MCMC algorithms are also presented via simulation examples. |
| Note: The following files were submitted by the author for peer review, but cannot be converted to PDF. You must view these files (e.g. movies) online. |  |
| fig2-3.wmf fig2-4.wmf <br> Table1.tex <br> Table2.tex <br> Table3.tex <br> Table4.tex <br> Table5.tex |  |

fig1-1.wmf
fig1-2.wmf
fig1-3.wmf fig1-4.wmf fig2-1.wmf fig2-2.wmf Paroli_Spezia2.tex

# BAYESIAN VARIABLE SELECTION IN MARKOV MIXTURE MODELS 

ROBERTA PAROLI<br>Dipartimento di Scienze Statistiche, Università Cattolica S.C., Milano

LUIGI SPEZIA<br>Dipartimento di Statistica, Università Ca' Foscari, Venezia

May 2007


#### Abstract

Bayesian methods for variable selection and model choice have become increasingly popular in recent years, due to advances in Markov chain Monte Carlo (MCMC) computational algorithms. Several methods have been proposed in literature in the case of linear and generalized linear models. In this paper we adapt some of the most popular algorithms to a class of non-linear and non-Gaussian time series models, i.e. the Markov mixture models (MMM). We also propose the "Metropolization" of the algorithm of Kuo and Mallick (1998), in order to tackle variable selection efficiently, both when the complexity of the model is high, as in MMM, and when the exogenous variables are strongly correlated. Numerical comparisons among the competing MCMC algorithms are also presented via simulation examples.


Keywords - Gibbs variable selection; Kuo-Mallick method; Metropolized-Kuo-Mallick method; Stochastic search variable selection.

## 1 Introduction

In recent years, several methods have been proposed in Bayesian literature for selecting exogenous variables in regression models: they include Stochastic Search Variable Selection (SSVS) by George and McCullogh (1993), the unconditional priors Gibbs sampler (KM) by Kuo and Mallick (1998)
and Gibbs Variable Selection (GVS) by Dellaportas et al. (2000). These methods have been also extended to generalized linear models (George et al. 1996), log-linear models (Ntzoufras et al., 2000) and multivariate regression models (Brown et al., 1998). Other authors provide alternative solutions to the variable selection problem: see, among others, Carlin and Chib (1995), Denison et al. (1998), Kohn et al. (2001), Nott and Green (2004), Nott and Kohn (2005), Schneider and Corcoran (2004) for some different approaches and recent developements. We focus our attention on the first three methods because they are the most popular and the most referred in recent literature.

The basic idea of these approaches is that the promising predictors are identified by their highest posterior probability, that is the "best" subset of predictors is that with the most frequent appearance in the sequence of the Markov chain Monte Carlo (MCMC) algorithm developed to evaluate the posterior probability. To give a simple idea, we consider the multiple linear regression model defined by the usual equation

$$
Y=\sum_{j=1}^{q} \beta_{j} X_{j}+E,
$$

where $Y$ is a dependent variable; $X_{1}, X_{2}, \ldots, X_{q}$ are the potential fixed exogenous variables or covariates, which are candidates for inclusion in the model, and $\beta_{1}, \ldots, \beta_{q}$ are the regression coefficients. As usual the r.v. $E$ is Normal with zero mean and variance $\sigma^{2}\left(\mathcal{N}\left(0, \sigma^{2}\right)\right)$. There are $2^{q}$ possible submodels to select, i.e. there are $2^{q}$ possible subsets of predictors. The selection of a subset is also equivalent to setting to 0 those regression coefficients that are not included in the final model, so a binary indicator vector $\gamma$ of dimension $q$ is added in the specification of the model and it represents which of the $q$ variables are included: any $\gamma_{j}(j=1, \ldots, q)$ assumes the value 1 if the variable $X_{j}$ is included, 0 if it is not included.

Let $f(\gamma)$ be the prior distribution of $\gamma$; the "best" subset of regressors is identified by analysing the marginal posterior distribution of $\gamma, f\left(\gamma \mid y^{n}\right) \propto f\left(y^{n} \mid \gamma\right) f(\gamma)$, which contains the information relevant to variable selection supported by the data $y^{n}=\left(y_{1}, y_{2}, \ldots, y_{n}\right)^{\prime}$ and by the prior. Intractable integrals are involved in the computation of this posterior distribution and in the evaluation of its mode, so an MCMC sampling procedure can be developed.

The differences among the above mentioned methods of variable selection depend on how the
vector $\gamma$ is involved in the model. In particular, in the SSVS approach the indicator vector $\gamma$ is involved through the prior of the regression coefficients; in the Kuo-Mallick approach $\gamma$ is part of the regression equation, while the GVS approach is a "hybrid" of SSVS and KM.

Our aim is twofold: first, we extend these three approaches to Markov switching autoregressive models with exogeneous variables (MSARX) and non-homogeneous hidden Markov models (NHHMM), which belong to the class of Markov mixture models (MMM), as defined by Chib (1996); second, we propose the "Metropolization" of the algorithm of Kuo and Mallick (1998), giving rise to Metropolized-Kuo-MallicK (MKMK) method, which increases the performance of the algorithm both when the complexity of the model is high, as in MMM, and when exogenous variables are strongly correlated. The correct selection of exogenous variables is due to the acceptance in block of both the coefficients and the indicators of the exogenous variables, because blocking increases the precision of the estimators and improves the chain mixing.

The paper is organized as follows. The four techniques for variable selection, i.e. KM, SSVS, GVS and MKMK, are applied to MSARX and to NHHMM in Section 2, while, in Section 3, several simulation examples are performed to illustrate the methods and to compare numerically the competing algorithms.

## 2 Variable selection in Markov mixture models

In this section we apply KM, SSVS, GVS and MKMK to the Markov mixtures models. The general MMM is analysed by Chib (1996), who describes it in terms of finite mixture distributions in which the component populations are selected according to contemporary states of an unobserved Markov process: at each time $t$, a realization of the Markov chain, also called regime, occurs hiddenly. The conditional density of the current observation of the time series, given the previous ones and the previous regime, is a finite mixture of densities, whose mixing distribution is the row of the transition matrix corresponding to the previous regime.

Two special MMM are considered here: Markov switching autoregressive models and hidden Markov models. Suppose we have a couple of discrete-time stochastic processes, one observed and one unobserved, or hidden: the realizations of the latter can be observed only through the
realizations of the former one. Suppose also the dynamics of the unobserved process affects the dynamics of the observed one. We model the hidden process as a finite-state Markov chain and assume that the observed process, given the Markov chain, is a sequence of conditional independent random variables, whose conditional distribution depends on the Markov chain only through its contemporary state. The conditional distributions belong to a single parametric family, usually Poisson or Normal. This class of models is called hidden Markov models (HMM). A hidden Markov chain may be assumed because it is either suggested by the physical nature of the observed phenomenon, such as in speech recognition (Elliott, Aggoun, Moore (1995)) and in biology for the analysis of the architecture of proteins and nucleic acids (Koski (2001)), or just for convenience in formulating the model, such as in overdispersed Poisson processes (MacDonald and Zucchini (1997)).

By contrast, in Markov switching autoregressive (MSAR) models the conditional independence condition is replaced by the autoregressive dependence condition and the conditional distributions are always Normals. We have that different autoregressions, each one depending on a latent regime, alternate according to the regime switching, which is driven by the Markov chain. These models have been introduced in econometric literature to study economic and financial time series with asymmetric cycles and changes in regime generated by a stochastic process (Hamilton (1993)). Krolzig (1997), Kim and Nelson (1999) and Franses and van Dijk (2000) provide generalizations and applications of this class of models to economic and financial time series, in which the hidden states represent different macroeconomic regimes.

### 2.1 Markov Switching Autoregressive models with exogenous variables

A Markov switching autoregressive model of order $m$ and $p$, i.e. $\operatorname{MSAR}(m, p)$, is defined as a bivariate discrete-time stochastic process $\left\{S_{t} ; Y_{t}\right\}$, where $\left\{S_{t}\right\}$ is a latent, or hidden, finite-state Markov chain and $\left\{Y_{t}\right\}$, given $\left\{S_{t}\right\}$, satisfies the order- $p$ dependence and the contemporary dependence conditions: we have a sequence of observed random variables $\left\{Y_{t}\right\}$ depending on the $p$ past observations, whose conditional distributions depend on $\left\{S_{t}\right\}$ only through the contemporary state of the Markov chain.

Let $\left\{S_{t}\right\}$ be a discrete-time, first-order, homogeneous, ergodic Markov chain on a finite statespace with cardinality $m$. The transition matrix of the process is $P=\left[p_{i, j}\right]$, where $p_{i, j}=P\left(S_{t}=\right.$ $j \mid S_{t-1}=i$, for any $i, j \in\{1, \ldots, m\}$ and any $t=2, \ldots, n$. The sequence of the states of the Markov chain is denoted with $s^{n}=\left(s_{1}, \ldots, s_{n}\right)^{\prime}$ and, for any $t=1, \ldots, n$, any $s_{t}$ has values in $\{1, \ldots, m\}$.

Hence, given the order- $p$ dependence and the contemporary dependence conditions, the equation describing $\operatorname{MSAR}(m, p)$ is

$$
Y_{t(i)}=\mu_{i}+\sum_{\tau=1}^{p} \varphi_{\tau(i)} y_{t-\tau}+E_{t(i)}
$$

where $Y_{t(i)}$ denotes the conditional variable $Y_{t}$ when $S_{t}=i$, for any $1 \leq t \leq n$ and for any $i \in\{1, \ldots, m\}$, and $E_{t(i)}$ denotes the Gaussian noise $E_{t}$ when $S_{t}=i$, with zero mean and variance $\sigma_{i}^{2}\left(E_{t(i)} \sim \mathcal{N}\left(0 ; \sigma_{i}^{2}\right)\right)$, for any $i \in\{1, \ldots, m\}$, with the discrete-time processes $\left\{E_{t}\right\}$, given $\left\{S_{t}\right\}$, satisfying the conditional independence and the contemporary dependence conditions. Any signal $\mu_{i}$, any variance $\sigma_{i}^{2}$ and any autoregressive coefficient $\varphi_{\tau(i)}$, for any $\tau=1, \ldots, p$, depend on the current state $i$ of the Markov chain, for any $i \in\{1, \ldots, m\}$.

If we join to the $\operatorname{MSAR}(m, p)$ model $q$ contemporary exogenous variables, we have the $\operatorname{MSARX}(m, p)$ model, defined by the equation

$$
\begin{equation*}
Y_{t(i)}=\mu_{i}+\sum_{\tau=1}^{p} \varphi_{\tau(i)} y_{t-\tau}+\sum_{j=1}^{q} \beta_{j(i)} X_{t j}+E_{t(i)} \tag{1}
\end{equation*}
$$

where $X_{t 1}, X_{t 2}, \ldots, X_{t q}$ are the potential fixed exogenous variables which are candidates for inclusion in the final model and $\beta_{i}=\left(\beta_{1(i)}, \ldots, \beta_{q(i)}\right)^{\prime}$ are their $q$ coefficients when the hidden chain visits state $i$; hence the observed process, at any time $t$, can be affected by different exogenous variables, according to the hidden state we assume in $t$. The matrix $n \times q$ of the observations of the exogenous variables will be denoted with $X=\left[x_{t j}\right]$, for $t=1, \ldots, n$ and $j=1, \ldots, q$. To automatically satisfy the constraint of stationarity of any AR subprocess on any vectors $\varphi_{i}=\left(\varphi_{1(i)}, \ldots, \varphi_{p(i)}\right)^{\prime}$, we reparametrize the model in terms of the log trasformation of the partial autocorrelations $r_{i}=\left(r_{1(i)}, \ldots, r_{p(i)}\right)^{\prime}$, according to Mariott et al. (1996).

