

# Efficient Estimation of Structural VARMA with Stochastic Volatility

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Macroeconomists and policy makers use time series models for structural inference as an alternative to DSGEs.

- robust to explicit theoretical modeling assumptions
- guide development of better DSGEs

Main instrument for inference is impulse response functions and forecast error variance decompositions.

- closely linked to the infinite MA representation (i.e. Wold representation) for multivariate time series data

Key challenge is identifying impulse response functions / structural shocks.

Identification is **not** only about **static orthogonal rotations**.

If errors are **Gaussian white noise**, then the characteristic roots of the infinite MA representation are **not identified**.

The **two** identification problems can be treated **separately**.

The equilibrium of a general class of DSGE models can be represented by a **finite order VARMA** (Fernandez-Villaverde, *et al.*, 2007, AER).

If the characteristic MA roots are all outside the unit circle, the representation is **fundamental**.

- then, we can obtain a VAR approximation by inverting the MA part and truncating lags (SVAR approach)

Many DSGEs (e.g. permanent income model of Hansen *et al.*, 1991) lead to clearly **non-fundamental** representations, where at least one root is inside the unit circle.

But theory is rarely precise on **which** non-fundamental representation (i.e. combination of roots inside the unit circle) is appropriate.

# Summarizing the Issue

Fundamental and non-fundamental representations are **observationally equivalent** under **Gaussian errors**.

Standard **SVAR** approach only identifies **fundamental** structural shocks, which is often inconsistent with theory.

Lippi and Reichlin (1994, JoE) argue IRFs and FEVDs should be based on **set identification**.

- difficult with SVARs because information about characteristic MA roots is lost when VAR lags are truncated
- even with VARMA, can be computationally infeasible

Observational equivalence can be eliminated with non-Gaussian errors.

If errors are **i.i.d. non-Gaussian**, then infinite MA representations are **unique** up to scaling, order of shocks, and time shifts (Chan *et al.*, 2006, Biometrika, Gouriéroux *et al.*, 2017, WP).

Therefore, specifying **i.i.d. non-Gaussian** errors simultaneously identifies both:

- 1 the fundamental or one of the non-fundamental representations,
- 2 the structural representation (i.e. orthogonal rotations of errors no longer observationally equivalent).

Could this be a viable solution?

# Problems With Non-Gaussian Errors

Specifying **i.i.d. non-Gaussian** errors entails a number of practical problems:

- 1 Different distributions lead to different representations being identified, so the question becomes **which** non-Gaussian distribution should be employed?
- 2 VARMA with non-Gaussian distribution are highly **non-linear** in parameters, and **no feasible** computation methodology is currently available.
- 3 **Statistically identified** structural shocks are difficult to interpret.

# Proposed Alternative

We argue that an alternative approach is to model errors as **conditionally Gaussian** by introducing **heteroskedastic** errors.

A large literature (Clark and Ravazzolo, 2014, JAE; Carriero *et al.*, 2016, JBES, Chan and Eisenstat, 2018, JAE) has consistently demonstrated that allowing for stochastic volatility in VARs is crucial in modeling macroeconomic time series.

Therefore, it is a natural extension to model VARMA with **Gaussian errors** and **stochastic volatility**.

We prove that under mild regularity conditions, we thus can obtain infinite MA representations that are **unique up to static orthogonal rotations**.



# Benefits of VARMA With Stochastic Volatility

The main advantages the proposed approach are:

- 1 conditionally Gaussian errors means **computation** is non-trivial but **feasible**: we develop an efficient MCMC algorithm in a Bayesian state-space setting,
- 2 characteristic MA **roots are identified**, but **not the structural representation**,
- 3 **structural shocks** can be identified using standard **theory-driven restrictions** as in the SVAR approach,
- 4 alternatives stochastic volatility specifications can be evaluated using Bayesian model comparison methods.

# Structural VARMA

Consider an  $n \times 1$  vector of observations  $\mathbf{y}_t$ , modeled as a VARMA( $p, q$ ) process:

$$\mathbf{B}(L)\mathbf{y}_t = \mathbf{A}(L)\boldsymbol{\varepsilon}_t, \quad \boldsymbol{\varepsilon}_t \sim \mathcal{N}(0, \mathbf{I}_n).$$

where  $\mathbf{A}(L)$  and  $\mathbf{B}(L)$  are polynomial matrices in the **back shift operator**  $L$  and  $\mathbf{B}_0 = \mathbf{I}_n$ .

