Normality tests for latent variables*

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Abstract

We exploit the rationale behind the Expectation Maximization algorithm to derive simple to implement and interpret LM normality tests for the innovations of the latent variables in linear state space models against generalized hyperbolic alternatives, including symmetric and asymmetric Student ts. We decompose our tests into third and fourth moment components, and obtain one-sided likelihood ratio analogues, whose asymptotic distribution we provide. When we apply our tests to a common trend model which combines the expenditure and income versions of US aggregate real output to improve its measurement, we reject normality if the sample period extends beyond the Great Moderation.

Keywords: Cointegration, Gross Domestic Product, Gross Domestic Income, Kurtosis, Kuhn-Tucker test, Skewness, Supremum test, Wiener-Kolmogorov-Kalman smoother.

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1 Introduction

Latent variable models that relate a set of observed variables to a meaningful set of unobserved influences are widely used in many applied fields. The list of empirical studies that make use of those models is vast. In this paper, we consider a classic application of signal extraction techniques whereby we obtain an improved aggregate (real) production series by combining its expenditure (GDP) and income (GDI) measures, which differ because they are constructed using largely independent data sources (see Landefeld, Seskin and Fraumeni (2008) for a review).

We will use this model in section 7, but in developing it, one particularly relevant decision we must make is the normality of the underlying variables, which implies the normality of the observed variables and justifies the use of the Kalman filter for inferring the true underlying output from its two measures. In contrast, if the innovations are not Gaussian, the Kalman filter only provides the best linear filter for the latent variable, which can be noticeably different from its conditional expectation. To illustrate this point, consider the simplest possible example in which a negatively skewed signal x is observed cloaked in some additive symmetric noise ϵ. As can be seen in Figure 1, the linear projection can display important biases relative to the conditional expectation of x given the observed series $y = x + \epsilon$. Intuitively, the conditional expectation takes into account that the asymmetry in x implies that large negative/positive realizations of y are more/less likely to result from the signal, while the linear projection assigns a constant fraction of y to x regardless.

The remarkable increase in computing power has made possible the implementation of simulation-based estimation and filtering techniques for non-Gaussian dynamic latent variable models (see e.g. Johannes and Polson (2009)). However, the majority of practitioners continue to rely on the Kalman filter, which is far simpler to implement and explain. Undoubtedly, those practitioners would benefit from the existence of diagnostics that could tell them the extent to which normality of the latent variables is at odds with the data. Although there are many readily available normality tests, they are designed to be directly applied to the observed variables in static models or their one-period ahead prediction errors in dynamic ones.

The objective of our paper is precisely to derive simple to implement and interpret tests for non-normality in all or a subset of (the innovations to) the state variables. We focus on Lagrange Multiplier (LM) tests, which only require estimation of the model under the null. As is well known, Likelihood ratio (LR), Wald and LM tests are asymptotically equivalent under

the null and sequences of local alternatives, and therefore they share their optimality properties. Aside from computational reasons, the advantage of LM tests is that rejections provide a clear indication of the specific directions along which modelling efforts should focus.

Nowadays, the computational advantages of LM tests might seem irrelevant, but in our case they are of first-order importance because the density function of the observed variables or their innovations is typically unknown when the distribution of the latent variables is not Gaussian, and in many cases it can only be approximated by simulation (see Durbin and Koopman (2012) for an extensive discussion in the context of dynamic models). As a result, the log-likelihood function under the alternative, its score and information matrix can seldom be obtained in closed form despite the fact that we can compute the true log-likelihood function under the Gaussian null. We overcome this stumbling block by exploiting what we call "the EM principle". Specifically, we generalize Louis' (1982) score formula in order to obtain the first derivatives of the log likelihood with respect to the shape parameters that characterize departures from normality. The EM algorithm studied in Dempster, Laird and Rubin (1977) is a well known procedure for obtaining maximum likelihood estimates in both static and dynamic latent variable models (see e.g. Rubin and Thayer (1982) or Watson and Engle (1983), respectively). However, to the best of our knowledge it has only been used for testing purposes by Fiorentini and Sentana (2015), who employ it to assess neglected serial dependence in non-Gaussian static factor models.

Our approach introduces a relatively minor complication: the influence functions that constitute the basis of our tests are serially correlated in dynamic models. In this regard, our methods are related to Bai and Ng (2005) and Bontemps and Meddahi (2005), who derive moment-based normality tests for a single observed variable or its innovations in potentially serially correlated contexts by relying on heteroskedastic and autocorrelation consistent estimators of the asymptotic variances. Nevertheless, we derive analytical expressions for the autocovariance matrices of the influence functions, which we would expect a priori to lead to more reliable finite sample sizes for our statistics than their non-parametric counterparts. For that reason, our approach is more closely related to Harvey and Koopman (1992), who apply standard univariate normality tests for observed variables to the smoothed values of the innovations in the underlying components of a univariate random walk plus noise model explicitly taking into account the serial correlation implied by the model for those estimates. Unlike us, though, none of those authors justify their procedures by appealing to the likelihood principle or consider multivariate models.

For most practical purposes, departures from normality can be attributed to two different sources: excess kurtosis and skewness. Although our EM-based LM approach can be applied far more generally, we follow Mencía and Sentana (2012) in considering Generalized Hyperbolic (GH) alternatives, which include the symmetric and asymmetric Student t, normal-gamma mixtures, hyperbolic, normal inverse Gaussian and symmetric and asymmetric Laplace distributions. The main advantage of these GH alternatives is that they lead to easy to interpret moment tests that focus on third and fourth moments. In particular, they coincide with the moments underlying the Jarque and Bera (1980) test in the univariate case. At the same time, the number of moments that are effectively tested in multivariate contexts is proportional to the number of series involved, unlike tests against Hermite expansions of the multivariate normal density, which suffer from the curse of dimensionality (see Amengual and Sentana (2015) for a comparison in the context of copulas). Importantly, we show that our tests are not affected by the sampling variability in the model parameters estimated under the null, so we can treat them as if they were known.

The rest of the paper is organized as follows. Section 2 describes the econometric model, as well as the GH alternatives. We derive our normality tests against the Student t first and the GH distribution later in sections 3 and 4, respectively. Then, in section 5 we illustrate our procedures with two popular examples, while in section 6 we discuss the results of our Monte Carlo experiments. Section 7 explores in detail the information about aggregate output in the GDP and GDI measures. Finally, we present our conclusions in section 8. The Supplemental Material contains proofs and provides additional results.

2 The model

2.1 Linear state-space models

A linear, time-invariant, parametric state-space model for a finite dimensional vector of N observed series, $\mathbf{y}_t$, can be recursively defined in the time domain by the system of stochastic difference equations

$$\mathbf{y}_t = \boldsymbol{\pi}(\boldsymbol{\theta}) + \mathbf{H}(\boldsymbol{\theta})\boldsymbol{\xi}_t \tag{1}$$

$$\boldsymbol{\xi}_t = \mathbf{F}(\boldsymbol{\theta})\boldsymbol{\xi}_{t-1} + \mathbf{M}(\boldsymbol{\theta})\boldsymbol{\varepsilon}_t \tag{2}$$

$$\boldsymbol{\varepsilon}_t|\mathcal{I}_{t-1}, \boldsymbol{\phi} \sim iid\ D(\mathbf{0}, \mathbf{I}_K, \boldsymbol{\varphi}) \tag{3}$$

where $\boldsymbol{\phi} = (\boldsymbol{\theta}', \boldsymbol{\varphi}')'$, $\boldsymbol{\theta} \in \Theta \subseteq \mathbb{R}^p$ is a vector of p first and second moment parameters, $\boldsymbol{\pi} : \Theta \to \mathbb{R}^N$ is the mean vector of the observed series, $\mathbf{H} : \Theta \to \mathbb{R}^{N \times M}$, $\mathbf{F} : \Theta \to \mathbb{R}^{M \times M}$ and $\mathbf{M} : \Theta \to \mathbb{R}^{M \times K}$ are matrix valued functions of coefficients, many of whose elements will typically be either 0 or 1, $\boldsymbol{\xi}_t$ is an M-dimensional vector of state variables, $\boldsymbol{\varepsilon}_t$ is a K-dimensional vector of standardized structural *iid* innovations driving those variables whose distribution depends on a vector of shape parameters $\boldsymbol{\varphi}$, and $\mathcal{I}_{t-1}$ is an information set that contains the values of $\mathbf{y}_t$ and $\boldsymbol{\xi}_t$ up to and including $t-1$.

We assume that $N \leq K \leq M$ to avoid dynamic singularities. We also assume that the model above is correctly specified, in the sense that there is some $\boldsymbol{\theta}$ for which (1) and (2) constitute the true data generating process of $\{\mathbf{y}_t, \boldsymbol{\xi}_t\}$. In this context, static models will be such that $\mathbf{F}(\boldsymbol{\theta}) = \mathbf{0}$ for all $\boldsymbol{\theta}$.

There are multiple alternative representations of state-space models,[1]but in this paper we follow the one in Harvey (1989), except that we have deliberately subsumed any possible error in the measurement equation (1) into the state vector so as to be able to test for normality not only in the minimal possible set of state variables but also in the measurement errors. For that reason, equations (1) and (2) closely resemble the usual state representation in the engineering literature, in which the elements of $\boldsymbol{\varepsilon}_t$ would be regarded as control variables (see Anderson and Moore (1979)). For ease of exposition, we do not look at models with exogenous regressors or those in which some of the system matrices are deterministic functions of time or observable predertermined variables.[2]

We assume without loss of generality that the columns of the matrix $\mathbf{M}(\boldsymbol{\theta})$ are linearly independent so that there are no redundant elements in $\boldsymbol{\varepsilon}_t$. Typically, $\mathbf{M}(\boldsymbol{\theta})$ will be a selection matrix whose columns are (proportional to) vectors of the M-imensional canonical basis, but in principle they could be different. As a result, we can uniquely recover $\boldsymbol{\varepsilon}_t$ from $\boldsymbol{\xi}_t$ as

$$\boldsymbol{\varepsilon}_t = \mathbf{M}^{+}(\boldsymbol{\theta})[\mathbf{I}_M - \mathbf{F}(\boldsymbol{\theta})L]\boldsymbol{\xi}_t, \tag{4}$$

where $\mathbf{M}^{+}(\boldsymbol{\theta}) = [\mathbf{M}'(\boldsymbol{\theta})\mathbf{M}(\boldsymbol{\theta})]^{-1}\mathbf{M}'(\boldsymbol{\theta})$ denotes the Moore-Penrose inverse of $\mathbf{M}(\boldsymbol{\theta})$. We also assume no linear combination of all the observables $\{\mathbf{y}_t\}$ has zero variance.

[1]For example, Durbin and Koopman (2012) shift the transition equation (2) forward by one period, as in Anderson and Moore (1979), and include measurement errors in (1), which they assume are orthogonal to the innovations in the state variables. On the other hand, Komunjer and Ng (2011) substitute the transition equation (2) into the measurement equation (1), thereby creating an alternative measurement equation whose innovations are perfectly correlated with the innovations in the transition equation.

[2]Minor changes to the testing procedures we propose will render them applicable to those situations.

Finally, we assume that the researcher makes sure that the model parameters $\boldsymbol{\theta}$ are identified before estimating the model, which often requires restrictions on the system matrices (see e.g. section 2.3 of Fiorentini, Galesi and Sentana (2016) and the references therein). These assumptions are satisfied in virtually all applications of state-space models.

2.2 Null and alternative hypotheses

In section 4 we derive computationally simple tests of the null hypothesis that the structural innovations are Gaussian against the alternative that they follow a member of the *GH* family of distributions introduced by Barndorff-Nielsen (1977) and studied in detail by Blæsild (1981). This is a rather flexible family of multivariate distributions that nests not only the normal and Student t but also many other examples such as the asymmetric Student t, the hyperbolic and normal inverse Gaussian distributions, as well as symmetric and asymmetric versions of the normal-gamma mixture and Laplace. As we mentioned in the introduction, the main advantages of these *GH* alternatives is that they lead to easy to interpret moment tests that focus on third and fourth moments, but in such a way that the number of conditions which are effectively tested is proportional to the number of series involved. Nevertheless, for clarity we first present the relevant results regarding testing multivariate normal versus multivariate Student t innovations in the next section, and then generalize them to the GH case.

In many applications, the researcher may only be interested in testing whether the source of non-normality comes from a subset of the underlying components, which have some meaningful interpretation. In our empirical application, for example, it matters whether the potential non-normality is a feature of the true GDP or its measurement errors. Given that we can always re-order the vector of structural innovations $\boldsymbol{\varepsilon}_t$ and postmultiply the matrix $\mathbf{M}(\boldsymbol{\theta})$ by a permutation matrix, without loss of generality we can assume that the non-Gaussian distribution is confined to the first $R \leq K$ innovations under the alternative. Henceforth, we refer to the relevant components as $\boldsymbol{\varepsilon}_t^{\mathrm{GH}} = \mathbf{S}_{RK}\boldsymbol{\varepsilon}_t$, with $\mathbf{S}_{RK} = (\mathbf{I}_R, \mathbf{0})$, and to the remaining ones as $\boldsymbol{\varepsilon}_t^{\mathrm{N}}$.

As a result, we explicitly consider the following alternative hypotheses:

1. The joint distribution of all structural innovations is *GH*: $H_J : \boldsymbol{\varepsilon}_t \sim GH_K(\eta, \psi, \boldsymbol{\beta})$;

2. The joint distribution of the first R structural innovations is *GH* while the rest are Gaussian: $H_S : \boldsymbol{\varepsilon}_t^{\mathrm{GH}} \sim GH_R(\eta, \psi, \boldsymbol{\beta})$, $\boldsymbol{\varepsilon}_t^{\mathrm{N}} \sim N_{K-R}(\mathbf{0}, \mathbf{I}_{K-R})$.[3]

[3] We might also envisage an alternative situation in which the elements of $\boldsymbol{\varepsilon}_t$ are cross-sectionally independent but non-Gaussian, see Almuzara, Amengual and Sentana (2017) for details.

3 Multivariate normal versus Student t innovations

The multivariate Student t distribution generalizes the multivariate normal distribution through a single additional parameter ν, which is usually known as the degrees of freedom. For convenience, we work with its reciprocal, η, so that Gaussianity requires $\eta \to 0^+$.

3.1 The score under Gaussianity

LM tests are usually obtained from the score associated to the (marginal) likelihood function of the observed variables, $f_{\mathbf{Y}}(\mathbf{Y}_T|\boldsymbol{\phi})$, with $\mathbf{Y}_T = vec(\mathbf{y}_1, \ldots, \mathbf{y}_T)$, evaluated under the Gaussian null. Unfortunately, the functional form of $f_{\mathbf{Y}}(\mathbf{Y}_T|\boldsymbol{\phi})$ is generally unknown under the alternative, and the same is true of its score vector evaluated under the null despite the fact that we can easily compute the Gaussian likelihood function. For that reason, we rely on a variant of Louis' (1982) score formula, which is based on the so-called "EM principle"; see also Ruud (1991) and Tanner (1996).

Initially, we assume $\boldsymbol{\theta}$ is fixed and known, and later on we consider the effect of estimation of mean and variance parameters. Formally, the EM principle applied to this context says

Proposition 1 *Let $f_{\mathbf{E}}(\mathbf{E}_T|\boldsymbol{\varphi})$ denote the density of $\mathbf{E}_T = vec(\boldsymbol{\varepsilon}_1, \ldots, \boldsymbol{\varepsilon}_T)$ with respect to Lebesgue measure on $\mathbb{R}^{KT}$, which is continuous in $\mathbf{E}_T$ and differentiable in $\boldsymbol{\varphi}$. Then,*

$$\frac{\partial \ln f_{\mathbf{Y}}(\mathbf{Y}_T|\boldsymbol{\phi})}{\partial \boldsymbol{\varphi}} = E\left[\frac{\partial \ln f_{\mathbf{E}}(\mathbf{E}_T|\boldsymbol{\varphi})}{\partial \boldsymbol{\varphi}} \middle| \mathbf{Y}_T, \boldsymbol{\phi} \right]. \tag{5}$$

Remark 1: The identity (5) is different from Louis' (1982) original formula, in that it applies when $K > N$ but only for a subset of the parameters.

Remark 2: Let $f_{(\mathbf{Y},\mathbf{E})}(\mathbf{Y}_T, \mathbf{E}_T|\boldsymbol{\phi})$ denote the joint likelihood function for both observed variables $\{\mathbf{y}_t\}$ and structural innovations $\{\boldsymbol{\varepsilon}_t\}$ of model (1)–(2) for a sample of size T. This joint density will necessarily be singular in linear state-space models because of the restrictions the observed variables $\mathbf{Y}_T$ place on the latent $\mathbf{E}_T$. The same comment applies to the conditional likelihood function of the latent variables given the observed ones, $f_{\mathbf{E}|\mathbf{Y}}(\mathbf{E}_T|\mathbf{Y}_T, \boldsymbol{\phi})$, which will usually be defined over a manifold of smaller dimension. Since the Kullback inequality implies that $E\left[\partial \ln f_{\mathbf{E}|\mathbf{Y}}(\mathbf{E}_T|\mathbf{Y}_T, \boldsymbol{\phi})/\partial \boldsymbol{\varphi} | \mathbf{Y}_T, \boldsymbol{\phi}\right] = \mathbf{0}$, it follows that we can obtain $\partial \ln f_{\mathbf{Y}}(\mathbf{Y}_T|\boldsymbol{\phi})/\partial \boldsymbol{\varphi}$ as the expected value of the unobservable score corresponding to $f_{(\mathbf{Y},\mathbf{E})}(\mathbf{Y}_T, \mathbf{E}_T|\boldsymbol{\phi})$ conditional on $\mathbf{Y}_T$ and $\boldsymbol{\phi}$. Therefore, an alternative formulation of (5) is

$$\frac{\partial \ln f_{\mathbf{Y}}(\mathbf{Y}_T|\boldsymbol{\phi})}{\partial \boldsymbol{\varphi}} = E\left[\frac{\partial \ln f_{(\mathbf{Y},\mathbf{E})}(\mathbf{Y}_T, \mathbf{E}_T|\boldsymbol{\varphi})}{\partial \boldsymbol{\varphi}} \middle| \mathbf{Y}_T, \boldsymbol{\phi} \right], \tag{6}$$

where

$$f_{(\mathbf{Y},\mathbf{E})}(\mathbf{Y}_T, \mathbf{E}_T|\boldsymbol{\phi}) = 1_{\{\mathbf{Y}_T = \mathbf{D}_T(\mathbf{E}_T,\boldsymbol{\theta})\}} f_{\mathbf{E}}(\mathbf{E}_T|\boldsymbol{\varphi}),$$

with $\mathbf{Y}_T = \mathbf{D}_T(\mathbf{E}_T, \boldsymbol{\theta})$ denoting the exact relationship between observed variables $\mathbf{Y}_T$ and innovations $\mathbf{E}_T$ implied by model (1)-(2) and $1_{\{.\}}$ the usual indicator function.

