# The Value of Structural Information in the VAR Model 

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#### Abstract

Economic policy decisions are often informed by empirical economic analysis. While the decision-maker is usually only interested in good estimates of outcomes, the analyst is interested in estimating the model. Accurate inference on the structural features of a model, such as cointegration, can improve policy analysis as it can improve estimation, inference and forecast efficiency from using that model. However, using a model does not guarantee good estimates of the object of interest and, as it assigns a probability of one to a model and zero to near-by models, takes extreme zero-one account of the 'weight of evidence' in the data and the researcher's uncertainty. By using the uncertainty associated with the structural features in a model set, one


obtains policy analysis that is not conditional on the structure of the model and can improve efficiency if the features are appropriately weighted. In this paper tools are presented to allow for unconditional inference on the vector autoregressive (VAR) model. In particular, we employ measures on manifolds to elicit priors on subspaces defined by particular features of the VAR model. The features considered are cointegration, exogeneity, deterministic processes and overidentification. Two applications - money demand in Australia, and a macroeconomic model of the UK proposed by Garratt, Lee, Persaran, and Shin (2002) are used to illustrate the feasibility of the proposed methods.

Key Words: Posterior probabilities; Laplace approximation; Structural modelling; Cointegration; Exogeneity; Model averaging.
JEL Codes: C11, C32, C52

## 1 Introduction.

An important function of empirical analysis is to provide information for decision making. This information is generally provided in the form of estimates of objects of interest such as forecasts of endogenous variables, effects of shocks measured by impulse response functions, probabilities, elasticities or distributions. In many cases, the decision maker is not directly interested in the underlying model used to produce such estimates, however, it is in the analyst's interest to detail how the results she provides rely upon the model. That is, the analyst, when providing the estimates of the objects of interest, must point out "This assumes that ..." Such restrictions upon the interpretation of the results do not aid the decision-maker in their task.

It is generally accepted, however, that to improve policy analysis it is important to have accurate inference on the support for the alternative models considered or to have such inference on the structural features of an encompassing model. As such, much effort is expended in investigating the empirical support for various economically and statistically plausible features. (If we condition upon particular features that are well supported by the data we can obtain efficiency gains in estimating parameters, in inference and in producing forecasts.) Examples of features of models that are of interest to analysts - but not necessarily decision-makers - include numbers of long run
relationships among variables, forms of these long run relationships, persistent and predictable long run behaviours of variables, short term behaviours, and the dimension of the system in variables or in parameters required for the problem of interest. Each of these features implies zero restrictions on particular parameters in a general model. If these features are supported by the data - and so are credible in the sense of Sims (1980) - and if they hold outside the sample, then imposing them can improve forecasts and inference, and hence policy suggestions. Unfortunately, the support in the data is often not clear or dogmatically for or against the restriction, and the researcher does not have strong prior belief in the restriction. It is, however, common to condition upon such features, effectively assigning a weight of one to the model implied by the restrictions being true and zero to all other plausible models. Even if the support is strongly for or against a particular restriction, with only slight support for the alternative unrestricted model, imposing the restriction ignores information from that less likely model which, if appropriately weighted, could improve forecasts.

There is therefore a conflict between the analyst's need to obtain the best model and the decision-maker's need for the least restrictive interpretation of the information provided by the analyst. As an alternative to conditioning on structural features, it is possible to improve policy analysis by presenting unconditional or averaged information. Gains in forecasting accuracy by simple averaging have been pioneered by Bates and Granger (1969) and discussed recently by Diebold and Lopez (1996), Newbold and Harvey (2001) and Terui and van Dijk (2002). Some explanation for this phenomenon in particular cases was provided by Hendry and Clements (2002). Alternatively, the weights can be determined to reflect the support for the model from which each estimate derives. This requires accurately reflecting the uncertainty associated with the structural features defining the model.

In this paper we present an approach for conducting unconditional inference on structural features of the cointegrating vector autoregressive model. We regard the restrictions on the general model implied by the structural features as producing a new model for comparison. The results will still be conditional upon the model set, but if this set covers a wide enough range of models, possibly those the analyst would have searched within otherwise, we see this as an improvement over conditional analysis. We work with models that nest within an encompassing model, however this is not a requirement as we take a Bayesian approach. We consider the joint probabilities of cointegration, overidentification, deterministic processes, and exogeneity. From
relationships among manifolds and orthogonal groups and their measures, we elicit measures on relevant subspaces of the parameter space. From these measures we develop prior distributions for elements of these subspaces as the parameter of interest. Thus we choose prior specification for models directly rather than on parameters that are subsequently restricted. Further, by enabling the expression of prior beliefs on parameters of interest, rather than on the instruments via which we obtain inference on that parameter of interest, we present a more coherent method of investigation.

The aim of this paper is to obtain unconditional policy analysis by which we mean we wish to obtain inference, estimates and forecasts from model averages in which the economically and econometrically important structural features may have weights other than zero or one. Examples of impulse responses are produced that derive from the unconditional, but correctly weighted model space.

The structure of the paper is as follows. In the Section 2 we introduce the general model of interest in this paper - the vector autoregressive model, the general structural features of interest, and the restrictions they imply. We demonstrate the approach with two applications: a model of Australian money demand; and a macroeconomic model of the UK economy proposed by Garratt, Lee, Persaran, and Shin (2002). These applications with the implied restrictions are outlined in Section 3. In Section 4 we present the priors we will be considering in the paper, the likelihood and a general expression for the posterior. The tools for inference in this paper, the Bayes factor and posterior probabilities, are introduced and expressions derived for the specific features of interest - impulse responses - in Section 5. The results of the application are presented in Section 6 and Section 7 concludes.

## 2 The Vector Autoregressive Model.

We work with the vector autoregressive model in the error correction form to simplify expressions of restrictions. The error correction model (ECM) of the $1 \times n$ vector time series process $y_{t}, t=1, \ldots, T$, conditioning on the $l$ observations $t=-l+1, \ldots, 0$, is

$$
\begin{align*}
\Delta y_{t} & =y_{t-1} \beta^{+} \alpha+d_{t} \mu+\Delta y_{t-1} \Gamma_{1}+\ldots+\Delta y_{t-l} \Gamma_{l}+\varepsilon_{t}  \tag{1}\\
& =y_{t-1} \beta^{+} \alpha+d_{t} \mu_{1} \alpha+d_{t} \mu_{2} \alpha_{\perp}+\Delta y_{t-1} \Gamma_{1}+\ldots+\Delta y_{t-l} \Gamma_{l}+\varepsilon_{t} \\
& =z_{1, t} \beta \alpha+z_{2, t} \Phi+\varepsilon_{t} \tag{2}
\end{align*}
$$

where $\Delta y_{t}=y_{t}-y_{t-1}, z_{1, t}=\left(d_{t}, y_{t-1}\right), z_{2, t}=\left(d_{t}, \Delta y_{t-1}, \ldots, \Delta y_{t-l}\right), \Phi=$ $\left(\alpha_{\perp}^{\prime} \mu_{2}^{\prime}, \Gamma_{1}^{\prime}, \ldots, \Gamma_{l}^{\prime}\right)^{\prime}$ and $\beta=\left(\mu_{1}^{\prime}, \beta^{+\prime}\right)^{\prime}$. The matrices $\beta^{+}$and $\alpha^{\prime}$ are $n \times r$ and assumed to have rank $r$, and if $r=n$ then $\beta^{+}=I_{n}$.

The following subsections define the restrictions of interest, combinations of which define different model features of interest which we may compare or weight using posterior probabilities.

As we consider a wide range of models in this paper, we will use a consistent notation to index each model to identify the cointegrating rank, the identifying restrictions, the form of exogeneity, and the deterministic processes in the model. We will denote the rank of a model by $r$, where $r=0,1, \ldots, n$. The particular identifying restrictions placed upon $\beta$ will be denoted by $o$, where $o=1, \ldots, \mathfrak{J}$ and $o=1$ will be understood to refer to the just identified model. Partitioning $y_{t}$ as $y_{t}=\left(y_{1, t} y_{2, t}\right)$ where $y_{1, t}$ is a $1 \times n_{1}$ vector, $n_{1} \geq r$, exogeneity of $y_{2, t}$ will be considered with respect to subsets of the parameters in the equation for $y_{1, t}$, where will use $\phi_{1}$ and $\phi_{2}$ to denote these subsets. The particular form of exogeneity restrictions in the model will be denoted by $e$, where $e=1, \ldots, 5$ and these refer respectively, to the model in which $y_{2, t}$ is: not exogenous with respect to $\phi_{1}$ or $\phi_{2}$; weakly exogenous with respect to $\phi_{1}$; strongly exogenous with respect to $\phi_{1}$; weakly exogenous with respect to $\phi_{2}$; and, strongly exogenous with respect to $\phi_{2}$. Finally, the particular form of deterministic processes will be denoted by $i$, where $i=1, \ldots, 5$ and these refer to the five models detailed in the subsection on deterministic processes below.

The vector identifying a particular model will therefore be $\omega=(r, o, e, i)$. For example, the least restricted model will be ( $n, 1,1,1$ ), while the most restricted model will be $(0, o, 5,5)$. Note that, as will become clear, there may be no sensible order to the models with $o>1$ by degree of restriction, and models with exogeneity $e=3$ and $e=4$ cannot be placed in a sensible order with respect to eachother. The models will be identified as $M_{\omega}$. When we are considering only a particular feature such as exogeneity, we will indicate this by referring to the model as $M_{(., ., e,)}$, and if we are conditioning upon a particular feature, such as rank, $M_{(e \mid r)}$. Where we have averaged across or marginalised with respect to the other features, we will indicate this by $M_{(r)}$, and the marginal likelihood for a model will be $m_{\omega}$.

Finally, we introduce the following terms to simplify the expressions in the posteriors. Let $\widetilde{z}_{t}=\left(\begin{array}{ll}z_{1, t} \beta & z_{2, t}\end{array}\right)$, and the $\left(r+k_{i}\right) \times n$ matrix $B=\left[\begin{array}{ll}\alpha^{\prime} & \Phi^{\prime}\end{array}\right]^{\prime}$.

We may now write the model as

$$
\begin{equation*}
\Delta y_{t}=\widetilde{z}_{t} B+\varepsilon \tag{3}
\end{equation*}
$$

### 2.1 Structural features

Within the model (1), a number of structural features are commonly of interest to economists and or econometricians. Here we detail five of these and the restrictions they imply for (1). To demonstrate we use a simple, and reasonably well understood example: money demand. The variables, all of which appear in logarithmic form, are defined as $y_{t}=\left(\begin{array}{ll}m_{t} & i n c_{t}\end{array}\right)$, where $m_{t}$ is the $\log$ measure of real money and $i n c_{t}$ is the measure of real income. The bivariate VAR has the form $y_{t}=d_{t} \mu+y_{t-1} \Pi_{1}+\ldots+y_{t-l-1} \Pi_{l+1}+\varepsilon_{t}$. Suppose we are interested in the one step ahead forecast of $m_{t}$ or the overall response path of $m_{t}$ to a shock in $i n c_{t-h}$ for $h=0,1,2, \ldots$.We are interested in estimation of the parameters determining the long and short run behaviour of $m_{t}$ and in forecasts of $m_{t}$, where the forecasts may also be over the long run, or both the long run and short run. Here we regard the long run as the equilibrium relationships to which the elements of $y_{t}$ would revert if all future errors were zero.