Assuming  $\mathbf{B}(z) \neq 0$  for all  $z \in \mathbb{C}$ ,  $|z| \leq 1$ , the infinite MA representation is

$$\mathbf{y}_t = \boldsymbol{\Psi}(L)\boldsymbol{\varepsilon}_t, \quad \boldsymbol{\Psi}(L) = \mathbf{B}(L)^{-1}\mathbf{A}(L). \quad (1)$$

The identification issue we focus on arises from the fact that there exist many  $\mathbf{D}(z)$  such that  $\mathbf{D}(z)\mathbf{D}(z^{-1})' = \mathbf{I}_n$  and

$$\mathbf{y}_t = \tilde{\boldsymbol{\Psi}}(L)\tilde{\boldsymbol{\varepsilon}}_t, \quad \tilde{\boldsymbol{\Psi}}(L) = \boldsymbol{\Psi}(L)\mathbf{D}(L), \quad \tilde{\boldsymbol{\varepsilon}}_t = \mathbf{D}(L^{-1})'\boldsymbol{\varepsilon}_t, \quad (2)$$

where  $L^{-1}$  is the **forward shift operator** and  $\tilde{\boldsymbol{\varepsilon}}_t \sim \mathcal{N}(0, \mathbf{I}_n)$ .

# Observational Equivalence

If  $\mathbf{D}(z) = \mathbf{D}$ ,  $\mathbf{D}\mathbf{D}' = \mathbf{I}_n$ , we obtain the standard case where (1) is **observationally equivalent** to (2) due static orthogonal rotations.

Even if  $\Psi_0 = \tilde{\Psi}_0$ , there exist  $\mathbf{D}(z)$  such that  $\det \Psi(z) = 0$  implies either  $\det \tilde{\Psi}(z) = 0$  or  $\tilde{\Psi}(z^{-1}) = 0$  and (2) is an observationally equivalent VARMA( $p, q$ ) to the one in (1).

- $\mathbf{D}(z)$  is a **Blaschke matrix** that “flips” some combination of roots of  $\det \Psi(z)$  inside /outside the unit circle
- there are up to  $2^{nq}$  countable, observationally equivalent representations generated this way

**SVAR** approach only considers the **fundamental** representation and ignores the (up to  $2^{nq} - 1$ ) non-fundamental ones.

- this is inconsistent with theory in many settings

# Proposed Solutions to Observational Equivalence

Lippi and Reichlin (1994, JoE) suggest set identification of IRFs and FEVDs based on all fundamental and non-fundamental representations.

- becomes computationally infeasible as  $n$ ,  $q$  increase
- inference typically imprecise

If structural shocks are **i.i.d. not Gaussian** then all observational equivalence vanishes: (1) can be made unique with trivial restrictions on  $\Psi_0$  (i.e., nonzero, non-decreasing diagonal elements).

If reduced form innovations are **conditionally Gaussian** with heteroskedastic errors, then the reduced form representation is identified up to static orthogonal rotations, and the structural representation is identified by theory-driven restrictions on  $\Psi_0$ .

# Multivariate Volatility Process

Let  $\{\Sigma_t : t \in \mathbb{Z}\}$  be a stochastic process where:

- each  $\Sigma_t$  a positive definite symmetric matrix,
- $\|\Sigma_t\| \leq \varsigma < \infty$  almost surely for all  $t \in \mathbb{Z}$ ,
- $\{\Sigma_t\}$  is a weakly stationary process with an absolutely integrable inverse spectral density.

Upper bound and stationarity are **not necessary** and can be replaced other regularity conditions.

- Upper bound allows for clearer proof and is not restrictive in practice since it can be arbitrarily large (we only need it to be finite).