Identifiability of the model is ensured by imposing increasing signals $\mu_{i}$, but the procedures we introduce can be easily adapted to models with any other type of constraint, provided the chosen constraint respects the geometry of the posterior distribution. An alternative option to
fixed constraints is postprocessing the MCMC outputs (Celeux et al. (2000)).
The unknown parameters of the model are $\left(\mu, \sigma^{2}, R, \beta, P\right)$ where $\mu$ is the vector of the $m$ signals $\mu_{i} ; \sigma^{2}$ is the vector of the $m$ variances $\sigma_{i}^{2} ; R$ is the matrix of the $m$ vectors $R_{i}$ of the $\log$ trasformations of the partial autocorrelations, i.e. $R=\left(R_{1}^{\prime}, \ldots, R_{i}^{\prime}, \ldots, R_{m}^{\prime}\right)^{\prime}$ with $R_{i}=$ $\left(R_{1(i)}, \ldots, R_{\tau(i)}, \ldots, R_{p(i)}\right)^{\prime}$ and $R_{\tau(i)}=\ln \left(\frac{1+r_{\tau(i)}}{1-r_{\tau(i)}}\right) ; \beta$ is the matrix of the $m$ vectors $\beta_{i}$, i.e. $\beta=\left(\beta_{1}^{\prime}, \ldots, \beta_{i}^{\prime}, \ldots, \beta_{m}^{\prime}\right)^{\prime} ; P$ is the matrix of the $m^{2}$ transition probabilities $p_{i, j}$; also the sequence of the hidden states $s^{n}$ is unknown and it will be estimated. The priors are so specified:

1) any parameter $\mu_{i}$ is Normal with known $\mu_{M}$ and $\sigma_{M}^{2}\left(\mu \sim \mathcal{N}\left(\mu_{M}, \sigma_{M}^{2}\right)\right.$, for any $\left.i=1, \ldots, m\right)$;
2) each entry of matrix $R$ is an independent Normal with known $\mu_{R}$ and $\sigma_{R}^{2}\left(R_{\tau(i)} \sim \mathcal{N}\left(\mu_{R} ; \sigma_{R}^{2}\right)\right.$, for any $i=1, \ldots, m$ and for any $\tau=1, \ldots, p)$;
3) each row of matrix $\beta$ is an independent multivariate Normal of dimension $q$ with known vector $\mu_{B}$ and matrix $\Sigma_{B}\left(\beta_{i} \sim \mathcal{N}_{q}\left(\mu_{B} ; \Sigma_{B}\right)\right.$, for any $\left.i=1, \ldots, m\right) ;$
4) any variance $\sigma_{i}^{2}$ is an independent Inverse Gamma with known $\nu$ and $\eta\left(\sigma_{i}^{2} \sim I G(\nu, \eta)\right.$, for any $i=1, \ldots, m) ;$
5) any row of $P, P_{i}=\left(p_{i, 1}, p_{i, 2}, \ldots, p_{i, m}\right)$, is an independent Dirichlet with parameter $\omega=$ $\left(\omega_{1}, \ldots, \omega_{m}\right)\left(P_{i \bullet} \sim \mathcal{D}(\omega)\right.$, for any $\left.i=1, \ldots, m\right)$.

The following subsections will be devoted to develope the MCMC variables selection procedures of KM, SSVS, GVS and MKMK for this type of models. We assume that the observed process can be influenced by different exogenous variables, depending on the state visited by the hidden Markov chain; this means that we have a multiple selection problem, consisting in the estimation of $m$ indicator vectors $\gamma_{i}=\left(\gamma_{1(i)}, \ldots, \gamma_{j(i)}, \ldots, \gamma_{q(i)}\right)^{\prime}$, with $i=1, \ldots, m$, which are entries of the indicator matrix $\gamma=\left(\gamma_{1}^{\prime}, \ldots, \gamma_{i}^{\prime}, \ldots, \gamma_{m}^{\prime}\right)^{\prime}$.

### 2.1.1 KM method

In the KM approach, we assume that binary indicators are associated to the parameters of any exogenous variables, for any state $i$, and that they are embedded in the model equation: model (1) becomes

$$
Y_{t(i)}=\mu_{i}+\sum_{\tau=1}^{p} \varphi_{\tau(i)} y_{t-\tau}+\sum_{j=1}^{q} \gamma_{j(i)} \beta_{j(i)} X_{t j}+E_{t(i)}
$$

The model also includes the prior distributions of the $\gamma_{j(i)}$ 's that are independent Bernoulli with known probability $\xi_{j}\left(\gamma_{j(i)} \sim B e\left(\xi_{j}\right)\right.$, for any $i$ and for any $\left.j=1, \ldots, q\right)$, where $\xi_{j}=P\left(\gamma_{j(i)}=\right.$ 1) $=1-P\left(\gamma_{j(i)}=0\right)$. If $\xi_{j}=0.5$, for any $j=1, \ldots, q$, we have the uniform or indifference prior $f(\gamma)=2^{-q}$, while if $\xi_{j}<0.5$ large models are penalized in favour of the parsimonious ones.

The MCMC scheme used to generate, at any iteration $k$, the sequence of $s^{n(k)}, \mu^{(k)}, \sigma^{2(k)}$, $R^{(k)}, \beta^{(k)}, \gamma^{(k)}, P^{(k)}$ of the KM algorithm for this class of models is Metropolis-within-Gibbs. The algorithm proposed here provides the identifiability constraint on the signals $\left(\mu_{i}^{(k)}<\mu_{j}^{(k)}\right.$, for any $k$ and for any $i, j=1, \ldots, m$, so that $i<j$ ), but it can be easily rearranged when another type of constraint is specified.
[step 1] The sequence $s^{n(k)}$ of hidden states is generated by the well-known forward filteringbackward sampling (ff-bs) algorithm (Chib (1996)).
[step 2] The $m$ parameters $\mu_{i}^{(k)}$, for any $i$, are independently generated from a Normal distribution with mean

$$
\frac{\sigma_{i}^{-2(k-1)} \sum_{\left\{t \geq 1: s_{t}^{(k)}=i\right\}}\left(y_{t}-\sum_{\tau=1}^{p} \varphi_{\tau(i)}^{(k-1)} y_{t-\tau}-\sum_{j=1}^{q} \gamma_{j(i)}^{(k-1)} \beta_{j(i)}^{(k-1)} x_{t, j}\right)+\mu_{M} \sigma_{M}^{-2}}{N_{i}^{(k)} \sigma_{i}^{-2(k-1)}+\sigma_{M}^{-2}}
$$

and variance

$$
\left(N_{i}^{(k)} \sigma_{i}^{-2(k-1)}+\sigma_{M}^{-2}\right)^{-1}
$$

where $N_{i}^{(k)}$ is the number of observations corresponding to the contemporary hidden state $i$ in the sequence $s^{n(k)}$ generated at step [1]. The entries of the vector $\mu^{(k)}$ must be in increasing order to satisfy the identifiability constraint, so we apply the constrained permutation sampling algorithm (Frühwirth-Schnatter (2001)) to order them.
[step 3] The parameters $R_{\tau(i)}^{(k)}$, for any $\tau=1, \ldots, p$ and any $i$, are independently generated from the random walk $R_{\tau(i)}^{(k)}=R_{\tau(i)}^{(k-1)}+U_{R}$, where $U_{R}$ is a univariate Gaussian noise with zero mean and constant variance, for any $k$. Then any vector $R_{i}^{(k)}$ is accepted with probability

$$
A\left(R_{i}^{(k-1)} ; R_{i}^{(k)}\right)=\min \left\{1 ; \frac{\pi\left(R_{i}^{(k)} \mid \mu^{(k)}, \sigma^{2(k-1)}, \beta^{(k-1)}, \gamma^{(k-1)}, s^{n(k)}, X, y^{n}\right)}{\pi\left(R_{i}^{(k-1)} \mid \mu^{(k)}, \sigma^{2(k-1)}, \beta^{(k-1)}, \gamma^{(k-1)}, s^{n(k)}, X, y^{n}\right)}\right\}
$$

for any $i$, where the numerator of the acceptance ratio is proportional to
$\exp \left\{-\frac{1}{2} \sigma_{i}^{-2(k-1)} \sum_{\left\{t \geq 1: s_{t}^{(k)}=i\right\}}\left(y_{t}-\mu_{i}^{(k)}-\sum_{\tau=1}^{p} \varphi_{\tau(i)}^{(k)} y_{t-\tau}-\sum_{j=1}^{q} \gamma_{j(i)}^{(k-1)} \beta_{j(i)}^{(k-1)} x_{t, j}\right)^{2}-\frac{1}{2} \sigma_{R}^{-2} \sum_{\tau=1}^{p}\left(R_{\tau(i)}^{(k)}-\mu_{R}\right)^{2}\right\}$
and the denominator is identical to the numerator except the introduction of $\varphi_{\tau(i)}^{(k-1)}$ and $R_{i}^{(k-1)}$. Notice that once the $R_{i}^{(k)}$, s have been generated, the new vectors $\varphi_{i}^{(k)}$ 's can be obtained by the inverse trasformation.
[step 4] The parameters $\beta_{i}^{(k)}$, for any $i$, are independently generated from a Normal distribution of dimension $q$ with mean vector

$$
\left(X^{* \prime} Q_{i}^{(k)} X^{*}+\Sigma_{B}^{-1}\right)^{-1}\left(X^{* \prime} Q_{i}^{(k)} \widehat{y}_{i}^{T(k)}+\Sigma_{B}^{-1} \mu_{B}\right)
$$

and covariance matrix

$$
\left(X^{* \prime} Q_{i}^{(k)} X^{*}+\Sigma_{B}^{-1}\right)^{-1}
$$

where $X^{*}$ is the $n \times q$ matrix with elements $x_{t j}^{*}=\gamma_{j(i)}^{(k-1)} x_{t j}$, for any $t=1, \ldots, n$ and $j=1, \ldots, q$; $Q_{i}^{(k)}$ is a $n \times n$ diagonal matrix whose $t-t h$ term on the diagonal is either $\sigma_{i}^{-2(k-1)}$, if $s_{t}^{(k)}=i$, or zero, otherwise; $\widehat{y}_{i}^{n}$ is a $n$-dimensional vector whose generic $t$-th element is either $y_{t}-\mu_{i}^{(k)}-\sum_{\tau=1}^{p} \varphi_{\tau(i)}^{(k)} y_{t-\tau}$, if $s_{t}^{(k)}=i$, or zero, otherwise.
[step 5] For any $i$-th row of the indicator matrix $\gamma^{(k)}$, every coefficient $\gamma_{j(i)}^{(k)}$ is independently generated from a Bernoulli distribution

$$
\gamma_{j(i)}^{(k)} \sim f\left(\gamma_{j(i)} \mid \mu^{(k)}, \sigma^{2(k-1)}, \beta^{(k)}, \gamma_{\backslash j(i)}^{(k-1)}, s^{n(k)}, X, y^{n}\right) \sim B e\left(\pi_{j(i)}^{(k)}\right)
$$

where $\gamma_{\backslash j(i)}^{(k-1)}=\left(\gamma_{1(i)}^{(k)}, \ldots, \gamma_{j-1(i)}^{(k)}, \gamma_{j+1(i)}^{(k-1)}, \ldots, \gamma_{q(i)}^{(k-1)}\right)^{\prime}$ is the vector of the entries of $\gamma_{(i)}^{(k-1)}$ in which the $j$-th entry is suppressed and the first $(j-1)$-th are updated by the new values. The parameter of the Bernoulli distribution is $\pi_{j(i)}^{(k)}=a_{j(i)}^{(k)} /\left(a_{j(i)}^{(k)}+b_{j(i)}^{(k)}\right)$, with

$$
\begin{aligned}
& a_{j(i)}^{(k)}=f\left(y^{n} \mid \mu^{(k)}, \sigma^{2(k-1)}, \beta^{(k)}, s^{n(k)}, X, \gamma_{\backslash j(i)}^{(k-1)}, \gamma_{j(i)}^{(k)}=1\right) f\left(\gamma_{\backslash j(i)}^{(k-1)}, \gamma_{j(i)}^{(k)}=1\right) \propto \\
& \propto \exp \left\{-\frac{1}{2} \sigma_{i}^{-2(k-1)} \sum_{\left\{t \geq 1: s_{t}^{(k)}=i\right\}}\left(y_{t}-\mu_{i}^{(k)}-\sum_{\tau=1}^{p} \varphi_{\tau(i)}^{(k)} y_{t-\tau}-\sum_{h=1}^{q} \theta_{h(i)}^{*(k)} x_{t h}\right)^{2}\right\} \xi_{j}
\end{aligned}
$$