### 2.1.1 Cointegration

As has been observed in empirical studies, many economic variables of interest are not stationary, yet economic theory, or empirical evidence, suggests stable long run relationships exist between these variables. The statistical theory of cointegration (Granger, 1983, and Engle and Granger, 1987), in which a set of nonstationary variables combine linearly to form stationary relationships, and the attendant Granger's representation theorem provide a useful specification to incorporate this economic behaviour into the error correction model and allows the separation of long run and short run behaviour. For cointegration analysis of (1), of interest is the coefficient matrix $\beta^{+}$ (and $\alpha$ ) which are of rank $r \leq n$. Of particular interest then, is $r$ as $(n-r)$ is the number of common stochastic trends in $y_{t}$, and $r$ is the number of $I(0)$ combinations of the element of $y_{t}$ extant. In the case $r<n$ and assuming for now $\mu_{1}=0, \beta^{+}$is the matrix of cointegration coefficients, $y_{t} \beta^{+}$are the stationary relations towards which the elements of $y_{t}$ are attracted, and $\alpha$ is the matrix of factor loading coefficients or adjustment coefficients determining the rate of adjustment of $y_{t}$ towards $y_{t} \beta^{+}$.

In the money demand example, $r \in[0,1,2]$. It is common to regard the money demand relation as the cointegrating relation between the integrated variables in $y_{t} \sim I(1)$, and supply is exogenous (see for example Johansen, 1995 and Funke, Hall and Beeby, 1997). That is, $\zeta_{t}=\beta_{1} m_{t}+\beta_{2} i n c_{t}=$ $y_{t} \beta^{+} \sim I(0), E\left(\zeta_{t}\right)=d_{t} \mu_{1}$ and possibly $\mu_{1} \neq 0$. Therefore, for the analysis to make sense, we require that cointegration should hold (and so $r=1$ ). In this case we would have the error correction representation for $y_{t}$ as

$$
\Delta y_{t}=y_{t-1} \beta^{+} \alpha+d_{t} \mu+\Delta y_{t-1} \Gamma_{1}+\ldots+\Delta y_{t-l} \Gamma_{l}+\varepsilon_{t}
$$

where $\beta^{+}=\left(\beta_{1}, \beta_{2}\right)^{\prime}$ and $\alpha=\left(\alpha_{1}, \alpha_{2}\right)$.

### 2.1.2 Exogeneity

As it is usually accepted in econometrics that there are benefits from parsimony, another important issue is the dimension of the system to be estimated in terms of the number of equations. Recall the partition $y_{t}=\left(y_{1, t} y_{2, t}\right)$. If the set of variables in $y_{2, t}$ can be treated as exogenous for inferential purposes, a partial system may be estimated in which no equations are estimated for these variables. This is essentially ignoring information that contributes nothing to the inference. As an example, it is not uncommon to assume that to estimate the income elasticity of money, a researcher would be interested in whether an equation for income need be estimated, or could this analysis be done with a single equation.

Under the condition of cointegration, the representation of the model in (1) will be useful for the analysis of exogeneity. Partition $\alpha=\left(\begin{array}{ll}\alpha_{1} & \alpha_{2}\end{array}\right)$ conformably with the dimensions of $y_{1, t}$ and $y_{2, t}$. In this article we consider weak exogeneity of $y_{2, t}$ with respect to the parameters influencing long run behaviour of $y_{1, t}, \phi_{1}=\left(\operatorname{vec}(\beta)^{\prime}, \operatorname{vec}\left(\alpha_{1}\right)^{\prime}\right)^{\prime}$. If our interest is in estimating or conducting inference on the subset of parameters $\phi_{1}$, it may not be necessary to estimate the full set of $n$ equations for $y_{t}$. That is, conditions may exist which allow us to condition on the variables $y_{2, t}$ and therefore only model the equations for $y_{1, t}$. This condition is that $y_{2, t}$ be weakly exogenous with respect to $\phi_{1}$. As shown in Urbain (1992) and Johansen (1992) inter alia, $y_{2, t}$ will be weakly exogenous with respect to $\phi_{1}$ if $\alpha_{2}=0$. To preserve the rank of $\alpha$ requires that $n_{1} \geq r$, which implies we cannot have more than $n-r$ variables weakly exogenous with respect to $\phi_{1}$. An important model in the
literature which relies upon this assumption is the triangular model (Phillips, 1991) used by Phillips (1994) in which $n_{1}=r$.

For a given cointegrating rank $r$, denote by $M_{(e \mid r)}$ the various models of exogeneity. The model with no exogeneity restrictions imposed is $e=1$ and the model with weak exogeneity of $y_{2, t}$ with respect to $\phi_{1}$ is $e=2$. Other forms of exogeneity include: strong exogeneity of $y_{2, t}$ with respect to the parameters influencing long run behaviour of $y_{1, t}, \phi_{1}(e=3)$; weak exogeneity of $y_{2, t}$ with respect to the parameters influencing long and short run behaviour of $y_{1, t}\left(\phi_{2}=\left(\phi_{1}^{\prime}, \operatorname{vec}\left(\Gamma_{11}\right)^{\prime}, \operatorname{vec}\left(\Gamma_{21}\right)^{\prime}\right)^{\prime}\right)(e=4) ;$ and strong exogeneity of $y_{2, t}$ with respect to the parameters influencing long and short run behaviour of $y_{1, t}(e=5)$. These imply further restrictions upon the parameters in (1) such as Granger noncausality, however we do not explore them here as the first case is sufficient to demonstrate the approach.

If we are interested in whether we may estimate the money demand equation $\zeta_{t}$ (and so estimate $\beta^{+}$) from a single equation for $m_{t}$, then this would require that the variables $i n c_{t}$ be weakly exogenous with respect to $\beta^{+}(e=2)$. $\alpha_{1}$ is the adjustment coefficient in the equation for $\Delta m_{t}$ and $\alpha_{2}$ is the same in the equation for $i n c_{t}$ such that these parameters determine the response in $y_{t}$ to a nonzero value of $\zeta_{t-1}$. Weak exogeneity of inc $c_{t}$ with respect to $\beta^{+}$ implies $\alpha_{2}=0$.

### 2.1.3 Overidentifying restrictions on the cointegrating vectors

As discussed in Garratt et al. (2002), when modelling economic systems, economic theory tends be more useful when it focuses upon the form of long run, or equilibrium, relationships between variables and leaves the short run relations unrestricted (see Sims 1980 for discussion about the dangers of imposing incredible restrictions on short run dynamics). This leads us to the consideration of the direction of the cointegrating space or the form of the cointegrating relations and to what valid linear restrictions can be imposed on $\beta$, as representing the long run relations. For money demand, the stability (in the sense that velocity is $I(0)$ but may have deterministic trends - we discuss this latter possibility in the following subsection) of the (log of the) inverse velocity of money, $\nu_{t}=m_{t}-i n c_{t}$ is an important issue for econometric analysis. Thus it would be sensible to allow this to be a long run relation such that $\zeta_{t}=\nu_{t}$ is another direction in the model set to be considered.

In both the classical and Bayesian approaches, to test the appropriateness of such restrictions and to estimate the restricted model, requires a specifi-
cation of the model subject to these restrictions. In the classical maximum likelihood approach, Johansen (1995) has provided methods for estimation with, and testing of, these restrictions. The three restrictions commonly investigated are presented in Johansen (1995, Chapter 5) as the following hypotheses.

$$
\begin{array}{ll}
(o=1) & \text { No restrictions upon } \beta . \\
(o=2) & H_{0}: \beta=H \psi
\end{array}
$$

where the dimensions of the respective matrices are: $H n \times s, \psi s \times r, r \leq s$. $(o=3) \quad H_{0}: \beta=(b \varphi)=(b b \perp \psi)$
where the dimensions of the respective matrices are: $b n \times s, b_{\perp} n \times(n-s)$, $\psi(n-s) \times(r-s), s \leq r$.
$(o=4) \quad H_{0}: \beta=\left(H_{1} \psi_{1}, H_{2} \psi_{2}, \ldots, H_{l} \psi_{l}\right)$
where the dimensions of the respective matrices are: $H_{i} n \times s_{i}, \psi_{i}$ is $s_{i} \times r_{i}$, $r_{i} \leq s_{i}, l \leq r, \sum_{i} r_{i}=r$.

The restriction in $o=1$ imposes no restriction on the space of $\beta$, in $o=2$ the cointegrating space is completely determined. The third restriction, $o=3$, restricts the cointegrating space to pass through a known vector or set of $s$ vectors, $b$, and the remaining $r-s$ vectors, $b_{\perp} \psi$, are unknown except that they are orthogonal to $b$, such that the space of $\beta$ is not completely known. The final hypothesis, $\mathfrak{J}=o=4$, generalizes the first two.

### 2.1.4 Deterministic terms

Economists are commonly interested in the presence or absence of deterministic processes in $y_{t}$ or $y_{t} \beta^{+}$. For both statistical and economic reasons, the persistent and predictable, or deterministic, component economic behaviour is important. Of interest are questions such as whether linear or quadratic drifts are present in $y_{t}$ and whether nonzero constant terms and deterministic trends are present in $y_{t} \beta^{+}$. For example, the velocity of money in many countries has not remained stable over the long run. For extended periods it has displayed what appears to be a clear trend. If we were to assume the velocity was an equilibrium or long run relation of interest, it would be important to allow for some trend in this relation. It is well known, however, that simplistic treatment of the deterministic terms by testing whether $\mu$ or some elements of $\mu$ are zero leads to the strange and unsatisfactory situation that very different trending behaviour is implied in the levels of the process for differing values of $r$. Therefore $\mu$ is decomposed into $\mu=\mu_{1} \alpha+\mu_{2} \alpha_{\perp}$ where $\mu_{1}=\mu \alpha^{\prime}\left(\alpha \alpha^{\prime}\right)^{-1}$ and $\mu_{2}=\mu \alpha_{\perp}^{\prime}\left(\alpha_{\perp} \alpha_{\perp}^{\prime}\right)^{-1}$ such that $\mu_{1}$ represents the
deterministic processes associated with $y_{t} \beta^{+}$and $\mu_{2}$ represents those for $y_{t}$ (see Johansen, 1995 Section 5.7 for further discussion).

Assuming $d_{t}=(1, t)$, then for each $j=1,2, d_{t} \mu_{j}=\mu_{j, \iota}+t \mu_{j, \delta}$. Although a wider range of models are clearly available, the five most commonly considered may be stated as follows, where $M_{r, i}$ is the $i^{\text {th }}$ model of deterministic terms at given rank $r$ :

$$
\begin{array}{ll}
M_{r, 1}: & d_{t} \mu=\mu_{1, \alpha} \alpha+\mu_{2, \iota} \alpha_{\perp}+\left(\mu_{1, \delta} \alpha+\mu_{2, \delta} \alpha_{\perp}\right) t \\
M_{r, 2} & : \\
d_{t} \mu=\mu_{1, \iota} \alpha+\mu_{2, \iota} \alpha_{\perp}+\mu_{1, \delta} \alpha t \\
M_{r, 3} & : \\
M_{r, 4} & : \\
d_{t} \mu=\mu_{1, \iota} \alpha+\mu_{2, \iota} \alpha_{\perp} \\
M_{r, 5} & : \\
d_{t} \mu=0
\end{array}
$$

## 3 Empirical example: The Garratt, Lee, Pesaran and Shin (2002) structural VAR model of the UK economy.

Garratt, Lee, Pesaran, and Shin (2002) provide an extensive model of the UK economy which focuses upon the long run relations, but incorporates useful short run restrictions to improve modelling. In their paper, Garratt et al. highlight two differences in their approach from other large models. First it is developed for a small open economy, and second it takes a new and practical approach to incorporating long run relations while leaving short run relations largely unrestricted. The variables in the econometric model are

$$
y_{t}=\left(r_{t}, w_{t}, \Delta p_{t}, p_{t}-p_{t}^{*}, e_{t}, h_{t}-w_{t}, r_{t}^{*}, w_{t}^{*}, p_{t}^{o}\right)
$$

where, in logarithms, $p_{t}^{o}$ is the price of oil, $w_{t}$ is UK real per capita GDP and $w_{t}^{*}$ is the foreign (OECD) real per capita GDP, $p_{t}$ is the UK producer price index, $p_{t}^{*}$ is foreign (OECD) producer prices, $e_{t}$ is the nominal Sterling effective exchange rate, $h_{t}$ UK real per capita M0 money stock, $r_{t}=0.25 \ln \left(1+R_{t} / 100\right)$ where $R_{t}$ is a function of 90 day interest rates and $r_{t}^{*}$ a similar function of the US, Germany, Japan and France 90 day rates.