# Gaussian Scale Mixture Process

The Gaussian scale mixture process (GSMP) is given by  $\{\mathbf{u}_t : t \in \mathbb{Z}\}$  such that

$$(\mathbf{u}_t | \{\boldsymbol{\Sigma}_t\}) \stackrel{iid}{\sim} \mathcal{N}(0, \boldsymbol{\Sigma}_t). \quad (3)$$

Marginally of  $\{\boldsymbol{\Sigma}_t\}$ ,  $\mathbf{u}_t$  and  $\mathbf{u}_{t-j}$  are **mean independent** for all  $j$ .

Marginally of  $\{\boldsymbol{\Sigma}_t\}$ ,  $\mathbf{u}_t$  and  $\mathbf{u}_{t-j}$  are **stochastically independent** if and only if  $\boldsymbol{\Sigma}_t$  and  $\boldsymbol{\Sigma}_{t-j}$  are stochastically independent.

For any non-singular  $\mathbf{D}$ ,  $\tilde{\mathbf{u}}_t = \mathbf{D}\mathbf{u}_t$  is also a GSMP.

# VARMA With GSMP Errors

Consider the **reduced form** VARMA( $p, q$ ):

$$\mathbf{y}_t = \mathbf{B}(L)^{-1} \boldsymbol{\Theta}(L) \mathbf{u}_t = \boldsymbol{\Phi}(L) \mathbf{u}_t, \quad (\mathbf{u}_t | \boldsymbol{\Sigma}_t) \sim \mathcal{N}(0, \boldsymbol{\Sigma}_t), \quad (4)$$

where  $\{\boldsymbol{\Sigma}_t\}$  satisfies the previous assumptions and  $\boldsymbol{\Theta}_0 = \mathbf{I}_n$  and  $\boldsymbol{\Phi}_0 = \mathbf{I}_n$ .

A **structural form** is obtained by:

$$\mathbf{y}_t = \mathbf{B}(L)^{-1} \boldsymbol{\Theta}(L) \boldsymbol{\Sigma}_t^{\frac{1}{2}} \mathbf{Q}_t \boldsymbol{\varepsilon}_t = \boldsymbol{\Psi}_t(L) \boldsymbol{\varepsilon}_t, \quad \boldsymbol{\varepsilon}_t \sim \mathcal{N}(0, \mathbf{I}_n),$$

where  $\mathbf{Q}_t \mathbf{Q}_t' = \mathbf{I}_n$ .

Does there exist an alternative VARMA( $p, q$ ) process

$$\tilde{\mathbf{y}}_t = \mathbf{B}(L)^{-1} \tilde{\boldsymbol{\Theta}}(L) \tilde{\mathbf{u}}_t = \tilde{\boldsymbol{\Phi}}(L) \tilde{\mathbf{u}}_t, \quad (\tilde{\mathbf{u}}_t | \tilde{\boldsymbol{\Sigma}}_t) \sim \mathcal{N}(0, \tilde{\boldsymbol{\Sigma}}_t), \quad (5)$$

such that  $\{\mathbf{y}_t\}$  and  $\{\tilde{\mathbf{y}}_t\}$  are **observationally equivalent**?

# Uniqueness Result

## Theorem

Assume  $\{\Sigma_t\}$  satisfies regularity conditions. The process  $\{\tilde{\mathbf{y}}_t\}$  given by (5) is observationally equivalent to the process  $\{\mathbf{y}_t\}$  given by (4) if and only if  $\{\tilde{\mathbf{u}}_t\} \stackrel{d}{=} \{\mathbf{u}_t\}$  and  $\tilde{\Phi}(L) = \Phi(L)$ , where  $\stackrel{d}{=}$  denotes equivalence in distribution.

Two comments regarding the proof:

- 1 Main challenge:  $\mathbf{u}_t$  not stochastically independent.
- 2 Key insight: if  $\tilde{\Phi}(z) \neq \Phi(z)$ , then  $\mathbf{D}(z) = \tilde{\Phi}(z)^{-1}\Phi(z)$  yields  $\tilde{\mathbf{u}}_t = \mathbf{D}(L)\mathbf{u}_t$  such that  $\tilde{\mathbf{u}}_t$  is not mean-independent of  $\tilde{\mathbf{u}}_{t-j}$  for some  $j$ ; therefore,  $\tilde{\mathbf{u}}_t$  is not a GSMP.



The data uniquely identifies either a fundamental or one of the non-fundamental representations.