and

$$
\begin{aligned}
& b_{j(i)}^{(k)}=f\left(y^{n} \mid \mu^{(k)}, \sigma^{2(k-1)}, \beta^{(k)}, s^{n(k)}, X, \gamma_{\backslash j(i)}^{(k-1)}, \gamma_{j(i)}^{(k)}=0\right) f\left(\gamma_{\backslash j(i)}^{(k-1)}, \gamma_{j(i)}^{(k)}=0\right) \propto \\
& \propto \exp \left\{-\frac{1}{2} \sigma_{i}^{-2(k-1)} \sum_{\left\{t \geq 1: s_{t}^{(k)}=i\right\}}\left(y_{t}-\mu_{i}^{(k)}-\sum_{\tau=1}^{p} \varphi_{\tau(i)}^{(k)} y_{t-\tau}-\sum_{h=1}^{q} \theta_{h(i)}^{* *(k)} x_{t h}\right)^{2}\right\}\left(1-\xi_{j}\right),
\end{aligned}
$$

where $f\left(y^{n} \mid \mu^{(k)}, \sigma^{2(k-1)}, \beta^{(k)}, s^{n(k)}, X, \gamma_{\backslash j(i)}^{(k-1)}, \gamma_{j(i)}^{(k)}=1\right)$ and $f\left(y^{n} \mid \mu^{(k)}, \sigma^{2(k-1)}, \beta^{(k)}, s^{n(k)}, X, \gamma_{\backslash j(i)}^{(k-1)}\right.$, $\gamma_{j(i)}^{(k)}=0$ ) are the likelihoods computed either with $\gamma_{j(i)}^{(k)}=1$ or $\gamma_{j(i)}^{(k)}=0$, with the $\theta_{h(i)}^{*(k)} \mathrm{s}$ entries of the vector $\theta_{(i)}^{(k)}=\left(\gamma_{1(i)}^{(k)} \beta_{1(i)}^{(k)}, \ldots, \gamma_{q(i)}^{(k-1)} \beta_{q(i)}^{(k)}\right)^{\prime}$ in which the $j$-th entry is replaced by $\beta_{j(i)}^{(k)}$, i.e. $\theta_{(i)}^{*(k)}=\left(\gamma_{1(i)}^{(k)} \beta_{1(i)}^{(k)}, \ldots, \gamma_{j-1(i)}^{(k)} \beta_{j-1(i)}^{(k)}, \beta_{j(i)}^{(k)}, \quad \gamma_{j+1(i)}^{(k-1)} \beta_{j+1(i)}^{(k)}, \ldots, \gamma_{q(i)}^{(k-1)} \beta_{q(i)}^{(k)}\right)^{\prime}$, and the $\theta_{h(i)}^{* *(k)}$ 's entries of the vector $\theta_{(i)}^{(k)}$ in which the $j$-th entry is replaced by 0 , i.e. $\theta_{(i)}^{* *(k)}=$ $\left(\gamma_{1(i)}^{(k)} \beta_{1(i)}^{(k)}, \ldots, \gamma_{j-1(i)}^{(k)} \beta_{j-1(i)}^{(k)}, 0, \gamma_{j+1(i)}^{(k-1)} \beta_{j+1(i)}^{(k)}, \ldots, \gamma_{q(i)}^{(k-1)} \beta_{q(i)}^{(k)}\right)^{\prime}$.
[step 6] The $m$ parameters $\sigma_{i}^{2(k)}$ are independently generated from an Inverse Gamma distribution with parameters

$$
\frac{N_{i}^{(k)}}{2}+\nu \quad \text { and } \quad \frac{1}{2} \sum_{\left\{t \geq 1: s_{t}^{(k)}=i\right\}}\left(y_{t}-\mu_{i}^{(k)}-\sum_{\tau=1}^{p} \varphi_{\tau(i)}^{(k)} y_{t-\tau}-\sum_{j=1}^{q} \gamma_{j(i)}^{(k)} \beta_{j(i)}^{(k)} x_{t, j}\right)^{2}+\eta
$$

[step 7] Each vector $P_{i \bullet}^{(k)}$ is independently generated from a Dirichlet distribution $\mathcal{D}\left(\omega+N_{i \bullet}^{(k)}\right)$, where $N_{i \bullet}^{(k)}=\left(N_{i, 1}^{(k)}, \ldots, N_{i, m}^{(k)}\right)$, with $N_{i, j}^{(k)}$ the number of couples of consecutive hidden states $i, j$ in the sequence $s^{n(k)}$, for any $i, j$.

### 2.1.2 SSVS method

In the SSVS method, the indicator vectors $\gamma_{i}$ are involved in the model through the prior of the coefficients $\beta_{j(i)}$, that are mixtures of two Normals with different variances

$$
\beta_{j(i)} \mid \gamma_{j(i)} \sim \gamma_{j(i)} \mathcal{N}\left(\mu_{\beta_{j}}, c_{j}^{2} \tau_{j}^{2}\right)+\left(1-\gamma_{j(i)}\right) \mathcal{N}\left(\mu_{\beta_{j}}, \tau_{j}^{2}\right)
$$

for specified $c_{j}^{2}$ and $\tau_{j}^{2}$ and $\mu_{\beta_{j}}$ the $j$-th entry of $\mu_{B}$. The prior distribution of any vector $\beta_{i}$, for any $i$, is a multivariate Normal of dimension $q$

$$
\beta_{i} \mid \gamma_{i} \sim \mathcal{N}_{q}\left(\mu_{B}, D_{\gamma(i)}\right)
$$

where $D_{\gamma(i)}=\operatorname{diag}\left[\left(\delta_{1(i)} \tau_{1}\right)^{2}, \ldots,\left(\delta_{q(i)} \tau_{q}\right)^{2}\right]$ with $\delta_{j(i)}=c_{j}$ if $\gamma_{j(i)}=1$ and $\delta_{j(i)}=1$ if $\gamma_{j(i)}=0$. The choice of the hyperparameters $c_{j}^{2}$ and $\tau_{j}^{2}$ and the interpretation of this prior is discussed in George and McCullogh (1993) and Ntzoufras et al. (2000). The hyperparameters $c_{j}$ and $\tau_{j}$ must be chosen carefully because the posterior probability is mostly dependent on them, given that the exogenous variables are never eliminated from the model space, but coefficients of the variables to be excluded are set close to zero with high probability. So it is recommended to run various
simulations with different values of $c_{j}$ and $\tau_{j}$ and to analyse how the posterior probability of any vector $\gamma_{i}$ is sensitive to the choice of these tuning factors.

The prior distributions of the remaining parameters are the same as for those in the KM method.
In the SSVS scheme used to generate the sequence of $s^{n(k)}, \mu^{(k)}, \sigma^{2(k)}, R^{(k)}, \beta^{(k)}, \gamma^{(k)}, P^{(k)}$, only the step [7] remains the same as for KM, since it does not depend on $\beta_{i}$ and $\gamma_{i}$; step [5] is completly different, while the remaining steps are identical to those in KM, except the presence of the $\gamma_{j(i)}$ 's that multiply the $\beta_{j(i)}$ 's. More in details, step [5] is so defined.
[step 5] For any $i$-th row of the indicator matrix $\gamma^{(k)}$, every coefficient $\gamma_{j(i)}^{(k)}$ is independently generated from a Bernoulli distribution $B e\left(a_{j(i)}^{(k)} /\left(a_{j(i)}^{(k)}+b_{j(i)}^{(k)}\right)\right)$, where

$$
a_{j(i)}^{(k)}=f\left(\beta_{i}^{(k)} \mid \gamma_{\backslash j(i)}^{(k-1)}, \gamma_{j(i)}^{(k)}=1\right) f\left(\gamma_{\backslash j(i)}^{(k-1)}, \gamma_{j(i)}^{(k)}=1\right) \propto f\left(\beta_{i}^{(k)} \mid \gamma_{\backslash j(i)}^{(k-1)}, \gamma_{j(i)}^{(k)}=1\right) \xi_{j}
$$

and

$$
b_{j(i)}^{(k)}=f\left(\beta_{i}^{(k)} \mid \gamma_{\backslash j(i)}^{(k-1)}, \gamma_{j(i)}^{(k)}=0\right) f\left(\gamma_{\backslash j(i)}^{(k-1)}, \gamma_{j(i)}^{(k)}=0\right) \propto f\left(\beta_{i}^{(k)} \mid \gamma_{\backslash j(i)}^{(k-1)}, \gamma_{j(i)}^{(k)}=0\right)\left(1-\xi_{j}\right)
$$

with $f\left(\beta_{i}^{(k)} \mid \gamma_{\backslash j(i)}^{(k-1)}, \gamma_{j(i)}^{(k)}=1\right)$ or $f\left(\beta_{i}^{(k)} \mid \gamma_{\backslash j(i)}^{(k-1)}, \gamma_{j(i)}^{(k)}=0\right)$ the density functions of the prior distribution of $\beta_{i} \mid \gamma_{i}$ when $\gamma_{j(i)}=1$ or $\gamma_{j(i)}=0$, respectively. Their explicit forms are

$$
f\left(\beta_{i}^{(k)} \mid \gamma_{\backslash j(i)}^{(k-1)}, \gamma_{j(i)}^{(k)}=1\right) \propto \exp \left\{-\frac{1}{2} \sum_{j \neq i}\left(\beta_{i}^{(k)}-\mu_{B}\right)^{\prime} D_{\gamma(i)}^{*(k)-1}\left(\beta_{i}^{(k)}-\mu_{B}\right)\right\}
$$

and

$$
f\left(\beta_{i}^{(k)} \mid \gamma_{\backslash j(i)}^{(k-1)}, \gamma_{j(i)}^{(k)}=0\right) \propto \exp \left\{-\frac{1}{2} \sum_{j \neq i}\left(\beta_{i}^{(k)}-\mu_{B}\right)^{\prime} D_{\gamma(i)}^{* *(k)-1}\left(\beta_{i}^{(k)}-\mu_{B}\right)\right\}
$$

where $D_{\gamma(i)}^{*(k)}$ and $D_{\gamma(i)}^{* *(k)}$ are the covariance matrix of the prior of $\beta_{i}^{(k)}$ in which the $j$-th entries of $D_{\gamma(i)}^{*(k)}$ are replaced by $\delta_{j(i)}=c_{j}\left(\right.$ which corresponds to $\left.\gamma_{j(i)}^{(k)}=1\right)$, while those of $D_{\gamma(i)}^{* *(k)}$ by $\delta_{j(i)}=1\left(\right.$ which corresponds to $\left.\gamma_{j(i)}^{(k)}=0\right)$.

### 2.1.3 GVS method

GVS is a mixed approach between the two previous methods. Here the indicator variable influences both the prior distribution of the regression coefficients, as in SSVS, and the model equation, as in KM. The prior distributions are the same as for SSVS. In the MCMC scheme, steps [1], [2], [3], [6] and [7] are identical to those in KM, while steps [4] and [5] are a combination of KM and SSVS.

In details, the elements of the parameters of the Bernoulli distributions of any $\gamma_{j(i)}^{(k)}$ are

$$
a_{j(i)}^{(k)} \propto f\left(y^{n} \mid \mu^{(k)}, \sigma^{2(k-1)}, \beta^{(k)}, s^{n(k)}, X, \gamma_{\backslash j(i)}^{(k-1)}, \gamma_{j(i)}^{(k)}=1\right) f\left(\beta_{i}^{(k)} \mid \gamma_{\backslash j(i)}^{(k-1)}, \gamma_{j(i)}^{(k)}=1\right) \xi_{j}
$$

and

$$
b_{j(i)}^{(k)} \propto f\left(y^{n} \mid \mu^{(k)}, \sigma^{2(k-1)}, \beta^{(k)}, s^{n(k)}, X, \gamma_{\backslash j(i)}^{(k-1)}, \gamma_{j(i)}^{(k)}=0\right) f\left(\beta_{i}^{(k)} \mid \gamma_{\backslash j(i)}^{(k-1)}, \gamma_{j(i)}^{(k)}=0\right)\left(1-\xi_{j}\right) .
$$

Notice that both the likelihoods $f\left(y^{n} \mid \cdot\right)$, as in KM, and the conditional priors $f\left(\beta_{i}^{(k)} \mid \cdot\right)$, as in SSVS, are present in the expressions of $a_{j(i)}^{(k)}$ and $b_{j(i)}^{(k)}$.