The long run relations which form the cointegrating relations, subject to all restrictions finally imposed as a result of the analysis by Garratt, et al.,
are

$$
\begin{aligned}
p_{t}-p_{t}^{*}-e_{t} & =u_{1, t} \\
r_{t}-r_{t}^{*} & =u_{2, t} \\
w_{t}-w_{t}^{*} & =u_{3, t} \\
r_{t}-\Delta p_{t} & =u_{4, t} \text { and } \\
\beta_{32}\left(h_{t}-w_{t}\right) & =\mu_{11, \delta} t+\beta_{22} r_{t}+u_{5, t}
\end{aligned}
$$

where the $u_{i, t}$ are $I(0)$ with unrestricted means. Assuming the rank $r=5$, these results suggest a cointegrating space spanning the space of the matrix $\beta=\left(\begin{array}{ll}H_{1} & \beta_{2}^{c}\end{array}\right)$ where

$$
\begin{aligned}
H_{1}^{\prime} & =\left[\begin{array}{llllllllll}
0 & 0 & 0 & 0 & 1 & -1 & 0 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 & 0 & 0 & 0 & -1 & 0 & 0 \\
0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 & -1 & 0 \\
0 & 1 & 0 & -1 & 0 & 0 & 0 & 0 & 0 & 0
\end{array}\right], \\
\beta_{2}^{c} & =H_{2} \varphi \\
H_{2}^{\prime} & =\left[\begin{array}{cccccccccc}
-1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & -1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 1 & 0 & 0 & 0
\end{array}\right] \\
\varphi^{\prime} & =\left(\begin{array}{llll}
\mu_{11, \delta} & \beta_{22} & \beta_{32}
\end{array}\right) .
\end{aligned}
$$

There are three parameters ${ }^{1}$ to be estimated in $\beta$. In their paper, Garratt et al. make oil prices strictly exogenous with respect to the rest of the system ${ }^{2}$. The parameterisation they use implies weak exogeneity of oil prices with respect to $\alpha$ and $\beta$. The restriction that there is no quadratic trend

[^0]in $y_{t}$ implies $\mu_{11, \delta} \alpha=0$. Further, the exclusion of a trend from all long run relations except the money-income relation, $u_{5, t}$, implies the restriction upon the first row of $\beta$ is $\mu_{1, \delta}=\left(0,0,0,0, \mu_{11, \delta}\right)$.

The combinations of restrictions implied by the above model can be denoted in the notation of Section 2 as $M_{\omega}$ with $\omega=(5,4,2,2)$, that is, the cointegrating rank is 5 , we employ the overidentifying restrictions on $\beta$ of type 4 , oil prices are weakly exogenous with respect to $\alpha$ and $\beta$, and there is no quadratic drift in $y_{t}$ but there may be a trend in $y_{t} \beta$.

The range of models we include in our model set are defined by $r \in$ $[0,1, \ldots, 9], e \in[1,2], o \in[1,4]$, and $i \in[1, \ldots, 5]$ for a total of 200 models. As a number of the models implied by combinations of these restrictions are either impossible or observationally equivalent, we need only estimate 87 models. We provide only a preliminary analysis which is not intended to be an alternative to the more complete classical analysis of Garratt, et al. A number of issues dealt with in their full classical analysis such as, for example, lag length determination are not taken into account in our study as they are beyond the scope of this paper.

## 4 Priors and posteriors.

In this section the forms of the priors and resultant posterior are presented. We restrict ourselves to flat priors where possible, although consideration is given to informative priors when discussing the parameters of interest. For the model in (3), assume the rows of the $T \times n$ matrix $\varepsilon=\left(\varepsilon_{1}^{\prime}, \varepsilon_{2}^{\prime}, \ldots, \varepsilon_{T}^{\prime}\right)^{\prime}$ are $\varepsilon_{t} \sim \operatorname{iid} N(0, \Sigma)$. The likelihood can then be written as

$$
\begin{equation*}
L(y \mid \Sigma, B, \beta, \omega, \widetilde{Z}) \propto|\Sigma|^{-\frac{T}{2}} \exp \left\{-\frac{1}{2} \operatorname{tr}\left(\Sigma^{-1} \varepsilon^{\prime} \varepsilon\right)\right\} \tag{4}
\end{equation*}
$$

### 4.1 The prior for $(\Sigma, B, \omega)$.

The priors for the elements of $\omega=(r, o, e, i)$ are not independent, as certain combinations are either impossible, meaningless (such as, for example, $r=0$ with $o=2$ ) or observationally equivalent to another combination (such as, for example, $r=0$ with $o=2$ or $r=n$ with $i=1$ or 2 ). However, after excluding these combinations we specify the remaining values of $\omega$ to be equally likely. This implies we use the prior for the rank $r$ as $p(r)=$ $(n+1)^{-1}$ and for the deterministic models $p(i \mid r)=1 / 5$ for $0<r<n$ and
$i \in[1,2,3,4,5], p(i \mid r=0)=1 / 3$ for $i \in[1,2,4]$ since at $r=0 i=2$ and $i=3$ are observationally equivalent as are $i=4$ and $i=5$, and $p(i \mid r=n)=1 / 3$ for $i \in[1,3,5]$ since at $r=n, i=1$ and $i=2$ are observationally equivalent as are $i=3$ and $i=4$. As oil prices are weakly exogenous with respect to $\phi_{1}$, we set $y_{2, t}=p_{t}^{o}$ and $y_{1, t}$ is the vector of remaining variables. The prior density for the states of exogeneity $e \in[1,2]$ is $p(e \mid r)=1 / 2$ for $r<n$ and for the states of overidentification of $\beta, o \in[1,4], p(o)=1 / 2$. The standard diffuse prior for $\Sigma, p(\Sigma) \propto|\Sigma|^{-(n+1) / 2}$ is used.

As $B$ changes dimensions across the different models of $\omega$ and each element of the matrix $B$ has the real line as its support, the Bayes factors for different models will not be well defined if an improper prior on $B$, such as $p(B \mid \beta, \omega) \propto 1$ were used. For discussion on this point see (among many others) Lindley (1957), Bartlett (1957), Jeffreys (1961) and more recently O'Hagan (1995). For this reason a weakly informative proper prior for $B$ must be used. We take the prior for $B$ conditional upon $(\Sigma, \beta, \omega)$ as normal with zero mean and covariance $\Sigma \quad\left(\widetilde{\beta}^{\prime} \underline{H} \widetilde{\beta}\right)^{-1}$ where $\underline{H}=0.1 I_{\left(r+k_{i}\right)}$ and

$$
\widetilde{\beta}=\left[\begin{array}{cc}
\beta & 0 \\
0 & I_{k_{i}}
\end{array}\right]
$$

such that $\widetilde{\beta}^{\prime} \underline{H} \widetilde{\beta}=0.1 I_{\left(r+k_{i}\right)}$.

### 4.2 Eliciting a prior on $\beta$.

In this section we outline earlier work in Bayesian cointegration analysis, focussing on problems addressed and limitations of these approaches as they relate to the aim of this paper. Then we present a general analysis of an alternative approach. For specific applications and a less technical outline of this approach we refer to Strachan and Inder (2003) and Strachan and van Dijk (2003).

Linear restrictions and the cointegrating space: It is well known that as $\beta$ and $\alpha$ appear as a product in (2), $r^{2}$ restrictions need to be imposed on the elements of $\beta$ and $\alpha$ to just identify these elements. These restrictions are commonly imposed upon $\beta$ by assuming $c \beta$ is invertible for known $(r \times n)$ matrix $c$ and the restricted $\beta$ to be estimated is $\bar{\beta}=\beta(c \beta)^{-1}$. The free elements are collected in $\bar{\beta}_{2}=c_{\perp} \bar{\beta}$ where $c_{\perp} d^{\prime}=0$. A common choice in theoretical work is $c=\left[\begin{array}{ll}I_{r} & 0\end{array}\right]$ such that $\bar{\beta}=\left[\begin{array}{ll}I_{r} & \bar{\beta}_{2}^{\prime}\end{array}\right]^{\prime}$. A prior is then specified
for $\bar{\beta}_{2}$ which is then estimated and often its value is interpreted.
There exist practical problems with incorrectly selecting $c$. The implications for classical analysis of this issue are discussed in Boswijk (1996) and Luukkonen, Ripatti and Saikkonen (1999) and in Bayesian analysis by Strachan (2003). In each of these papers examples are provided which demonstrate the importance of correctly determining $c$.

Assuming that $c$ is known, Kleibergen and van Dijk $(1994,1998)$ and Bauwens and Lubrano (1996) detail remaining pathologies and features which complicate analysis associated with the posterior for $\bar{\beta}_{2}$ with a flat prior. Kleibergen and van Dijk (1994) demonstrate how a variable addition specification - which would provide a natural way of performing inference on $r$ by nesting the reduced rank model within a full rank model - results in an improper posterior distribution at reduced ranks, thus precluding inference. For the non-nested reduced rank model, as in (2), Kleibergen and van Dijk (1994) outline the additional issue of local nonidentification which manifests itself in the likelihood and results in asymptotes in the marginal posterior distributions, nonexistence of moments of $\bar{\beta}_{2}$, and precludes the use of MCMC due to reducibility of the chain. As a solution they propose using the Jeffreys prior as the behaviour of this prior in problem areas of the support offsets the problematic behaviour of the likelihood. Kleibergen and van Dijk (1998) and Kleibergen and Paap (2002) use a singular value decomposition to nest the rank $r<n$ model within the rank $n$ model. Importantly, they include in the posterior the Jacobian for the transformation from the full rank model to the parameters of the reduced rank model into the posterior. In this specification, the Jacobian behaves in a similar way to the Jeffreys prior in the problem areas of the support, however this approach allows freer expression of prior beliefs than the Jeffreys prior. Use of the Jeffreys prior or the singular value decomposition avoid the issue of local nonidentification, result in proper posteriors and allow use of MCMC, however the posterior again has no moments of $\bar{\beta}_{2}$.

Bauwens and Lubrano (1996) begin with the reduced rank model and provide a study of the posterior distribution of $\bar{\beta}_{2}$. They use the results for the $1-1$ poly - $t$ density of Drèze (1976) to show the posterior has no moments due to a deficiency of degrees of freedom. Similar results have been shown for the simultaneous equations model (Drèze 1976, Kleibergen and van Dijk 1998). Nonexistence of moments is not commonly a concern for estimation as modal estimates exist as alternative estimates of location. However, as the kernel of the 1-1 poly $-t$ is a ratio of the kernels of two student $-t$ densities,
the posterior may be bimodal - with the modes sometimes well apart from eachother - making it difficult to both locate the global mode and bringing into question the interpretation of the mode as a measure of location.