VARs with time-varying volatilities (i.e. Sims and Zha, 2006, AER; Chan and Eisenstat, 2018, JAE) are **over-identified** because they enforce **fundamentalness**.

- this holds even for models with log-volatilities modeled as random walks (i.e. non-stationary volatility process)
- can be shown to hold for TVP-VARs with time-varying coefficients (though not a trivial extension)

We can **test** if the fundamentalness restriction is **supported by the data**.

We consider a Bayesian state-space model given by:

$$\begin{aligned} \mathbf{B}_0 \mathbf{y}_t &= \boldsymbol{\mu} + \mathbf{B}_1 \mathbf{y}_{t-1} + \cdots + \mathbf{B}_p \mathbf{y}_{t-p} \\ &\quad + \boldsymbol{\Theta}_1 \mathbf{u}_{t-1} + \cdots + \boldsymbol{\Theta}_q \mathbf{u}_{t-q}, \quad (\mathbf{u}_t \mid \boldsymbol{\Sigma}_t) \sim \mathcal{N}(0, \boldsymbol{\Sigma}_t), \\ (\boldsymbol{\Sigma}_t^{-1} \mid \boldsymbol{\Sigma}_{t-1}) &\sim \mathcal{IW} \left( \nu, \frac{1}{\nu} \boldsymbol{\Sigma}_{t-1}^{-\delta} \right), \end{aligned}$$

where  $\nu > n$  and  $|\delta| < 1$ .

$\mathcal{W}(\cdot)$  denotes the Wishart distribution; the state equation defines the Wishart process of Philipov and Glickman (2006, JBES).

Adding the constraint  $\|\boldsymbol{\Sigma}_t\| < \varsigma$  ensures that this volatility process satisfies the regularity conditions and has nice properties.

We set  $\mathbf{B}_0$  to be a lower triangular matrix, with ones on the diagonals.

The free elements of  $\mathbf{B}_l$  ( $l = 0, \dots, p$ ) and  $\Theta_l$  ( $l = 1, \dots, q$ ) are modeled with **Stochastic Search Variable Selection** priors (George *et al.*, 2008, JoE):

$$\begin{aligned}(B_{l,ij} | \gamma_{l,ij}^B) &\sim (1 - \gamma_{l,ij}^B) \mathcal{N}(0, \tau_{0,l,ij}^2) + \gamma_{l,ij}^B \mathcal{N}(0, \tau_{1,l,ij}^2), \\ (\Theta_{l,ij} | \gamma_{l,ij}^\Theta) &\sim (1 - \gamma_{l,ij}^\Theta) \mathcal{N}(0, \tau_{0,l,ij}^2) + \gamma_{l,ij}^\Theta \mathcal{N}(0, \tau_{1,l,ij}^2),\end{aligned}$$

where  $\gamma_{l,ij}^B \in \{0, 1\}$ ,  $\gamma_{l,ij}^\Theta \in \{0, 1\}$  and  $\tau_{0,l,ij}^2 \ll \tau_{1,l,ij}^2$ .

This approximates the *Echelon form* for unique VARMA specifications (see Chan *et al.*, 2016, JoE).

- It is needed because  $\mathbf{B}(L)$  and  $\Theta(L)$  have redundant coefficients when, e.g.,  $[\mathbf{B}_p; \Theta_q]$  has row rank less than  $n$ .



# Log Likelihood

The log likelihood for the first  $t$  observations is given by:

$$\ln p(\mathbf{y}_1, \dots, \mathbf{y}_t | \cdot) = -\frac{1}{2} \ln(2\pi) - \frac{1}{2} |\tilde{\Theta}_{1:t} \tilde{\Sigma}_{1:t} \tilde{\Theta}'_{1:t}| - \frac{1}{2} \tilde{\mathbf{v}}'_{1:t} \left( \tilde{\Theta}_{1:t} \tilde{\Sigma}_{1:t} \tilde{\Theta}'_{1:t} \right)^{-1} \tilde{\mathbf{v}}_{1:t}. \quad (6)$$

The full log likelihood is obtained by replacing  $t$  with  $T$ .