Importantly, while the mean and variance parameters enter in the indicator function, the shape parameters do not, so that the right hand side of (6) is well defined.[4]

Remark 3: One must be careful in applying the "EM principle" to the score with respect to $\boldsymbol{\theta}$. We deal with this situation in Proposition 3b and 6b (see the proof of Lemma 7 in the Supplemental Material A)

In the case of Student t innovations, we can use the expression provided in Fiorentini, Sentana and Calzolari (2003) for the score with respect to η under the Gaussian null:

$$\frac{\partial \ln f_{(\mathbf{Y},\mathbf{E})}(\mathbf{Y}_T, \mathbf{E}_T|\boldsymbol{\phi})}{\partial \eta} = \frac{R(R+2)}{4} - \frac{R+2}{2}\varsigma_t^{\mathrm{GH}} + \frac{1}{4}(\varsigma_t^{\mathrm{GH}})^2, \tag{7}$$

where $\varsigma_t^{\mathrm{GH}} = \boldsymbol{\varepsilon}_t^{\mathrm{GH}\prime}\boldsymbol{\varepsilon}_t^{\mathrm{GH}}$ and $\boldsymbol{\varepsilon}_t^{\mathrm{GH}} = \mathbf{S}_{RK}\boldsymbol{\varepsilon}_t$. Thus, we can regard (7) as the M-step in Louis' (1982) formula (5). Next, we can apply the E-step by taking expectations. Specifically, if $\boldsymbol{\varepsilon}_{t|T}(\boldsymbol{\theta})$ denotes the Kalman smoothed values of the innovations at t given $\mathbf{Y}_T$, which contains past, present and future values of the observed series, and $\boldsymbol{\Omega}_{t|T}(\boldsymbol{\theta})$ the corresponding mean-square error, we have that $\boldsymbol{\varepsilon}_t|\,\mathbf{Y}_T, \boldsymbol{\theta} \sim N[\boldsymbol{\varepsilon}_{t|T}(\boldsymbol{\theta}), \boldsymbol{\Omega}_{t|T}(\boldsymbol{\theta})]$ under the null of normality, so that

$$\frac{\partial \ln f_{\mathbf{Y}}(\mathbf{Y}_T|\boldsymbol{\phi})}{\partial \eta} = \frac{R(R+2)}{4} - \frac{R+2}{2}E\left[\varsigma_t^{\mathrm{GH}}\middle|\,\mathbf{Y}_T, \boldsymbol{\theta}\right] + \frac{1}{4}E\left[(\varsigma_t^{\mathrm{GH}})^2\middle|\,\mathbf{Y}_T, \boldsymbol{\theta}\right]$$

only involves the computation of $E\left[\varsigma_t^{\mathrm{GH}}\middle|\,\mathbf{Y}_T, \boldsymbol{\theta}\right]$ and $E\left[(\varsigma_t^{\mathrm{GH}})^2\middle|\,\mathbf{Y}_T, \boldsymbol{\theta}\right]$, whose expressions we derive in the main appendix. Thus, we can show that:

Proposition 2 *The score of the Student t log-likelihood with respect to the shape parameter η when $\eta = 0$ is given by*

$$\bar{s}_{\mathrm{k}t|T}(\boldsymbol{\theta}) = \frac{1}{T}\sum_{t=1}^{T} s_{\mathrm{k}t|T}(\boldsymbol{\theta}) = \frac{1}{T}\sum_{t=1}^{T} \mathbf{b}'_{\mathrm{k}t|T}(\boldsymbol{\theta})\mathbf{m}_{\mathrm{k}t|T}(\boldsymbol{\theta}),$$

where $\mathbf{m}_{\mathrm{k}t|T}(\boldsymbol{\theta}) = [1, \mathbf{m}'_{2t|T}(\boldsymbol{\theta}), \mathbf{m}'_{4t|T}(\boldsymbol{\theta})]'$, $\mathbf{b}_{\mathrm{k}t|T}(\boldsymbol{\theta}) = [b_{0t|T}(\boldsymbol{\theta}), \mathbf{b}'_{2t|T}(\boldsymbol{\theta}), \mathbf{b}'_{4t|T}(\boldsymbol{\theta})]'$,

$$\begin{aligned} \mathbf{m}_{2t|T}(\boldsymbol{\theta}) &= vec[\boldsymbol{\varepsilon}_{t|T}(\boldsymbol{\theta})\boldsymbol{\varepsilon}_{t|T}(\boldsymbol{\theta})'], \\ \mathbf{m}_{4t|T}(\boldsymbol{\theta}) &= vec\{[\boldsymbol{\varepsilon}_{t|T}(\boldsymbol{\theta}) \odot \boldsymbol{\varepsilon}_{t|T}(\boldsymbol{\theta})][\boldsymbol{\varepsilon}_{t|T}(\boldsymbol{\theta}) \odot \boldsymbol{\varepsilon}_{t|T}(\boldsymbol{\theta})]'\}. \end{aligned} \tag{8}$$

[4] We thank a referee for pointing this out.

and

$$b_{0t|T}(\boldsymbol{\theta}) = c_0 + \{c_1 + c_2 tr[\boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})]\} tr[\boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})] + 2c_2 tr\{[\boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})]^2\},$$
$$\mathbf{b}_{2t|T}(\boldsymbol{\theta}) = \{c_1 + 2c_2 tr[\boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})]\} (\mathbf{S}'_{RK} \otimes \mathbf{S}'_{RK}) vec(\mathbf{I}_R) + 4c_2 (\mathbf{S}'_{RK} \otimes \mathbf{S}'_{RK}) vec[\boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})],$$
$$\mathbf{b}_{4t|T}(\boldsymbol{\theta}) = c_2 (\mathbf{S}'_{RK} \otimes \mathbf{S}'_{RK}) \boldsymbol{\ell}_{R^2},$$

with $c_0 = R(R+2)/4$, $c_1 = -(R+2)/2$, $c_2 = 1/4$ *and* $\boldsymbol{\ell}_H$ *a vector of* H *ones.*

3.2 Asymptotic covariance matrix of the score under Gaussianity

As is well known, the Kalman smoothed process $\boldsymbol{\varepsilon}_{t|T}(\boldsymbol{\theta})$ will typically be serially correlated in spite of $\boldsymbol{\varepsilon}_t$ being *iid*. Consequently, the same will be true of $s_{\mathrm{k}t|T}(\boldsymbol{\theta})$. In addition, the autocovariances of $\boldsymbol{\varepsilon}_{t|T}(\boldsymbol{\theta})$ change with both t and T. Nevertheless, we show in Supplemental Material B that it suffices to compute the autocovariances of powers of $\boldsymbol{\varepsilon}_{t|\infty}(\boldsymbol{\theta})$, which is the Wiener-Kolmogorov filter of $\boldsymbol{\varepsilon}_t$ based on a double-infinite sample of the observable vector $\mathbf{y}_t$, for the purposes of obtaining the asymptotic variance of $\sqrt{T}\bar{s}_{\mathrm{k}t|T}(\boldsymbol{\theta})$. For that reason, we define $\mathbf{m}_{jt}(\boldsymbol{\theta})$ as $\mathbf{m}_{jt|T}(\boldsymbol{\theta})$ in (8) with $\boldsymbol{\varepsilon}_{t|\infty}(\boldsymbol{\theta})$ in place of $\boldsymbol{\varepsilon}_{t|T}(\boldsymbol{\theta})$, $\mathbf{b}_j(\boldsymbol{\theta})$ as $\mathbf{b}_{jt|T}(\boldsymbol{\theta})$ in Proposition 2 with $\boldsymbol{\Omega}_\infty(\boldsymbol{\theta})$ replacing $\boldsymbol{\Omega}_{t|T}(\boldsymbol{\theta})$,[5]and $\bar{s}_{\mathrm{k}T}(\boldsymbol{\theta})$ as the associated average score.

In practice, however, we do not generally know $\boldsymbol{\theta}$. Therefore, we need to obtain the asymptotic covariance matrix of $\sqrt{T}\bar{s}_{\mathrm{k}T}(\hat{\boldsymbol{\theta}}_T)$, where $\hat{\boldsymbol{\theta}}_T$ is the Gaussian Maximum Likelihood estimator of $\boldsymbol{\theta}$, which is the efficient estimator under the null. Importantly, the second part of the following proposition shows that the sampling variability of the Gaussian ML estimators of $\boldsymbol{\theta}$ does not affect the asymptotic variance of the test:

Proposition 3 *Under the null hypothesis of Gaussian innovations:*

a) $\lim_{T\to\infty} V[\sqrt{T}\bar{s}_{\mathrm{k}T}(\boldsymbol{\theta})|\boldsymbol{\theta}] = \mathbf{b}'_4(\boldsymbol{\theta})\boldsymbol{\kappa}_4(\boldsymbol{\theta})\mathbf{b}_4(\boldsymbol{\theta}) - \mathbf{b}'_2(\boldsymbol{\theta})\boldsymbol{\kappa}_2(\boldsymbol{\theta})\mathbf{b}_2(\boldsymbol{\theta}) = \mathcal{C}_{\mathrm{k}}(\boldsymbol{\theta})$,

where

$$\boldsymbol{\kappa}_i(\boldsymbol{\theta}) = \sum_{j=-\infty}^{\infty} cov[\mathbf{m}_{it}(\boldsymbol{\theta}), \mathbf{m}_{it-j}(\boldsymbol{\theta})], \tag{9}$$

denotes the autocovariance generating function of $\mathbf{m}_{it}(\boldsymbol{\theta})$ *evaluated at one.*

b) $$\lim_{T\to\infty} cov[\sqrt{T}\bar{s}_{\mathrm{k}T}(\boldsymbol{\theta}), \sqrt{T}\bar{\mathbf{s}}_{\mathrm{MV}T}(\boldsymbol{\theta})|\boldsymbol{\theta}] = \mathbf{0},$$

where $\bar{\mathbf{s}}_{\mathrm{MV}T}(\boldsymbol{\theta})$ *denotes the average Gaussian score with respect to the conditional mean and variance parameters* $\boldsymbol{\theta}$*.*

[5]Under the usual controllability and observability conditions (see e.g. Harvey (1989)), which we assume henceforth, $\boldsymbol{\Omega}_{t|\infty}(\boldsymbol{\theta})$ will not depend on t in steady state, so we can write $\boldsymbol{\Omega}_\infty(\boldsymbol{\theta}) = \boldsymbol{\Omega}_{t|\infty}(\boldsymbol{\theta})$.

3.3 The test statistic

We can easily compute an LM test for multivariate normality versus multivariate Student t distributed innovations on the basis of the value of the score of the log-likelihood function corresponding to η evaluated at the Gaussian ML estimates $\hat{\boldsymbol{\phi}}_T = (\hat{\boldsymbol{\theta}}_T', \mathbf{0}')'$.

Proposition 4 *The LM test of normality against a multivariate Student t can be expressed as:*

$$LM_T^{Student}(\hat{\boldsymbol{\theta}}_T) = T\frac{\bar{s}^2_{\mathrm{k}t|T}(\hat{\boldsymbol{\theta}}_T)}{\mathcal{C}_{\mathrm{k}}(\hat{\boldsymbol{\theta}}_T)},$$

which is asymptotically distributed as a χ^2_1 *under the null.*

The fact that $\eta = 0$ lies at the boundary of the admissible parameter space invalidates the usual distribution of the LR and Wald tests, which under the null will be a 50:50 mixture of χ^2_0 (=0 with probability 1) and χ^2_1. Although the distribution of the LM test statistic remains valid, intuition suggests that the one-sided nature of the alternative hypothesis should be taken into account to obtain a more powerful test. For that reason, we follow Fiorentini, Sentana and Calzolari (2003) in using the Kühn-Tucker (KT) multiplier test introduced by Gouriéroux, Holly and Monfort (1980) instead, which is equivalent in large samples to the LR and Wald tests. Thus, we would reject H_0 at the $100\varkappa\%$ significance level if the average score with respect to η evaluated under the Gaussian null is strictly positive *and* the LM statistic exceeds the $100(1-2\varkappa)$ percentile of a χ^2_1 distribution.[6]In this respect, it is important to mention that when there is a single restriction, as in our case, those one-sided tests would be asymptotically locally most powerful.

4 Multivariate normal versus *GH* innovations

4.1 The *GH* as a location-scale mixture of normals

We can gain some intuition about the *GH* distribution by considering Blaesild's (1981) interpretation as a location-scale mixture of normals in which the mixing variable is a Generalized

[6]Intuituitively, under the null of normality

$$\sqrt{T}\frac{\bar{s}_{\mathrm{k}t|T}(\hat{\boldsymbol{\theta}}_T)}{\sqrt{\mathcal{C}_{\mathrm{k}}(\hat{\boldsymbol{\theta}}_T)}}$$

will be asymptotically distributed as a standard normal. Therefore, the one-sided nature of the alternative hypothesis implies that the relevant critical value for size $\varkappa$ is given by the $(1-\varkappa)^{th}$ quantile of a standard normal insttead of the usual $(1-\varkappa/2)^{th}$ one.

Inverse Gaussian (GIG). Specifically, if $\boldsymbol{\varepsilon}$ is a GH vector, then it can be expressed as

$$\boldsymbol{\varepsilon} = \boldsymbol{\alpha} + \boldsymbol{\Upsilon}\boldsymbol{\beta}\zeta^{-1} + \zeta^{-\frac{1}{2}}\boldsymbol{\Upsilon}^{\frac{1}{2}}\boldsymbol{\varepsilon}^{\circ}, \tag{10}$$

where $\boldsymbol{\alpha}, \boldsymbol{\beta} \in \mathbb{R}^K$, $\boldsymbol{\Upsilon}$ is a symmetric positive definite matrix of order K, $\boldsymbol{\varepsilon}^{\circ} \sim N(\mathbf{0}, \mathbf{I}_K)$ and the positive mixing variable ζ is an independent GIG with parameters $-\nu$, γ and δ, or $\zeta \sim GIG(-\nu, \gamma, \delta)$ for short, where $\nu \in \mathbb{R}$ and $\gamma, \delta \in \mathbb{R}^+$ (see Jørgensen (1982) and Johnson, Kotz, & Balakrishnan (1994) for further details). Obviously, the distribution of $\boldsymbol{\varepsilon}$ becomes a simple scale mixture of normals, and thereby spherical, when $\boldsymbol{\beta}$ is zero. By restricting $\boldsymbol{\alpha}$ and $\boldsymbol{\Upsilon}$, Mencía and Sentana (2012) derive a standardized version of the GH distribution with zero mean and identity covariance matrix, which therefore depends exclusively on three shape parameters that we can set to zero under normality: $\boldsymbol{\beta}$, which introduces asymmetries, $\eta = -.5\nu^{-1}$ and $\psi = (1+\gamma)^{-1}$, whose product $\tau = \eta\psi$ effectively controls excess kurtosis in the vicinity of the Gaussian null.

4.2 The score under Gaussianity

As in section 3, there is no analytical expression for the log-likelihood function under the alternative, so once again we resort to the generalized Louis' (1982) formula. However, we face two additional difficulties. First, there are three different paths along which a symmetric GH distribution converges to a Gaussian distribution. Fortunately, Mencía and Sentana (2012) showed that the score of the relevant kurtosis parameter evaluated under the null of normality is proportional along those three paths to the score with respect to $\tau = \eta\psi$ evaluated at $\tau = 0$. Second, $\boldsymbol{\beta}$ vanishes from the log-likelihood function as $\tau \to 0$.

One standard solution in the literature to deal with testing situations with underidentified parameters under the null involves fixing those parameters to some arbitrary values, and then computing the appropriate test statistic for the chosen values. To apply this idea to the LM test in our context, we need:

Proposition 5 *The score of the asymmetric GH with respect to the parameter τ when $\tau = 0$ for fixed values of the skewness parameters $\boldsymbol{\beta}$ is given by*

$$\begin{aligned} \bar{s}_{\mathrm{GH}T}(\boldsymbol{\theta}, \boldsymbol{\beta}) &= \frac{1}{T}\sum_{t=1}^{T}\left[s_{\mathrm{kt}|T}(\boldsymbol{\theta}) + \boldsymbol{\beta}'\mathbf{s}_{\mathrm{st}|T}(\boldsymbol{\theta})\right], \\ \mathbf{s}_{\mathrm{st}|T}(\boldsymbol{\theta}) &= \mathbf{b}'_{\mathrm{st}|T}(\boldsymbol{\theta})\mathbf{m}_{\mathrm{st}|T}(\boldsymbol{\theta}), \end{aligned} \tag{11}$$

where $\mathbf{m}_{\mathrm{s}t|T}(\boldsymbol{\theta}) = [\mathbf{m}'_{1t|T}(\boldsymbol{\theta}), \mathbf{m}'_{3t|T}(\boldsymbol{\theta})]'$, $\mathbf{b}_{\mathrm{s}t|T}(\boldsymbol{\theta}) = [\mathbf{b}'_{1t|T}(\boldsymbol{\theta}), \mathbf{b}'_{3t|T}(\boldsymbol{\theta})]'$,

$$\begin{aligned}\mathbf{m}_{1t|T}(\boldsymbol{\theta}) &= \boldsymbol{\varepsilon}_{t|T}(\boldsymbol{\theta}),\\ \mathbf{m}_{3t|T}(\boldsymbol{\theta}) &= vec\{\boldsymbol{\varepsilon}_{t|T}(\boldsymbol{\theta})[\boldsymbol{\varepsilon}_{t|T}(\boldsymbol{\theta}) \odot \boldsymbol{\varepsilon}_{t|T}(\boldsymbol{\theta})]'\},\end{aligned}$$

and

$$\begin{aligned}\mathbf{b}_{1t|T}(\boldsymbol{\theta}) &= [c_3 + tr(\boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta}))]\mathbf{S}'_{RK} + 2\mathbf{S}'_{RK}\boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta}),\\ \mathbf{b}_{3t|T}(\boldsymbol{\theta}) &= \mathbf{S}'_{RK}\boldsymbol{\ell}_R \otimes \mathbf{S}'_{RK},\end{aligned}$$

with $c_3 = -(R+2)$ *and* $\boldsymbol{\ell}_H$ *a vector of* H *ones.*

This result provides an intuitive interpretation for $s_{\mathrm{GH}t|T}(\boldsymbol{\theta}, \boldsymbol{\beta})$ as a linear combination of a kurtosis component, $s_{\mathrm{k}t|T}(\boldsymbol{\theta})$, and R skewness components, $\mathbf{s}_{\mathrm{s}t|T}(\boldsymbol{\theta})$.

4.3 Asymptotic covariance matrix of the score under Gaussianity

If we denote by $\bar{s}_{\mathrm{s}T}(\boldsymbol{\theta})$ the average score with $\boldsymbol{\varepsilon}_{t|T}(\boldsymbol{\theta})$ and $\boldsymbol{\Omega}_{t|T}(\boldsymbol{\theta})$ replaced by $\boldsymbol{\varepsilon}_{t|\infty}(\boldsymbol{\theta})$ and $\boldsymbol{\Omega}_{\infty}(\boldsymbol{\theta})$, respectively, arguments analogous to those in section 3.2 allow us to prove the following result:

Proposition 6 *Under the null hypothesis of Gaussian innovations:*
a) $\sqrt{T}\bar{s}_{\mathrm{k}T}(\boldsymbol{\theta})$ *and* $\sqrt{T}\bar{\mathbf{s}}_{\mathrm{s}T}(\boldsymbol{\theta})$ *are asymptotically independent, and*

$$\lim_{T\to\infty} V[\sqrt{T}\bar{\mathbf{s}}_{\mathrm{s}T}(\boldsymbol{\theta})|\boldsymbol{\theta}] = \mathbf{b}'_3(\boldsymbol{\theta})\boldsymbol{\kappa}_3(\boldsymbol{\theta})\mathbf{b}_3(\boldsymbol{\theta}) - \mathbf{b}'_1(\boldsymbol{\theta})\boldsymbol{\kappa}_1(\boldsymbol{\theta})\mathbf{b}_1(\boldsymbol{\theta}) = \mathcal{C}_{\mathrm{s}}(\boldsymbol{\theta}),$$

with $\boldsymbol{\kappa}_i(\boldsymbol{\theta})$ *defined in (9).*

b) $$\lim_{T\to\infty} cov[\sqrt{T}\bar{\mathbf{s}}_{\mathrm{s}T}(\boldsymbol{\theta}), \sqrt{T}\bar{\mathbf{s}}_{\mathrm{MV}T}(\boldsymbol{\theta})|\boldsymbol{\theta}] = \mathbf{0}.$$

The second part of this proposition, combined with the second part of Proposition 3, implies that the scores of the conditional mean and variance parameters $\boldsymbol{\theta}$ and the scores of the shape parameters $\boldsymbol{\varphi}$ are asymptotically independent under the null of Gaussianity, so that we need not worry about parameter uncertainty, at least in large samples. Interestingly, this implication is closely related to Proposition 3 in Fiorentini and Sentana (2007), which contains an analogous result for multivariate, dynamic location-scale models with non-Gaussian innovations. It is also related to Bontemps and Meddahi (2005), who show that univariate normality tests based on third and higher order Hermite polynomials are insensitive to parameter uncertainty too.

In Supplemental Material C, we provide a numerically reliable algorithm for computing the asymptotic covariance matrices $\mathcal{C}_{\mathrm{s}}(\boldsymbol{\theta})$ and $\mathcal{C}_{\mathrm{k}}(\boldsymbol{\theta})$ for any state space model.