### 2.1.4 MKMK method

When the complexity of the model is high, as for MSARX and NHHMM, or when exogenous variables are strongly correlated, the performances of KM, SSVS and GVS can be improved by introducing Metropolized-Kuo-MallicK method (MKMK), based on the acceptance in block of the coefficients $\beta_{i}$ 's and the indicators of the exogenous variables $\gamma_{i}$ 's: Liu, Wong, Kong (1994) showed that the precision of the estimators increases when the unknown parameters are gathered in separated blocks, while Gilks and Roberts (1996) showed that blocking improves the chain mixing.

In the basic MCMC algorithm only one step must be modified w.r.t. KM: for any state $i$, the coefficients of the exogenous variables and the associated dummy indicators are sampled jointly and updated in block within a Metropolis step. For MSARX models, the pair of vectors $\left(\beta_{i}^{(k)} ; \gamma_{i}^{(k)}\right)$ must be accepted in block, so steps [4] and [5] of Subsection 2.1.1 are joined in only one step and modified as follows. We assume a random walk proposal for any $\beta_{j(i)}^{(k)}$ and an independent Bernoulli proposal for any $\gamma_{j(i)}^{(k)}$.

For any $i(i=1, \ldots, m)$, first, parameters $\beta_{j(i)}^{(k)}$, for any $j=1, \ldots, q$, are independently generated from the random walk $\beta_{j(i)}^{(k)}=\beta_{j(i)}^{(k-1)}+U_{B}$, where $U_{B}$ is a Gaussian noise with zero mean and constant variance, for any $k$; then, every parameter $\gamma_{j(i)}^{(k)}$, for any $j=1, \ldots, q$, is independently generated from a Bernoulli distribution with probability $\pi_{j(i)}^{(k)}=a_{j(i)}^{(k)} /\left(a_{j(i)}^{(k)}+b_{j(i)}^{(k)}\right)$, where $a_{j(i)}^{(k)}$ and $b_{j(i)}^{(k)}$ are the same as in step [5] of Subsection 2.1.1; finally the pair of vectors $\left(\beta_{i}^{(k)} ; \gamma_{i}^{(k)}\right)$ is
accepted in block with probability

$$
\begin{aligned}
& A\left[\left(\beta_{i}^{(k-1)} ; \gamma_{i}^{(k-1)}\right) ;\left(\beta_{i}^{(k)} ; \gamma_{i}^{(k)}\right)\right]=\min \left\{1 ; \frac{\pi\left(\beta_{i}^{(k)} ; \gamma_{i}^{(k)}\right) q\left(\beta_{i}^{(k-1)} ; \gamma_{i}^{(k-1)} \mid \beta_{i}^{(k)} ; \gamma_{i}^{(k)}\right)}{\pi\left(\beta_{i}^{(k-1)} ; \gamma_{i}^{(k-1)}\right) q\left(\beta_{i}^{(k)} ; \gamma_{i}^{(k)} \mid \beta_{i}^{(k-1)} ; \gamma_{i}^{(k-1)}\right)}\right\}= \\
& =\min \left\{1 ; \frac{f\left(y^{n} \mid \beta_{i}^{(k)} ; \gamma_{i}^{(k)}\right) p\left(\gamma_{i}^{(k)}\right) p\left(\beta_{i}^{(k)}\right) q\left(\gamma_{i}^{(k-1)} \mid \beta_{i}^{(k-1)} ; \beta_{i}^{(k)} ; \gamma_{i}^{(k)}\right)}{f\left(y^{n} \mid \beta_{i}^{(k-1)} ; \gamma_{i}^{(k-1)}\right) p\left(\gamma_{i}^{(k-1)}\right) p\left(\beta_{i}^{(k-1)}\right) q\left(\gamma_{i}^{(k)} \mid \beta_{i}^{(k)} ; \beta_{i}^{(k-1)} ; \gamma_{i}^{(k-1)}\right)}\right\},
\end{aligned}
$$

by the factorization of the proposal density, i.e.

$$
q\left(\beta_{i}^{(k)} ; \gamma_{i}^{(k)} \mid \beta_{i}^{(k-1)} ; \gamma_{i}^{(k-1)}\right)=q\left(\gamma_{i}^{(k)} \mid \beta_{i}^{(k)} ; \beta_{i}^{(k-1)} ; \gamma_{i}^{(k-1)}\right) q\left(\beta_{i}^{(k)} \mid \beta_{i}^{(k-1)} ; \gamma_{i}^{(k-1)}\right)
$$

by the independence of $\beta_{i}^{(k)}$ on $\gamma_{i}^{(k-1)}$ and cancelling the ratio $q\left(\beta_{i}^{(k-1)} \mid \beta_{i}^{(k)}\right) / q\left(\beta_{i}^{(k)} \mid \beta_{i}^{(k-1)}\right)$ for the symmetry of the proposal distribution.

### 2.2 Non-homogeneous hidden Markov models

In the MSARX models of the previous section, the Markov chain of the latent process is supposed homogeneous, that is the transition probabilities do not depend on time $t$. In the most general case, the Markov chain can be non-homogeneous, that is the transition probabilities are time-varying, by assuming that their dynamics depend on exogenous variables. If we consider a HMM with a nonhomogeneous Markov chain, we obtain a Non-Homogeneous Hidden Markov model (NHHMM). It is a discrete-time stochastic processes $\left\{S_{t}, Y_{t}\right\}$ where $\left\{S_{t}\right\}$ is now a non-homogeneous Markov chain whose transition probabilities $p_{i, j}^{t}=P\left(S_{t}=j \mid S_{t-1}=i\right)$ vary at any time $t$ by depending on exogenous variables. We suppose that they can be expressed as logistic functions of $q$ exogenous variables. Let $X_{t}$ be a $(q+1)$-dimensional vector of the observations of the exogenous deterministic variables at time $t, X_{t}=\left(1, x_{t, 1}, \ldots, x_{t, q}\right)^{\prime}$ and $\alpha_{i, j}$ be a $(q+1)$-dimensional vector of parameters, $\alpha_{i, j}=\left(\alpha_{0(i, j)}, \alpha_{1(i, j)}, \ldots, \alpha_{q(i, j)}\right)^{\prime}$, for any $i, j \in\{1, \ldots, m\}$, if $i \neq j$, and a $(q+1)$-dimensional vector of zeros, if $i=j$. For any $t=2, \ldots, n$, the transition probabilities can be expressed as

$$
\begin{align*}
& \operatorname{logit}\left(p_{i, j}^{t}\right)=\ln \left(p_{i, j}^{t} / p_{i, i}^{t}\right)=\sum_{h=0}^{q} \alpha_{h(i, j)} x_{t h}=X_{t}^{\prime} \alpha_{i, j} \\
& p_{i, j}^{t}=\left(\exp \left(X_{t}^{\prime} \alpha_{i, j}\right)\right) /\left(1+\sum_{j \neq i} \exp \left(X_{t}^{\prime} \alpha_{i, j}\right)\right), \tag{2}
\end{align*}
$$

for any $i, j \in\{1, \ldots, m\}$ and let $P_{t}=\left[p_{i, j}^{t}\right]$ be the $m \times m$ transition matrix of the chain, for any $t=2, \ldots, n$.

Process $\left\{Y_{t}\right\}$, given process $\left\{S_{t}\right\}$, is an observed sequence of conditionally independent random variables, whose conditional distributions depend on $\left\{S_{t}\right\}$ only through the contemporary $S_{t}$. Here we assume the probability density function of any $Y_{t}$, given $S_{t}$, is Gaussian; hence we have Gaussian NHHMM. The model equation is

$$
Y_{t(i)}=\mu_{i}+E_{t(i)}
$$

where $E_{t(i)}$ denotes the Gaussian random variable $E_{t}$, when $S_{t}=i$, with zero mean and variance $\sigma_{i}^{2}\left(E_{t(i)} \sim \mathcal{N}\left(0 ; \sigma_{i}^{2}\right)\right)$, for any $i$, with the discrete-time process $\left\{E_{t}\right\}$, given $\left\{S_{t}\right\}$, satisfying the conditional independence and the contemporary dependence conditions; so $Y_{t(i)} \sim \mathcal{N}\left(\mu_{i} ; \sigma_{i}^{2}\right)$, for any $i$. Notice that the inferential tools and the variable selection machinery we are going to describe can also be applied when another conditional distribution, both discrete and continuous, is hypothesized.

Identifiability of the model is again ensured by imposing increasing signals $\mu_{i}$. The unknown parameters and the latent data of the NHHMM are $\mu, \sigma^{2}, \alpha, s^{n}$, where $\mu$ is the vector of the $m$ signals $\mu_{i}, \sigma^{2}$ is the vector of the $m$ variances $\sigma_{i}^{2}, \alpha$ is the $m \times m$ matrix of the vectors $\alpha_{i, j}$ and $s^{n}$ is the sequence of the hidden states.

The following prior assumptions are stated in analogy with those in Subsection 2.1:

1) the prior distribution of any parameter $\mu_{i}$ is an independent Normal with known $\mu_{M}$ and $\sigma_{M}^{2}$ $\left(\mu_{i} \sim \mathcal{N}\left(\mu_{M}, \sigma_{M}^{2}\right)\right.$, for any $\left.i=1, \ldots, m\right) ;$
2) the prior distribution of any vector $\alpha_{i, j}$ is an independent multivariate Normal of dimension $(q+1)$ with known vector $\mu_{A}$ and matrix $\Sigma_{A}\left(\alpha_{i, j} \sim \mathcal{N}_{q+1}\left(\mu_{A} ; \Sigma_{A}\right)\right.$, for any $i, j=1, \ldots, m$, with $i \neq j) ;$
3) the prior distribution of any variance $\sigma_{i}^{2}$ is assumed to be an independent Inverse Gamma $\left(\sigma_{i}^{2} \sim I G(\nu, \eta)\right.$, for any $\left.i=1, \ldots, m\right)$.

### 2.2.1 KM method

Let $\alpha_{i}=\left(\alpha_{0(i)}, \alpha_{1(i)}, \ldots, \alpha_{q(i)}\right)^{\prime}, i=1, \ldots, m$, be the $i$-th row of matrix $\alpha$; we associate to any entry of matrix $\alpha$, i.e. the vector of parameters $\alpha_{i, j}$, a vector of binary indicators $\gamma_{(i)}=$ $\left(1, \gamma_{1(i)}, \ldots, \gamma_{q(i)}\right)^{\prime}$ of dimension $q+1$, which regulates the inclusion of the exogenous variables
when the hidden chain visits state $i$ at time $t-1$; finally, let $\gamma$ be the $m \times(q+1)$ matrix of any vector $\gamma_{(i)}$, i.e. $\gamma=\left(\gamma_{(1)}^{\prime}, \ldots, \gamma_{(m)}^{\prime}\right)^{\prime}$. Expression (2) becomes

$$
\begin{align*}
& \operatorname{logit}\left(p_{i, j}^{t}\right)=\sum_{h=0}^{q} \gamma_{h(i)} \alpha_{h(i, j)} x_{t h}=X_{t}^{\prime} \theta_{(i, j)} \\
& p_{i, j}^{t}=\exp \left(\sum_{h=0}^{q} \gamma_{h(i)} \alpha_{h(i, j)} x_{t h}\right) /\left(1+\sum_{j \neq i} \exp \left(\sum_{h=0}^{q} \gamma_{h(i)} \alpha_{h(i, j)} x_{t h}\right)\right)=  \tag{3}\\
& =\exp \left(X_{t}^{\prime} \theta_{(i, j)}\right) /\left(1+\sum_{j \neq i} \exp \left(X_{t}^{\prime} \theta_{(i, j)}\right)\right),
\end{align*}
$$

with $\theta_{(i, j)}=\left(\alpha_{0(i, j)}, \gamma_{1(i)} \alpha_{1(i, j)}, \ldots, \gamma_{q(i)} \alpha_{q(i, j)}\right)^{\prime}$.
Notice that we assume the transition from state $i$ to any other state $j$ is regulated by the same set of covariates, hence the same vector $\gamma_{(i)}$ is associated to any $\left(\alpha_{0(i, j)}, \alpha_{1(i, j)}, \ldots, \alpha_{q(i, j)}\right)^{\prime}$ $(j=1, \ldots, m)$ and consequently to $P_{i \bullet}$.