The term  $\left( \tilde{\Theta}_{1:t} \tilde{\Sigma}_{1:t} \tilde{\Theta}'_{1:t} \right)^{-1}$  complicates computation because it is nonlinear in  $\Theta_l$  and  $\Sigma_t$ , but the quadratic term is **not expensive** to compute.

If the  $\Theta(L)$  was **fundamental**, we could set  $\Sigma_{1-q} = \Sigma_0 = 0$  and simplify computation (e.g. Chan, 2013, JoE).

- We cannot take this approach here because  $\Theta(L)$  is possibly **non-fundamental**; we instead treat  $\Sigma_{1-q}, \dots, \Sigma_0$  as parameters.

# Parameters and Priors

The parameters in this specification are the free elements in  $\mathbf{B}_0$ , along with  $\mathbf{B}_1, \dots, \mathbf{B}_p$ ,  $\boldsymbol{\mu}$ ,  $\boldsymbol{\Theta}_1, \dots, \boldsymbol{\Theta}_q$ ,  $\boldsymbol{\Sigma}_{1-q}, \dots, \boldsymbol{\Sigma}_0$ ,  $\{\gamma_{l,ij}^B\}$ ,  $\{\gamma_{l,ij}^\Theta\}$ ,  $\nu$  and  $\delta$ .

Along with the priors already discussed, we assign the following:

$$\Pr(\gamma_{l,ij}^B = 0) = 0.5,$$

$$\Pr(\gamma_{l,ij}^\Theta = 0) = 0.5,$$

$$\boldsymbol{\mu} \sim \mathcal{N}(0, 10),$$

$$\boldsymbol{\Sigma}_{q-1}^{-1} \sim \mathcal{W}(n+1, \mathbf{I}_n),$$

$$\nu - n \sim \mathcal{G}(1, 1),$$

$$\delta \sim \mathcal{U}(-1, 1),$$

where  $\mathcal{G}$  denotes the **Gamma** distribution and  $\mathcal{U}$  the **Uniform** distribution.

The MCMC algorithm involves iterating seven sampling steps:

- 1  $(\boldsymbol{\mu}, \mathbf{B}_0, \dots, \mathbf{B}_p \mid \boldsymbol{\Theta}_1, \dots, \boldsymbol{\Theta}_q, \{\boldsymbol{\Sigma}_t\}, \{\gamma_{l,ij}^B\})$ ,
- 2  $(\boldsymbol{\Theta}_1, \dots, \boldsymbol{\Theta}_q \mid \boldsymbol{\mu}, \mathbf{B}_0, \dots, \mathbf{B}_p, \{\boldsymbol{\Sigma}_t\}, \{\gamma_{l,ij}^\Theta\})$ ,
- 3  $(\boldsymbol{\Sigma}_{1-q}, \dots, \boldsymbol{\Sigma}_T \mid \boldsymbol{\mu}, \mathbf{B}_0, \dots, \mathbf{B}_p, \boldsymbol{\Theta}_1, \dots, \boldsymbol{\Theta}_q, \mathbf{S}, \nu, \delta)$ ,
- 4  $(\gamma_{l,ij}^B \mid B_{l,ij})$  for  $l = 0, \dots, p$ ,  $i, j = 1, \dots, n$  (and  $i > j$  for  $l = 0$ ),
- 5  $(\gamma_{l,ij}^\Theta \mid \Theta_{l,ij})$  for  $l = 1, \dots, q$ ,  $i, j = 1, \dots, n$ ,
- 6  $(\nu \mid \{\boldsymbol{\Sigma}_t\}, \delta)$ ,
- 7  $(\delta \mid \{\boldsymbol{\Sigma}_t\}, \nu)$ .

# Sampling Implementation

Step 1 entails standard sampling from the multivariate Gaussian distribution.

Steps 4-5 are also standard in the SSVS literature.

Steps 6-7 are provided in Philipov and Glickman (2006, JBES).

Step 3 is efficiently implemented using the [particle Gibbs with ancestry sampling](#) as in Lindsten *et al.* (2014, JMLR).

Step 2 is the most challenging: we implement an [independence MH sampler](#) based on the [Whittle likelihood](#) approach of Dahlhaus (2002, AoS).



# Whittle Likelihood

Let  $\check{\Theta}_k = \sum_{l=0}^q \Theta_l e^{il\lambda_k}$ , where  $i = \sqrt{-1}$  and  $\lambda_k = 2\pi k/T$ .