4.4 The test statistic

If we combine Propositions 5 and 6, we can easily show that the LM test statistic for a given value of $\boldsymbol{\beta}$ will be given by

$$LM_T^{GH}(\hat{\boldsymbol{\theta}}_T, \boldsymbol{\beta}) = \frac{T}{\mathcal{C}_{\mathrm{k}}(\hat{\boldsymbol{\theta}}_T) + \boldsymbol{\beta}'\mathcal{C}_{\mathrm{s}}(\hat{\boldsymbol{\theta}}_T)\boldsymbol{\beta}} \left\{ \frac{1}{T} \sum_{t=1}^{T} \left[s_{\mathrm{k}t|T}(\hat{\boldsymbol{\theta}}_T) + \boldsymbol{\beta}' \mathbf{s}_{\mathrm{s}t|T}(\hat{\boldsymbol{\theta}}_T) \right] \right\}^2 ,$$

which will also follow an asymptotic χ^2_1 distribution under H_0.

But since it is often unclear what value of $\boldsymbol{\beta}$ to choose, we prefer a second approach, which consists in computing the LM test for all possible values of $\boldsymbol{\beta}$ and then taking the supremum over those parameter values. Remarkably, we can maximize $LM_T^{GH}(\boldsymbol{\theta}, \boldsymbol{\beta})$ with respect to $\boldsymbol{\beta}$ in closed form, and also obtain the asymptotic distribution of the resulting sup test statistic:

Proposition 7 *The supremum with respect to $\boldsymbol{\beta}$ of the LM tests based on (11) is equal to*

$$\sup_{\boldsymbol{\beta}} LM_T^{GH}(\hat{\boldsymbol{\theta}}_T, \boldsymbol{\beta}) = LM_T^{Student}(\hat{\boldsymbol{\theta}}_T) + T \left[\frac{1}{T} \sum_{t=1}^{T} \mathbf{s}_{\mathrm{s}t|T}(\hat{\boldsymbol{\theta}}_T) \right]' \mathcal{C}_{\mathrm{s}}^{-1}(\hat{\boldsymbol{\theta}}_T) \left[\frac{1}{T} \sum_{t=1}^{T} \mathbf{s}_{\mathrm{s}t|T}(\hat{\boldsymbol{\theta}}_T) \right],$$

which is asymptotically distributed as a χ^2_{R+1} under the null.

Given that $s_{\mathrm{k}t|T}(\boldsymbol{\theta})$ is asymptotically orthogonal to the other R moment conditions in $\mathbf{s}_{\mathrm{s}t|T}(\boldsymbol{\theta})$ from the first part of Proposition 6, we can conduct a partially one-sided test by combining the KT one-sided version of the symmetric GH test and the moment test based on $\mathbf{s}_{\mathrm{s}t|T}(\boldsymbol{\theta})$, which should be equivalent in large samples to the corresponding LR test (see Proposition 6 in Mencía and Sentana (2012) for a more formal argument). The asymptotic distribution of the joint test under the null will be a 50:50 mixture of χ^2_R and χ^2_{R+1}, whose p-values are the equally weighted average of those two χ^2 p-values.

5 Further discussion

In section 7 we will use our methods for improving GDP measurement. But since they apply far more generally, in this section we first describe how to implement our testing procedures in a generic model. Next, we illustrate them with two popular textbook examples: a static factor model and the so-called local-level dynamic model. Finally, we use these examples to explicitly compare our proposed testing procedures to previous suggestions.

5.1 Practical implementation

Assume the researcher has already (i) specified the model and (ii) computed if necessary the Gaussian maximum likelihood estimates, $\hat{\boldsymbol{\theta}}_T$.

STEP 1: *Compute the smoothed influence functions.* Propositions 2 and 5 provide the explicit expressions required to compute the contribution to the score of each smoothed innovation, namely $s_{\mathrm{k}t|T}(\boldsymbol{\theta})$ and $\mathbf{s}_{\mathrm{s}t|T}(\boldsymbol{\theta})$, respectively. One simply needs as inputs the smoothed innovations $\boldsymbol{\varepsilon}_{t|T}(\boldsymbol{\theta})$, which are required to compute $\mathbf{m}_{jt|T}(\boldsymbol{\theta})$, for $j =, 1..., 4$, and the mean-square error for the vector of innovations being tested, $\boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})$, which are necessary to compute the vector of coefficients $\mathbf{b}_{jt|T}(\boldsymbol{\theta})$, for $j = 0, ..., 4$. Importantly, both of these quantities can be easily obtained from a standard Kalman filter-smoother.

STEP 2: *Obtain the asymptotic variance of the test statistics.* Although this can be done in different ways, in what follows we describe a numerically reliable and computationally efficient algorithm that avoids simulations.

STEP 2.1: *Obtain the* VARMA *representation of the Wiener-Kolmogorov filter for the innovations.* It turns out that the Wiener-Kolmogorov filter always has a finite-order VARMA representation with scalar autoregressive part for all the models in this paper. This feature follows from the fact that their autocovariance generating functions are rational polynomials. Specifically, there exist positive integers p and q, a set of scalars $\phi_1, ..., \phi_p \in \mathbb{R}$ and a set of matrices $\boldsymbol{\Theta}_0, \boldsymbol{\Theta}_1, ..., \boldsymbol{\Theta}_q \in \mathbb{R}^{(M+K)\times K}$ such that

$$(1 - \phi_1 L - \cdots - \phi_p L^p) \begin{pmatrix} \hat{\boldsymbol{\xi}}_{t-1|\infty} \\ \hat{\boldsymbol{\varepsilon}}_{t|\infty} \end{pmatrix} = (\boldsymbol{\Theta}_0 + \boldsymbol{\Theta}_1 L + \cdots + \boldsymbol{\Theta}_q L^q)\boldsymbol{\varepsilon}_t.$$

This representation, in which importantly the VAR component is scalar, is useful to the extent that the coefficients $\phi_1, ..., \phi_p$ and matrices $\boldsymbol{\Theta}_0, \boldsymbol{\Theta}_1, ..., \boldsymbol{\Theta}_q$ can be obtained in terms of the parametrization of the state matrices $\mathbf{H}$, $\mathbf{F}$ and $\mathbf{M}$. This can be done using symbolic software such as Mathematica. We refer the reader to Lemma 4.

STEP 2.2: *Compute the autocovariance function implied by the Wiener-Kolmogorov filter of the innovations.* To do so, consider a VARMA process with scalar VAR part for a K_x-dimensional process $\mathbf{x}_t$,

$$\phi(L)\mathbf{x}_t = \boldsymbol{\Theta}(L)\mathbf{u}_t$$

where $\phi(z) = 1 - \phi_1 z - \cdots - \phi_p z^p$ and $\boldsymbol{\Theta}(z) = \boldsymbol{\Theta}_0 + \boldsymbol{\Theta}_1 z + \ldots \boldsymbol{\Theta}_q z^q$. In Section C of the Supplemental Material we provide a detailed algorithm to compute the autocovariance function

of $\mathbf{x}_t$, whence we can compute the autocovariances of $(\hat{\boldsymbol{\xi}}'_{t-1|\infty}, \hat{\boldsymbol{\varepsilon}}'_{t|\infty})'$.

STEP 2.3: *Compute the expressions that appear in Propositions 3 and 6.* To do so, e.g. one can obtain the autocovariance function of $\mathbf{m}_{h,t|\infty}(\boldsymbol{\theta})$ for $h = 1, \ldots, 4$ from the expressions in i), ii), iii) and iv) in the proof of Proposition *6*. Next, one can cumulate the autocovariance matrices of $\mathbf{m}_{h,t|\infty}(\boldsymbol{\theta})$ for $h = 1 \ldots 4$ until some convergence criterion is satisfied. This gives a numerical approximation to $\boldsymbol{\kappa}_h(\hat{\boldsymbol{\theta}}_T)$ for $h = 1, \ldots, 4$. Finally, one computes $\mathbf{b}_h(\hat{\boldsymbol{\theta}}_T)$, which only requires knowledge of the contemporaneous covariance matrix of the Wiener-Kolmogorov filter because $\boldsymbol{\Omega}_\infty = \mathbf{I}_K - \boldsymbol{\Gamma}_0$.

Codes for all the different steps above, as well as detailed derivations for the expressions in STEP 2 are available upon request.

5.2 Static factor models

We start by considering a single factor version of a traditional (i.e. static, conditionally homoskedastic and exact) factor model. Specifically,

$$\begin{array}{rcl} \mathbf{y}_t & = & \boldsymbol{\pi} + \mathbf{c} f_t + \mathbf{v}_t, \\ \left(\begin{array}{c} f_t \\ \mathbf{v}_t \end{array} \right) \Big| \mathcal{I}_{t-1}, \boldsymbol{\phi} & \sim & iid\ D \left[\left(\begin{array}{c} 0 \\ \mathbf{0} \end{array} \right), \left(\begin{array}{cc} 1 & \mathbf{0} \\ \mathbf{0} & \boldsymbol{\Gamma} \end{array} \right), \boldsymbol{\eta} \right], \end{array} \tag{12}$$

where $\mathbf{y}_t$ is an $N \times 1$ vector of observable variables with constant conditional mean $\boldsymbol{\pi}$, f_t is an unobserved common factor, whose constant variance we have normalized to 1 to avoid the usual scale indeterminacy, $\mathbf{c}$ is the $N \times 1$ vector of factor loadings, $\mathbf{v}_t$ is an $N \times 1$ vector of idiosyncratic noises, which are conditionally orthogonal to f_t, $\boldsymbol{\Gamma}$ is an $N \times N$ diagonal positive definite matrix of constant idiosyncratic variances, and $\boldsymbol{\theta} = (\boldsymbol{\pi}', \mathbf{c}', \boldsymbol{\gamma}')'$, with $\boldsymbol{\gamma} = vecd(\boldsymbol{\Gamma})$.

We can easily express model (12) as in (1)–(2) with $\boldsymbol{\xi}_t = (f_t, \mathbf{v}'_t)'$, $\mathbf{H}(\boldsymbol{\theta}) = (\mathbf{c}, \mathbf{I}_N)$, $\mathbf{F}(\boldsymbol{\theta}) = \mathbf{0}$,

$$\mathbf{M}(\boldsymbol{\theta}) = \left(\begin{array}{cc} 1 & \mathbf{0} \\ \mathbf{0} & diag^{1/2}(\boldsymbol{\gamma}) \end{array} \right)$$

and $\boldsymbol{\varepsilon}_t = (f_t, \mathbf{v}_t^{*\prime})'$, where $\mathbf{v}_t^* = \boldsymbol{\Gamma}^{-1/2}\mathbf{v}_t$. Note that this specification trivially implies that

$$\mathbf{y}_t | \mathcal{I}_{t-1}, \boldsymbol{\phi} \sim iid\ D^*[\boldsymbol{\pi}, \boldsymbol{\Sigma}(\boldsymbol{\theta}), \boldsymbol{\varphi}], \quad \text{with} \quad \boldsymbol{\Sigma}(\boldsymbol{\theta}) = \mathbf{c}\mathbf{c}' + \boldsymbol{\Gamma}.$$

While the normality of $\boldsymbol{\xi}_t$ implies the normality of $\mathbf{y}_t$, in principle the distribution of $\mathbf{y}_t$ and $\boldsymbol{\xi}_t$ will be different under the alternative.

Letting $\mathbf{D}(\boldsymbol{\theta}) = \mathbf{H}(\boldsymbol{\theta})\mathbf{M}(\boldsymbol{\theta})$, we can show that

$$\boldsymbol{\varepsilon}_{t|\infty}(\boldsymbol{\theta}) = \boldsymbol{\varepsilon}_{t|t}(\boldsymbol{\theta}) = \mathbf{D}'(\boldsymbol{\theta})[\mathbf{D}(\boldsymbol{\theta})\mathbf{D}'(\boldsymbol{\theta})]^{-1}\mathbf{D}(\boldsymbol{\theta})(\mathbf{y}_t - \boldsymbol{\pi}),$$

so like in any other static model, $\boldsymbol{\varepsilon}_{t|\infty}(\boldsymbol{\theta})$ will be white noise, with covariance matrix

$$\boldsymbol{\Gamma}(\mathbf{0}) = \mathbf{D}'(\boldsymbol{\theta})[\mathbf{D}(\boldsymbol{\theta})\mathbf{D}'(\boldsymbol{\theta})]^{-1}\mathbf{D}(\boldsymbol{\theta}).$$

In addition,

$$\boldsymbol{\Omega}_{t|T}(\boldsymbol{\theta}) = \boldsymbol{\Omega}_{t|\infty}(\boldsymbol{\theta}) = \boldsymbol{\Omega}_{\infty}(\boldsymbol{\theta}) = \mathbf{I}_K - \mathbf{D}'(\boldsymbol{\theta})[\mathbf{D}(\boldsymbol{\theta})\mathbf{D}'(\boldsymbol{\theta})]^{-1}\mathbf{D}(\boldsymbol{\theta}),$$

which has rank N rather than $N+1$, so that the conditional density will be degenerate. Hence, we will have that under the null,

$$\boldsymbol{\varepsilon}_t|\mathbf{Y}_T, \boldsymbol{\theta} \sim N[\boldsymbol{\varepsilon}_{t|t}(\boldsymbol{\theta}), \boldsymbol{\Omega}_{\infty}(\boldsymbol{\theta})],$$

which contains all the information we need to compute the normality tests.

To provide some intuition, though, it is convenient to focus on tests that look exclusively at the common factor. If we could observe f_t, then we could write the joint log-likelihood function of $\mathbf{y}_t$ and f_t as the sum of the marginal log-likelihood function of f_t and the log-likelihood function of $\mathbf{y}_t$ conditional on f_t, which would coincide with the marginal log-likelihood function of the idiosyncratic terms $\mathbf{v}_t$. If we maintained the assumption that this conditional distribution was Gaussian, and confined the non-normality to the marginal distribution of f_t, the results in Mencía and Sentana (2012) would imply that the LM test of the null hypothesis that f_t is Gaussian versus the alternative that it follows an asymmetric Student t would be based on the following influence conditions:

$$\left.\begin{array}{c} H_3(f_t) = f_t^3 - 3f_t, \\ H_4(f_t) = f_t^4 - 6f_t^2 + 3, \end{array}\right\} \tag{13}$$

which coincide with the third and fourth Hermite polynomials for f_t underlying the usual Jarque and Bera (1980) test.

Unfortunately, f_t is unknown. But we can easily compute the expected values of these expressions conditional on $\mathbf{y}_t$, which under normality are simple functions of

$$f_{t|t}(\boldsymbol{\theta}) = E(f_t|\mathbf{y}_t) = \omega_f(\boldsymbol{\theta})\mathbf{c}'\boldsymbol{\Gamma}^{-1}(\mathbf{y}_t - \boldsymbol{\pi})$$

and

$$\omega_f(\boldsymbol{\theta}) = V(f_t|\mathbf{y}_t) = \frac{1}{\mathbf{c}'\boldsymbol{\Gamma}^{-1}\mathbf{c} + 1}.$$

In particular, we can show that the expected values of the elements of (13) are proportional to $H_3[f_{t|t}(\boldsymbol{\theta})/\sqrt{1-\omega_f(\boldsymbol{\theta})}]$ and $H_4[f_{t|t}(\boldsymbol{\theta})/\sqrt{1-\omega_f(\boldsymbol{\theta})}]$, respectively, where $V[f_{t|t}(\boldsymbol{\theta})] = 1 - \omega_f(\boldsymbol{\theta})$

by virtue of the fact that

$$V(f_t) = E[V(f_t|\mathbf{y}_t)] + V[E(f_t|\mathbf{y}_t)].$$

Somewhat remarkably, therefore, the LM test for the normality of the latent common factor will numerically coincide with the usual LM test for the normality of its best estimator in the mean square error sense. Obviously, analogous calculations apply to each element of $\mathbf{v}_t$.

5.3 The local-level model

Consider now the random walk plus noise model studied in Harvey and Koopman (1992):

$$\begin{aligned} y_t &= \pi + x_t + v_t \\ x_t &= x_{t-1} + f_t \\ \left(\begin{array}{c} f_t \\ v_t \end{array} \right) \Big| \mathcal{I}_{t-1}, \boldsymbol{\phi} &\sim iid\ D\left[\left(\begin{array}{c} 0 \\ 0 \end{array} \right), \left(\begin{array}{cc} \sigma_f^2 & 0 \\ 0 & \sigma_v^2 \end{array} \right), \boldsymbol{\varphi} \right], \end{aligned}$$

where x_t is the "signal" component, v_t the orthogonal "non-signal" component, and $\boldsymbol{\theta}$ refers to the model parameters that characterize the autocovariance structure of the observed series. Once again, we can easily express this model as in (1)–(2) with $\boldsymbol{\xi}_t = (f_t, v_t')'$, $\mathbf{H}(\boldsymbol{\theta}) = (1, 1)$,

$$\mathbf{F}(\boldsymbol{\theta}) = \left(\begin{array}{cc} \alpha & 0 \\ 0 & 0 \end{array} \right), \quad \mathbf{M}(\boldsymbol{\theta}) = \left(\begin{array}{cc} \sigma_f^2 & 0 \\ 0 & \sigma_v^2 \end{array} \right)$$

and $\boldsymbol{\varepsilon}_t = (f_t^*, v_t^*)'$, where $f_t^* = \sigma_f^{-1} f_t$ and $v_t^* = \sigma_v^{-1} v_t$.

Since there are only two shocks, we could look at (i) a test of joint normality, (ii) a test of normality of the "signal" with the maintained hypothesis of normality for the "non-signal", and (iii) vice versa.

For the sake of brevity, let us focus on the non-signal component. Proposition 5 implies that for symmetric Student t alternatives, the score with respect to the reciprocal of the degrees of freedom parameter evaluated under the null will be given by

$$E\left[s_{\eta t}^{S_v}(\boldsymbol{\theta}, \mathbf{0}) | \mathbf{Y}_T \right] = \frac{1}{2}\sqrt{\frac{3}{2}} \left[1 - \omega_{\mathrm{v}t|T}(\boldsymbol{\theta}) \right]^2 - \sqrt{\frac{3}{2}} \left[1 - \omega_{\mathrm{v}t|T}(\boldsymbol{\theta}) \right] v_{t|T}^{*2}(\boldsymbol{\theta}) + \frac{1}{2}\sqrt{\frac{1}{6}} v_{t|T}^{*4}(\boldsymbol{\theta}). \quad (14)$$

But the optimality of the Wiener-Kolmogorov-Kalman filter under Gaussianity implies that

$$V(v_t^*) = V[v_{t|T}^*(\boldsymbol{\theta})] + V[v_t^* - v_{t|T}^*(\boldsymbol{\theta})],$$

which in turns means that

$$V[v_{t|T}^*(\boldsymbol{\theta})] = 1 - \omega_{\mathrm{v}t|T}(\boldsymbol{\theta}).$$

Hence, expression (14) is proportional to the fourth order Hermite polynomial of the standardized variable $v^*_{t|T}(\boldsymbol{\theta})/\sqrt{1-\omega_{\mathrm{v}t|T}(\boldsymbol{\theta})}$. Therefore, for this model our proposed LM test also yields exactly the same influence function as an LM test of normal versus Student t that would treat $v^*_{t|T}(\boldsymbol{\theta})$ as an *iid* series. Unlike in the static model considered in section 5.2, though, the elements of (14) are serially correlated.

5.4 Comparison with alternative approaches

5.4.1 Univariate tests applied to the smoothed innovations

As we mentioned in the introduction, Harvey and Koopman (1992) applied standard univariate normality tests for observed variables to the smoothed values of the innovations in the underlying components of a local level model explicitly taking into account the serial correlation in those filtered estimates implied by the model.

Their asymmetry test is based on the skewness coefficient

$$sk_{\varepsilon_i} = m_{\varepsilon_i 3}/m_{\varepsilon_i 2}^{3/2},$$

where

$$m_{\varepsilon_i j} = T^{-1}\sum_{t=1}^{T}(\varepsilon_{it|T} - \bar{\varepsilon})^j$$

is the j^{th} centred sample moment of the smoothed innovations of either the signal ($i=1$) or the noise ($i=2$). Under normality, the asymptotic variance of sk_{ε_i} will be given by $\zeta_{\varepsilon_i}(\boldsymbol{\theta},3)$, where

$$\zeta_{\varepsilon_i}(\boldsymbol{\theta},\lambda) = \lambda!\sum_{j=-\infty}^{\infty}[\rho_{\varepsilon_i}(j)]^\lambda$$

provides the sum of powers of the autocorrelations, which are the autocorrelations of the powers of the original Gaussian series (see Lomnicki (1961)).

Similarly, their excess kurtosis test is based on the sample excess kurtosis coefficient

$$k_{\varepsilon_i} = m_{\varepsilon_i 4}/m_{\varepsilon_i 2}^2 - 3,$$

whose asymptotic variance under normality will be given by $\zeta_{\varepsilon_i}(\boldsymbol{\theta},4)$.