The prior distribution of any $\gamma_{h(i)}$ is an independent Bernoulli with known probability $\xi_{h}$ $\left(\gamma_{h(i)} \sim B e\left(\xi_{h}\right)\right.$, for any $i=1, \ldots, m$ and any $\left.h=1, \ldots, q\right)$.

The steps of the MCMC algorithm used to simulate the sequence of $s^{n(k)}, \mu^{(k)}, \alpha^{(k)}, \gamma^{(k)}, \sigma^{2(k)}$ are the following.
[step 1] The sequence of hidden states $s^{n(k)}$ is still generated by the ff-bs algorithm.
[step 2] The $m$ parameters $\mu_{i}^{(k)}$, for any $i$, are independently generated from a Normal distribution with mean

$$
\frac{\sigma_{i}^{-2(k-1)} \sum_{\left\{t \geq 1: s_{t}^{(k)}=i\right\}} y_{t}+\mu_{M} \sigma_{M}^{-2}}{N_{i}^{(k)} \sigma_{i}^{-2(k-1)}+\sigma_{M}^{-2}}
$$

and variance

$$
\left(N_{i}^{(k)} \sigma_{i}^{-2(k-1)}+\sigma_{M}^{-2}\right)^{-1}
$$

with the entries of the vector $\mu^{(k)}$ in increasing order to satisfy the identifiability constraint, so we apply the constrained permutation sampling algorithm (Frühwirth-Schnatter (2001)) to order them.
[step 3] Any vector $\alpha_{i, j}^{(k)}$, for any $i, j$, with $i \neq j$, is independently generated from a random walk $\alpha_{i, j}^{(k)}=\alpha_{i, j}^{(k-1)}+U_{A}$, where $U_{A}$ is a multivariate Gaussian noise with zero mean and constant
precision matrix, for any $k$. Then any matrix $\alpha_{i}=\left(\alpha_{i, 1}^{\prime}, \ldots, \alpha_{i, m}^{\prime}\right)^{\prime}$ is accepted with probability

$$
A\left(\alpha_{i}^{(k-1)} ; \alpha_{i}^{(k)}\right)=\min \left\{1 ; \frac{\pi\left(\alpha_{i}^{(k)} \mid s^{n(k)}, \gamma_{(i)}^{(k-1)}, X\right)}{\pi\left(\alpha_{i}^{(k-1)} \mid s^{n(k)}, \gamma_{(i)}^{(k-1)}, X\right)}\right\}
$$

where the numerator of the acceptance ratio is proportional to

$$
\prod_{\left\{t \geq 2: s_{t-1}^{(k)}=i\right\}} p_{s_{t-1}^{(k)}, s_{t}^{(k)}}^{t(k)} \exp \left\{-\frac{1}{2} \sum_{j \neq i}\left(\theta_{(i, j)}^{(k)}-\mu_{A}\right)^{\prime} \Sigma_{A}^{-1}\left(\theta_{(i, j)}^{(k)}-\mu_{A}\right)\right\}
$$

and the denominator is identical to the numerator except the introduction of $\alpha_{i}^{(k-1)}$ and $p_{i, j}^{t(k-1)}$, with $p_{i, j}^{t(k)}$ and $p_{i, j}^{t(k-1)}$ functions of $\alpha_{i}^{(k)}$ and $\alpha_{i}^{(k-1)}$, respectively, and $X=\left(X_{1}^{\prime}, \ldots, X_{t}^{\prime}, \ldots, X_{n}^{\prime}\right)^{\prime}$. [step 4] Any coefficient $\gamma_{h(i)}^{(k)}(h=1, \ldots, q ; i=1, \ldots, m)$ is independently generated from a Bernoulli distribution $\operatorname{Be}\left(a_{j(i)}^{(k)} /\left(a_{j(i)}^{(k)}+b_{j(i)}^{(k)}\right)\right)$, with

$$
a_{h(i)}^{(k)} \propto \prod_{\left\{t \geq 2: s_{t-1}^{(k)}=i\right\}} \frac{\exp \left(X_{t}^{\prime} \theta_{\left(i, s_{t}^{(k)}\right)}^{*(k)}\right)}{1+\sum_{j \neq i} \exp \left(X^{\prime} \theta_{(i, j)}^{*(k)}\right)} \xi_{h}
$$

and

$$
b_{h(i)}^{(k)} \propto \prod_{\left\{t \geq 2: s_{t-1}^{(k)}=i\right\}} \frac{\exp \left(X_{t}^{\prime} \theta_{\left(i, s_{t}^{(k)}\right)}^{* *(k)}\right)}{1+\sum_{j \neq i} \exp \left(X_{t}^{\prime} \theta_{(i, j)}^{* *(k)}\right)}\left(1-\xi_{h}\right)
$$

where $\theta_{(i, j)}^{*(k)}$ is the vector $\theta_{(i, j)}^{(k)}$ in which the $h$-th entry is replaced by $\alpha_{h(i, j)}^{(k)}$, i.e. $\theta_{(i, j)}^{*(k)}=\left(\alpha_{0(i, j)}^{(k)}, \ldots\right.$, $\left.\gamma_{h-1(i)}^{(k)} \alpha_{h-1(i, j)}^{(k)}, \alpha_{h(i, j)}^{(k)}, \gamma_{h+1(i)}^{(k-1)} \alpha_{h+1(i, j)}^{(k)}, \ldots, \gamma_{q(i)}^{(k-1)} \alpha_{q(i, j)}^{(k) \prime}\right)^{\prime}$, while in $\theta_{(i, j)}^{* *(k)}$ the $h$-th entry of $\theta_{(i, j)}^{(k)}$ is replaced by 0 , i.e. $\theta_{(i, j)}^{* *(k)}=\left(\alpha_{0(i, j)}^{(k)}, \ldots \gamma_{h-1(i)}^{(k)} \alpha_{h-1(i, j)}^{(k)}, 0, \gamma_{h+1(i)}^{(k-1)} \alpha_{h+1(i, j)}^{(k)}, \ldots, \gamma_{q(i)}^{(k-1)} \alpha_{q(i, j)}^{(k)}\right)^{\prime}$. [step 5] The $m$ parameters $\sigma_{i}^{2(k)}$ are generated from independent Inverse Gamma distributions with parameters

$$
\frac{N_{i}^{(k)}}{2}+\nu \quad \text { and } \quad \frac{1}{2} \sum_{\left\{t \geq 1: s_{t}^{(k)}=i\right\}}\left(y_{t}-\mu_{i}^{(k)}\right)^{2}+\eta
$$

### 2.2.2 SSVS method

In the SSVS method, the indicator variable is involved in the model through the conditional priors of the coefficients $\alpha_{i, j}$ that are mixtures of two Normals with different variances: i.e. any vector $\alpha_{i, j}$, given $\gamma_{(i)}$, is an independent multivariate Normal of dimension $(q+1)$

$$
\alpha_{i, j} \mid \gamma_{(i)} \sim \mathcal{N}_{q+1}\left(\mu_{A}, D_{\gamma(i)}\right)
$$

where $D_{\gamma(i)}=\operatorname{diag}\left[1,\left(\delta_{1(i)} \tau_{1}\right)^{2}, \ldots,\left(\delta_{q(i)} \tau_{q}\right)^{2}\right]$ with $\delta_{h(i)}=c_{h}$ if $\gamma_{h(i)}=1$ and $\delta_{h(i)}=1$ if $\gamma_{h(i)}=0$, for specified $c_{h}$ and $\tau_{h}$.

The prior distributions of the remaining parameters are the same of those elicited for the KM method; the expression of the logit trasformation of the $p_{i, j}^{t}$ 's is still the (2).

In the SSVS scheme used to generate the sequence $s^{n(k)}, \mu^{(k)}, \sigma^{2(k)}, \alpha^{(k)}, \gamma^{(k)}$, steps [2] and [5] remain the same as for KM, since they do not depend on $\alpha_{(i)}$ and $\gamma_{(i)}$, steps [3] and [4] change totally and step [1] is identical to that in KM except the presence of the $\gamma_{(i)}$ 's that multiply the $\alpha_{i, j}$ 's.
[step 3] Any vector $\alpha_{i, j}^{(k)}$, for any $i, j$, with $i \neq j$, is independently generated from a multivariate random walk $\alpha_{i, j}^{(k)}=\alpha_{i, j}^{(k-1)}+U_{A}$. Then any matrix $\alpha_{i}^{(k)}=\left(\alpha_{i, 1}^{(k) \prime}, \ldots, \alpha_{i, m}^{(k) \prime}\right)^{\prime}$ is accepted with probability

$$
A\left(\alpha_{i}^{(k-1)} ; \alpha_{i}^{(k)}\right)=\min \left\{1 ; \frac{\pi\left(\alpha_{i}^{(k)} \mid s^{n(k)}, X\right)}{\pi\left(\alpha_{i}^{(k-1)} \mid s^{n(k)}, X\right)}\right\}
$$

where the numerator of the acceptance ratio is proportional to

$$
\prod_{\left\{t \geq 2: s_{t-1}^{(k)}=i\right\}} p_{s_{t-1}^{(k)}, s_{t}^{(k)}}^{t(k)} \cdot \exp \left\{-\frac{1}{2} \sum_{j \neq i}\left(\alpha_{i, j}^{(k)}-\mu_{A}\right)^{\prime} D_{\gamma(i)}^{(k)-1}\left(\alpha_{i, j}^{(k)}-\mu_{A}\right)\right\}
$$

and the denominator is identical to the numerator except the introduction of $\alpha_{i}^{(k-1)}$ and $p_{i, j}^{t(k-1)}$, with $p_{i, j}^{t(k)}$ and $p_{i, j}^{t(k-1)}$ functions of $\alpha_{i}^{(k)}$ and $\alpha_{i}^{(k-1)}$, respectively, and $X=\left(X_{1}^{\prime}, \ldots, X_{t}^{\prime}, \ldots, X_{n}^{\prime}\right)^{\prime}$. [step 4] For every hidden state $i$, any coefficient $\gamma_{h(i)}^{(k)}$ is independently generated, for $h=1, \ldots, q$, from a Bernoulli distribution $B e\left(a_{j(i)}^{(k)} /\left(a_{j(i)}^{(k)}+b_{j(i)}^{(k)}\right)\right)$, with

$$
a_{h(i)}^{(k)} \propto f\left(\alpha_{i}^{(k)} \mid \gamma_{\backslash h(i)}^{(k-1)}, \gamma_{h(i)}^{(k)}=1\right) \xi_{h} \quad \text { and } \quad b_{h(i)}^{(k)} \propto f\left(\alpha_{i}^{(k)} \mid \gamma_{\backslash h(i)}^{(k-1)}, \gamma_{h(i)}^{(k)}=0\right)\left(1-\xi_{h}\right),
$$

where $f\left(\alpha_{i}^{(k)} \mid \gamma_{\backslash h(i)}^{(k-1)}, \gamma_{h(i)}^{(k)}=1\right)$ or $f\left(\alpha_{i}^{(k)} \mid \gamma_{\backslash h(i)}^{(k-1)}, \gamma_{h(i)}^{(k)}=0\right)$ are the density functions of the $(q+1)$ variate Normal prior distribution of $\alpha_{i, j} \mid \gamma_{(i)}$ when $\gamma_{h(i)}=1$ or $\gamma_{h(i)}=0$, respectively, that is

$$
f\left(\alpha_{i}^{(k)} \mid \gamma_{\backslash h(i)}^{(k-1)}, \gamma_{h(i)}^{(k)}=1\right) \propto \exp \left\{-\frac{1}{2} \sum_{j \neq i}\left(\alpha_{i, j}^{(k)}-\mu_{A}\right)^{\prime} D_{\gamma(i)}^{*(k)-1}\left(\alpha_{i, j}^{(k)}-\mu_{A}\right)\right\}
$$

and

$$
f\left(\alpha_{i}^{(k)} \mid \gamma_{\backslash h(i)}^{(k-1)}, \gamma_{h(i)}^{(k)}=0\right) \propto \exp \left\{-\frac{1}{2} \sum_{j \neq i}\left(\alpha_{i, j}^{(k)}-\mu_{A}\right)^{\prime} D_{\gamma(i)}^{* *(k)-1}\left(\alpha_{i, j}^{(k)}-\mu_{A}\right)\right\},
$$

where $D_{\gamma(i)}^{*(k)}$ and $D_{\gamma(i)}^{* *(k)}$ are the covariance matrix of the prior of $\alpha_{i, j}^{(k)}$ in which the $h$-th entries $\delta_{h(i)}$ are replaced by $c_{h}\left(\right.$ which corresponds to $\left.\gamma_{h(i)}^{(k)}=1\right)$, and by $1\left(\right.$ which corresponds to $\left.\gamma_{h(i)}^{(k)}=0\right)$, respectively.