Exogeneity is a commonly employed restriction and is important in our application. For our application in which we combine restrictions to define new models, we have the additional problem that the posterior for $\bar{\beta}_{2}$ is improper when exogeneity is imposed. As there are no published references to this results an outline of the result is provided in Appendix 1. Nonexistence of moments or an improper posterior are significant issues as they imply we know a priori any estimate of an object of interest, $g\left(\bar{\beta}_{2}\right)$ - obtained by averaging across the set of models - will not exist (or be infinite) if exogeneity is imposed or if $g\left(\bar{\beta}_{2}\right)$ is a convex or linear function of $\bar{\beta}_{2}$.

Further, it is clear from the discussion on the prior for $B$ that a flat prior on $\bar{\beta}_{2}$ cannot be employed to obtain posterior probabilities for $\omega$, since the dimensions of $\bar{\beta}_{2}$ depend upon $\omega$. As argued in the introduction, an advantage of the Bayesian approach is the ability to explicitly incorporate prior beliefs into the analysis. A flat improper prior is generally intended to reflect ignorance about the parameter of interest, therefore the above issues with the posterior at least, may be resolved by relinquishing this option and making use of an informative prior on $\bar{\beta}_{2}$. For example, a student-t prior may be used, or inequality restrictions - such as a marginal propensity to consume between zero and one - are often useful. Priors such as the Jeffreys prior have been proposed which may resolve some of the above problems, however their application is often complicated and one is suspicious that such priors are advocated more as a fix to a problem in the likelihood and less as a representation of prior beliefs. Further the Jeffreys prior does not allow for model averaging. Therefore, to preserve the options of both informative and uninformative priors, to preserve the function of the prior as a representation of prior beliefs, to simplify the application and estimation, and as we do not see $\bar{\beta}_{2}$ as the parameter of interest, we diverge at this point from much of the earlier literature in both specifying our parameter of interest and eliciting an uninformative prior on that parameter.

The parameter of interest: Here we explain the above comment regarding 'the parameter of interest' and implications of using $\bar{\beta}_{2}$. We denote the space spanned by a matrix $A$ by $s p(A)$. In cointegration analysis it is not the values of the elements of $\beta$ that are the object of interest, rather it is the space spanned by $\beta, \mathfrak{p}=s p(\beta)$, and this space is in fact all we are able to uniquely estimate. The parameter $\mathfrak{p}$ is an $r$-dimensional hyperplane in $R^{n}$
containing the origin and as such is an element of the Grassman manifold $G_{r, n-r}$ (James, 1954), $\mathfrak{p} \in G_{r, n-r}$. Before we derive the priors for $\mathfrak{p}$ we briefly comment on the relationship between priors for $\bar{\beta}_{2}$ and $\mathfrak{p}$. First we must introduce some notation for matrix spaces and measures on these spaces. For a more intuitive discussion of these concepts see Strachan and Inder (2003).

The $r \times r$ orthogonal matrix $C$ is an element of the orthogonal group of $r \times r$ orthogonal matrices denoted by $O(r)=\left\{C(r \times r): C^{\prime} C=I_{r}\right\}$, that is $C \in O(r)$. The $n \times r$ semi-orthogonal matrix $V$ is an element of the Stiefel manifold denoted by $V_{r, n}=\left\{V(n \times r): V^{\prime} V=I_{r}\right\}$, that is $V \in V_{r, n}$. As the vectors of any $V$ are linearly independent (since they are orthogonal) the columns of $V$ define a plane, $\mathfrak{p}$, which is an element of the $(n-r) r$ dimensional Grassman manifold, that is $\mathfrak{p}=s p(V) \in G_{r, n-r}$. The cointegrating space for an $n$ dimensional system with cointegrating rank $r$ is an example of an element of $G_{r, n-r}$. Finally, let the $j^{\text {th }}$ largest eigenvalue of the matrix $A$ be denoted $\lambda_{j}(A)$.

As discussed in James (1954), the invariant measures on the orthogonal group, the Stiefel manifold and the Grassman manifold are defined in exterior product differential forms (for measures on the orthogonal group and the Stiefel manifold, see also Muirhead 1982, Ch. 2). For brevity we denote these measures as follows. For a $(n \times n)$ orthogonal matrix $\left[b_{1}, b_{2}, \ldots, b_{n}\right] \in O(n)$ where $b_{i}$ is a unit $n$-vector such that $\beta=\left[b_{1}, b_{2}, \ldots, b_{r}\right] \in V_{r, n}, r<n$, the measure on the orthogonal group $O(n)$ is denoted $d v_{n}^{n} \equiv \Lambda_{i=1}^{n} \Lambda_{j=i+1}^{n} b_{j}^{\prime} d b_{i}$, the measure on the Stiefel manifold $V_{r, n}$ is denoted $d v_{r}^{n} \equiv \Lambda_{i=1}^{r} \Lambda_{j=i+1}^{n} b_{j}^{\prime} d b_{i}$, and the the measure on the Grassman manifold $G_{r, n-r}$ is denoted $d g_{r}^{n} \equiv$ $\Lambda_{i=1}^{r} \Lambda_{j=n-r+1}^{n} b_{j}^{\prime} d b_{i}$. These measures are invariant (to left and right orthogonal translations).

Theorem 1 The Jacobian for the transformation from $\mathfrak{p} \in G_{r, n-r}$ to vec $\left(\bar{\beta}_{2}\right) \in$ $R^{(n-r) r}$ is defined by

$$
\begin{equation*}
\underline{d g_{r}^{n}}=\pi^{-(n-r) r} \Pi_{j=1}^{r} \frac{\Gamma[(n+1-j) / 2]}{\Gamma[(r+1-j) / 2]}\left|I_{r}+\bar{\beta}_{2}^{\prime} \bar{\beta}_{2}\right|^{-n / 2}\left(d \bar{\beta}_{2}\right) \tag{5}
\end{equation*}
$$

where $\Gamma(q)=\int_{0}^{\infty} u^{q-1} e^{-u} d u$ for $q>0$. The underscore denotes the normalised measure such that $\int_{G_{r, n-r}} d g_{r}^{n}=1$.

Proof. In deriving the invariant measure on the Grassman manifold, James (1954) presents a relationship between an element of the Stiefel manifold, $V \in V_{r, n}$, and element of the Grassman manifold, $\mathfrak{p}=s p(\beta) \in G_{r, n-r}$
where the $r$-frame $\beta \in V_{r, n}$ and an element of the orthogonal group, $C \in$ $O(r) . \beta$ has a particular (fixed) orientation in $\mathfrak{p}$ such that it has only $(n-r) r$ free elements. Thus as $\mathfrak{p}$ is permitted to vary over all of $G_{r, n-r}, \beta$ is not free to vary over all of $V_{r, n}$. For $\mathfrak{p}=s p(V), V$ is determined uniquely given $\mathfrak{p}$ and orientation of $V$ in $\mathfrak{p}$ by $C \in O(r)$, such that $V=\beta C$. Note that as $\mathfrak{p}$ is permitted to vary over all of $G_{r, n-r}, V$ is free to vary over all of $V_{r, n}$. The resulting relationship between the measures is

$$
\begin{align*}
d v_{r}^{n} & =d g_{r}^{n} d v_{r}^{r} \\
\text { or } \underline{d v_{r}^{n}} & =\underline{d g_{r}^{n}} \underline{d v_{r}^{r}} . \tag{6}
\end{align*}
$$

James ${ }^{3}$ obtains the volume of $G_{r, n-r}$ as

$$
\begin{align*}
\int_{G_{r, n-r}} d g_{r}^{n} & =\frac{\int_{V_{r, n}} d v_{r}^{n}}{\int_{O(r)} d v_{r}^{r}} \\
& =\pi^{(n-r) r} \Pi_{j=1}^{r} \frac{\Gamma[(r+1-j) / 2]}{\Gamma[(n+1-j) / 2]} \tag{7}
\end{align*}
$$

Since the polynomial term accompanying the exterior product of the differential forms is equivalent to the Jacobian for the transformation (Muirhead 1982, Theorem 2.1.1), we can see from the expression (6) that the Jacobian for the transformation $V$ to $(\beta, C)$ is one.

Next consider the transformation from $V \in V_{r, n}$, to $\bar{\beta}_{2} \in R^{(n-r) r}$ and $C \in O(r)$ presented by Phillips (1994, Lemma 5.2 and see also Chikuse, 1998) and reproduced here:

$$
V=\left[c^{\prime}+c_{\perp}^{\prime} \bar{\beta}_{2}\right]\left[I_{r}+\bar{\beta}_{2}^{\prime} \bar{\beta}_{2}\right]^{-1 / 2} C .
$$

The differential form for this transformation is

$$
\begin{equation*}
\frac{d v_{r}^{n}}{}=\pi^{-(n-r) r} \Pi_{j=1}^{r} \frac{\Gamma[(n+1-j) / 2]}{\Gamma[(r+1-j) / 2]}\left|I_{r}+\bar{\beta}_{2}^{\prime} \bar{\beta}_{2}\right|^{-n / 2} d \bar{\beta}_{2}\left(\underline{d v_{r}^{r}}\right) \tag{8}
\end{equation*}
$$

(Phillips, 1994).
Equating (6) and (8) gives the result. Another, slightly more general proof for the same result is presented in Chikuse (1998).

[^1]Thus while a uniform distribution on $G_{r, n-r}$ implies a uniform distribution on $V_{r, n}$, this uniform distribution on $G_{r, n-r}$ implies a Cauchy distribution for $\bar{\beta}_{2}$. This last result was also derived by a very different approach by Villani (2000) for the case where $c=\left[\begin{array}{ll}I_{r} & 0\end{array}\right]$, although it holds for general $c$.

This transformation of the measure is relevant in both Bayesian and classical applications. As discussed in Phillips (1994), the form in (8) which introduces Cauchy tails into the distribution for $\bar{\beta}_{2}$ explains why applying linear restrictions to the maximum likelihood estimator of Johansen, $\widehat{\beta}=\left[\begin{array}{ll}\widehat{\beta}_{1}^{\prime} & \widehat{\beta}_{2}^{\prime}\end{array}\right]^{\prime}$ results in an estimator, $\stackrel{\rightharpoonup}{\beta}=\widehat{\beta}_{2} \widehat{\beta}_{1}^{-1}$, which is occasionally unreliable. The finite sample distribution for $\widehat{\bar{\beta}}$ has Cauchy tails and this Cauchy behaviour is a direct result of imposing the linear restrictions. This form also provides an alternative explanation for the rather similar but Bayesian results of Bauwens and Lubrano (1996). They show posterior Cauchy tail behaviour of the Bayesian estimator of $\bar{\beta}=\beta_{2} \beta_{1}^{-1}$ where no (additional) prior information on the cointegrating space is employed, although they use a 1-1 poly-t argument to find this result. Similar results can be found for the simultaneous equations model in Kleibergen and van Dijk (1998) and Drèze (1976).

Generally, estimating the cointegrating space using linear identifying restrictions will result in Cauchy tail behaviour unless there are other terms - such as prior information - offsetting the effect of this transformation. As one example of this effect of prior information, Bauwens and Lubrano (1996) show that overidentifying restrictions - which therefore reduce the number of free parameters to be estimated and, importantly, restrict the range of $\mathfrak{p}$ within $G_{r, n-r}$ - will result in a posterior with as many moments as overidentifying restrictions.