It is interesting to compare these tests to our LM tests based on Propositions 4 and 7. The procedures proposed by Harvey and Koopman (1992) can be regarded as moment tests of

$$\begin{aligned} E[f^{*3}_{t|T}(\boldsymbol{\theta})] &= 0, \qquad E[f^{*4}_{t|T}(\boldsymbol{\theta}) - 3] = 0, \\ E[v^{*3}_{t|T}(\boldsymbol{\theta})] &= 0, \qquad E[v^{*4}_{t|T}(\boldsymbol{\theta}) - 3] = 0, \end{aligned}$$

where $f^*_{t|T}(\boldsymbol{\theta})$ and $v^*_{t|T}(\boldsymbol{\theta})$ are standardized smoothed innovations. Thus, the main difference is that they look at third and fourth moments, while we use the log-likelihood scores, which are proportional to the third and fourth Hermite polynomials. The main advantage of the latter is that they are not affected by the sampling variability in $\hat{\boldsymbol{\theta}}_T$, as we have shown in Propositions 3 and 6. Nevertheless, Harvey and Koopman (1992) indicate that their tests are also asymptotically insensitive to parameter uncertainty when the standardization of $f^*_{t|T}(\boldsymbol{\theta})$ and $v^*_{t|T}(\boldsymbol{\theta})$ relies on sample moments (see also Bontemps and Meddahi (2005)).[7] In fact, we can show that their tests and ours are asymptotically equivalent under the null hypothesis in the local level model in section 5.3.

5.4.2 Reduced form tests

Assuming covariance stationarity, possibly after some suitable transformation, we can find the autocorrelation structure of the observed series generated by (1)-(2), as well as the corresponding Wold representation, which will typically resemble a VARMA model, with potentially long but finite AR and MA orders, but restricted coefficient matrices because $M \geq N$.

As a result, we will be able to write

$$[\mathbf{y}_t - \boldsymbol{\pi}(\boldsymbol{\theta})] = \sum_{j=1}^{p_y} \mathbf{A}_j(\boldsymbol{\theta})[\mathbf{y}_{t-j} - \boldsymbol{\pi}(\boldsymbol{\theta})] + \mathbf{w}_t + \sum_{j=1}^{q_y} \mathbf{B}_j(\boldsymbol{\theta})\mathbf{w}_{t-j},$$

where $\mathbf{w}_t$ is a serially uncorrelated sequence, linearly unpredictable on the basis of lagged values of $\mathbf{y}_t$. In fact, assuming that the Wold representation is strictly invertible,

$$\mathbf{w}_t = \left[\mathbf{I}_N + \sum_{j=1}^{q_y} \mathbf{B}_j(\boldsymbol{\theta})L^j\right]^{-1} \left[\mathbf{I}_N - \sum_{j=1}^{q_y} \mathbf{A}_j(\boldsymbol{\theta})L^j\right] [\mathbf{y}_t - \boldsymbol{\pi}(\boldsymbol{\theta})]. \tag{15}$$

This relationship is the basis for the comparison of our tests, which target the components in $\boldsymbol{\varepsilon}_t$ directly, to existing tests, which target $\mathbf{w}_t$ instead. If $\boldsymbol{\varepsilon}_t|\mathcal{I}_{t-1}$ is *iid* normal, then $\mathbf{y}_t$ will be a Gaussian process, and therefore $\mathbf{w}_t|\mathcal{I}_{t-1}$ will be *iid* normal too. As a result, checking the normality of the latter provides an indirect way of checking the normality of the former. Nevertheless, if some elements of $\boldsymbol{\varepsilon}_t$ are not normal, then the conditional distribution of the reduced form innovations will typically be extremely complicated, especially taking into account that they are unlikely to follow a martingale difference sequence in dynamic contexts.[8]The problem is that the conditional mean of the observed variables given their past alone will no longer be given

[7] In that regard, the situation seems analogous to the Jarque and Bera (1980) tests, whose distribution is insensitive to parameter uncertainty for many models (see Fiorentini, Sentana and Calzolari (2004)).

[8] Although we would expect it to be closer to a normal than $\boldsymbol{\varepsilon}_t$ because of the averaging implicit in (15).

by the one-period ahead linear prediction generated by the Kalman filter recursions, $\mathbf{y}_{t|t-1}(\boldsymbol{\theta})$. Similarly, the conditional variance will not usually coincide with the associated mean-square error matrix $\boldsymbol{\Sigma}_{t|t-1}(\boldsymbol{\theta})$.

Still, it may be worth considering tests against the following alternative model

$$\mathbf{y}_t|\mathbf{y}_{t-1},...,\mathbf{y}_1,\boldsymbol{\phi} \sim GH[\mathbf{y}_{t|t-1}(\boldsymbol{\theta}),\boldsymbol{\Sigma}_{t|t-1}(\boldsymbol{\theta}),\eta,\psi,\boldsymbol{\beta}],$$

which maintains the assumption that the conditional mean and variance coincide with their values under normality, but allows for a non-Gaussian distribution. The assumption that the distribution of $\mathbf{y}_t$ conditional on $\mathbf{Y}_{t-1}$ is GH but with a mean vector and covariance matrix given by the usual Gaussian Kalman filter recursions may be regarded as a way of constructing a convenient auxiliary model that coincides with the model of interest for $\boldsymbol{\varphi} = \mathbf{0}$, but whose log-likelihood function and score we can obtain in closed form for every possible value of $\boldsymbol{\theta}$ when $\boldsymbol{\varphi} \neq \mathbf{0}$. The pay-off is that the resulting model falls within the framework studied by Mencía and Sentana (2012). Specifically, if we define the standardized reduced form innovations as

$$\mathbf{w}^*_{t|t-1}(\boldsymbol{\theta}) = \boldsymbol{\Sigma}^{-\frac{1}{2}}_{t|t-1}(\boldsymbol{\theta})[\mathbf{y}_t - \mathbf{y}_{t|t-1}(\boldsymbol{\theta})],$$

and their (square) Euclidean norm as

$$\varsigma_{t|t-1}(\boldsymbol{\theta}) = \mathbf{w}^*_{t|t-1}(\boldsymbol{\theta})'\mathbf{w}^*_{t|t-1}(\boldsymbol{\theta}) = [\mathbf{y}_t - \mathbf{y}_{t|t-1}(\boldsymbol{\theta})]'\boldsymbol{\Sigma}^{-1}_{t|t-1}(\boldsymbol{\theta})[\mathbf{y}_t - \mathbf{y}_{t|t-1}(\boldsymbol{\theta})],$$

we can write the influence functions underling their test as

$$\begin{aligned} s^{MS}_{\mathrm{kt}|t-1}(\boldsymbol{\theta}) &= \frac{1}{4}\varsigma^2_{t|t-1}(\boldsymbol{\theta}) - \frac{N+2}{2}\varsigma_{t|t-1}(\boldsymbol{\theta}) + \frac{N(N+2)}{4}, \\ \mathbf{s}^{MS}_{\mathrm{st}|t-1}(\boldsymbol{\theta}) &= \boldsymbol{\Sigma}^{\frac{1}{2}}_{t|t-1}(\theta)\mathbf{w}^*_{t|t-1}(\boldsymbol{\theta})\left[\varsigma_{t|t-1}(\boldsymbol{\theta}) - (N+2)\right]. \end{aligned}$$

Propositions 3 and 5 in Mencía and Sentana (2012) provide expressions for the asymptotic covariance matrix of the sample average of those influence functions in terms of $\boldsymbol{\Sigma}(\boldsymbol{\theta}) = V(\mathbf{w}_t)$, which typically coincides with the steady state value of $\boldsymbol{\Sigma}_{t|t-1}(\boldsymbol{\theta})$ (see footnote 4).

6 Monte Carlo simulations

In this section, we study the finite sample size and power properties of the testing procedures discussed above by means of several extensive Monte Carlo exercises. We do so in the context of three different models:[9]

[9] Results for a trivariate version of the static factor model (12) can be found in section 4.2.2 of Almuzara, Amengual and Sentana (2017).

1. the cointegrated single factor model we use in our empirical application in section 7,

2. the illustrative local level model in section 5.3, and

3. a multivariate version of this local level model in which there is a single integrated common trend, but the number of observed series is 10, each of which containing an *iid* idiosyncratic component.

6.1 Simulation and estimation details

We assess the power properties of our tests by generating non-Gaussian data in three alternative designs:

1. All structural innovations are jointly GH: $\boldsymbol{\varepsilon}_t \sim GH(\eta, \psi, \boldsymbol{\beta})$ (alternative J);

2. The distribution of the innovations to the signal component is GH while the idiosyncratic shocks are Gaussian: $f_t \sim GH(\eta, \psi, \beta)$, $\mathbf{v}_t \sim N(\mathbf{0}, \mathbf{I}_N)$ (alternative S_f);

3. The joint distribution of the innovations to the idiosyncratic variables is GH while the common component is Gaussian: $\mathbf{v}_t \sim GH(\eta, \psi, \boldsymbol{\beta})$, $f_t \sim N(0, 1)$ (alternative $S_{\mathbf{v}}$).

We consider two examples of GH distributions: a symmetric Student t with 8 degrees of freedom and an asymmetric Student t with 8 degrees of freedom and skewness vector $\boldsymbol{\beta} = -\boldsymbol{\ell}_{K\times 1}$. Thus, we end up with a total of seven different specifications for $\boldsymbol{\varepsilon}_t$, including the Gaussian null. For each distributional assumption, we generate 10,000 samples of size T exploiting the location-scale mixture of normal representation of the GH distribution we discussed in section 4.1.

We use standard MATLAB routines for estimation. In the case of the local-level model, we rely on its IMA$(1,1)$ reduced form representation to improve the computational efficiency of the algorithm. Finally, we compute the asymptotic variances of the test statistics by truncating the infinite sum in expression (9) when the additional terms lead to increments lower than 10^{-5}.[10]

Given that in all the models we observe a "pile-up" problem, whereby the fraction of negative values of the average kurtosis scores exceeds 50% under the null, we employ a parametric bootstrap procedure based on 10,000 simulated samples. In this way, we can automatically compute size-adjusted rejection rates, as forcefully argued by Horowitz and Savin (2000). Importantly, our bootstrap procedure does not exploit the asymptotic orthogonality of the scores

[10] In Supplemental Appendix E we report analogous results but using a HAC estimator to compute the asymptotic variances of the influence functions underlying our test statistics. As expected, the results are far less reliable than when we use the theoretical expressions.

between mean and variance parameters on the one hand and shape parameters on the other in Propositions 3 and 6. On the contrary, it explicitly takes into account the sensitivity of the critical values to the estimated values of $\boldsymbol{\theta}$ in order not to rule out higher order refinements (see Appendix D.1 in Amengual and Sentana (2015) for details).

In all the tables, the row labels H_J, H_{S_f}, and $H_{S_{\mathbf{v}}}$ refer to the score tests in Propositions 4 and 7 corresponding to the J, S_f and $S_{\mathbf{v}}$ alternative hypotheses, while Red denotes the reduced form tests discussed in section 5.4.2. For each of those labels, Kt and Sk refer to the kurtosis and skewness components of the corresponding test statistics, while GH indicates the sum of the two.

6.2 Small sample properties

6.2.1 Cointegrated dynamic factor model

We simulate data from the model (16) that we use in our empirical application, with $\rho_x = .5$, $\rho_{\epsilon_E} = .2$, $\rho_{\epsilon_I} = .8$, $\sigma_f^2 = 1$ and $\sigma_{v_i}^2$ chosen such that $q_E = 2$ and $q_I = .5$, where $q_i = \sigma_x^2/\sigma_{\epsilon_i}^2 = [\sigma_f^2(1-\rho_{\epsilon_i}^2)]/[(1-\rho_x^2)\sigma_{v_i}^2]$ represents the signal-to-noise ratio for y_{it} for $i = E, I$.

Panels A of Tables 1 and 2 report rejection rates under the null at the 1%, 5% and 10% levels for $T = 100$ and $T = 250$, respectively, which roughly correspond to the sample sizes in our empirical application in section 7. The results make clear that the parametric bootstrap works remarkably well for both sample sizes.[11]

Panels B of the same tables report the rejection rates at the 5% level of the tests under each of different alternative hypotheses that we consider. As expected, the most powerful test for any given alternative is typically the score test we have designed against that particular alternative. In that regard, we find that while the reduced form tests have non-trivial power, especially under alternative J, they are clearly dominated by the tests aimed at the structural innovations.

6.2.2 Univariate local level model

Table 3 contains the results for samples of size $T = 250$ of the local-level model in section 5.3 in which the signal-to-noise ratio $q = \sigma_f^2/\sigma_v^2$ is set to 2, as in Harvey and Koopman (1992). For comparison purposes, we also include their original tests.

Our results confirm the asymptotic equivalence between their tests and the less powerful two-sided versions of ours (not reported). More generally, we essentially reach the same conclusions

[11] Given the number of Monte Carlo replications, the 95% asymptotic confidence intervals for the rejection probabilities under the null are (.80,1.20), (4.57,5.43) and (9.41,10.59) at the 1, 5 and 10% levels.

for size and power as in the previous two examples.

6.2.3 Multivariate local level model

To assess the performance of our tests when the cross-section dimension is moderately large, in Table 4 we provide results for a ten-variate model with a single common trend and uncorrelated idiosyncratic terms. Specifically, we assume $\boldsymbol{\pi} = \mathbf{0}$, $\mathbf{c} = \boldsymbol{\ell}_{10}$ and $\boldsymbol{\gamma} = q^{-1}\boldsymbol{\ell}_{10}$, where the signal-to-noise ratio q is set to 2, as in the univariate version. We also maintain $T = 250$. Once again, we reach analogous conclusions for size and power as in the other three examples. The main difference is that rejection rates are almost 100% under S_J and $S_{\mathbf{v}}$ because the number of non-normal innovations is substantially larger than in the univariate case. Moreover, the precision with which the common factor is filtered is much higher than in the previous example because, *ceteris paribus*, the increase in the cross-sectional dimension N increases the observability of the latent variables. As a result, we obtain rejection rates close to the nominal ones in cases in which the mantained assumption of normality is indeed satisfied.

7 Inferring real output from GDP and GDI

7.1 The model

As we mentioned in the introduction, in theory the expenditure (GDP) and income (GDI) measures of output should be equal, but they differ because they are calculated from different sources. Traditionally, the difference between the two, officially known as the "statistical discrepancy" (see Grimm (2007)), was regarded by many academic economists as a curiosity in the US National Input and Product Accounts (NIPA) elaborated by the Bureau of Economic Analysis (BEA) of the Department of Commerce. However, the Great Recession substantially renewed interest in the possibility of obtaining more reliable GDP growth figures by combining the two measures (see e.g. Nalewaik (2010, 2011), Greenaway-McGrevy (2011) and especially Aruoba et al (2016), which provides the background for the Philadelphia Fed GDPplus measure). Some national statistical offices compute a simple equally weighted average of the different aggregate series, and in fact, BEA began providing this average in 2015. More sophisticated combination methods would give higher weights to the more precise GDP measures, as argued by Stone, Champernowne and Meade (1942) (see Weale (1992) for an account of the earlier literature).

As emphasized by Smith, Weale and Satchell (1998), though, dynamic considerations also matter because the contemporaneously filtered GDP series and its succesive updates as future

data becomes available will depend on the specification of the underlying processes. The secular growth in GDP and GDI has understandably led all previous studies to apply a signal-extraction framework to their growth rates, but doing so rules out by construction the possibility of saying anything about the level of U.S. output, which is of considerable interest on its own. In addition, taken literally, the absence of cointegration between the expenditure and income measures, with cointegrating vector (1,-1), implies an implausible diverging statistical discrepancy. Figure 2a contains the temporal evolution of the US quarterly (log) GDP and GDI series between 1984Q3 and 2015Q2, with shaded areas indicating NBER recessions. Although the two series differ, their (1,-1) cointegration relationship is evident. In turn, Figure 2b shows that their first differences are also highly correlated, but with a rich dynamic bivariate structure. Finally, Figure 2c makes clear that the statistical discrepancy is a persistent but stationary series whose movements are unrelated to the business cycle.

In view of the previous considerations, we prefer to formulate and estimate a model with covariance stationary measurement errors and an integrated common trend in the (log) levels of the two output measures.[12]Specifically, if y_{Et} and y_{It} denote (log) GDP and GDI, respectively, the model that we consider is

$$\begin{aligned}
\begin{pmatrix} y_{Et} \\ y_{It} \end{pmatrix} &= \begin{pmatrix} 1 \\ 1 \end{pmatrix} x_t + \begin{pmatrix} \epsilon_{Et} \\ \epsilon_{It} \end{pmatrix} \qquad (16)\\
(1-\rho_x L)(\Delta x_t - \mu) &= f_t \\
(1-\rho_{\epsilon_E} L)(\epsilon_{Et} - \delta/2) &= v_{Et} \\
(1-\rho_{\epsilon_I} L)(\epsilon_{It} + \delta/2) &= v_{It} \\
\begin{pmatrix} f_t \\ v_{Et} \\ v_{It} \end{pmatrix} \Bigg| \mathcal{I}_{t-1}, \boldsymbol{\phi} &\sim iid\ D\left[\begin{pmatrix} 0 \\ 0 \\ 0 \end{pmatrix}, \begin{pmatrix} \sigma_f^2 & 0 & 0 \\ 0 & \sigma_{v_E}^2 & 0 \\ 0 & 0 & \sigma_{v_I}^2 \end{pmatrix}, \boldsymbol{\varphi} \right],
\end{aligned}$$

where x_t is the "true GDP" common factor, whose rate of growth follows an AR(1) process with mean μ, autoregressive coefficient ρ_x and innovation variance σ_f^2, while ϵ_{Et} and ϵ_{It} are the measurement errors in the (log) expenditure and income measures, respectively, which follow covariance stationary AR(1) processes with unconditional means $\pm\delta/2$, autoregressive coefficients

[12] Arguably, a sufficiently flexible specification for the measurement errors in first-diferences by means of high-order autoregressive moving average processes may mitigate and, eventually, eliminate the consequences of ignoring cointegration, at least for the growth rates (see Almuzara, Fiorentini and Sentana (2018)).

ρ_{ϵ_E} and ρ_{ϵ_I}, and innovation variances $\sigma^2_{v_E}$ and $\sigma^2_{v_I}$.[13] Our specification of the serial correlation structure of the latent series follows from the empirical analysis in earlier versions of Fiorentini and Sentana (2017), who found evidence in favour of AR(1) processes for both the first difference of the common factor and the levels of the measurement errors. Importantly, our model allows for systematic biases in the measurement errors through δ, the difference between those biases determining the mean of the statistical discrepancy while their levels fixing the initial conditions.[14]

7.2 Estimation under the null and normality tests

We initially estimate the model using data from 1984Q3 to 2007Q2. We chose the final date to exclude the Great Recession from the sample. As for the start date, it marks the beginning of the so-called Great Moderation, as in Nalewaik (2010). We estimate the model in the time domain on the basis of the bivariate Gaussian likelihood of the stationarity-inducing transformation $\Delta y_{Et} + \Delta y_{It}$ and $y_{Et} - y_{It}$, systematically exploring its surface to make sure that we have found the global maximum. Panel A of Table 5 presents the estimates of the model parameters and their corresponding standard errors obtained from the asymptotic information matrix, which we compute using its frequency domain closed-form expression. As expected, we find that the growth rate of the "true" aggregate real output series is reasonably persistent. Our estimates also suggest that GDP provides a better measure of output than GDI, in the sense that GDP measurement errors have both a smaller autoregressive coefficient –in absolute value– and a smaller variance parameter. Indeed, the negative serial correlation coefficient for the GDP measurement error implies a tendency to compensate prior measurement errors, while the highly persistent GDI measurement error indicates that the difference between the growth rates of GDI and the true output measure are close to white noise.