### 2.2.3 GVS method

As in the previous sections, the GVS method is a mixed approach between the two previous methods: the indicator variable influences both the prior distributions of the coefficients of the exogenous variables, as in SSVS, and the equation of the model (3), as in KM.

The prior distributions are the same of those in SSVS; in the MCMC scheme, steps [1], [2], [5] are identical to those in KM, step [3] is identical to that in SSVS, while step [4] is a combination of KM and SSVS. In details, the elements of the parameters of the Bernoulli distributions of any $\gamma_{j(i)}^{(k)}$ are

$$
a_{h(i)}^{(k)} \propto f\left(s^{n} \mid \alpha_{i}^{(k)}, X, \gamma_{\backslash h(i)}^{(k-1)}, \gamma_{h(i)}^{(k)}=1\right) f\left(\alpha_{i}^{(k)} \mid \gamma_{\backslash h(i)}^{(k-1)}, \gamma_{h(i)}^{(k)}=1\right) \xi_{h}
$$

and

$$
b_{h(i)}^{(k)} \propto f\left(s^{n} \mid \alpha_{i}^{(k)}, X, \gamma_{\backslash h(i)}^{(k-1)}, \gamma_{h(i)}^{(k)}=0\right) f\left(\alpha_{i}^{(k)} \mid \gamma_{\backslash h(i)}^{(k-1)}, \gamma_{h(i)}^{(k)}=0\right)\left(1-\xi_{h}\right)
$$

Notice that both the densities $f\left(s^{n(k)} \mid \cdot\right)$, as in KM, and the conditional priors $f\left(\alpha_{i}^{(k)} \mid \cdot\right)$, as in SSVS, are present in the expressions of $a_{h(i)}^{(k)}$ and $b_{h(i)}^{(k)}$.

### 2.2.4 MKMK method

For NHHMM the pair of vectors $\left(\alpha_{i}^{(k)} ; \gamma_{(i)}^{(k)}\right)$ must be accepted in block, so steps [3] and [4] of Subsection 2.2 .1 are joined in only one step and modified as follows. We assume the following proposals: a multivariate random walk for any $\alpha_{i, j}^{(k)}$ and an independent Bernoulli for any entry of the vector $\gamma_{(i)}^{(k)}$.

For any $i$, first, parameters $\alpha_{i, j}^{(k)}(j=1, \ldots, m$, with $i \neq j)$ are independently generated from the multivariate random walk $\alpha_{i, j}^{(k)}=\alpha_{i, j}^{(k-1)}+U_{A}$, then, every parameter $\gamma_{h(i)}^{(k)}$, for any $h=1, \ldots, q$, is independently generated from a Bernoulli distribution with probability $\pi_{h(i)}^{(k)}=$ $a_{h(i)}^{(k)} /\left(a_{h(i)}^{(k)}+b_{h(i)}^{(k)}\right)$, where $a_{h(i)}^{(k)}$ and $b_{h(i)}^{(k)}$ are the same as in step [4] of Subsection 2.2.1; finally
the pair of vectors $\left(\alpha_{i}^{(k)} ; \gamma_{(i)}^{(k)}\right)$ is accepted in block with probability

$$
\begin{aligned}
& A\left[\left(\alpha_{i}^{(k-1)} ; \gamma_{i}^{(k-1)}\right) ;\left(\alpha_{i}^{(k)} ; \gamma_{i}^{(k)}\right)\right]=\min \left\{1 ; \frac{\pi\left(\alpha_{i}^{(k)} ; \gamma_{i}^{(k)}\right) q\left(\alpha_{i}^{(k-1)} ; \gamma_{i}^{(k-1)} \mid \alpha_{i}^{(k)} ; \gamma_{i}^{(k)}\right)}{\pi\left(\alpha_{i}^{(k-1)} ; \gamma_{i}^{(k-1)}\right) q\left(\alpha_{i}^{(k)} ; \gamma_{i}^{(k)} \mid \alpha_{i}^{(k-1)} ; \gamma_{i}^{(k-1)}\right)}\right\}= \\
& =\min \left\{1 ; \frac{f\left(y^{n} \mid \alpha_{i}^{(k)} ; \gamma_{i}^{(k)}\right) p\left(\gamma_{i}^{(k)}\right) p\left(\alpha_{i}^{(k)}\right) q\left(\gamma_{i}^{(k-1)} \mid \alpha_{i}^{(k-1)} ; \alpha_{i}^{(k)} ; \gamma_{i}^{(k)}\right)}{f\left(y^{n} \mid \alpha_{i}^{(k-1)} ; \gamma_{i}^{(k-1)}\right) p\left(\gamma_{i}^{(k-1)}\right) p\left(\alpha_{i}^{(k-1)}\right) q\left(\gamma_{i}^{(k)} \mid \alpha_{i}^{(k)} ; \alpha_{i}^{(k-1)} ; \gamma_{i}^{(k-1)}\right)}\right\}
\end{aligned}
$$

by the factorization of the proposal density, by the independence of $\alpha_{i}^{(k)}$ on $\gamma_{i}^{(k-1)}$ and cancelling the ratio $q\left(\alpha_{i}^{(k-1)} \mid \alpha_{i}^{(k)}\right) / q\left(\alpha_{i}^{(k)} \mid \alpha_{i}^{(k-1)}\right)$ for the symmetry of the proposal distribution.

## 3 Simulation examples

### 3.1 MSARX models

We generate time series of length $n=500$, for each pair ( $m, p$ ) with $m=2,3$ and $p=1,2,3$ and different coefficient vectors $\beta_{i}$ and $\gamma_{i}$ for any states $i=1, \ldots, m$. The data generation details are the following: we obtain the exogenous variables $X_{1}, \ldots, X_{5}$ as a multivariate Normal, with independent components or with low or high correlations between them. Specifically, for the case of correlated variables, we set to non zero the correlations between variables $\left(X_{2}, X_{4}\right)$ and $\left(X_{1}, X_{5}\right)$, while the correlations between the other variables are random values in modulo less then the previous ones. With these values we explore several cases of correlation between variables which are included in or excluded from the true model. For any state $i=1, \ldots, m$ and $\tau=1, \ldots, p$, the values $r_{\tau(i)}$ 's are randomly generated from uniform distributions between -1 and +1 and then transformed to obtain the autoregressive coefficients $\varphi_{\tau(i)}$ 's. The remaining parameters are set as follows:

- for $m=2$, we suppose that the true model contains the covariates $X_{4}$ and $X_{5}$ in state 1, and $X_{2}, X_{3}, X_{4}$ in state 2 . The corresponding coefficients and the values of the other parameters are:

| state | $\mu_{i}$ | $\sigma_{i}^{2}$ | $\beta_{j(i)}$ |  |  |  |  |  | $\gamma_{j(i)}$ |  |  |  |  | $P$ |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | 1 | 1 | 0 | 0 | 0 | 1 | 1.2 | 0 | 0 | 0 | 1 | 1 | 0.6 | 0.4 |  |  |
| 2 | 3 | 0.8 | 0 | 1 | 0.8 | 1 | 0 | 0 | 1 | 1 | 1 | 0 | 0.2 | 0.8 |  |  |

- for $m=3$, we suppose that in states 1 and 2 there are the same covariates of the case $m=2$, while in state 3 the covariates in the model are $X_{2}, X_{4}$. The corresponding coefficients and the
true values of the parameters are:

| state | $\mu_{i}$ | $\sigma_{i}^{2}$ | $\beta_{j(i)}$ |  |  |  |  | $\gamma_{j(i)}$ |  |  |  | $P$ |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | 1 | 1 | 0 | 0 | 0 | 1 | 1.2 | 0 | 0 | 0 | 1 | 1 | 0.6 | 0.3 | 0.1 |
| 2 | 3 | 0.8 | 0 | 1 | 0.8 | 1 | 0 | 0 | 1 | 1 | 1 | 0 | 0.05 | 0.8 | 0.15 |
| 3 | 5 | 1.25 | 0 | 1.2 | 0 | 0.8 | 0 | 0 | 1 | 0 | 1 | 0 | 0.1 | 0.2 | 0.7 |

Notice that the signals $\mu_{i}$ are in increasing order to ensure the identifiability.
For all the analysed models, we use independent $\mathcal{N}(0,10)$ priors for any parameter $\mu_{i}$ and $R_{\tau(i)}$ (for any $i=1, \ldots, m$, and $\left.\tau=1, \ldots, p\right)$; independent $I G(0.5,0.5)$ for any variances $\sigma_{i}^{2}$ and indifference priors for any $\gamma_{j(i)}$ with $\xi_{j}=0.5$. The priors of any row of the transition matrix $P_{i}$ are independent Dirichlet with parameter $\omega=\left(\omega_{1}, \ldots, \omega_{k}, \ldots, \omega_{m}\right)$, and $\omega_{k}=m \cdot I(i=k)+0.6 \cdot I(i \neq$ $k)$, for any $i, k=1, \ldots, m$, where $I(\cdot)$ is the indicator function. By this prior the probability of persistence is greater than the probability of transition: the probabilty of persistence is about 0.7 and it slowly decreases as the number of states increases. Finally, in KM the priors of any $\beta_{j(i)}$ are independent $\mathcal{N}(0,10)$, for $i=1, \ldots, m$ and $j=1, \ldots, q$, while in SSVS and GVS the conditional priors of $\beta_{i} \mid \gamma_{i}$ are $\mathcal{N}_{5}\left(0, D_{\gamma(i)}\right)$, where $D_{\gamma(i)}$ is defined through the choice of the tuning parameters $c_{j}$ and $\tau_{j}$ : by the results of several simulations of the full models and following the suggestions by George and McCullogh (1993), we use $c_{j}=10$ and $\tau_{j}=0.3$, for any $j=1, \ldots, q$.

Tables $1,2,3,4$ show the results of the posterior model probabilities, estimated through the relative frequencies, of the highest frequency models, for all the variable selection methods we analysed, both in the case of independent covariates and in some special cases of correlation (i.e. $0.3,0.7,0.9)$. The results are based on 300,000 iterations for all the algorithms, after a burn-in period of 100,000 iterations. All the Markov chain are initialized at the full model, while the starting points of the parameters are generated at random from the priors.