Let  $\tilde{\mathbf{V}}_k = (\mathbf{v}_1 e^{-i\lambda_k}, \dots, \mathbf{v}_T e^{-iT\lambda_k})$  and  $\text{vec}(\tilde{\mathbf{W}}) = \mathbf{\Omega} \text{vec}(\check{\Theta}_k^{-1} \tilde{\mathbf{V}}_k)$ , where

$$\mathbf{\Omega} = \begin{pmatrix} \Sigma_1^{-1} & \Sigma_1^{-1} & \cdots & \cdots & \Sigma_{\lfloor \frac{1+T}{2} \rfloor}^{-1} \\ \Sigma_1^{-1} & \ddots & \ddots & \ddots & \vdots \\ \vdots & \ddots & \Sigma_{\lfloor \frac{r+c}{2} \rfloor}^{-1} & \ddots & \vdots \\ \vdots & \ddots & \ddots & \ddots & \vdots \\ \Sigma_{\lfloor \frac{1+T}{2} \rfloor}^{-1} & \cdots & \cdots & \cdots & \Sigma_T^{-1} \end{pmatrix}.$$

$\mathbf{\Omega}$  is  $nT \times nT$ , with  $r$  denoting the row,  $c$  the column and  $\lfloor \cdot \rfloor$  the floor operator.

# Whittle Likelihood

The Whittle **likelihood** is given by:

$$\ln p_W(\mathbf{v}; \boldsymbol{\theta}) = -\frac{1}{2T^2} \sum_{k=0}^{T-1} \left( (2n) \ln(2\pi) \sum_{t=1}^T |\check{\boldsymbol{\Theta}}_k \boldsymbol{\Sigma}_t \check{\boldsymbol{\Theta}}_k^*| + \frac{1}{2\pi} \text{vec}(\widetilde{\mathbf{W}}_k)^* \text{vec}(\widetilde{\mathbf{V}}_k) \right) \quad (7)$$

The **score** is given by:

$$\frac{d \ln p_W(\mathbf{v}; \boldsymbol{\theta})}{d\boldsymbol{\Theta}'_l} = \frac{1}{T} \sum_{k=0}^{T-1} \text{Re} \left[ \check{\boldsymbol{\Theta}}_k^{-1} + \frac{1}{2\pi T} \check{\boldsymbol{\Theta}}_k^{-1} \widetilde{\mathbf{V}}_k \widetilde{\mathbf{W}}_k^* \check{\boldsymbol{\Theta}}_k^{-1} \right]. \quad (8)$$

Both are relatively **easy to evaluate**, which means the Whittle likelihood is **easy to maximize**.

Dahlhaus (2002, AoS) proved that if  $\hat{\boldsymbol{\theta}}$  maximizes  $p_W(\mathbf{v}; \boldsymbol{\theta})$  then it **converges** to the MLE of the actual likelihood  $p(\mathbf{v}; \boldsymbol{\theta})$ , and it is **asymptotically efficient**.

Consequently, we set the proposal density to

$$\boldsymbol{\theta}^c \sim \mathcal{N} \left( \hat{\boldsymbol{\theta}}, -\hat{\mathbf{H}}_{\boldsymbol{\theta}}^{-1} \right),$$

where  $\hat{\boldsymbol{\theta}}$  maximizes  $\ln p(\boldsymbol{\theta}) + \ln p_W(\mathbf{v}; \boldsymbol{\theta})$  and

$$\hat{\mathbf{H}}_{\boldsymbol{\theta}} = \frac{d^2 \ln p(\boldsymbol{\theta})}{d\boldsymbol{\theta} d\boldsymbol{\theta}'} + \frac{d^2 \ln p_W(\mathbf{v}; \boldsymbol{\theta})}{d\boldsymbol{\theta} d\boldsymbol{\theta}'}$$

**Initial Monte Carlo** results that the algorithm is reasonably efficient with VARMA(2, 2) and up to  $n = 4$  variables.

**Extensive Monte Carlo** exercise is under way.

**Real data application** will re-examine SVARs for effects of **monetary policy shocks** and **news shocks**.