[13] In terms of the formulation (1)–(2), we have that $\boldsymbol{\pi}(\boldsymbol{\theta}) = (\delta/2, -\delta/2)'$, $\boldsymbol{\xi}_t = (1, x_t, x_{t-1}, \epsilon_{Et}, \epsilon_{It})'$,

$$\mathbf{H}(\boldsymbol{\theta}) = \begin{pmatrix} 0 & 1 & 0 & 1 & 0 \\ 0 & 1 & 0 & 0 & 1 \end{pmatrix}, \mathbf{F}(\boldsymbol{\theta}) = \begin{pmatrix} 1 & 0 & 0 & 0 & 0 \\ \mu(1-\rho_x) & 1+\rho_x & -\rho_x & 0 & 0 \\ 0 & 1 & 0 & 0 & 0 \\ 0 & 0 & 0 & \rho_{\epsilon_E} & 0 \\ 0 & 0 & 0 & 0 & \rho_{\epsilon_I} \end{pmatrix}, \mathbf{M}(\boldsymbol{\theta}) = \begin{pmatrix} 0 & 0 & 0 \\ \sigma_f & 0 & 0 \\ 0 & 0 & 0 \\ 0 & \sigma_{v_E} & 0 \\ 0 & 0 & \sigma_{v_I} \end{pmatrix}$$

and $\boldsymbol{\varepsilon}_t = (f_t, v_{Et}, v_{It})'$.

[14] For identification purposes, though, we assume without loss of generality that the magnitude of those biases is the same for the two output series. We also assume that the two measurement errors are uncorrelated, which guarantees the non-parametric identification of the signal from the noise (see Almuzara, Fiorentini and Sentana (2018) for further details). The fact that the two measures of output are obtained from independent sources provides some plausibility to this assumption (but see Aruoba et al (2016)).

In turn, the normality tests reported in Panel B of Table 5 suggest that the soothing effects of the so–called Great Moderation propagated beyond second moments because the normality of the innovations to the underlying GDP growth rates is not rejected at conventional levels. On the other hand, we clearly reject the null of Gaussian innovations in the measurement errors. In fact, we reject not only when we use the joint test but also when we look at the skewness and kurtosis components separately. In contrast, the bivariate normality test of the reduced form innovations fails to reject its null hypothesis, which confirms the power advantages of looking at the structural innovations we documented in section 6.

To gain some further insight, in Figure 3 we plot the temporal evolution of the smoothed innovations (top panels), as well as the influence functions underlying the kurtosis tests (middle panels) and skewness tests (bottom panels) for both common factor (left panels) and measurement errors (right panels). Panels 3d and 3f indicate that an unusual measurement issue in both series around the first quarter of 2000 leads to the rejection of the Gaussian null for the measurement errors.

In Table 6 we present analogous results for a slightly larger sample that includes the Great Recession (1984Q3-2015Q2). As can be seen from Panel A, there are no dramatic changes in the parameter estimates, except perhaps for a higher persistence in the common factor, whose innovations have an unsurprisingly larger variance too. Nevertheless, the smoothed series are almost identical over the overlapping period. Figure 4 presents the evolution of the two output measures and our smoothed estimate in the period surrounding the Great Recession. As can be seen, GDP kept increasing over the entire 2007 while GDI began to show early warning signs of stagnation one year before. In the fourth quarter of 2008, though, both series experimented a dramatic drop, with GDI recovering slightly earlier than GDP. Our estimate tends to closely follow the GDP series, but taking into account the differing behavior of GDI around the turning points.

The large fall in output experienced in 2008Q4 implies that we also reject the normality of the common factor over this extended period. In that regard, we would like to emphasize that plots of the influence functions $s_{\mathrm{k},t|T}(\boldsymbol{\theta})$ and $\mathbf{s}_{\mathrm{s},t|T}(\boldsymbol{\theta})$'s seem to be more informative than plots of the smoothed innovations for the purposes of detecting non-normality. For example, Figure 5, which is entirely analogous to Figure 3 but including the Great Recession, confirms that 2008Q4 has a huge impact on the skewness and kurtosis scores of the common factor, resulting in a

strong rejection of the null.

Nevertheless, if we take a longer historical perspective, and start our sample soon after the Treasury - Federal Reserve Accord whereby the Fed stopped its wartime pegging of interest rates, the Great Recession is no longer an isolated outlier. There are several other periods, including the turbulences in the late 70's, early 80's, in which the normality of the "true GDP" innovations is clearly rejected (see Supplemental Material F for details).

7.3 The model under the alternative

Given those rejections, the natural next step is to estimate the parameters and obtain smoothed versions of the latent variables under the alternative distributions that we have considered. In view of the fact that the rejection of the null comes from both skweness and kurtosis, we consider an asymmetric Student t, a popular member of the asymmetric GH distribution, as DGP for the structural innovations. To estimate the model, we rely on a Metropolis-within-Gibbs algorithm which exploits the interpretation of the asymmetric Student t as a location-scale mixture of normals in (10). We estimate this model with 500,000 draws for the parameters and 250,000 for the latent variables, which correspond to 1 in 20 and 1 in 40 of the 10^7 original simulations (see Supplemental Material D for further details on the posterior simulator).

For the sake of brevity, we focus on the shape parameters, which are reported in Figure 6, with the left and right panels corresponding to the posterior distributions for the samples 1984Q3-2007Q2 and 1984Q3-2015Q2, respectively. Interestingly, when we exclude the Great Recession from the sample, the 95% credible intervals of all the skewness parameters include the origin. In the longer sample, in contrast, the asymmetry coefficient of the latent "true GDP" series becomes statistically significantly different from zero, which is in line with the evidence obtained from our proposed score test statistics in the previous section. Similarly, there is a shift in the mode and median of the reciprocal of the degrees of freedom (top panels) towards a lower number when we use the longer sample. The results in Figure F2 in the Supplemental Material F confirm the agreement between our proposed tests and the posterior intervals.

Finally, in Figure 7 we compare $\Delta x_{t|T}$ under the null and under asymmetric t innovations. In order to account for parameter uncertainty in both models, we also estimate the Gaussian specification using a simplified version of the MCMC algorithm which imposes $\eta = 0$ but uses the same number of draws. The top panel (Figure 7a) reports the median of the posteriors, while the bottom one (Figure 7b) reports the centered 95% error bands, computed by substracting the

median from the quantiles 97.5% and 2.5%. As can be seen, the median values are quite similar across distributions, but the drop in 2008Q4 seems to be sharper under asymmetric Student t innovations. Perhaps more interestingly, while we find that the asymmetric t seems to generate narrower (wider) intervals on the right (left) of the distribution in normal times, their magnitudes increase substantially during the Great Recession, exacerbating the asymmetry of the error band too. Importantly, this pattern starts to appear –albeit moderately– a few quarters before 2008Q4. In contrast, the Gaussian error bands are symmetric and almost constant irrespective of whether the economy is in a recession or not.

8 Conclusions

We exploit the EM principle to derive simple to implement and interpret LM-type tests of normality in all or a subset of the innovations to the latent variables in state space models against *GH* alternatives, which include the symmetric and asymmetric Student t, together with many other popular distributions. We decompose our tests into third and fourth moment components, and obtain one-sided LR analogues, whose asymptotic distribution we provide.

We perform a Monte Carlo study of the finite sample size and power of our procedures, explicitly comparing them to previously proposed tests. For all the models that we consider, our results detect a pile-up problem, whereby the fraction of negative values of the average kurtosis scores exceeds 50% under the null. For that reason, we employ a parametric bootstrap procedure, which improves the reliability of our tests under the null. In terms of power, we find that the most powerful test for any given alternative is usually the score test we have designed against it. We also find that while the tests that are based on the reduced form innovations have non-trivial power, they are clearly dominated by our proposed tests, which aim at the structural innovations.

When we apply our tests to a common trend model which combines the levels of the expenditure and income versions of US aggregate real output to improve its measurement, we reject normality of the innovations to the true GDP if the sample span extends beyond the Great Moderation (1984Q3-2007Q2). In contrast, the GDP/GDI measurement errors seem to be non-normal regardless of the period. For that reason, we develop a non-linear, simulation-based filtering procedure that improves over the Kalman filter, and highlights the importance of taking non-normality into account during turbulent periods such as the Great Recession.

From a methodological point of view, our EM-based approach can be successfully used in cross-sectional contexts too. In particular, it is straightforward to employ it for proving that many of the diagnostics suggested by Pagan and Vella (1989) for Tobit models do indeed coincide with the LM tests against specific alternatives in Chesher and Irish (1987) and Gouriéroux et al (1987). While the linearity implicit in (1)-(2) helps us obtain closed-form expressions for all the relevant quantities, it is not a requirement for applying our methodology in different contexts. Analyzing other latent variable models in which non-Gaussianity might be relevant constitutes a very interesting avenue for future research.

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Table 1: Monte Carlo rejection rates (in %) under null and alternative hypotheses for the bivariate cointegrated, dynamic single factor model ($T = 100$)

		Panel A: Null hypothesis			Panel B: Alternative hypotheses (5%) Student t			asymmetric Student t		
		1%	5%	10%	J	S_f	$S_{\mathbf{v}}$	J	S_f	$S_{\mathbf{v}}$
H_J	Kt	1.15	4.72	9.43	55.73	6.72	44.09	71.44	12.64	55.04
	Sk	1.00	4.92	10.30	31.77	6.79	25.05	68.09	17.31	50.62
	GH	1.02	4.67	9.79	48.13	6.88	37.11	74.04	16.26	57.02
H_{S_f}	Kt	0.94	4.71	9.60	19.54	13.83	6.70	39.00	26.72	13.38
	Sk	0.91	4.69	9.79	13.03	10.07	6.11	33.83	29.56	10.26
	GH	0.95	4.69	9.65	18.22	12.90	6.40	39.76	30.13	13.08
$H_{S_{\mathbf{v}}}$	Kt	1.08	4.75	9.70	48.35	4.84	46.40	58.30	5.02	55.61
	Sk	1.09	4.87	9.92	27.60	5.29	27.15	51.41	6.30	54.84
	GH	1.04	4.83	9.94	42.96	5.14	41.71	61.15	5.74	60.98
Red	Kt	1.04	4.76	9.58	53.15	7.71	37.89	70.70	14.98	48.17
	Sk	0.88	4.61	8.91	24.33	5.02	21.65	31.23	5.30	23.59
	GH	0.99	4.45	9.02	47.45	6.79	34.36	65.45	12.37	44.22

Notes: Results based on 10,000 samples of size $T = 100$ from model (16) with $\rho_x = .5$, $\rho_{\epsilon_E} = .2$, $\rho_{\epsilon_I} = .8$, $\sigma_f^2 = 1$ and $\sigma_{v_i}^2$ chosen such that $q_E = 2$ and $q_I = .5$, where $q_i = \sigma_x^2/\sigma_{\epsilon_i}^2$ represents the signal-to-noise ratio for y_{it} for $i = E, I$. The column labels J, S_f, $S_{\mathbf{v}}$ refer to the alternative $\boldsymbol{\varepsilon}_t \sim GH(\eta, \psi, \boldsymbol{\beta})$ (i.e. $R = 3$), $f_t \sim GH(\eta, \psi, \beta)$, $\mathbf{v}_t \sim N(\mathbf{0}, \mathbf{I}_N)$ ($R = 1$) and $\mathbf{v}_t \sim GH(\eta, \psi, \boldsymbol{\beta})$, $f_t \sim N(0, 1)$ ($R = 2$), respectively. The row labels H_J, H_{S_f}, and $H_{S_{\mathbf{v}}}$ refer to the score tests in Propositions 4 and 7 corresponding to the J, S_f, and $S_{\mathbf{v}}$ alternative hypotheses, while Red denotes the reduced form tests discussed in section 5.4.2. In Panel B, Student t refers to the DGP for the GH being symmetric Student t with 8 degrees of freedom and, analogously, asymmetric Student t to the asymmetric Student t with 8 degrees of freedom and skewness vector $\boldsymbol{\beta} = -\boldsymbol{\ell}_R$. For each of those labels, Kt and Sk refer to the kurtosis and skewness components of the corresponding test statistics, while GH indicates the sum of the two.

Table 2: Monte Carlo rejection rates (in %) under null and alternative hypotheses for the bivariate cointegrated, dynamic single factor model ($T = 250$)

		Panel A: Null hypothesis			Panel B: Alternative hypotheses (5%)					
					Student t			asymmetric Student t		
		1%	5%	10%	J	S_f	$S_{\mathbf{v}}$	J	S_f	$S_{\mathbf{v}}$
H_J	Kt	0.83	4.67	9.72	88.54	9.89	76.00	96.80	23.30	86.98
	Sk	1.02	5.33	10.19	42.42	8.77	33.85	95.50	36.18	82.65
	GH	0.98	4.99	9.85	80.82	9.73	66.07	98.55	34.51	90.56
H_{S_f}	Kt	1.07	4.81	9.79	34.44	22.74	8.27	64.40	48.53	22.74
	Sk	1.11	5.25	10.04	17.07	12.33	6.45	55.84	58.49	16.27
	GH	1.09	5.08	10.09	31.41	20.69	7.86	67.19	59.01	22.76
$H_{S_{\mathbf{v}}}$	Kt	0.86	4.78	9.78	81.86	5.60	79.33	91.86	6.87	88.03
	Sk	1.15	5.21	10.15	35.49	6.07	35.22	83.47	8.32	86.65
	GH	1.03	4.89	9.83	74.06	5.83	72.00	93.88	7.99	92.91
Red	Kt	0.93	4.68	9.61	85.85	11.43	66.66	96.22	27.25	80.75
	Sk	1.22	5.15	10.72	31.06	5.41	27.85	41.49	6.24	31.54
	GH	0.98	4.71	9.96	80.97	9.57	60.67	94.33	23.22	76.20

Notes: Results based on 10,000 samples of size $T = 250$ from model (16) with $\rho_x = .5$, $\rho_{\epsilon_E} = .2$, $\rho_{\epsilon_I} = .8$, $\sigma_f^2 = 1$ and $\sigma_{v_i}^2$ chosen such that $q_E = 2$ and $q_I = .5$, where $q_i = \sigma_x^2/\sigma_{\epsilon_i}^2$ represents the signal-to-noise ratio for y_{it} for $i = E, I$. The column labels J, S_f, $S_{\mathbf{v}}$ refer to the alternative $\boldsymbol{\varepsilon}_t \sim GH(\eta, \psi, \boldsymbol{\beta})$ (i.e. $R = 3$), $f_t \sim GH(\eta, \psi, \beta)$, $\mathbf{v}_t \sim N(\mathbf{0}, \mathbf{I}_N)$ ($R = 1$) and $\mathbf{v}_t \sim GH(\eta, \psi, \boldsymbol{\beta})$, $f_t \sim N(0, 1)$ ($R = 2$), respectively. The row labels H_J, H_{S_f}, and $H_{S_{\mathbf{v}}}$ refer to the score tests in Propositions 4 and 7 corresponding to the J, S_f, and $S_{\mathbf{v}}$ alternative hypotheses, while Red denotes the reduced form tests discussed in section 5.4.2. In Panel B, Student t refers to the DGP for the GH being symmetric Student t with 8 degrees of freedom and, analogously, asymmetric Student t to the asymmetric Student t with 8 degrees of freedom and skewness vector $\boldsymbol{\beta} = -\boldsymbol{\ell}_R$. For each of those labels, Kt and Sk refer to the kurtosis and skewness components of the corresponding test statistics, while GH indicates the sum of the two.

Table 3: Monte Carlo rejection rates (in %) under the null and alternative hypotheses for the local-level model

		Panel A: Null hypothesis			Panel B: Alternative hypotheses (5%) Student t			asymmetric Student t		
		1%	5%	10%	J	S_f	S_v	J	S_f	S_v
H_J	Kt	1.15	5.15	10.07	56.63	25.12	13.82	90.53	53.15	33.40
	Sk	1.06	5.20	10.26	24.27	13.33	8.87	95.14	63.39	36.51
	GH	1.19	5.02	10.33	48.81	21.71	12.22	95.64	64.17	39.28
H_{S_f}	Kt	1.14	5.23	10.69	47.35	29.64	7.80	83.55	59.20	16.38
	Sk	1.03	4.82	10.22	19.81	13.77	5.81	88.63	68.30	8.68
	GH	1.22	5.17	9.94	42.65	26.21	7.24	90.45	69.28	15.45
H_{S_v}	Kt	1.03	4.72	9.93	40.70	11.13	18.34	82.43	26.60	41.91
	Sk	1.05	4.89	9.92	14.67	6.47	9.49	72.92	8.18	43.37
	GH	1.04	4.70	9.84	35.77	9.94	15.97	84.85	22.82	47.82
Red	Kt	1.08	5.37	10.30	55.48	25.49	11.25	89.72	54.63	27.29
	Sk	1.17	4.99	10.04	22.31	13.11	6.83	94.90	63.05	16.58
	GH	1.20	5.22	10.09	49.66	22.93	10.34	95.58	64.10	26.14
HK_f	Kt	1.14	5.49	10.68	43.99	26.97	7.33	82.00	56.92	15.26
	Sk	1.04	4.83	10.19	19.82	13.75	5.79	88.67	68.30	8.73
	GH	1.22	5.23	9.96	41.95	25.67	7.06	90.29	69.15	15.06
HK_v	Kt	1.03	4.41	10.33	36.81	9.64	16.18	80.21	24.26	39.17
	Sk	1.05	4.89	9.99	14.66	6.51	9.49	72.91	8.19	43.38
	GH	1.05	4.81	9.98	35.25	9.70	15.51	84.54	22.41	47.29

Notes: Results based on 10,000 samples of size $T = 250$ from the local-level model discussed in section 5.3 in which the signal-to-noise ratio $q = \sigma_f^2/\sigma_v^2$ is set to 2. The column labels J, S_f, $S_{\mathbf{v}}$ refer to the alternative $\boldsymbol{\varepsilon}_t \sim GH(\eta, \psi, \boldsymbol{\beta})$ (i.e. $R = 2$), $f_t \sim GH(\eta, \psi, \beta)$, $v_t \sim N(0, 1)$ ($R = 1$) and $v_t \sim GH(\eta, \psi, \beta)$, $f_t \sim N(0, 1)$ ($R = 1$), respectively. The row labels H_J, H_{S_f}, and H_{S_v} refer to the score tests in Propositions 4 and 7 corresponding to the J, S_f, and S_v alternative hypotheses, Red denotes the reduced form tests discussed in section 5.4.2, while HK denotes the original Harvey and Koopman (1992) tests discussed in section 5.4.1. In Panel B, Student t refers to the DGP for the GH being symmetric Student t with 8 degrees of freedom and, analogously, asymmetric Student t to the asymmetric Student t with 8 degrees of freedom and skewness vector $\boldsymbol{\beta} = -\boldsymbol{\ell}_R$. For each of those labels, Kt and Sk refer to the kurtosis and skewness components of the corresponding test statistics, while GH indicates the sum of the two.

Table 4: Monte Carlo rejection rates (in %) under null and alternative hypotheses for the multivariate local-level model

		Panel A: Null hypothesis			Panel B: Alternative hypotheses (5%) Student t			asymmetric Student t		
		1%	5%	10%	J	S_f	$S_{\mathbf{v}}$	J	S_f	$S_{\mathbf{v}}$
H_J	Kt	0.91	4.79	9.51	100.00	13.63	100.00	100.00	34.12	100.00
	Sk	1.04	5.29	10.53	98.31	9.91	96.99	99.97	44.21	99.12
	GH	1.17	4.91	10.46	100.00	10.24	100.00	100.00	45.78	100.00
H_{S_f}	Kt	0.96	5.16	10.20	67.44	64.58	4.98	84.78	86.99	8.44
	Sk	1.08	5.36	9.60	28.13	58.96	5.10	86.72	95.72	5.88
	GH	0.95	5.60	10.14	62.40	59.98	5.44	90.96	95.81	7.95
$H_{S_{\mathbf{v}}}$	Kt	0.67	5.05	9.83	100.00	5.35	100.00	100.00	5.34	100.00
	Sk	0.87	5.27	10.13	97.70	4.98	97.48	99.90	5.26	99.48
	GH	0.93	5.33	10.06	100.00	4.90	100.00	100.00	5.29	100.00
Red	Kt	0.95	4.86	9.45	100.00	14.48	100.00	100.00	36.05	100.00
	Sk	1.10	5.18	10.23	98.18	10.84	96.63	99.93	46.46	98.15
	GH	1.10	4.91	10.46	100.00	11.10	100.00	100.00	48.75	100.00

Notes: Results based on 10,000 samples of size $T = 250$ from a 10-variate version of the local-level model with $\boldsymbol{\pi} = \mathbf{0}$, $\mathbf{c} = \boldsymbol{\ell}_{10}$ and $\boldsymbol{\gamma} = q^{-1}\boldsymbol{\ell}_{10}$, where q reflects the signal-to-noise ratio, which we set to 2. The column labels J, S_f, $S_{\mathbf{v}}$ refer to the alternative $\boldsymbol{\varepsilon}_t \sim GH(\eta, \psi, \boldsymbol{\beta})$ (i.e. $R = 11$), $f_t \sim GH(\eta, \psi, \beta)$, $\mathbf{v}_t \sim N(\mathbf{0}, \mathbf{I}_N)$ ($R = 1$) and $\mathbf{v}_t \sim GH(\eta, \psi, \boldsymbol{\beta})$, $f_t \sim N(0, 1)$ ($R = 10$), respectively. The row labels H_J, H_{S_f}, and $H_{S_{\mathbf{v}}}$ refer to the score tests in Propositions 4 and 7 corresponding to the J, S_f, and $S_{\mathbf{v}}$ alternative hypotheses. In Panel B, Student t refers to the DGP for the GH being symmetric Student t with 8 degrees of freedom and, analogously, asymmetric Student t to the asymmetric Student t with 8 degrees of freedom and skewness vector $\boldsymbol{\beta} = -\boldsymbol{\ell}_R$. For each of those labels, Kt and Sk refer to the kurtosis and skewness components of the corresponding test statistics, while GH indicates the sum of the two.