The Jacobian defined by (5) implies that a flat prior on $\mathfrak{p}$ is informative with respect to $\beta_{2}$ and vice versa. This leads us to consider the implications of a flat prior on $\beta_{2}$ for the prior on $\mathfrak{p}$.

Theorem 2 The Jacobian for the transformation from $\bar{\beta}_{2} \in R^{(n-r) r}$ to $\mathfrak{p} \in$ $G_{r, n-r}$ is defined by

$$
\begin{align*}
\left(d \bar{\beta}_{2}\right) & =\pi^{(n-r) r} \Pi_{j=1}^{r} \frac{\Gamma[(r+1-j) / 2]}{\Gamma[(n+1-j) / 2]}\left|I_{r}+(c \beta)^{\prime-1} \beta^{\prime} c_{\perp}^{\prime} c_{\perp} \beta(c \beta)^{-1}\right|^{n / 2}\left(d g_{r}^{n}\right) \\
& =J d g_{r}^{n} . \tag{9}
\end{align*}
$$

Proof. Invert (9) and replace $\bar{\beta}_{2}$ by $c_{\perp} \beta(c \beta)^{-1}$.
A common justification for the linear restrictions is that an economist will usually have some idea about which variables will enter the cointegrating relations and so she chooses $c$ to select the rows of coefficients most likely to be nonzero - more generally linearly independent from eachother - and then normalise on these coefficients. This is a necessary assumption to ensure $(c \beta)^{-1}$ exists. As the next theorem shows, using these linear restrictions, however, has the unexpected and undesirable result that the Jacobian for $\bar{\beta}_{2} \rightarrow \mathfrak{p}$ places more weight in the direction where the coefficients thought most likely to be different from zero are, in fact, zero (or linearly dependent).

Theorem 3 Given $r$, use of the normalisation $\bar{\beta}_{2}={\underset{c}{\perp}}^{\beta}(c \beta)^{-1}$ results in a transformation of measures for the transformation $\bar{\beta}_{2} \in R^{(n-r) r} \rightarrow \mathfrak{p} \in$ $G_{r, n-r}$ that places infinite mass in the region of null space of c relative to the complement of this region.

Proof. Let $\rho_{c_{\perp}}$ be the plane defined by the null space of $c$. Define a ball, $\mathfrak{B}$, of fixed diameter, $d$, around $\rho_{c_{\perp}}$ and let $N_{0}=\mathfrak{B} \cap G_{r, n-r}$ and $N=G_{r, n-r}-N_{0}$. Since for $d>0, \int_{N} J d g_{r}^{n}$ is finite whereas $\int_{N_{0}} J d g_{r}^{n}=\infty$, we have

$$
\frac{\int_{N_{0}} J d g_{r}^{n}}{\int_{N} J d g_{r}^{n}}=\infty
$$

To summarise, normalisation of $\beta$ by choice of $c$ with a flat prior on $\bar{\beta}_{2}$ implies infinite prior odds against this normalisation.

To demonstrate this result, consider a $n$-dimensional system for $y=$ $\left(x^{\prime}, z^{\prime}\right)^{\prime}$ where $x$ is a $r$ vector. To implement linear restrictions a normalisation must begin by first choosing $c$. Suppose it is believed that if a cointegrating relationship exists then it will most likely involve the elements of $x$ in linearly independent relations. That is in $y \beta=x \beta_{1}+z \beta_{2} \sim I(0)$, $\operatorname{det}\left(\beta_{1}\right)$ is believed far from zero making it safe to normalise on $\beta_{1}$, and so choose $c=\left[\begin{array}{ll}I_{r} & 0\end{array}\right]$ and estimate $\bar{\beta}_{2}=c_{\perp} \beta(c \beta)^{-1}$.

From (9) we see as $\mathfrak{p}=s p(\beta) \rightarrow s p(c), c_{\perp} \beta \rightarrow 0_{(n-r) \times r}$ and $c \beta \rightarrow O(r)$ and $J \rightarrow 1$. However, as vectors in $\beta$ approach the null space of $c$, that is $\operatorname{det}(c \beta) \rightarrow 0$, then $(c \beta)^{-1} \rightarrow \infty$, and thus $J \rightarrow \infty$. As a result the prior will more heavily weight regions where $\operatorname{det}(c \beta)=\operatorname{det}\left(\beta_{1}\right) \approx 0$, contrary to the intention of the economist. As a trivial example, consider our money demand study with $r=1$ and $\zeta_{t}=\beta_{1} m_{t}+\beta_{2} i n c_{t}$. If we believe money is
most likely to enter the cointegrating relation, we would choose $c=(1,0)$ as we believe $\beta_{1} \neq 0$. Yet the Jacobian places infinite weight in the region $\beta_{1}=0$ excluding $m_{t}$ from the cointegrating relation.

A uniform prior on the cointegrating space: There is clearly a need to consider a new approach to eliciting priors for $\beta$. We wish to avoid the problems outlined above deriving from the use of linear restrictions with normalisation to identify the elements of $\beta$ and the subsequent treatment of $\bar{\beta}_{2}$ as the parameter of interest. Our recommendation is, if the economist wishes to incorporate prior beliefs about the cointegrating relations, these should be expressed in the prior distribution for the cointegrating space.

As we have claimed the cointegrating space to be the parameter of interest, rather than $\bar{\beta}_{2}$, we propose working directly with $\mathfrak{p}=s p(\beta)$ avoiding the linear restrictions and normalisation. Initially we present a distribution and identifying restrictions for $\beta$ from the form of the uniform distribution for $\mathfrak{p}$ over $G_{r, n-r}$ using the results of James (1954) (see also Strachan and Inder, 2003). The identifying restrictions on $\beta$ follow naturally from this approach. This prior has the form

$$
\begin{equation*}
p(\beta)=\frac{1}{\int_{G_{r, n-r}} d g_{r}^{n}} \tag{10}
\end{equation*}
$$

where $\beta$ is the $r$-frame with fixed orientation in $\mathfrak{p}$. In the proof of Theorem 1 , the measure on $G_{r, n-r}$ used in the above expression is derived from its relationship with the spaces $V_{r, n}$ and $O(r)$.

To avoid using linear restrictions with a normalisation to identify $\beta$ it is necessary to find an alternative set of restrictions that do not require knowledge of $c$ and which avoid the issues associated with the posterior for $\bar{\beta}_{2}$. Fortunately the definition (6) and the discussion in the proof of Theorem 1 provide a natural solution to this question. That is use $\beta \in V_{r, n_{i}}$ which implies $r(r+1) / 2$ restrictions. The dimension of the Grassman manifold is only $(n-r) r$ while the dimension of the Stiefel manifold $V_{r, n}$ is $n r-r(r+1) / 2$, which exceeds that of $G_{r, n-r}$ by $r(r-1) / 2$. In (6), these remaining restrictions come from the orientation of $\beta$ in $\mathfrak{p}$ by $C \in O(r)$. The prior, the posterior (as is made clear later) and the differential form for $\beta$ are all invariant to translations of the form $\beta \rightarrow \beta H, H \in O(r)$. Therefore it is possible to work directly with $\beta$ as an element of the Stiefel manifold and adjust the integrals with respect to $\beta$ by $\left(\int_{O(r)} d v_{r}^{r}\right)^{-1}$ as shown in (7). Note that these identifying restrictions do not distort the weight on the space of the parameter of interest, $\mathfrak{p}$, and it is never necessary to actually specify the
orientation of $\beta$ in $\mathfrak{p}$.
Thus, contrary to the situation when using linear identifying restrictions, we are able to employ innocuous identifying restrictions, place a prior directly on the parameter of interest and, as we show below, we achieve a better behaved posterior about which we know much more. Before we discuss the posterior, however, we extend this approach to informative priors on the cointegrating space.

An informative prior on the cointegrating space: If an economist believes a parameter is likely to have a particular value, to incorporate this prior belief she places more prior mass around this likely point. When considering the cointegrating space $\mathfrak{p}$, we will denote our desired location or the likely value as $\mathfrak{p}^{H}=s p(H \kappa)$ (as in the Garrett et al. case) where $H \in V_{s, n}$ is a known $n \times s(s \geq r)$ matrix, $H_{\perp} \in V_{n-s, n}$ its orthogonal complement and $\kappa$ is an $s \times r$ full rank $r$ matrix. To obtain $H$ in $V_{s, n}$, first specify the general matrix $H^{g}$ with the desired coefficient values. One might consider as an example the matrix $H_{2}$ presented in Section 3. Next map this to $V_{r, n}$ by the transformation $H=H^{g}\left(H^{g \prime} H^{g}\right)^{-1 / 2}$.

At the extreme, a dogmatic prior for $\mathfrak{p}$ could be specified by letting $\beta=$ $H \kappa V, V \in O(r)$. Next define $\kappa V=V_{\kappa} \in V_{r, s}$ and specify the prior in (10) for $V_{\kappa}$. This resulting prior which assigns unit probability mass to $\mathfrak{p}=\mathfrak{p}^{H}$.

Next we specify an informative, nondogmatic, prior for $\mathfrak{p}$ centered at $\mathfrak{p}=\mathfrak{p}^{H}$ but with positive mass elsewhere in $G_{r, n-r}$.

Let the random scalar $\tau$ have $E(\tau)=0$ and $E\left(\tau^{2}\right)=\sigma^{2}$. The value of $\sigma$ will control the tightness of the prior density around $\mathfrak{p}^{H}$. Next construct

$$
\begin{aligned}
P_{\tau} & =H H^{\prime}+H_{\perp} H_{\perp}^{\prime} \tau \\
& =\left[H H_{\perp}\right]\left[\begin{array}{cc}
I_{r} & 0 \\
0 & I_{n-r} \tau
\end{array}\right]\left[\begin{array}{c}
H^{\prime} \\
H_{\perp}^{\prime}
\end{array}\right]
\end{aligned}
$$

and let the elements of the $n \times r$ matrix $Z$ be independently distributed as standard normal, $N(0,1)$. The matrix $X=P_{\tau} Z$ can be decomposed as $X=\beta^{*} \kappa$ where $\beta^{*} \in V_{r, n}$ and $\kappa$ is an $r \times r$ upper triangular matrix. For $\tau \neq 0$ and $|\tau|<\infty$, the space of $\beta^{*}, \mathfrak{p}=s p\left(\beta^{*}\right)$, is a direct weighted sum of the spaces $\mathfrak{p}^{H}$ and $\mathfrak{p}^{H_{\perp}}$ with the weight determined by $\tau$.

At $\tau=0$ and $\tau= \pm \infty, \mathfrak{p}$ is respectively $\mathfrak{p}^{H}$ and $\mathfrak{p}^{H_{\perp}}$. It is for this reason that we chose $E(\tau)=0$ such that with respect to $\tau$, the space will on average be $\mathfrak{p}^{H}$. One choice for $\tau$ is $N(0,1)$ and the form of the resultant density for
$\beta^{*}$ and the hyperparameter $\tau$ is

$$
\begin{equation*}
p(\tau, \beta)=\tau^{-(n-r) r} \exp \left\{-\frac{\tau^{2}}{2}\right\}\left|\beta^{\prime} P_{\tau^{2}}^{-1} \beta\right|^{-n / 2} \mathfrak{c}_{r} \tag{11}
\end{equation*}
$$

where $\mathfrak{c}_{r}=2^{-r-1 / 2} \pi^{r(r-1) / 4-(n+1) r / 2} \Pi_{j=1}^{r} \Gamma[(n+1-j) / 2]$. This prior treats the point $\mathfrak{p}^{H_{\perp}}$, which occurs at $\tau=\infty$, as an improbable (practically impossible) event regardless of the choice of $\sigma$. This is desirable since at $\tau=\infty$ the dimension of the cointegrating space, $\operatorname{dim}(\mathfrak{p})$, would become $\operatorname{dim}\left(\mathfrak{p}^{H_{\perp}}\right)=\min (p-r, r)$ rather than $r$.