## TABLES 1, 2, 3, 4 ABOUT HERE

The results show that in general all the four methods predict the right MSARX model, choosing that model which includes the right covariates. In the case of independent covariates the model probability is high for each pair of $m$ and $p$, while as the correlation between covariates increases the
model probability decreases till to assume very small values (i.e. $7 \%$, for the case of $\operatorname{MSARX}(2,1)$, in Table 4, for the SSVS method).

It can be pointed out also that all the sampling schemes explore local modes of the posterior distribution because different results are obtained in different runs of the algorithms, although all runs are able to identify the true model.

It can be seen that SSVS generally performs worse than the other methods, but reasonably well, and only in the case of independent covariates the model probabilities are higher than $60 \%$, while in the case of high correlation they are always less than $30 \%$.

On the other hand, GVS seems to be the best method, since in all the analysed cases, the true model is identified with probability greater than $60 \%$, sometimes exceeding the value $90 \%$, even if the choice of the tuning factors can represent a serious drawback of this method, as in SSVS. Finally, the performances of KM and MKMK are almost equivalent and their model probabilities, in some cases, tend to the values reached by GVS.

In Figure 1 the MCMC evolution of the corresponding ergodic posterior probabilities are plotted in some special cases, at every $50-\mathrm{th}$ iteration after the burn-in. We dealt with the convergence of the MCMC algorithms by using many plots of output values and formal convergence tests (i.e. Gamerman (1997)).

## FIGURE 1 ABOUT HERE

We compared also the performances of all methods in problems with more data, for example with $n=1000$, or with a large number of potential covariates, $q=30$ or 50 . As we would expect, the results are similar to those presented here and, again, GVS seems to perform better than other methods. Moreover, as the number of data increases the values of the highest posterior model probability increases and the convergence is reached more fastly.

### 3.2 NHHMM

The simulation examples for evaluating the performance of the previous methods are based on time series of length $n=500$, generated from NHHMM $(m)$, for each values of $m=2,3$, with
$q=5$ potential covariates and different coefficients $\alpha_{h(i, j)}$ of the logistic trasformation in (2), for $i=1, \ldots, m ; j=1, \ldots, m$ and $h=0,1, \ldots, q$.

The exogenous variables $X_{1}, \ldots, X_{5}$ are obtained again as multivariate Normal; for the cases of correlated covariates we suppose that the pairs of variables $\left(X_{2}, X_{4}\right)$ and $\left(X_{1}, X_{5}\right)$ have a nonzero fixed correlation, while the correlations between other variables are zero or random values, in modulo less than the previous ones. The parameters values are taken as follows:

- for $m=2$, we suppose that the true model contains covariates $X_{2}, X_{3}$, in state 1 and $X_{2}, X_{3}$, $X_{4}$, in state 2 ; the values of the parameters are:

| state | $\mu_{i}$ | $\sigma_{i}^{2}$ | $\alpha_{h(i, j)}$ |  |  |  |  |  | $\gamma_{h(i)}$ |  |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | 1 | 1 |  | $\alpha_{(1,2)}$ | 2 | 0 | -2 | 2 | 0 | 0 | 1 | 0 | 1 | 1 | 0 | 0 |
| 2 | 3 | 0.8 | $\alpha_{(2,1)}$ | 2 | 0 | -2 | 2 | 2 | 0 | 1 | 0 | 1 | 1 | 1 | 0 |  |

- for $m=3$, we suppose in states 1 and 2 the model contains the same covariates as in case $m=2$, while in the state 3 the covariates in the model are $X_{1}$ and $X_{3}$; the values of the parameters are:

| state | $\mu_{i}$ | $\sigma_{i}^{2}$ | $\alpha_{h(i, j)}$ |  |  |  |  |  |  |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | 1 | 1 | $\alpha_{(1,2)}$ | 2 | 0 | -2 | 2 | 0 | 0 | 1 | 0 | 1 | 1 | 0 | 0 |
| 2 | 3 | 0.8 | $\alpha_{(1,3)}$ | 0.5 | 0 | 1 | 0.5 | 0 | 0 | 1 | 0 | 1 | 1 | 0 | 0 |
| 3 | 5 | 1.25 | $\alpha_{(2,1)}$ | 2 | 0 | -2 | 2 | 2 | 0 | 1 | 0 | 1 | 1 | 1 | 0 |
|  | $\alpha_{(2,3)}$ | 1 | 0 | 0.5 | -0.5 | 1 | 0 | 1 | 0 | 1 | 1 | 1 | 0 |  |  |
|  | $\alpha_{(3,1)}$ | 2 | -2 | 0 | 2 | 0 | 0 | 1 | 1 | 0 | 1 | 0 | 0 |  |  |
|  | $\alpha_{(3,2)}$ | 2 | -2 | 0 | 2 | 0 | 0 | 1 | 1 | 0 | 1 | 0 | 0 |  |  |

We use independent $\mathcal{N}(0,10)$ priors for any parameter $\mu_{i}$; independent $I G(0.5,0.5)$ for any variance $\sigma_{i}^{2}$ and indifference priors for any $\gamma_{h(i)}$, with $\xi_{h}=0.5$, for any $h=1, \ldots, q$. For KM the prior of any vector $\alpha_{i, j}$ (for any $i=1, \ldots, m$, and $j=1, \ldots, m$ ) is a multivariate Normal $N_{6}\left(0, \Sigma_{A}\right)$, with $\Sigma_{A}=10 * I$ (where $I$ is the identity matrix), while, for SSVS and GVS, the conditional priors of vectors $\alpha_{i, j} \mid \gamma_{(i)}$ are $\mathcal{N}_{6}\left(0, D_{\gamma(i)}\right)$, where $D_{\gamma(i)}$ is defined throught the choice of the tuning parameters $c_{h}$ and $\tau_{h}$, which are still set $c_{h}=10$ and $\tau_{h}=0.3$, for any $i=1, \ldots, m$ and $h=1, \ldots, 5$. Again, all the Markov chain are initialized at the full model and the starting values of the parameters are generated at random from the priors.

In Table 5 the highest posterior model probabilities are diplayed, both in the case of independent covariates and in some special cases of correlation (i.e. $0.3,0.7,0.9$ ). The results are based on 300,000 iterations, after a burn-in period of 100,000 iterations.

## TABLE 5 ABOUT HERE

For KM, SSVS, GVS, the results are very similar to those described for the MSARX models, even if the values of the posterior model probablilities are much smaller: this means that the performance of the methods decreases when the complexity of the model increases $(m>2)$, expecially also when exogenous variables are strongly correlated. The algorithms in fact visit a larger set of promising models that incorrectly include other variables or exclude the correct variables: we can see that, in the case of independence, the true models have posterior model probabilities less than $30 \%$ or, when high correlation exists ( 0.7 or 0.9 ), less than $10 \%$. In these last cases (marked with $\left(^{*}\right)$ in Table 5) the algorithms do not identify the correct model and the associated probabilities are not the highest.

By contrast, MKMK works extremly well and always identify the right set of covariates with high probablilities, both when the correlations we fixed are null and when they are positive.

In Figure 2 the MCMC evolution of the corresponding ergodic posterior probabilities are plotted in some special cases, at every $50-$ th iteration after the burn-in. The convergence of the samplers has been checked through plots and traditional convergence tests (i.e. Gamer$\operatorname{man}(1997))$.

## FIGURE 2 ABOUT HERE

The comparison among the results of Table 6 and those of the previous tables shows that our method is very satisfactory: the values of the posterior model probabilities are always very high, also in the critical cases of NHHMM.

## 4 Conclusions

Bayesian variable selection is an active research area in recent years. In this paper we developed some alternative Bayesian procedures, based on Markov chain Monte Carlo algorithms, for variable
selection in special Markov mixture models, i.e. the Markov switching autoregressive models with exogenous variables (MSARX) and the non-homogeneous hidden Markov models (NHHMM). We focused our attention on three methods based on the analysis of the highest model posterior probability: the Stochastic Search Variable Selection (SSVS) by George and McCullogh (1993), the unconditional priors Gibbs sampler (KM) by Kuo and Mallick (1998) and the Gibbs Variable selection (GVS) by Dellaportas et al. (2000). These methods have been initially proposed for linear models and so we decided to explore their performances when they are applied to highly complex stochastic processes.

Simulation results show that GVS selects the true model with an excellent relative frequency, both when the exogenous variables are independent and when they are correlated. By contrast, SSVS and KM visit the right models with a much lower probability and sometimes they are not able to choose the best ones.

The selection procedures of SSVS and GVS provides the employment of some tuning factors which, in our modelling, are hard to identify; hence, our goal was the improvement of the performances of KM, which is extremely straightforward, since it requires only to specify the prior and no pilot runs are required to define the hyperparameters.

Our method is based on the scheme of KM, which is improved by accepting in block, through a Metropolis step, both the coefficients and the indicators of the exogenous variables; so we decided to call it Metropolized-Kuo-MallicK (MKMK). MKMK works very well and better than the other methods, expecially in the NHHMM case, because always selects the right model with a high model posterior probability. The correct selection of exogenous variables is due to the acceptance in block of both the coefficients and the indicators of the exogenous variables, because blocking increases the precision of the estimators and improves the chain mixing.

In Paroli and Spezia (2005), we successfully applied MKMK to a real data problem in the context of non-homogeneous Markov mixtures of periodic autoregressions, that is MSARX models with a Markov-dependent periodic component, all driven by NHHMM. We observed a series of hourly mean concentrations of sulphur dioxide simultaneously to six meteorological variables (wind speed, temperature, rain, solar radiation, relative humidity, pressure) and MKMK allows us to select the
exogenous variables that, for any state of the hidden Markov chain, can influence the dynamics of the observed process and/or the latent one.

## References

Brown, P.J. Vannucci, M., Fearn, T. (1998) Multivariate Bayesian variable selection and prediction. Journal of the Royal Statistical Society, Series B, 60, 627-641.

Carlin, B.P. and Chib, S. (1995) Bayesian model choice via Markov chain Monte Carlo. Journal of the Royal Statistical Society, Series B, 57, 473-484.

Celeux G., Hurn M., Robert C.P. (2000) Computational and Inferential Difficulties With Mixture Posterior Distributions. Journal of the American Statistical Association, 95, 957-970.

Chib, S. (1996) Calculating posterior distributions and modal estimates in Markov mixture models. Journal of Econometrics, 75, 79-97.

Dellaportas, P., Forster, J.J., Ntzoufras, I. (2000) Bayesian Variable Selection Using the Gibbs Sampler. In Generalized Linear Models: A Bayesian Perspective (eds D.K. Dey, S.K. Ghosh, B.K. Mallick), 273-286. Marcel Dekker, New York.

Denison, D.G.T., Mallick, B.K., Smith, A.F.M. (1998) Automatic Bayesian curve fitting. Journal of the Royal Statistical Society, Series B, 60, 333-350.

Elliott R.J., Aggoun L., Moore J.B. (1995) Hidden Markov Models: Estimation and Control. Springer, New York.

Franses P.H. and van Dijk D. (2000) Nonlinear Time Series Models in Empirical Finance. Cambridge University Press, Cambridge.

Frühwirth-Schnatter, S. (2001) Markov Chain Monte Carlo Estimation of Classical and Dynamic Switching and Mixture Models. Journal of the American Statistical Association, 96, 194-209.

Gamerman, D. (1997) Markov Chain Monte Carlo: Stochastic simulation for Bayesian inference. Chapman \& Hall, London.

George, E.I. and McCullogh, R.E. (1993) Variables selection via Gibbs-Sampling. Journal of the American Statistical Association, 88, 881-889.

George, E.I., McCullogh, R.E., Tsay, R. (1996) Two Approaches to Bayesian Model Selection
with Applications. In Bayesian Analysis in Statistics and Econometrics: Essays in Honor of Arnold Zellner (eds D.A. Berry, K.M. Chaloner, J.K. Geweke), 339-348. John Wiley \& Sons, New York.

Gilks W.R. and Roberts G.O. (1996) Strategies for improving MCMC. In Markov Chain Monte Carlo in Prectice (eds W.R. Gilks, S. Richardson, D.J. Spieghelhalter), 89-114. Chapman \& Hall, London.