Table 5: Parameter estimates and normality tests during Great Moderation

Panel A: ML estimates

Param.	estimate	std. err.
μ	0.765	0.330
δ	0.181	0.040
ρ_x	0.536	0.105
ρ_{ϵ_E}	-0.672	0.152
ρ_{ϵ_I}	0.940	0.036
σ_f^2	0.135	0.027
$\sigma_{v_E}^2$	0.010	0.005
$\sigma_{v_I}^2$	0.153	0.025

Panel B: Normality tests

		statistic	p-value
H_{S_f}	Kt	0.646	0.211
	Sk	1.540	0.215
	GH	2.186	0.237
$H_{S_\mathbf{v}}$	Kt	5.901	0.008
	Sk	7.914	0.019
	GH	13.815	0.002
Red	Kt	1.585	0.104
	Sk	1.478	0.478
	GH	3.063	0.299

Notes: Data: Quarterly real GDP and GDI from 1984Q3 to 2007Q2. Model: Bivariate cointegrated, dynamic single factor model (16); see section 7 for parameter definitions. In Panel A, estimates are Gaussian ML of the bivariate Gaussian likelihood of the stationary transformation $\Delta y_{Et} + \Delta y_{It}$ and $y_{Et} - y_{It}$ in the time domain. Standard errors are obtained from the asymptotic information matrix, which is computed using its frequency domain closed-form expression. In Panel B, the row labels H_{S_f} and $H_{S_\mathbf{v}}$ refer to the score tests in Propositions 4 and 7 corresponding to the S_f ($R = 1$) and $S_\mathbf{v}$ ($R = 2$) alternative hypotheses, respectively, while Red denotes the reduced form tests discussed in section 5.4.2. For each of those labels, Kt and Sk refer to the kurtosis and skewness components of the corresponding test statistics, while GH indicates the sum of the two.

Table 6: Parameter estimates and normality tests during the Great moderation and the Great Recession

Panel A: ML estimates

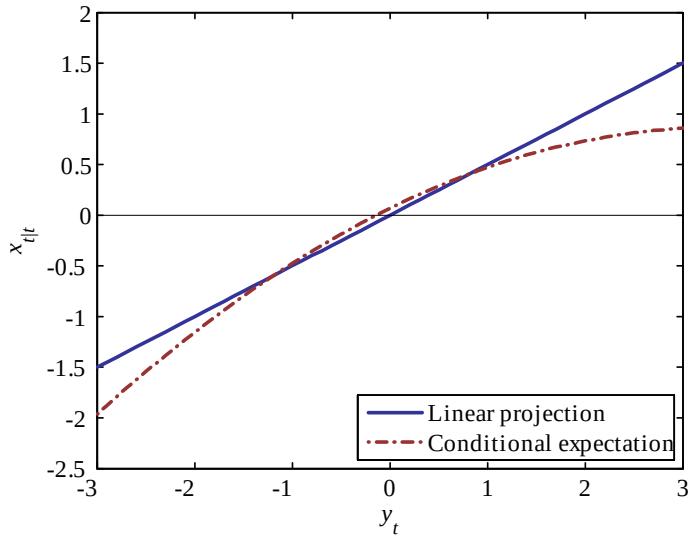
Param.	estimate	std. err.
μ	0.642	0.196
δ	0.033	0.036
ρ_x	0.643	0.080
ρ_{ϵ_E}	-0.384	0.204
ρ_{ϵ_I}	0.938	0.032
σ^2_f	0.169	0.031
$\sigma^2_{v_E}$	0.022	0.010
$\sigma^2_{v_I}$	0.150	0.023

Panel B: Normality tests

		statistic	p-value
H_{S_f}	Kt	64.691	0.000
	Sk	22.542	0.000
	GH	87.233	0.000
$H_{S_\mathbf{v}}$	Kt	8.210	0.002
	Sk	4.398	0.111
	GH	12.607	0.004
Red	Kt	20.828	0.000
	Sk	7.818	0.020
	GH	28.645	0.000

Notes: Data: Quarterly real GDP and GDI from 1984Q3 to 2015Q2. Model: Bivariate cointegrated, dynamic single factor model (16); see section 7 for parameter definitions. In Panel A, estimates are Gaussian ML of the bivariate Gaussian likelihood of the stationary transformation $\Delta y_{Et} + \Delta y_{It}$ and $y_{Et} - y_{It}$ in the time domain. Standard errors are obtained from the asymptotic information matrix, which is computed using its frequency domain closed-form expression. In Panel B, the row labels H_{S_f} and $H_{S_\mathbf{v}}$ refer to the score tests in Propositions 4 and 7 corresponding to the S_f ($R = 1$) and $S_\mathbf{v}$ ($R = 2$) alternative hypotheses, respectively, while Red denotes the reduced form tests discussed in section 5.4.2. For each of those labels, Kt and Sk refer to the kurtosis and skewness components of the corresponding test statistics, while GH indicates the sum of the two.

Figure 1: Linear projection versus conditional expectation in a non-Gaussian univariate static factor model



Notes: The observed variable is $y_t = \frac{1}{2}x_t + \frac{\sqrt{3}}{2}\epsilon_t$. We assume that the joint distribution of x_t and ϵ_t is asymmetric Student t with zero mean, identity covariance matrix, 8 degrees of freedom and skewness vector parameter $\mathbf{b} = (-1, 0)'$. Given that the joint distribution of y_t and x_t will also be an asymmetric Student t, we can use the expressions in Mencía (2012) to compute the conditional expectation of x_t given y_t.

Figure 2: Expenditure (GDP) and income (GDI) measures of real output

Figure 2a: Quarterly real (log) GDP and (log) GDI

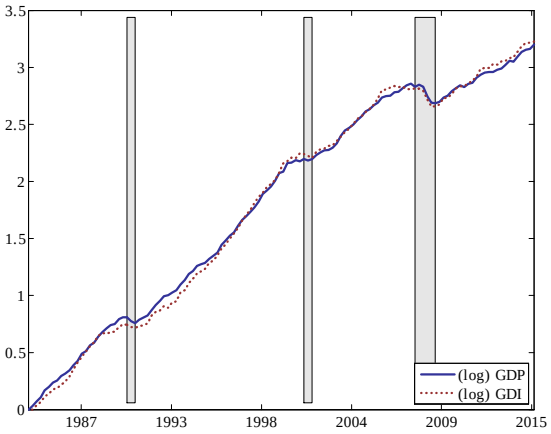


Figure 2b: Quarterly real GDP and GDI growth

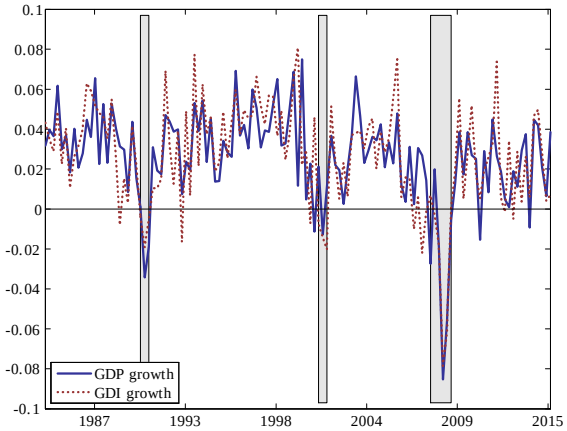
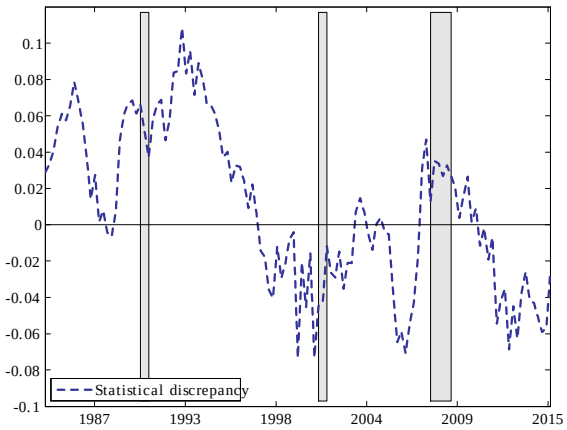


Figure 2c: Statistical discrepancy



Notes: Data: Quarterly real GDP and GDI from 1984Q3 to 2015Q2. Statistical discrepancy is defined as $\log(GDP) - \log(GDI)$. Shaded areas represent NBER recessions.

Figure 3: Smoothed innovations and influence functions for the kurtosis and skewness tests: Sample 1984Q3 to 2007Q2.

Figure 3a: Smoothed innovations for the underlying factor

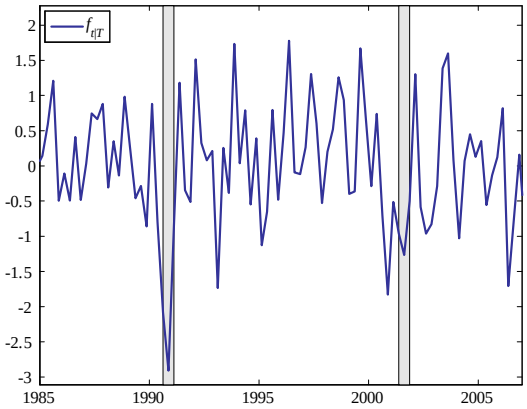


Figure 3b: Smoothed innovations for the measurement errors

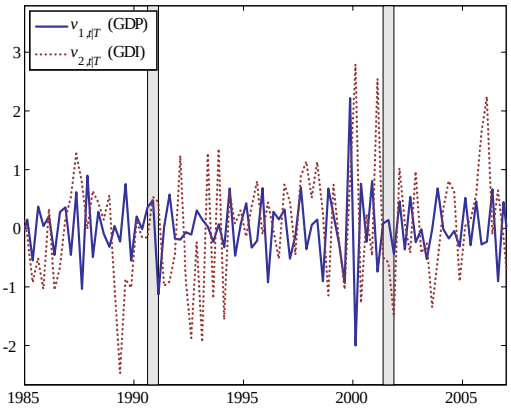


Figure 3c: Influence functions for the underlying factor (kurtosis)

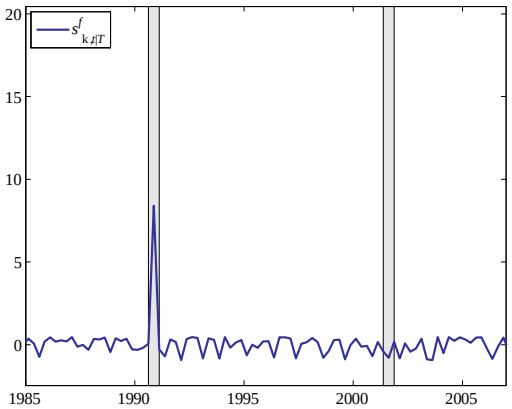


Figure 3d: Influence functions for the measurement errors (kurtosis)

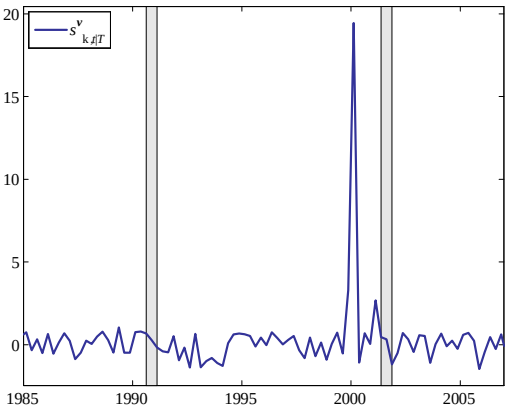


Figure 3e: Influence functions for the underlying factor (skewness)

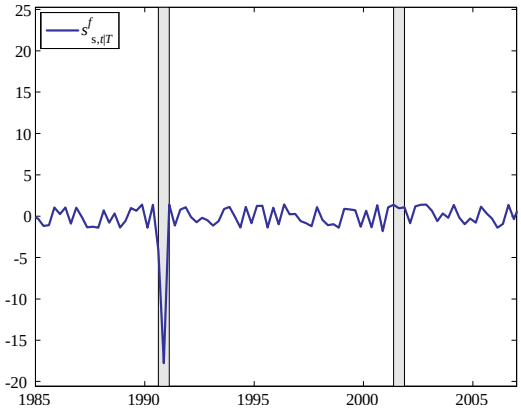
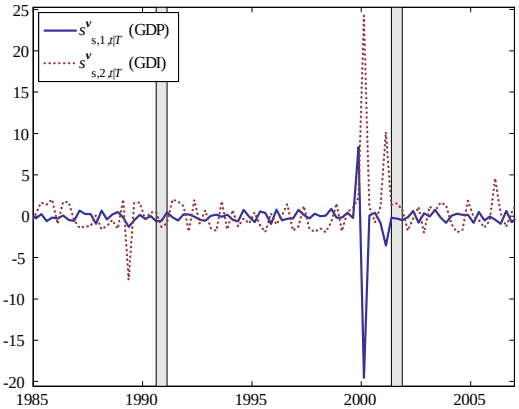
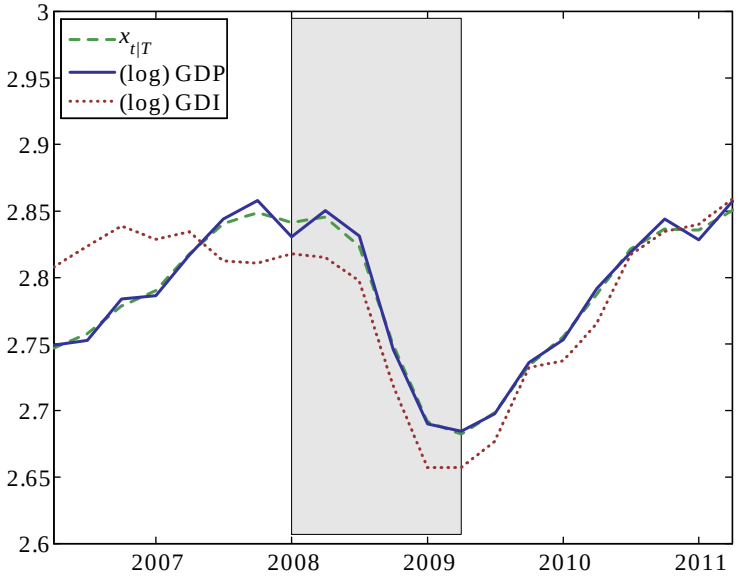


Figure 3f: Influence functions for the measurement errors (skewness)



Notes: Smoothed innovations and influence functions were obtained from fitting the bivariate cointegrated, dynamic single factor model (16) to the quarterly real GDP and GDI from 1984Q3 to 2007Q2; see Table 4 for parameter estimates. Shaded areas represent NBER recessions.

Figure 4: GDP, GDI and smoothed estimate of real output around the Great Recession



Notes: The smoothed estimate $x_{t|T}$ was obtained from fitting the bivariate cointegrated, dynamic single factor model (16) to the quarterly real GDP and GDI from 1984Q3 to 2015Q2; see Table 5 for parameter estimates. The shaded area represents the NBER recession.

Figure 5: Smoothed innovations and influence functions for the kurtosis and skewness tests: Sample 1984Q3 to 2015Q2.

Figure 5a: Smoothed innovations for the underlying factor

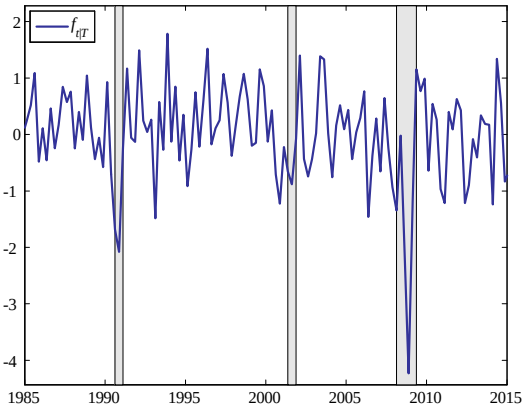


Figure 5b: Smoothed innovations for the measurement errors

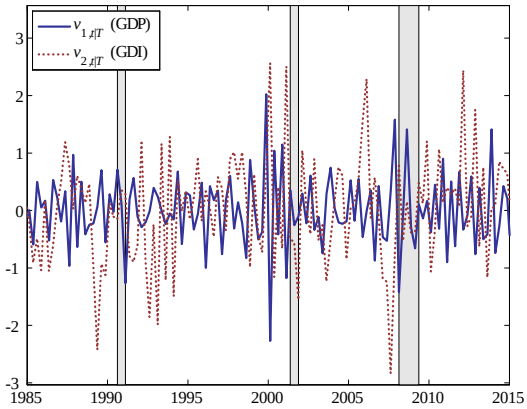


Figure 5c: Influence functions for the underlying factor (kurtosis)

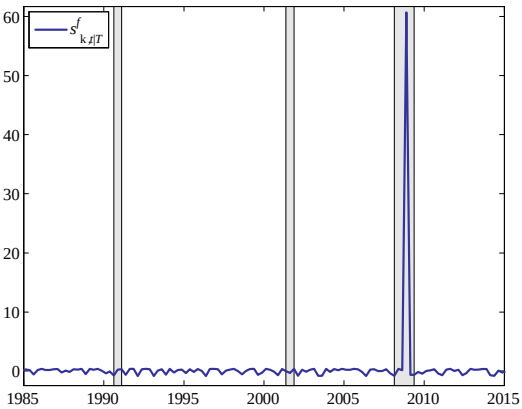


Figure 5d: Influence functions for the measurement errors (kurtosis)

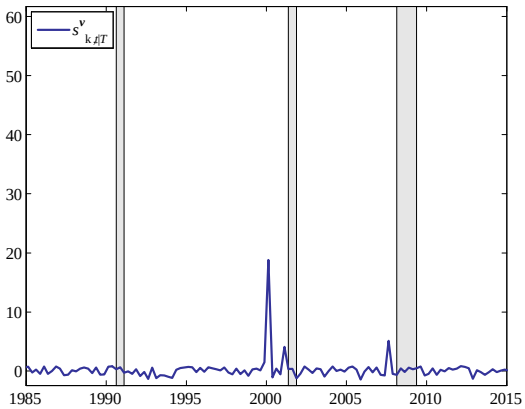


Figure 5e: Influence functions for the underlying factor (skewness)

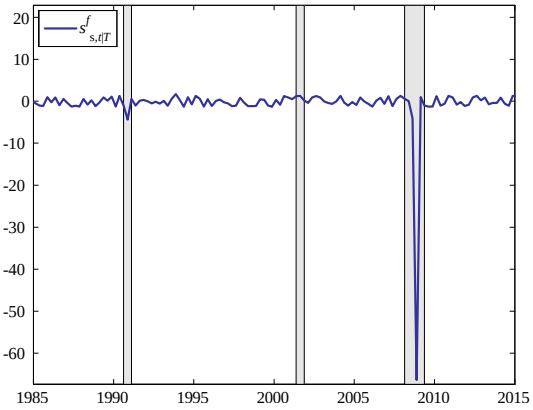
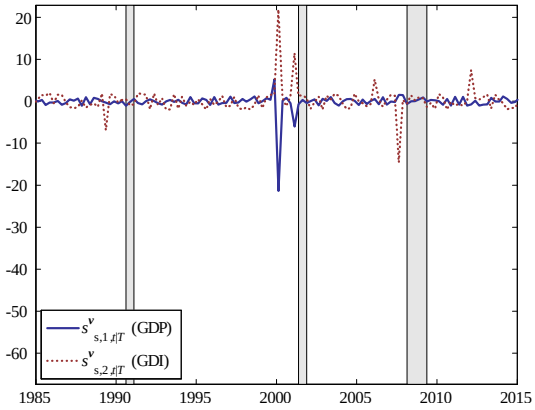


Figure 5f: Influence functions for the measurement errors (skewness)



Notes: Smoothed innovations and influence functions were obtained from fitting the bivariate cointegrated, dynamic single factor model (16) to the quarterly real GDP and GDI from 1984Q3 to 2015Q2; see Table 5 for parameter estimates. Shaded areas represent NBER recessions.