As an alternative, if the researcher would prefer to assign more weight in the direction of $\mathfrak{p}^{H_{\perp}}$ but preserve $\operatorname{dim}(\mathfrak{p})=r$ with probability one, she may choose $P_{\tau}=H H^{\prime}\left(1-\tau^{2}\right)^{1 / 2}+H_{\perp} H_{\perp}^{\prime} \tau$ with $\tau \in[-1,1]$. Again the choice of $E(\tau)=0$ would make sense and $E\left(\tau^{2}\right)=\sigma^{2}$ controls the tightness of the density around $\mathfrak{p}^{H}$. A possible choice of a distribution for $\eta=\tau+1$ may be Beta over $\eta \in[0,2]$ which allows some mass to be distributed around $\mathfrak{p}^{H_{\perp}}$ by appropriate choice of parameter values.

### 4.3 The posteriors.

Using the priors specified above, the general form of the posterior is then

$$
\begin{aligned}
p(B, \Sigma, \beta, r, i \mid y) \propto & p(\beta)|\Sigma|^{-\left(T+n+k_{i}+r+1\right) / 2} \\
& \left.\times \exp \left\{-\frac{1}{2} \operatorname{tr} \Sigma^{-1}\left[T S+(B-\widetilde{B})^{\prime} V(B-\widetilde{B})\right]\right\} 12\right) \\
& \times(2 \pi)^{-n\left(k_{i}+r\right) / 2} 100^{n\left(k_{i}+r\right) / 2} \\
= & k(B, \Sigma, \beta, \omega \mid y)
\end{aligned}
$$

where $S=S_{00}-S_{01} \beta\left(\beta^{\prime} S_{11} \beta\right)^{-1} \beta^{\prime} S_{10}, \widetilde{B}=\left[\begin{array}{cc}\widetilde{\alpha}^{\prime} & \widetilde{\Phi}^{\prime}\end{array}\right]^{\prime}, \widetilde{\alpha}=\left(\beta^{\prime} S_{11} \beta\right)^{-1} \beta^{\prime} S_{10}$, $\widetilde{\Phi}=S_{22}^{-1} S_{20}$, and $V=\widetilde{\beta}^{\prime}\left(\Sigma_{t=1}^{T} z_{t}^{\prime} z_{t}+\underline{H}\right) \widetilde{\beta}$ where $z_{t}=\left(z_{1, t} z_{2, t}\right)$. The values for the $S_{i j}$ are defined as

$$
\begin{array}{rlr}
T M_{i j} & =h_{i j}+\Sigma_{t=1}^{T} z_{i, t}^{\prime} z_{j, t} \quad \text { for } i \text { and } j=1,2, \\
h_{i j} & =0 \text { if } i \neq j \text { and } h_{i i}=0.01 I, \\
T M_{20} & =\Sigma_{t=1}^{T} z_{2, t}^{\prime} \Delta y_{t}, \quad T M_{10}=\Sigma_{t=1}^{T} z_{1, t}^{\prime} \Delta y_{t}, \\
T M_{00} & =\Sigma_{t=1}^{T} \Delta y_{t}^{\prime} \Delta y_{t} \quad \text { and so }
\end{array}
$$

$$
S_{i j}=M_{i j}-M_{i 2} M_{22}^{-1} M_{2 j} \quad \text { for } i j=0,1,2
$$

except $i=j=2$ where

$$
\begin{aligned}
& S_{22}=M_{22}-M_{21} M_{11}^{-1} M_{12} \text { and } \\
& S_{20}=M_{20}-M_{21} M_{11}^{-1} M_{10} .
\end{aligned}
$$

For later use we also define $D_{0}=D_{1}-D_{2}, D_{1}=S_{11}$ and $D_{2}=S_{01} S_{11}^{-1} S_{10}$.
For $B \in R^{\left(k_{i}+r\right) n}$ and $\Sigma$ positive definite (denoted $\Sigma>0$ ), to estimate the relevant Bayes factors, $B_{j l}=\frac{m_{j}}{m_{l}}$, for the models of interest, estimates of the marginal likelihoods, e.g.

$$
\begin{equation*}
m_{j}=\sum_{\omega} \int_{R^{\left(k_{i}+r\right)_{n}}} \int_{\Sigma>0} \int_{G_{r, n-r}} k_{\theta}(B, \Sigma, \beta, \omega \mid y)\left(d g_{r}^{n}\right)(d \Sigma)(d B), \tag{13}
\end{equation*}
$$

are required. To perform the integration in (13) of $\theta=(\Sigma, B, \beta)$, we first analytically integrate (12) with respect to $(\Sigma, B)$ as these parameters have conditional posteriors of standard form. This integration gives us the following.

Theorem 4 The marginal posterior for $(\beta, \omega)$ is

$$
\begin{equation*}
p(\beta, \omega \mid y) \propto \mathfrak{g}_{\omega}\left|S_{00}\right|^{-T / 2}\left|M_{22}\right|^{-n / 2}\left|\beta^{\prime} D_{0} \beta\right|^{-T / 2}\left|\beta^{\prime} D_{1} \beta\right|^{(T-n) / 2} p(\beta) \tag{14}
\end{equation*}
$$

where in this case $\mathfrak{g}_{\omega}=T^{-n r / 2} \pi^{-\left(n_{i}-r\right) r / 2} 100^{n\left(k_{i}+r\right) / 2}$.
Proof. See, for example, Zellner (1971)
Remark: It is from the expression (14) that we see that not only is $d g_{r}^{n}$ invariant to $\beta \rightarrow \beta C$ for $C \in O(r)$, but so is $k(\beta)$ and thus the posterior.

Next we need to integrate (14) with respect to $\beta$ to obtain the posterior for $\omega$. Since $\mathfrak{g}_{\omega}$ is finite for the class of priors considered, that the Bayes factor is finite requires the integral with respect to $\beta$ to be finite. The following are some general results with respect to this integral.

Theorem 5 The marginal posterior density for $\beta$ conditional upon $\omega$ has the same form for each model considered:

$$
\begin{align*}
p(\beta \mid \omega, y) & \propto\left|\beta^{\prime} D_{0} \beta\right|^{-T / 2}\left|\beta^{\prime} D_{1} \beta\right|^{(T-n) / 2}  \tag{15}\\
& =k_{\beta}(\beta)
\end{align*}
$$

where $k_{\beta}(\beta)=\left|\beta^{\prime} D_{0} \beta\right|^{-T / 2}\left|\beta^{\prime} D_{1} \beta\right|^{(T-n) / 2}$.

Theorem 6 The marginal posterior density for $\beta$ conditional upon $(r, i)$ in (15) is proper and all finite moments exist.

Proof. Denote by $b_{i j}$ any element of $\beta$. The proof follows from the result that the integral

$$
M_{\beta}=\int_{V_{r, n}}\left|b_{i j}\right|^{m} k_{\beta}(\beta) d v_{r}^{n}
$$

for $m=0,1,2, \ldots$ is bounded above almost everywhere by the finite integral $M \int_{-1}^{1}\left|b_{i j}\right|^{m} d b_{i j}$. As the elements of $\beta, b_{i j}$, have compact support, it is only necessary for this proof to show that $k_{\beta}(\beta) d v_{r}^{n}$ is bounded above almost everywhere by some finite constant function over $V_{r, n}$ (note the adjustment to the integral over $G_{r, n-r}$ simply requires division by the finite volume of $O(r)$, thus we only need consider the integral over $V_{r, n}$ ). As demonstrated in the proof to Theorem 1 in Section (4.2), $d g_{r}^{n}$ is integrable and therefore bounded above almost everywhere by some finite constant, $M_{1}$.

The eigenvalues $\lambda_{j}\left(D_{l}\right)$ for $l=0,1$, will be positive and finite with probability one. By the Poincaré separation theorem, since $\beta \in V_{r, n}$, then

$$
\Pi_{j=1}^{r} \lambda_{n-r+j}\left(D_{l}\right) \leq\left|\beta^{\prime} D_{1} \beta\right| \leq \Pi_{j=1}^{r} \lambda_{j}\left(D_{l}\right)
$$

and so $k_{\beta}(\beta)$ is bounded above (and below) by some positive finite constant, $M_{2}$. Thus $k_{\beta}(\beta) d g_{r}^{n}$ has a finite upper bound, $M=M_{1} M_{2}$. With the compact support for $b_{i j}$, these conditions are sufficient to ensure the posterior for $\beta$ will be proper and all finite moments exist (see Billingsley 1979, pp. 174 and 180).

The importance of Theorem 6 becomes evident when we consider that economic objects of interest to decision-makers are often linear or convex functions of the cointegrating vectors. As discussed in a previous section, with linear identifying restrictions expectations of such objects are not defined unless overidentifying restrictions are imposed or an informative prior is used. Further, the result in Theorem 6 holds even when exogeneity is imposed - again in contrast to when linear identifying restrictions are used.

To obtain the posterior distribution of $\omega=(r, o, e, i), p(\omega \mid y)$, it is necessary to integrate (14) with respect to $\beta$ and so obtain an expression for

$$
\begin{equation*}
p(\omega \mid y)=\int p(\beta, \omega \mid y) d g_{r}^{n} \tag{16}
\end{equation*}
$$

The marginal density of $\beta$ conditional on $\omega$ in (15) is not of standard form. Although one may exist, we do not currently know of a simple, general analytical solution for $\mathfrak{c}_{\omega}=\int_{V_{r, n}} k_{\beta}(\beta) d g_{r}^{n}$ and so we estimate $\boldsymbol{c}_{\omega}$.

Two possible approaches to estimating $\mathfrak{c}_{\omega}$ are either to use Markov Chain Monte Carlo (MCMC) methods or to use deterministic methods to approximate the integral. Kleibergen and van Dijk (1998) develop a MCMC scheme in the simultaneous equations model and Kleibergen and Paap (2002) extend this to the cointegrating error correction model. Bauwens and Lubrano (1996) demonstrate an alternative approach. In each of these applications a method is presented to evaluate integrals using MCMC when $\beta$ has been identified using linear restrictions rather than those used in this paper. Strachan (2003) demonstrates the MCMC approach when $\beta$ has been identified using restrictions related to those of the ML estimator of Johansen (1992). An approach commonly used in classical work to approximate integrals over $V_{r, n}$, is to use the Laplace approximation which is computationally much faster than MCMC. Strachan and Inder (2003) present the Laplace approximation to (16).

The Laplace approximation is a second order asymptotic approximation to the marginal likelihood. There is an alternative, simpler, first order asymptotic approximation to the marginal likelihood which assumes dominance by the likelihood. That is, we may treat the Bayesian information criteria of Schwarz (1978) (BIC) as an asymptotic approximation to $-T / 2$ times the log marginal likelihood, $\boldsymbol{c}_{\omega}$, for each model. Thus we are able to obtain estimates of the posterior probabilities of the models. In the Applications section we employ both the Laplace and the BIC approaches.