Hamilton J.D. (1993) Estimation, Inference and Forecasting of Time Series Subject to Changes in Regime. In Handbook of Statistics, vol. 11 (eds G.S. Maddala, C.R. Rao, H.D. Vinod), 231-259. North-Holland, Amsterdam.

Kim C.J. and Nelson C.R. (1999) State-Space Models with Regime Switching: Classical and Gibbs-Sampling Approaches with Applications. The MIT Press, Cambridge.

Kohn, R., Smith, M., Chan, D. (2001) Nonparametric regression using linear combinations of basis functions. Statistics and Computing, 11, 313-322.

Koski T. (2001) Hidden Markov Models for Bioinformatics. Kluwer Academic Publishers, Dordrecht.

Krolzig H.-M. (1997) Markov-Switching Vector Autoregressions: Modelling, Statistical Inference and Applications to Business Cycle Analysis. Springer, Berlin.

Kuo, L. and Mallick, B. (1998) Variable Selection for Regression Models. Sankhya, Series B, 60, 65-81.

Liu J.S., Wong W.H., Kong A. (1994) Covariance structure of the Gibbs sampler with applications to the comparisons of the estimators and augmentation scheme. Biometrika, 81, 27-40.

MacDonald I.L. and Zucchini W. (1997) Hidden Markov and Other Models for Discrete-valued Time Series. Chapman \& Hall, London.

Mariott, J., Ravishanker, N., Gelfand, A., Pai, J. (1996) Bayesian Analysis of ARMA Processes: Complete Sampling-Based Inference under Exact Likelihoods. In Bayesian Analysis in Statistics and Econometrics: Essays in Honor of Arnold Zellner (eds D.A. Berry, K.M. Chaloner, J.K. Geweke), 243-256. John Wiley \& Sons, New York.

Nott, D.J. and Green, P.J. (2004) Bayesian variable selection and the Swendsen-Wang Algo-
rithm. Journal of Computational and Graphical Statistics, 13, 1-17.
Nott, D.J. and Kohn, R. (2005) Adaptive sampling for Bayesian variable selection. Biometrika, 92, 747-763.

Ntzoufras, I., Forster, J.J., Dellaportas, P. (2000) Stochastic Search Variable Selection for Loglinear Models. Journal of Statistical Computation and Simulation, 68, 23-37.

Paroli, R. and Spezia, L. (2005) Non-Homogeneous Markov Mixture of Periodic Autoregressions and Sulphur Dioxide Concentrations. Serie E.P. n. 126, Istituto di Statistica, Università Cattolica S.C., Milano.

Schneider, U. and Corcoran, J.N. (2004) Perfect Sampling for Bayesian Variable Selection in a Linear Regression Model. Journal of Statistical Planning and Inference, 126, 153-171.


Table 1: Highest posterior model probabilities for independent covariates

|  | corr $=0.3$ |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | $\mathrm{m}=2$ |  |  | $\mathrm{m}=3$ |  |  |  |
|  | state | 1 | 2 | state | 1 | 2 | 3 |
|  | model var. | 4,5 | 2,3,4 | model var. | 4,5 | 2,3,4 | 2,4 |
| $\mathrm{p}=1$ | KM | 0.61 | 0.76 | KM | 0.66 | 0.73 | 0.60 |
|  | SSVS | 0.53 | 0.35 | SSVS | 0.55 | 0.46 | 0.24 |
|  | GVS | 0.92 | 0.95 | GVS | 0.77 | 0.94 | 0.75 |
|  | MKMK | 0.78 | 0.82 | MKMK | 0.46 | 0.73 | 0.61 |
| $\mathrm{p}=2$ | KM | 0.62 | 0.77 | KM | 0.75 | 0.78 | 0.54 |
|  | SSVS | 0.53 | 0.44 | SSVS | 0.33 | 0.38 | 0.20 |
|  | GVS | 0.89 | 0.94 | GVS | 0.85 | 0.93 | 0.91 |
|  | MKMK | 0.65 | 0.60 | MKMK | 0.33 | 0.72 | 0.55 |
| $\mathrm{p}=3$ | KM | 0.75 | 0.87 | KM | 0.60 | 0.64 | 0.64 |
|  | SSVS | 0.52 | 0.39 | SSVS | 0.21 | 0.20 | 0.22 |
|  | GVS | 0.66 | 0.92 | GVS | 0.83 | 0.93 | 0.83 |
| MKMK |  | 0.60 | 0.75 | MKMK | 0.54 | 0.81 | 0.55 |

Table 2: Highest posterior model probabilities for correlated covariates: $\operatorname{corr}\left(\mathrm{X}_{2}, \mathrm{X}_{4}\right)=\operatorname{corr}\left(\mathrm{X}_{1}, \mathrm{X}_{5}\right)=0.3$

|  | $\mathbf{m}=\mathbf{2}$ |  |  |  |  |  |  |  |
| :--- | :--- | :--- | :--- | :--- | :---: | :---: | :---: | :---: |
|  | corr=0.7 |  |  |  |  |  |  |  |
|  | state | 1 | 2 | state | 1 | 2 | 3 |  |
|  | model var. | 4,5 | $2,3,4$ | model var. | 4,5 | $2,3,4$ | 2,4 |  |
| $\mathbf{p = 1}$ | KM | 0.68 | 0.74 | KM | 0.71 | 0.76 | 0.33 |  |
|  | SSVS | 0.36 | 0.31 | SSVS | 0.37 | 0.11 | 0.25 |  |
|  | GVS | 0.67 | 0.89 | GVS | 0.74 | 0.78 | 0.43 |  |
| $\mathbf{p = \mathbf { 2 }}$ | KM | 0.63 | 0.81 | KM | 0.36 | 0.60 | 0.70 |  |
|  | SSVS | 0.45 | 0.30 | SSVS | 0.28 | 0.10 | 0.11 |  |
|  | GVS | 0.77 | 0.88 | GVS | 0.86 | 0.87 | 0.83 |  |
|  | MKMK | 0.60 | 0.85 | MKMK | 0.40 | 0.58 | 0.71 |  |
| $\mathbf{p = 3}$ | KM | 0.41 | 0.59 | KM | 0.58 | 0.41 | 0.61 |  |
|  | SSVS | 0.36 | 0.23 | SSVS | 0.26 | 0.07 | 0.18 |  |
|  | GVS | 0.87 | 0.85 | GVS | 0.81 | 0.82 | 0.84 |  |
|  | MKMK | 0.55 | 0.61 | MKMK | 0.61 | 0.62 | 0.60 |  |

[^0]|  | $\mathbf{m}=\mathbf{2}$ |  |  |  |  |  |  |  |
| :--- | :--- | :---: | :---: | :--- | :---: | :---: | :---: | :---: |
|  | corr=0.9 |  |  |  |  |  |  |  |
|  | state | 1 | 2 | state | 1 | 2 | 3 |  |
|  | model var. | 4,5 | $2,3,4$ | model var. | 4,5 | $2,3,4$ | 2,4 |  |
| $\mathbf{p = 1}$ | KM | 0.65 | 0.25 | KM | 0.32 | 0.33 | 0.42 |  |
|  | SSVS | 0.18 | 0.07 | SSVS | 0.27 | 0.04 | 0.26 |  |
|  | GVS | 0.70 | 0.91 | GVS | 0.86 | 0.80 | 0.70 |  |
| $\mathbf{p = \mathbf { 2 }}$ | KM | 0.60 | 0.21 | KM | 0.35 | 0.31 | 0.34 |  |
|  | SSVS | 0.26 | 0.31 | SSVS | 0.24 | 0.19 | 0.20 |  |
|  | GVS | 0.87 | 0.89 | GVS | 0.83 | 0.61 | 0.91 |  |
|  | MKMK | 0.40 | 0.42 | MKMK | 0.19 | 0.38 | 0.27 |  |
| $\mathbf{p}=\mathbf{3}$ | KM | 0.36 | 0.46 | KM | 0.37 | 0.41 | 0.50 |  |
|  | SSVS | 0.24 | 0.18 | SSVS | 0.18 | 0.17 | 0.23 |  |
|  | GVS | 0.68 | 0.69 | GVS | 0.85 | 0.75 | 0.89 |  |
|  | MKMK | 0.43 | 0.35 | MKMK | 0.33 | 0.51 | 0.39 |  |

[^1]|  |  | independent |  |  |  | corr $=0.3$ |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | model var. | KM | SSVS | GVS | MKMK | KM | SSVS | GVS | MKMK |
| $\mathrm{m}=2$ |  |  |  |  |  |  |  |  |  |
| state 1 | 2,3 | 0.25 | 0.05 | 0.11 | 0.75 | 0.46 | 0.05 | 0.16 | 0.71 |
| state 2 | 2,3,4 | 0.19 | 0.06 | 0.27 | 0.85 | 0.30 | 0.04 | 0.26 | 0.88 |
| $\mathrm{m}=3$ |  |  |  |  |  |  |  |  |  |
| state 1 | 2,3 | 0.09 | 0.03 | 0.11 | 0.87 | 0.08 | 0.03 | 0.09 | 0.90 |
| state 2 | 2,3,4 | 0.08 | 0.04 | 0.10 | 0.79 | 0.10 | 0.04 | 0.26 | 0.95 |
| state 3 | 1,3 | 0.05 | 0.03 | 0.03 | 0.63 | 0.04 | 0.03 | 0.04 | 0.70 |
|  |  | corr $=0.7$ |  |  |  | corr $=0.9$ |  |  |  |
|  | model var. | KM | SSVS | GVS | MKMK | KM | SSVS | GVS | MKMK |
| $\mathrm{m}=2$ |  |  |  |  |  |  |  |  |  |
| state 1 | 2,3 | 0.41 | 0.04 | 0.18 | 0.74 | 0.49 | 0.04 | 0.08 | 0.68 |
| state 2 | 2,3,4 | 0.34 | 0.06 | 0.25 | 0.85 | 0.48 | 0.04 | 0.17 | 0.67 |
| $\mathrm{m}=3$ |  |  |  |  |  |  |  |  |  |
| state 1 | 2,3 | 0.07 ${ }^{(*)}$ | $0.04{ }^{(*)}$ | 0.13 | 0.64 | 0.02 ${ }^{(*)}$ | $0.02^{(*)}$ | 0.05 | 0.27 |
| state 2 | 2,3,4 | 0.08 ${ }^{(*)}$ | $0.04{ }^{(*)}$ | 0.12 | 0.44 | 0.01 ${ }^{(*)}$ | $0.04{ }^{(*)}$ | 0.09 | 0.81 |
| state 3 | 1,3 | 0.05 ${ }^{(*)}$ | $0.03{ }^{(*)}$ | 0.10 | 0.70 | $0.03{ }^{(*)}$ | $0.02^{(*)}$ | 0.02 | 0.20 |

Table 5: Highest posterior model probabilities for $\operatorname{NHHMM}(\mathrm{m})$ with independent
or correlated covariates $-\left(^{*}\right)=$ not the highest probability


Figure 1: Ergodic highest posterior model probabilities for the simulated datasets of $\operatorname{MSARX}(2,1)$, in state 1, for KM (black line); SSVS (grey line); GVS (bold line); MKMK(double black line)

(independent covariates)

(covariates with corr $=0.7$ )

(covariates with corr=0.3)

(covariates with corr $=0.9$ )

Figure 2: Ergodic highest posterior model probabilities for the simulated datasets of NHHMM(2),
in state 1, for KM (black line); SSVS (grey line); GVS (bold line); MKMK (double black line)


[^0]:    Table 3: Highest posterior model probabilities for correlated covariates: $\operatorname{corr}\left(\mathrm{X}_{2}, \mathrm{X}_{4}\right)=\operatorname{corr}\left(\mathrm{X}_{1}, \mathrm{X}_{5}\right)=0.7$

[^1]:    Table 4: Highest posterior model probabilities for correlated covariates: $\operatorname{corr}\left(\mathrm{X}_{2}, \mathrm{X}_{4}\right)=\operatorname{corr}\left(\mathrm{X}_{1}, \mathrm{X}_{5}\right)=0.9$