Figure 6: Posterior densities of shape parameters under the asymmetric Student t alternative

Figure 6a: η, 1984Q3 to 2007Q2

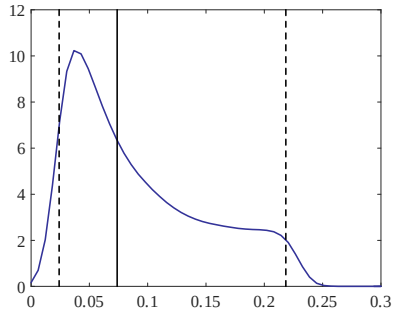


Figure 6b: η, 1984Q3 to 2015Q2

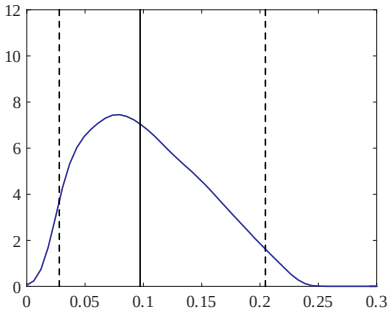


Figure 6c: β_x, 1984Q3 to 2007Q2

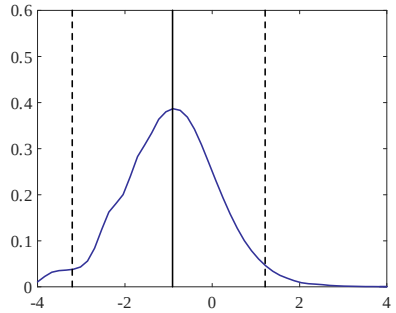


Figure 6d: β_x, 1984Q3 to 2015Q2

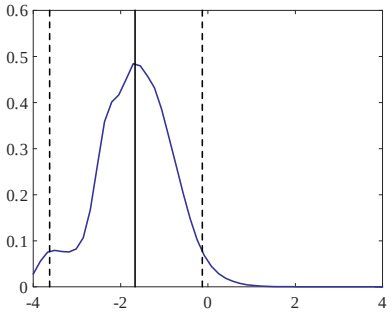


Figure 6e: β_{v_E}, 1984Q3 to 2007Q2

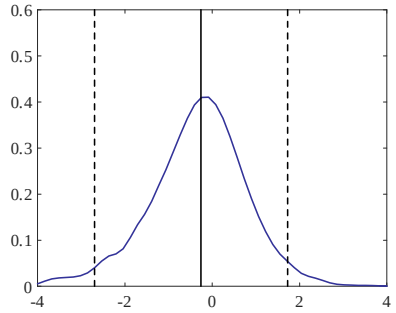


Figure 6f: β_{v_E}, 1984Q3 to 2015Q2

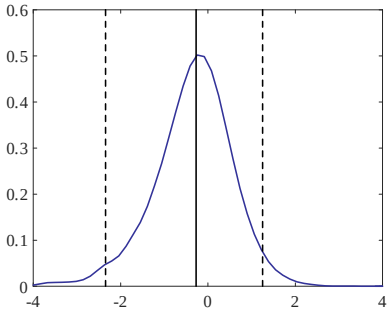


Figure 6g: β_{v_I}, 1984Q3 to 2007Q2

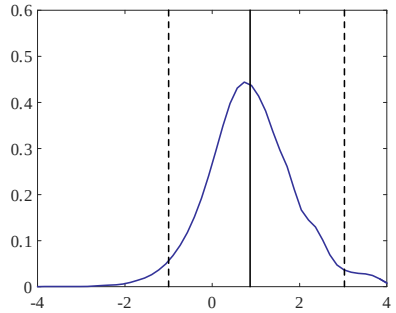
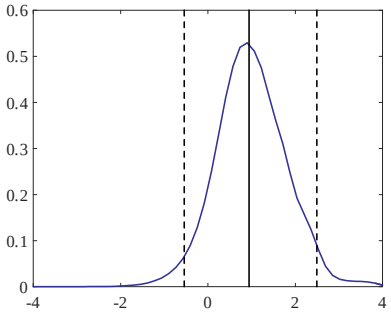


Figure 6h: β_{v_I}, 1984Q3 to 2015Q2



Notes: Data: Quarterly real GDP and GDI from 1984Q3 to 2007Q2 (2015Q2) in left (right) panels. Model: Bivariate cointegrated, dynamic single factor model (16) with multivariate asymmetric Student t innovations; see Section 7 for parameter definitions. η refers to the reciprocal of degrees of freedom while β_x (β_{v_E}) [β_{v_I}] refers to the skewness parameter of the "true GDP" (expenditure) [income] measure. Solid vertical lines refer to the median values while dashed lines report the 2.5% and 97.5% quantiles.

Figure 7: Smoothed "true GDP" growth under Gaussian and asymmetric Student t innovations

Figure 7a: Posterior median of $\Delta x_{t|T}$

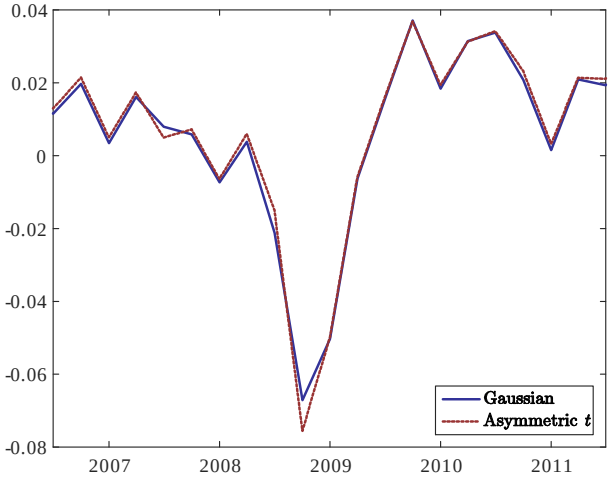


Figure 7b: Posterior 95% error bands for $\Delta x_{t|T}$

Notes: Data: Quarterly real GDP and GDI from 1984Q3 to 2015Q2. Model: Bivariate cointegrated, dynamic single factor model (16) with multivariate asymmetric Student t innovations; see Section 7 for parameter definitions. Results are based on 25,000 draws from the posterior simulator. Error bands refer to the 2.5% and 97.5% quantiles from which the median values where substracted.

Supplemental Material for

Normality tests for latent variables

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A Proofs

Lemmata

Lemma 1 *Let $\mathbf{z} \sim N(\boldsymbol{\mu}, \boldsymbol{\Sigma})$ be an n_z-dimensional real Gaussian random vector. Then,*
i) Expectation of second powers:

$$E(\mathbf{z}\mathbf{z}') = \boldsymbol{\mu}\boldsymbol{\mu}' + \boldsymbol{\Sigma},$$

ii) Expectation of third powers:

$$E[\mathbf{z}(\mathbf{z} \odot \mathbf{z})'] = \boldsymbol{\mu}(\boldsymbol{\mu} \odot \boldsymbol{\mu})' + 2(\boldsymbol{\Sigma} \odot \boldsymbol{\ell}_{n_z}\boldsymbol{\mu}') + \boldsymbol{\mu} vecd'(\boldsymbol{\Sigma}),$$

iii) Expectation of fourth powers:

$$\begin{aligned} E\left[(\mathbf{z} \odot \mathbf{z})(\mathbf{z} \odot \mathbf{z})'\right] = {} & (\boldsymbol{\mu} \odot \boldsymbol{\mu})(\boldsymbol{\mu} \odot \boldsymbol{\mu})' + 2(\boldsymbol{\Sigma} \odot \boldsymbol{\Sigma}) + vecd(\boldsymbol{\Sigma})vecd'(\boldsymbol{\Sigma}) \\ & + 4(\boldsymbol{\Sigma} \odot \boldsymbol{\mu}\boldsymbol{\mu}') + vecd(\boldsymbol{\mu}\boldsymbol{\mu}')vecd'(\boldsymbol{\Sigma}) + vecd(\boldsymbol{\Sigma})vecd'(\boldsymbol{\mu}\boldsymbol{\mu}'), \end{aligned}$$

where $\odot$ denotes the Hadamard (or elementwise) product, $vecd(\cdot)$ is the operator which stacks the diagonal elements of a square matrix in vector form and $\boldsymbol{\ell}_{n_z}$ is a vector of n_z ones.

Proof. The proof is tedious but straightforward. □

Lemma 2 *Define $\mathbf{m}_h : \mathbb{R}^{n_1} \times \mathbb{R}^{n_2} \to \mathbb{R}^{n_1 \times n_2}$ for $n_1, n_2 \in \mathbb{Z}_{++}$ and $h \in \{2, 3, 4\}$ as*

$$\begin{aligned} \mathbf{m}_2(\mathbf{w}_1, \mathbf{w}_2) &= vec(\mathbf{w}_1\mathbf{w}_2'), \\ \mathbf{m}_3(\mathbf{w}_1, \mathbf{w}_2) &= vec[\mathbf{w}_1(\mathbf{w}_2 \odot \mathbf{w}_2)'], \\ \mathbf{m}_4(\mathbf{w}_1, \mathbf{w}_2) &= vec[(\mathbf{w}_1 \odot \mathbf{w}_1)(\mathbf{w}_2 \odot \mathbf{w}_2)'], \end{aligned}$$

where $\mathbf{w}_1 \in \mathbb{R}^{n_1}$, $\mathbf{w}_2 \in \mathbb{R}^{n_2}$, and $vec(\cdot)$ is the vectorization (by columns) operator. Consider the real Gaussian random vector

$$\begin{pmatrix} \mathbf{x} \\ \mathbf{y} \\ \mathbf{z} \end{pmatrix} \sim N\left[\begin{pmatrix} \boldsymbol{\mu}_x \\ \boldsymbol{\mu}_y \\ \boldsymbol{\mu}_z \end{pmatrix}, \begin{pmatrix} \boldsymbol{\Sigma}_{xx} & \boldsymbol{\Sigma}_{xy} & \boldsymbol{\Sigma}_{xz} \\ \boldsymbol{\Sigma}_{xy}' & \boldsymbol{\Sigma}_{yy} & \boldsymbol{\Sigma}_{yz} \\ \boldsymbol{\Sigma}_{xz}' & \boldsymbol{\Sigma}_{yz}' & \boldsymbol{\Sigma}_{zz} \end{pmatrix}\right]$$

where $\mathbf{x}$ is n_x-dimensional, $\mathbf{y}$ is n_y-dimensional, and $\mathbf{z}$ is n_z-dimensional. Then,
i) Covariance with the first power:

$$\begin{aligned} cov\left[\mathbf{x}, \mathbf{m}_2(\mathbf{y}, \mathbf{z})\right] &= \mathbf{0}, \\ cov\left[\mathbf{x}, \mathbf{m}_3(\mathbf{y}, \mathbf{z})\right] &= 2\left[\boldsymbol{\ell}_{n_x} \otimes vec'(\boldsymbol{\Sigma}_{yz})\right] \odot (\boldsymbol{\Sigma}_{xz} \otimes \boldsymbol{\ell}_{n_y}') \\ &\quad + \left[vecd'(\boldsymbol{\Sigma}_{zz}) \otimes \mathbf{1}_{n_x \times n_y}\right] \odot \left(\boldsymbol{\ell}_{n_z}' \otimes \boldsymbol{\Sigma}_{xy}\right), \\ cov\left[\mathbf{x}, \mathbf{m}_4(\mathbf{y}, \mathbf{z})\right] &= \mathbf{0}, \end{aligned}$$

ii) Covariance with the second power:

$$\begin{aligned} cov\left[\mathbf{m}_2(\mathbf{x}, \mathbf{x}), \mathbf{m}_2(\mathbf{y}, \mathbf{z})\right] &= (\mathbf{1}_{n_x \times n_z} \otimes \boldsymbol{\Sigma}_{xy}) \odot \left(\boldsymbol{\Sigma}_{xz} \otimes \mathbf{1}_{n_x \times n_y}\right) \\ &\quad + (\boldsymbol{\ell}_{n_x} \otimes \boldsymbol{\Sigma}_{xz} \otimes \boldsymbol{\ell}_{n_y}') \odot \left(\boldsymbol{\ell}_{n_x}' \otimes \boldsymbol{\Sigma}_{xy} \otimes \boldsymbol{\ell}_{n_z}\right), \\ cov\left[\mathbf{m}_2(\mathbf{x}, \mathbf{x}), \mathbf{m}_3(\mathbf{y}, \mathbf{z})\right] &= \mathbf{0}, \end{aligned}$$

$$cov\left[\mathbf{m}_2(\mathbf{x},\mathbf{x}),\mathbf{m}_4(\mathbf{y},\mathbf{z})\right] = 4\left[\boldsymbol{\ell}_{n_x^2}\otimes vec'(\boldsymbol{\Sigma}_{yz})\right]\odot cov\left[\mathbf{m}_2(\mathbf{x},\mathbf{x}),\mathbf{m}_2(\mathbf{y},\mathbf{z})\right]$$
$$+2\left[\boldsymbol{\ell}_{n_x^2}\otimes\boldsymbol{\ell}'_{n_z}\otimes vecd'(\boldsymbol{\Sigma}_{yy})\right]\odot\left(\boldsymbol{\ell}_{n_x}\otimes\boldsymbol{\Sigma}_{xz}\otimes\boldsymbol{\ell}'_{n_y}\right)\odot\left(\boldsymbol{\Sigma}_{xz}\otimes\mathbf{1}_{n_x\times n_y}\right)$$
$$+2[\boldsymbol{\ell}_{n_x^2}\otimes vecd'(\boldsymbol{\Sigma}_{zz})\otimes\boldsymbol{\ell}'_{n_y}]\odot(\mathbf{1}_{n_x\times n_z}\otimes\boldsymbol{\Sigma}_{xy})\odot(\boldsymbol{\ell}'_{n_y}\otimes\boldsymbol{\Sigma}_{xy}\otimes\boldsymbol{\ell}_{n_x}),$$

iii) Covariance with the third power:

$$cov\left[\mathbf{m}_3(\mathbf{x},\mathbf{x}),\mathbf{m}_3(\mathbf{y},\mathbf{z})\right] = [vecd(\boldsymbol{\Sigma}_{xx})\otimes\boldsymbol{\ell}_{n_x}\otimes\boldsymbol{\ell}'_{n_yn_z}]\odot\{\boldsymbol{\ell}_{n_x}\otimes cov\left[\mathbf{x},\mathbf{m}_3(\mathbf{y},\mathbf{z})\right]\}$$
$$+2\left(\mathbf{1}_{n_x\times n_z}\otimes\boldsymbol{\Sigma}_{xy}\right)\odot\left[(\boldsymbol{\Sigma}_{xz}\odot\boldsymbol{\Sigma}_{xz})\otimes\mathbf{1}_{n_x\times n_y}\right]$$
$$+2[vec(\boldsymbol{\Sigma}_{xx})\otimes\boldsymbol{\ell}'_{n_yn_z}]\odot[\boldsymbol{\ell}_{n_x^2}\otimes vecd'(\boldsymbol{\Sigma}_{zz})\otimes\boldsymbol{\ell}'_{n_y}]\odot(\boldsymbol{\ell}'_{n_y}\otimes\boldsymbol{\Sigma}_{xy}\otimes\boldsymbol{\ell}_{n_x})$$
$$+4[vec(\boldsymbol{\Sigma}_{xx})\otimes\boldsymbol{\ell}'_{n_yn_z}]\odot\left[\boldsymbol{\ell}_{n_x^2}\otimes vec'(\boldsymbol{\Sigma}_{yz})\right]\odot\left(\boldsymbol{\Sigma}_{xz}\otimes\mathbf{1}_{n_x\times n_y}\right)$$
$$+4(\boldsymbol{\ell}_{n_x}\otimes\boldsymbol{\Sigma}_{xz}\otimes\boldsymbol{\ell}'_{n_y})\odot\left(\boldsymbol{\ell}'_{n_z}\otimes\boldsymbol{\Sigma}_{xy}\otimes\boldsymbol{\ell}_{n_x}\right)\odot\left(\boldsymbol{\Sigma}_{xz}\otimes\mathbf{1}_{n_x\times n_y}\right),$$
$$cov\left[\mathbf{m}_3(\mathbf{x},\mathbf{x}),\mathbf{m}_4(\mathbf{y},\mathbf{z})\right] = \mathbf{0},$$

iv) Covariance with the fourth power:

$$cov\left[\mathbf{m}_4(\mathbf{x},\mathbf{x}),\mathbf{m}_4(\mathbf{y},\mathbf{z})\right] = 4cov\left[\mathbf{m}_2(\mathbf{x},\mathbf{x}),\mathbf{m}_2(\mathbf{y},\mathbf{z})\right]\odot cov\left[\mathbf{m}_2(\mathbf{x},\mathbf{x}),\mathbf{m}_2(\mathbf{y},\mathbf{z})\right]$$
$$+4[vec(\boldsymbol{\Sigma}_{xx})\otimes\boldsymbol{\ell}'_{n_yn_z}]\odot cov\left[\mathbf{m}_2(\mathbf{x},\mathbf{x}),\mathbf{m}_4(\mathbf{y},\mathbf{z})\right]$$
$$+2[\boldsymbol{\ell}_{n_x}\otimes vecd(\boldsymbol{\Sigma}_{xx})\otimes\boldsymbol{\ell}'_{n_yn_z}]\odot\left[\boldsymbol{\ell}_{n_x^2}\otimes\boldsymbol{\ell}_{n_z}\otimes vecd'(\boldsymbol{\Sigma}_{yy})\right]\odot\left(\boldsymbol{\Sigma}_{xz}\otimes\mathbf{1}_{n_x\times n_y}\right)\odot\left(\boldsymbol{\Sigma}_{xz}\otimes\mathbf{1}_{n_x\times n_y}\right)$$
$$+2[\boldsymbol{\ell}_{n_x}\otimes vecd(\boldsymbol{\Sigma}_{xx})\otimes\boldsymbol{\ell}'_{n_yn_z}]\odot\left[\boldsymbol{\ell}_{n_x^2}\otimes vecd'(\boldsymbol{\Sigma}_{zz}\otimes\boldsymbol{\ell}_{n_y})\right]\odot\left(\boldsymbol{\ell}'_{n_z}\otimes\boldsymbol{\Sigma}_{xy}\otimes\boldsymbol{\ell}_{n_x}\right)\odot\left(\boldsymbol{\ell}'_{n_z}\otimes\boldsymbol{\Sigma}_{xy}\otimes\boldsymbol{\ell}_{n_x}\right)$$
$$+2[vecd(\boldsymbol{\Sigma}_{xx})\otimes\boldsymbol{\ell}_{n_x}\otimes\boldsymbol{\ell}'_{n_yn_z}]\odot\left[\boldsymbol{\ell}_{n_x^2}\otimes\boldsymbol{\ell}_{n_z}\otimes vecd'(\boldsymbol{\Sigma}_{yy})\right]\odot(\boldsymbol{\ell}_{n_x}\otimes\boldsymbol{\Sigma}_{xz}\otimes\boldsymbol{\ell}'_{n_y})\odot(\boldsymbol{\ell}_{n_x}\otimes\boldsymbol{\Sigma}_{xz}\otimes\boldsymbol{\ell}'_{n_y})$$
$$+2[vecd(\boldsymbol{\Sigma}_{xx})\otimes\boldsymbol{\ell}_{n_x}\otimes\boldsymbol{\ell}'_{n_yn_z}]\odot\left[\boldsymbol{\ell}_{n_x^2}\otimes vecd'(\boldsymbol{\Sigma}_{zz}\otimes\boldsymbol{\ell}_{n_y})\right]\odot(\mathbf{1}_{n_x\times n_z}\otimes\boldsymbol{\Sigma}_{xy})\odot(\mathbf{1}_{n_x\times n_z}\otimes\boldsymbol{\Sigma}_{xy})$$
$$+8[\boldsymbol{\ell}_{n_x}\otimes vecd(\boldsymbol{\Sigma}_{xx})\otimes\boldsymbol{\ell}'_{n_yn_z}]\odot\left[\boldsymbol{\ell}_{n_x^2}\otimes vec'(\boldsymbol{\Sigma}_{yz})\right]\odot\left(\boldsymbol{\ell}'_{n_z}\otimes\boldsymbol{\Sigma}_{xy}\otimes\boldsymbol{\ell}_{n_x}\right)\odot\left(\boldsymbol{\Sigma}_{xz}\otimes\mathbf{1}_{n_x\times n_y}\right)$$
$$+8[vecd(\boldsymbol{\Sigma}_{xx})\otimes\boldsymbol{\ell}_{n_x}\otimes\boldsymbol{\ell}'_{n_yn_z}]\odot\left[\boldsymbol{\ell}_{n_x^2}\otimes vec'(\boldsymbol{\Sigma}_{yz})\right]\odot(\mathbf{1}_{n_x\times n_z}\otimes\boldsymbol{\Sigma}_{xy})\odot(\boldsymbol{\ell}_{n_x}\otimes\boldsymbol{\Sigma}_{xz}\otimes\boldsymbol{\ell}'_{n_y})$$
$$+8\left(\boldsymbol{\Sigma}_{xy}\otimes\mathbf{1}_{n_x\times n_z}\right)\odot\left(\mathbf{1}_{n_x\times n_y}\otimes\boldsymbol{\Sigma}_{xz}\right)\odot\left(\boldsymbol{\ell}'_{n_z}\otimes\boldsymbol{\Sigma}_{xy}\otimes\boldsymbol{\ell}_{n_x}\right)\odot(\boldsymbol{\ell}_{n_x}\otimes\boldsymbol{\Sigma}_{xz}\otimes\boldsymbol{\ell}'_{n_y}),$$

where $\otimes$ denotes Kronecker product and $\mathbf{1}_{n_1\times n_2}$ denotes a matrix of ones of dimension $n_1\times n_2$.