As we wish to obtain estimates of economic objects of interest averaged across models we need to be able to obtain draws of $\beta$ from the posterior. The next subsection outlines an approach to obtaining MCMC draws from the posterior with the uniform prior used in this paper.

Obtaining MCMC draws from the posterior with an uninformative prior on the cointegrating space: As demonstrated in Strachan and Inder (2003), the mode of the marginal posterior for $\beta, \widetilde{\beta}$, is relatively straight forward to obtain. Denote this point by $H$. This gives us a method of developing a candidate density for the posterior with mass in the same location as the posterior by using an approach similar to that used to develop the informative prior in the previous section.

First, specify a distribution for $\tau$ which includes specifying $\sigma$. Let $H=\widetilde{\beta}$. Take a draw of $\tau$ and construct $P_{\tau}$, then draw $Z$ from the multivariate standard normal. Next, construct $X=P_{\tau} Z$ and then $\beta^{*}$ from the decomposition $X=\beta^{*} \kappa . \beta^{*}$ is then a draw from the candidate density for $\mathfrak{p}$ with location
$\mathfrak{p}^{H}=s p(\widetilde{\beta})$. Each of these steps is explained in the previous section.
Acceptance for a Metropolis Hastings scheme or weighting in an Importance sampling method will be determined by a function the ratio of the posterior to the candidate - an example of which is provided in (11) with $H=\widetilde{\beta}$.

Finally, a word about dispersion. For the candidate density to more closely match the posterior in form, the value of $\sigma$ could be calibrated to a desired level of dispersion, preferably to match that of the posterior. This can be achieved by using (MC)MC draws and the span variation measure (sv) of Villani (2000). This measure of variation can be used to express the degree of variation in a distribution as a proportion of the variation under the uniform distribution. The uniform distribution is an appropriate reference as it implies equal variation in every direction.

## 5 Applications

In this section we provide some preliminary results for the applications to two economic models. The first is relatively simple and involves only rank and exogeneity restrictions in a small model. The second involves rank, exogeneity, trend and overidentifying restrictions.

### 5.1 Australian money demand.

We consider a simple study of Australian money demand. The variables, all of which appear in logarithmic form, are defined as

$$
y_{t}=\left(\begin{array}{lll}
m_{t} & p_{t} & i n c_{t}
\end{array}\right),
$$

where $m_{t}$ is the measure of money - either M1, M3 or broad money (BM), $p_{t}$ is the price level, such that $m_{t}-p_{t}$ measures real money, and $i n c_{t}$ is real gross national income. The data are quarterly observations from September 1976 to December 2002 and were sourced from the web site of The Australian Bureau of Statistics, specifically tables D03, G09 and G02.

We are interested in estimation of the parameters determining the long and short run behaviour of $m_{t}$ and in forecasts of $m_{t}$, where the forecasts may be over the long run, or both the long run and short run. Here we do not regard the long run as the true long run path of $m_{t}$, rather the equilibrium
relationships to which the elements of $y_{t}$ would revert if all future errors were zero.

If we are interested in whether we may estimate the money demand equation

$$
z_{t}=\beta_{1} m_{t}+\beta_{2} p_{t}+\beta_{3} i n c_{t}=y_{t} \beta
$$

from a single equation for $m_{t}$, then this would require that the variables ( $p_{t} \quad$ inc $c_{t}$ ) be weakly exogenous with respect to $\beta$. As is commonly done (see for example Johansen, 1995 and Funke, Hall and Beeby, 1997), we regard the money demand relation as the cointegrating relation between the variables in $y_{t}$. Therefore, for the analysis to make sense, we require in addition that cointegration hold (and so $r=1$ or 2 ). In terms of the posterior probabilities, these joint conditions imply $(e=2, r=1)$. The same condition is said to hold for strong exogeneity of ( $p_{t} \quad$ inc $c_{t}$ ) with respect to $\beta$.

We can determine the values of other conditional probabilities as well. Note that strong exogeneity implies weak exogeneity, and exogeneity with respect to $\phi_{2}$ implies exogeneity with respect to $\phi_{1}$. Further, for ( $p_{t}$ inc $c_{t}$ ) to be exogenous with respect to $\beta$, this implies $\alpha_{2}=0$ and so the rank of $\Pi$ can be at most 1. These results imply that the probabilities of exogeneity conditional on all ranks above 1 will be zero.

The posterior probabilities of the ranks are reported in Table 1. These indicate that there is strong support for the requirement of $r=1$ with some support for $r=0$.

|  | $\widehat{p}(r \mid y)$ |  |  |
| :---: | :---: | :---: | :---: |
| $r$ | M1 | M3 | BM |
| 0 | 0.268 | 0.588 | 0.316 |
| 1 | 0.732 | 0.412 | 0.684 |
| 2 | 0.000 | 0.000 | 0.000 |
| 3 | 0.000 | 0.000 | 0.000 |

Table 1: Posterior probabilities of the ranks for money demand study.
For each money series ( $p_{t} \quad i n c_{t}$ ) is weakly exogenous with respect to the parameters determining the long run behaviour of money $\left(\phi_{1}\right)$ with probability one. This implies the figures in Table 1 are the marginal probabilities of the ranks, the joint probabilities of the ranks and weak exogeneity and also the conditional probabilities of the ranks given weak exogeneity $p(r \mid y)=p(r, e=2 \mid y)=p(r \mid e=2, y)$.

There is clearly strong support for weak exogeneity of prices and income with respect to the long run parameters, $\left(\beta, \alpha_{1}\right)$, in the regression equation for M1 and broad money and these features imply the rank of the system is one. Thus single equation estimation of money demand relations for Australian data is an appropriate method. On the question of whether the velocity of money appears to be a stable relation, we found no support for this model within the model set and this agreed with the classical results for this data set.

### 5.2 Structural model of the UK economy

Analysing their macroeconomic model within (1), Garratt, et al. find support for $r=5$ using Johansen's trace test. They also find support for the overidentifying restrictions and trend restrictions using a log-likelihood ratio test, where they used bootstrap estimates of critical values. They do not appear to test support for the weak exogeneity of oil prices. Below we present the posterior probabilities of the various models (zeros or near zeros are suppressed or omitted) where $e=2$ implies weak exogeneity of oil prices.

|  | $\widehat{p}(r, i, e=1 \mid y)$ |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $r$ | $i=1$ | $i=2$ | $i=3$ | $i=4$ | $i=5$ |  |
| 0 |  |  |  |  |  |  |
| 1 |  | 0.0101 | 0.0008 | 0.0324 | 0.0128 |  |
| 2 |  | 0.0004 | 0.0001 | 0.0040 | 0.0190 |  |
|  | $\widehat{p}(r, i, e=2 \mid y)$ |  |  |  |  |  |
| $r$ | $i=1$ | $i=2$ | $i=3$ | $i=4$ | $i=5$ |  |
| 0 |  |  |  |  |  |  |
| 1 |  | 0.0213 | 0.0085 | 0.2507 | 0.1262 |  |
| 2 |  | 0.0001 | 0.0002 | 0.0756 | 0.4376 |  |
| 3 |  |  |  |  | 0.0001 |  |

Marginal Probabilities

| $\widehat{p}(i \mid y)$ | $i=1$ | $i=2$ | $i=3$ | $i=4$ | $i=5$ |
| :---: | :---: | :---: | :---: | :---: | :---: |
|  | 0.0000 | 0.0320 | 0.0096 | 0.3626 | 0.5958 |
| $\widehat{p}(r \mid y)$ | $r=0$ | $r=1$ | $r=2$ | $r=3$ | $r=4$ |
|  |  | 0.4628 | 0.5371 | 0.0001 |  |

The posterior probabilities for the rank suggest support for a rank of one or two with $P(r=1 \mid y)=0.4628$ and $P(r=2 \mid y)=0.5371$. We also find
marginal probabilities of no deterministic processes $(i=5)$ of 0.5958 and of an intercept in the cointegrating relations $(i=4)$ of 0.3626 . The posterior probability that the oil prices are weakly exogenous is 0.9203 providing strong support for this restriction. The combined restrictions of overidentification, exogeneity, four stochastic trends and a linear trend in the long run moneyincome relation had a joint probability of effectively zero within this model set.

With the overidentifying restrictions, the only coefficients to be estimated in the long run relations, ignoring the intercepts, are in the money market equilibrium condition given by

$$
h_{t}-y_{t}=\mu_{21, \tau} t+\beta_{22} r_{t}+u_{4, t}
$$

Estimating the coefficients in this relation subject to the restrictions proposed by Garratt, et al., we obtain

$$
h_{t}-y_{t}=-0.0070 t-43.2148 r_{t}+u_{4, t}
$$

which compares with the classical estimate of Garratt et al. of

$$
h_{t}-y_{t}=-0.0073 t-56.0975 r_{t}+u_{4, t} .
$$

Both results suggest a downward trend in the money-income ratio which may be attributed to technological innovations in the finance sector (Garratt, et al. 2002).

Although there is a clear modal model, $M_{(r, o, e, i)}=M_{(2,1,2,5)}$, there is just as clearly some support for nearby models such as $M_{(1,1,2,4)}$ and $M_{(1,1,2,5)}$. We would like to incorporate the information value of these models for decision making and one way to achieve this is through averaging the economic object of interest. As an example of an averaged output which can be used as an input for decision making, Figure 1 presents the higher posterior density regions (hpds) for the impulse response function over 60 months for a response in relative UK prices, $p_{t}-p_{t}^{*}$, to a shock in oil prices, $p_{t}^{o}$. This output is averaged across all models and was produced from 100,000 draws from the full posterior. The intervals plot the boundaries of the $20 \%, 40 \%, 60 \%$ and $80 \%$ hpds. The UK during the period of the sample was a net oil exporter and we see the effect of this reflected in the figure as the distribution of the response path indicates initially that the rest of the world experiences a larger response to an oil price shock than the UK, after which the UK appears to
catch up slightly. However, the greater impact on world prices relative to UK prices seems to persist as after 60 months the path is centred around a slightly negative mode just above negative $1 \%$. This is not a surprising result given the likely exchange rate adjustment in the pound.

$$
* * * * * * * * * * \text { Figure } 1 \text { around here }{ }^{* * * * * * * * * *}
$$

It should be pointed out that these intervals are not comparable with the usual classical confidence intervals as they incorporate variable uncertainty, parameter uncertainty and model uncertainty. With this extra uncertainty it is sensible then that the intervals containing a given mass will be wider and the mass in any particular region does not have the same interpretation. Trimming the model set of unreasonable models would likely produce smaller intervals. However, the results we present are more informative on the question 'What will happen to relative prices in the UK if there is an oil price shock?' as they do not require the addendum: '... if this model and these parameter values are correct?'.

Figure 2 plots the hpds for the impulse response function over 60 months for a response in UK inflation, $\Delta p_{t}$, to a shock in oil prices, $p_{t}^{o}$, again produced from 100,000 draws from the posterior. The median response after 60 months shows a moderate increase in the level of inflation of around $2.5 \%$ and so the median impulse response is about where we would expect it and the $20 \%$ and $40 \%$ hpds are reasonable.

An interesting feature of both figures are the long tails at low lags. This tail behaviour is due entirely to the set of 40 models (out of 97 models) in which oil prices are not constrained to be weakly exogenous. Although these models are given a small (but not negligible) posterior probability (around $8 \%$ ), their implied response paths are so extreme that they have a noticeable influence upon the marginal distribution of the response.

It is to demonstrate this rather strange behaviour that we have reported the results using the BIC approximation to the posterior probabilities. The same plots of the hpds for the impulse response paths when we used the Laplace approximation or the MCMC estimation do not demonstrate such an extreme diversion in the tail and look similar to what we obtain if we use BIC but exclude the models in which oil prices are not exogenous $(e=1)$. The
reason for this is that the Laplace and MCMC methods tend to concentrate the mass of the density for the models on fewer models and attribute no mass to the models with $e=1$. The behaviour in Figure 2 demonstrates the risks of conditioning on particular models, but also the risks - also inherent in our approach - of not using a sufficiently well considered model set.

## 6 Conclusion.

In this paper we have presented an approach to obtaining inference on the structural features of the vector autoregressive model that are of interest to researchers and for policy analysis. This approach allows the incorporation of uncertainty about the 'true state of nature' into the conduct of policy analysis by producing output averaged across models rather than output conditional upon a particular model. The output produced this way allows policy recommendations to be made that are not conditional on a particular model, and thus this model averaging approach provides an important alternative to the more commonly used model selection approach. Specifically we provide techniques for estimating marginal likelihoods for models of cointegration, deterministic processes, exogeneity, and overidentifying restrictions upon the cointegrating space. These estimates are derived using a mixture of analytical integration and MCMC or asymptotic approximations to integrals. Two applications of these tools are provided. First for a simple example of a model of Australian money demand and, second, a more complete macroeconomic model of the UK proposed by Garratt, et al..

Very natural extensions of our approach are to include inequality conditions in the parameter space of the structural VAR or forms of nonlinearity in the model itself. For instance, in using a SVAR for business cycle analysis one may use prior information on the length and amplitude of the period of oscillation. An example of a possible nonlinear structure that may prove useful is presented in Paap and van Dijk (2003). Systematic use of inequality conditions and nonlinearity implies a more intense use of MCMC algorithms.

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## 9 Appendix

### 9.1 Posterior distribution of $\beta_{2}$ given exogeneity.

In this section we show the marginal likelihoods are not well defined for $\beta_{2}$ when weak exogeneity is imposed. The following results apply for a wide class of priors. To consider weak exogeneity with respect to $\beta$, we partition the
matrix $\alpha$ as $\alpha=\left(\begin{array}{cc}\alpha_{1} & \alpha_{2}\end{array}\right)$ such that the exogeneity restriction is implied by $\alpha_{2}=0$ and derive the marginal distribution of $\left(\alpha_{2}, \beta\right)$. Next we set $\alpha_{2}=0$ in $p\left(\alpha_{2}, \beta \mid y\right)$. If $\int p\left(\beta \mid \alpha_{2}=0, y\right)\left(\beta^{\prime} d \beta\right)=\infty$, then the posterior does not integrate to a finite constant and Bayes factors are not defined. Thus by demonstrating that the above integral is not finite when linear restrictions are imposed on $\beta$, such that $\beta=\left[I_{r} \beta_{2}^{\prime}\right]^{\prime}$ and $\beta_{2} \in R^{(n-r) r}$, we show the marginal likelihoods are not finite.

The marginal, joint posterior distribution for $(\alpha, \beta)$ given $r=p_{2}$, is

$$
p(\alpha, \beta \mid r, y) \propto\left|T S+(\alpha-\widehat{\alpha})^{\prime} \beta^{\prime} S_{11} \beta(\alpha-\widehat{\alpha})\right|^{-(\nu+r) / 2}
$$

such that

$$
\begin{aligned}
p(\alpha, \beta \mid r, y) \propto & \left|T S+(\alpha-\widehat{\alpha})^{\prime} \beta^{\prime} S_{11} \beta(\alpha-\widehat{\alpha})\right|^{-(\nu+r) / 2} \\
= & \left|\left(\beta^{\prime} S_{11} \beta\right)^{-1}+T^{-1}(\alpha-\widehat{\alpha}) S^{-1}(\alpha-\widehat{\alpha})^{\prime}\right|^{-(\nu+r) / 2} \\
& \times|T S|^{-(\nu+r) / 2}\left|\beta^{\prime} S_{11} \beta\right|^{-(\nu+r) / 2}
\end{aligned}
$$

Let $\sigma_{22}$ denote the last $p_{2}$ rows and columns of $T S$ and partition $T S$ as

$$
T S=\left[\begin{array}{ll}
S_{11} & S_{12} \\
S_{21} & \sigma_{22}
\end{array}\right]
$$

Next, denote the $p_{2} \times p_{2}$ matrix made up of the last $p_{2}$ rows and columns of $S_{00}$ by $S_{00,22}$, and note that $\sigma_{22}=S_{00,22}-\widehat{\alpha}_{2}^{\prime} \beta^{\prime} S_{11} \beta \widehat{\alpha}_{2}$. Next, we integrate with respect to $a_{1}$. The conditional distribution of $a \mid \beta$ is

$$
\begin{aligned}
p(\alpha \mid \beta, y) \propto & \left|\left(\beta^{\prime} S_{11} \beta\right)^{-1}+T^{-1}(\alpha-\widehat{\alpha}) S^{-1}(\alpha-\widehat{\alpha})^{\prime}\right|^{-(\nu+r) / 2} \\
= & \left|\left(\beta^{\prime} S_{11} \beta\right)^{-1}+g_{1}\left(\alpha_{2}\right)+g_{2}\left(\alpha_{1}\right)\right|^{-(\nu+r) / 2} \\
& \text { where } \\
g_{1}\left(\alpha_{2}\right)= & \left(\alpha_{2}-\widehat{\alpha}_{2}\right) \sigma_{22}^{-1}\left(\alpha_{2}-\widehat{\alpha}_{2}\right)^{\prime} \\
g_{2}\left(a_{1}\right)= & \left(\alpha_{1}-\widehat{\Delta}\right)\left(S_{11}-S_{12} \sigma_{22}^{-1} S_{21}\right)^{-1}\left(\alpha_{1}-\widehat{\Delta}\right)^{\prime} .
\end{aligned}
$$

Integrating with respect to $\alpha_{1}$ gives us the marginal distribution of $\left(\alpha_{2}, \beta\right)$ as

$$
\begin{aligned}
p\left(\alpha_{2}, \beta \mid r, y\right) \propto & \left|\sigma_{22}+\left(\alpha_{2}-\widehat{\alpha}_{2}\right)^{\prime} \beta^{\prime} S_{11} \beta\left(\alpha_{2}-\widehat{\alpha}_{2}\right)\right|^{-\left(\nu-p+r+p_{2}\right) / 2} \\
& \times\left|\beta^{\prime} S_{11} \beta\right|^{-\left(p-p_{2}\right) 2}\left|\sigma_{22}\right|^{\left(\nu-p+p_{2}\right) / 2}|S|^{-\nu / 2} .
\end{aligned}
$$

Since $S_{00,22}=\sigma_{22}+\widehat{\alpha}_{2}^{\prime} \beta^{\prime} S_{11} \beta \widehat{\alpha}_{2}=\sigma_{22}+S_{01,2} \beta\left(\beta^{\prime} S_{11} \beta\right)^{-1} \beta^{\prime} S_{10,2}$, then evaluating this expression at $\alpha_{2}=0$ and rearranging we have

$$
p\left(\beta \mid \alpha_{2}=0, r, y\right) \propto\left|\beta^{\prime} D \beta\right|^{-\nu / 2}\left|\beta^{\prime} D_{0,2} \beta\right|^{\left(\nu-p+p_{2}\right) / 2}
$$

where $D_{0,2}=S_{11}-S_{10,2} S_{00,22}^{-1} S_{01,2}, S_{10,2}=S_{01,2}^{\prime}$ is the last $p_{2}$ rows of $S_{10}$. If we partition $D$ and $D_{0,2}$ conformably as

$$
D=\left[\begin{array}{cc}
D_{11} & D_{12} \\
D_{21} & d
\end{array}\right] \text { and } D_{0,2}=\left[\begin{array}{cc}
\Delta_{11} & \Delta_{12} \\
\Delta_{21} & \delta
\end{array}\right]
$$

use the linear restrictions $\beta=\left[\begin{array}{ll}\beta_{2}^{\prime} & I_{r}\end{array}\right]_{\widetilde{\prime}}^{\prime}$, then let $d_{s}=d-D_{21} D_{11}^{-1} D_{12}$, $\delta_{s}=\delta-\Delta_{21} \Delta_{11}^{-1} \Delta_{12}, \widehat{\beta}_{2}=D_{11}^{-1} D_{12}$ and $\widetilde{\beta}_{2}=\Delta_{11}^{-1} \Delta_{12}$,

$$
\begin{aligned}
p\left(\beta_{2} \mid \alpha_{2}=0, r, y\right) \propto & \left|d_{s}+\left(\beta_{2}-\widehat{\beta}_{2}\right)^{\prime} D_{11}\left(\beta_{2}-\widehat{\beta}_{2}\right)\right|^{-l_{0}} \\
& \times\left|\delta_{s}+\left(\beta_{2}-\widetilde{\beta}_{2}\right)^{\prime} \Delta_{11}\left(\beta_{2}-\widetilde{\beta}_{2}\right)\right|^{l_{1}}
\end{aligned}
$$

Thus we have the $1-1$ poly-t form for the posterior of $\beta \mid \alpha_{2}=0$. As the posterior is integrable only if $2\left(l_{0}-l_{1}\right)-(p-r)>0$. In this case, then, since $p_{2}=r$

$$
\gamma=2\left(l_{0}-l_{1}\right)-(p-r)=\nu-\nu+p-p_{2}-p+r=0
$$

and the posterior is clearly not integrable. Note that is is possible to take $p_{2}>r$ provided $p_{1}>p-r$. In this case $\gamma=r-p_{2}<0$, again producing an improper posterior.

Taking strong exogeneity with respect to $\beta$ will result in $p_{2}$ being replaced by $k_{2}=p_{2}+l p$ giving

$$
\begin{aligned}
2\left(l_{0}-l_{1}\right)-(p-r) & =\nu-\nu+p-k_{2}-p+r \\
& =-p_{2}-l p+r<0
\end{aligned}
$$

and the posterior is not proper in any situation.

## 10 Figures.



Figure 1: Higher posterior density regions for the impulse response of relative UK prices $\left(p_{t}-p_{t}^{*}\right)$ to a shock in oil prices. The $x$-axis spans zero to sixty months.


Figure 2: Higher posterior density regions for the impulse response of relative UK inflation $\left(\Delta p_{t}\right)$ to a shock in oil prices. The $x$-axis spans zero to sixty months.


[^0]:    ${ }^{1}$ Note that we do not use linear identifying restrictions (or normalisation) for the vector $\beta_{2}^{c}$, in which coefficients must be estimated. Instead, as discussed below, we identify $\varphi$ by nonlinear restrictions of the form $\varphi^{\prime} \varphi=1$. We do this to simplify estimation, and to avoid the potential problem that the posterior may have no moments and possibly be improper, particularly when we impose exogeneity.
    ${ }^{2}$ The concept of strict exogeneity has been criticised (Engle, Hendry and Richard 1983 and Hendry 1995) for introducing ambiguity of interpretation. The concepts of weak, strong and super exogeneity do, however, have clear interpretations and implications. Therefore, it is fortunate that in making oil prices strictly exogenous, Garratt et al. in fact make them weakly exogenous with respect to $\beta^{+}$and $\alpha$. The weak exogeneity of oil prices implies $\alpha_{2}=0$.

[^1]:    ${ }^{3}$ We note that the sums, $\Sigma$, in (5.23) of James (1954) should be products, $\Pi$.