Proof. Again, the proof is tedious but straightforward. □

Lemma 3 *Consider the model (1)-(2) where $\boldsymbol{\varepsilon}_t = (\boldsymbol{\varepsilon}_t^{\mathrm{GH}\prime},\boldsymbol{\varepsilon}_t^{\mathrm{N}\prime})'$, with $\boldsymbol{\varepsilon}_t^{\mathrm{GH}}\sim GH_R(\eta,\psi,\boldsymbol{\beta})$ and $\boldsymbol{\varepsilon}_t^{\mathrm{N}}\sim N(\mathbf{0};\mathbf{I}_{K-R})$. Let $\varsigma_t^{\mathrm{GH}} = \boldsymbol{\varepsilon}_t^{\mathrm{GH}\prime}\boldsymbol{\varepsilon}_t^{\mathrm{GH}}$ and*

$$s_{\mathrm{k}t} = c_0 + c_1\varsigma_t^{\mathrm{GH}} + c_2(\varsigma_t^{\mathrm{GH}})^2,$$
$$\mathbf{s}_{\mathrm{s}t} = \boldsymbol{\varepsilon}_t^{\mathrm{GH}}(c_3+\varsigma_t^{\mathrm{GH}}),$$
$$s_{\mathrm{GH}t} = s_{\mathrm{k}t} + \boldsymbol{\beta}'\mathbf{s}_{\mathrm{s}t},$$

where $c_0 = R(R+2)/4$, $c_1 = -(R+2)/2$, $c_2 = 1/4$, and $c_3 = -(R+2)$. Then,
i) For any $\boldsymbol{\beta}\in\mathbb{R}^R$ and $\psi>0$,

$$\lim_{\eta\to0^+}\frac{1}{T}\frac{\partial\ln f(\mathbf{Y}_T,\mathbf{E}_T|\boldsymbol{\phi})}{\partial\eta} = -\lim_{\eta\to0^-}\frac{1}{T}\frac{\partial\ln f(\mathbf{Y}_T,\mathbf{E}_T|\boldsymbol{\phi})}{\partial\eta} = \frac{1}{T}\sum_{t=1}^{T}s_{\mathrm{GH}t},\ \textit{and}$$
$$\lim_{\eta\to0^\pm}\frac{1}{T}\frac{\partial\ln f(\mathbf{Y}_T,\mathbf{E}_T|\boldsymbol{\phi})}{\partial\psi} = 0.$$

ii) For any $\boldsymbol{\beta} \in \mathbb{R}^R$ *and* $\eta \in \mathbb{R}$,

$$\lim_{\psi \to 0^+} \frac{1}{T} \frac{\partial \ln f(\mathbf{Y}_T, \mathbf{E}_T | \boldsymbol{\phi})}{\partial \eta} = 0, \text{ and } \lim_{\psi \to 0^+} \frac{2}{T} \frac{\partial \ln f(\mathbf{Y}_T, \mathbf{E}_T | \boldsymbol{\phi})}{\partial \psi} = \frac{1}{T} \sum_{t=1}^{T} s_{\mathrm{GH}t}.$$

Proof. See Mencía and Sentana (2012). □

Lemma 4 *Consider the model (1)-(2) where* $\{\boldsymbol{\varepsilon}_t\}_{t=-\infty}^{\infty}$ *is white noise with identity covariance matrix. Further, assume that all the eigenvalues of* $\mathbf{F}$ *are inside the unit circle. If we observe the double-infinite sequence* $\mathbf{Y}_\infty = \{\mathbf{y}_t\}_{t=-\infty}^{\infty}$*, then the linear projection*

$$\begin{pmatrix} \hat{\boldsymbol{\xi}}_{t-1|\infty} \\ \hat{\boldsymbol{\varepsilon}}_{t|\infty} \end{pmatrix} = \mathcal{P}\left[\left. \begin{pmatrix} \boldsymbol{\xi}_{t-1} \\ \boldsymbol{\varepsilon}_t \end{pmatrix} \right| \mathbf{Y}_\infty \right] = \begin{bmatrix} \boldsymbol{\Psi}(L) \\ \boldsymbol{\Upsilon}(L) \end{bmatrix} \mathbf{y}_t,$$

where $\boldsymbol{\Psi}$ *and* $\boldsymbol{\Upsilon}$ *are absolutely summable two-sided filters in the lag operator* L*, will be given by*

$$\begin{bmatrix} \boldsymbol{\Psi}(z) \\ \boldsymbol{\Upsilon}(z) \end{bmatrix} = \begin{bmatrix} z\mathbf{F}^{-1}(z)\mathbf{M} \\ \mathbf{I}_K \end{bmatrix} \mathbf{D}'(z^{-1}) \left[\mathbf{D}(z)\mathbf{D}'(z^{-1}) \right]^{-1},$$

where

$$\mathbf{F}^{-1}(L) = (\mathbf{I}_M - \mathbf{F}L)^{-1} = \sum_{j=0}^{\infty} \mathbf{F}^j L^j \text{ and } \mathbf{D}(L) = \mathbf{H}\mathbf{F}^{-1}(L)\mathbf{M} = \sum_{j=0}^{\infty} \mathbf{D}_j L^j$$

with $\mathbf{D}_j = \mathbf{H}\mathbf{F}^j\mathbf{M}$ *for all* j*.*

Proof. Given that $\mathbf{y}_t = \mathbf{D}(L)\boldsymbol{\varepsilon}_t$, the joint autocovariance generating function for $(\mathbf{y}_t', \boldsymbol{\varepsilon}_t')'$ is easily seen to be

$$\mathbf{G}(z) = \begin{bmatrix} \mathbf{G}_{yy}(z) & \mathbf{G}_{y\varepsilon}(z) \\ \mathbf{G}_{\varepsilon y}(z) & \mathbf{G}_{\varepsilon\varepsilon}(z) \end{bmatrix} = \begin{bmatrix} \mathbf{D}(z)\mathbf{D}'(z^{-1}) & \mathbf{D}(z) \\ \mathbf{D}'(z^{-1}) & \mathbf{I}_K \end{bmatrix}$$

for any $z \in \mathbb{C}$. The Wiener-Kolmogorov filter for $\boldsymbol{\varepsilon}_t$ is given by

$$\hat{\boldsymbol{\varepsilon}}_{t|\infty} = \mathbf{G}_{\varepsilon y}(L)\mathbf{G}_{yy}^{-1}(L)\mathbf{y}_t = \mathbf{D}'(L^{-1}) \left[\mathbf{D}(L)\mathbf{D}'(L^{-1}) \right]^{-1} \mathbf{y}_t$$

It is then easily checked that for every t, $\hat{\boldsymbol{\varepsilon}}_{t|\infty}$ is well-defined as a mean-square limit under the assumptions of the Lemma. Moreover, because

$$\boldsymbol{\xi}_{t-1} = L\mathbf{F}^{-1}(L)\mathbf{M}\boldsymbol{\varepsilon}_t,$$

the filter for $\boldsymbol{\xi}_{t-1}$ follows from the filter for $\boldsymbol{\varepsilon}_t$, so it is also well-defined. □

Lemma 5 *Consider the model (1)-(2). The score of the asymmetric GH with respect to the parameter* τ *when* $\tau = 0$ *for fixed values of the skewness parameters* $\boldsymbol{\beta}$ *is given by*

$$\begin{aligned} \bar{s}_{\mathrm{GH}T}(\boldsymbol{\theta}, \boldsymbol{\beta}) &= \frac{1}{T} \sum_{t=1}^{T} \left[s_{\mathrm{k}t|T}(\boldsymbol{\theta}) + \boldsymbol{\beta}' \mathbf{s}_{\mathrm{s}t|T}(\boldsymbol{\theta}) \right], \\ s_{\mathrm{k}t|T}(\boldsymbol{\theta}) &= \mathbf{b}'_{\mathrm{k}t|T}(\boldsymbol{\theta}) \mathbf{m}_{\mathrm{k}t|T}(\boldsymbol{\theta}), \\ \mathbf{s}_{\mathrm{s}t|T}(\boldsymbol{\theta}) &= \mathbf{b}'_{\mathrm{s}t|T}(\boldsymbol{\theta}) \mathbf{m}_{\mathrm{s}t|T}(\boldsymbol{\theta}), \end{aligned}$$

where

$$\mathbf{m}_{\mathrm{k}t|T}(\boldsymbol{\theta}) = \begin{pmatrix} 1 \\ \mathbf{m}_{2t|T}(\boldsymbol{\theta}) \\ \mathbf{m}_{4t|T}(\boldsymbol{\theta}) \end{pmatrix}, \quad \mathbf{b}_{\mathrm{k}t|T}(\boldsymbol{\theta}) = \begin{pmatrix} b_{0t|T}(\boldsymbol{\theta}) \\ \mathbf{b}_{2t|T}(\boldsymbol{\theta}) \\ \mathbf{b}_{4t|T}(\boldsymbol{\theta}) \end{pmatrix},$$

$$\mathbf{m}_{\mathrm{s}t|T}(\boldsymbol{\theta}) = \begin{pmatrix} \mathbf{m}_{1t|T}(\boldsymbol{\theta}) \\ \mathbf{m}_{3t|T}(\boldsymbol{\theta}) \end{pmatrix}, \quad \mathbf{b}_{\mathrm{s}t|T}(\boldsymbol{\theta}) = \begin{pmatrix} \mathbf{b}_{1t|T}(\boldsymbol{\theta}) \\ \mathbf{b}_{3t|T}(\boldsymbol{\theta}) \end{pmatrix},$$

$$\begin{aligned}
b_{0t|T}(\boldsymbol{\theta}) &= c_0 + \{c_1 + c_2 tr[\boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})]\} tr[\boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})] + 2c_2 tr\{[\boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})]^2\}, \\
\mathbf{b}_{1t|T}(\boldsymbol{\theta}) &= [c_3 + tr(\boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})]\mathbf{S}'_{RK} + 2\mathbf{S}'_{RK}\boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta}), \\
\mathbf{b}_{2t|T}(\boldsymbol{\theta}) &= \{c_1 + 2c_2 tr[\boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})]\} \left(\mathbf{S}'_{RK} \otimes \mathbf{S}'_{RK}\right) vec(\mathbf{I}_R) + 4c_2 \left(\mathbf{S}'_{RK} \otimes \mathbf{S}'_{RK}\right) vec[\boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})], \\
\mathbf{b}_{3t|T}(\boldsymbol{\theta}) &= \mathbf{S}'_{RK}\boldsymbol{\ell}_R \otimes \mathbf{S}'_{RK}, \\
\mathbf{b}_{4t|T}(\boldsymbol{\theta}) &= c_2 \left(\mathbf{S}'_{RK} \otimes \mathbf{S}'_{RK}\right) \boldsymbol{\ell}_{R^2},
\end{aligned}$$

with $c_0 = R(R+2)/4$, $c_1 = -(R+2)/2$, $c_2 = 1/4$, $c_3 = -(R+2)$ *and* $\boldsymbol{\ell}_H$ *a vector of* H *ones.*

Proof. From Lemma 3, we can obtain the expression for the score with respect to τ for a fixed value of the skewness parameter vector $\boldsymbol{\beta}$, $s_{\mathrm{GH}t} = s_{\mathrm{k}t} + \boldsymbol{\beta}'\mathbf{s}_{\mathrm{s}t}$, which corresponds to the M-step of the EM algorithm. Next, we can apply the E-step to each of the components separately.

As for $s_{\mathrm{k}t}$, we have that $\boldsymbol{\varepsilon}_t | \mathbf{Y}_T, \boldsymbol{\theta} \sim N[\boldsymbol{\varepsilon}_{t|T}(\boldsymbol{\theta}), \boldsymbol{\Omega}_{t|T}(\boldsymbol{\theta})]$ under the null of normality, so that

$$s_{\mathrm{k}t|T}(\boldsymbol{\theta}) = c_0 + c_1 E\left[\varsigma^{\mathrm{GH}}_t \middle| \mathbf{Y}_T, \boldsymbol{\theta}\right] + c_2 E\left[(\varsigma^{\mathrm{GH}}_t)^2 \middle| \mathbf{Y}_T, \boldsymbol{\theta}\right]$$

involves the computation of $E[\varsigma_t | \mathbf{Y}_T, \boldsymbol{\theta}]$ and $E\left[\varsigma^2_t \middle| \mathbf{Y}_T, \boldsymbol{\theta}\right]$. To compute the first expectation, we can write

$$\begin{aligned}
E\left[\varsigma^{\mathrm{GH}}_t \middle| \mathbf{Y}_T\right] &= E\left[\boldsymbol{\varepsilon}^{\mathrm{GH}\prime}_t \boldsymbol{\varepsilon}^{\mathrm{GH}}_t \middle| \mathbf{Y}_T, \boldsymbol{\theta}\right] \\
&= tr\left\{E\left[\boldsymbol{\varepsilon}^{\mathrm{GH}}_t \boldsymbol{\varepsilon}^{\mathrm{GH}\prime}_t | \mathbf{Y}_T, \boldsymbol{\theta}\right]\right\} \\
&= tr[\boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})] + vec(\mathbf{I}_R)' vec[\boldsymbol{\varepsilon}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta}) \boldsymbol{\varepsilon}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})'],
\end{aligned}$$

where the first equality follows from the fact that $tr(A'B) = tr(BA')$, and the second one from Lemma 1.i. As for the second expectation,

$$\begin{aligned}
E\left[(\varsigma^{\mathrm{GH}}_t)^2 | \mathbf{Y}_T, \boldsymbol{\theta}\right] &= E\{\left[\boldsymbol{\varepsilon}^{\mathrm{GH}}_t \odot \boldsymbol{\varepsilon}^{\mathrm{GH}}_t\right]' \mathbf{1}_{R\times R} \left[\boldsymbol{\varepsilon}^{\mathrm{GH}}_t \odot \boldsymbol{\varepsilon}^{\mathrm{GH}}_t\right] | \mathbf{Y}_T, \boldsymbol{\theta}\} \\
&= tr[\mathbf{1}_{R\times R} E\{\left[\boldsymbol{\varepsilon}^{\mathrm{GH}}_t \odot \boldsymbol{\varepsilon}^{\mathrm{GH}}_t\right] \left[\boldsymbol{\varepsilon}^{\mathrm{GH}}_t \odot \boldsymbol{\varepsilon}^{\mathrm{GH}}_t\right]' | \mathbf{Y}_T, \boldsymbol{\theta}\}] \\
&= 2\boldsymbol{\ell}'_{R^2} vec[\boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta}) \odot \boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})] \\
&+ \boldsymbol{\ell}'_{R^2} vec\{vecd[\boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})] vecd[\boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})]'\} \\
&+ 4\boldsymbol{\ell}'_{R^2} vec[\boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta}) \odot \boldsymbol{\varepsilon}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta}) \boldsymbol{\varepsilon}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})'] \\
&+ \boldsymbol{\ell}'_{R^2} vec\{vecd[(\boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})] vecd'[\boldsymbol{\varepsilon}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta}) \boldsymbol{\varepsilon}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})']\} \\
&+ \boldsymbol{\ell}'_{R^2} vec\{vecd[\boldsymbol{\varepsilon}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta}) \boldsymbol{\varepsilon}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})'] vecd'[\boldsymbol{\Omega}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})]\} \\
&+ \boldsymbol{\ell}'_{R^2} vec\{[\boldsymbol{\varepsilon}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta}) \odot \boldsymbol{\varepsilon}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})] [\boldsymbol{\varepsilon}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta}) \odot \boldsymbol{\varepsilon}^{\mathrm{GH}}_{t|T}(\boldsymbol{\theta})]'\},
\end{aligned}$$

where the first equality is a rewriting of $(\varsigma_t^{\text{GH}})^2$, the second one follows from the aforementioned property of the trace, and the third one from Lemma 1.iii. Finally, to obtain the expression for $s_{\text{k},t|T}(\boldsymbol{\theta})$, we have made use of the following identities:

$$\begin{aligned}
\boldsymbol{\ell}'_{R^2} vec[\boldsymbol{\Omega}_{t|T}^{\text{GH}}(\boldsymbol{\theta}) \odot \boldsymbol{\Omega}_{t|T}^{\text{GH}}(\boldsymbol{\theta})] &= vec'[\boldsymbol{\Omega}_{t|T}^{\text{GH}}(\boldsymbol{\theta})] vec[\boldsymbol{\Omega}_{t|T}^{\text{GH}}(\boldsymbol{\theta})] \\
&= tr[\boldsymbol{\Omega}_{t|T}^{\text{GH}}(\boldsymbol{\theta}) \boldsymbol{\Omega}_{t|T}^{\text{GH}}(\boldsymbol{\theta})] = tr\{[\boldsymbol{\Omega}_{t|T}^{\text{GH}}(\boldsymbol{\theta})]^2\} \\
\boldsymbol{\ell}'_{R^2} vec\{vecd[\boldsymbol{\Omega}_{t|T}^{\text{GH}}(\boldsymbol{\theta})] vecd'[\boldsymbol{\Omega}_{t|T}^{\text{GH}}(\boldsymbol{\theta})]\} &= tr^2[\boldsymbol{\Omega}_{t|T}^{\text{GH}}(\boldsymbol{\theta})] \\
\boldsymbol{\ell}'_{R^2} vec\{\boldsymbol{\Omega}_{t|T}^{\text{GH}}(\boldsymbol{\theta}) \odot [\boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta}) \boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta})']\} &= vec'[\boldsymbol{\Omega}_{t|T}^{\text{GH}}(\boldsymbol{\theta})] vec[\boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta}) \boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta})'] \\
\boldsymbol{\ell}'_{R^2} vec\{vecd[\boldsymbol{\Omega}_{t|T}^{\text{GH}}(\boldsymbol{\theta})] vecd'[\boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta}) \boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta})']\} &= tr[\boldsymbol{\Omega}_{t|T}^{\text{GH}}(\boldsymbol{\theta})] vec'(\mathbf{I}_R) vec[\boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta}) \boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta})'],
\end{aligned}$$

together with

$$\begin{aligned}
vec[\boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta}) \boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta})'] &= (\mathbf{S}_{RK} \otimes \mathbf{S}_{RK})\, \mathbf{m}_{2,t|T}(\boldsymbol{\theta}) \\
vec\{[\boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta}) \odot \boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta})][\boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta}) \odot \boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta})]'\} &= (\mathbf{S}_{RK} \otimes \mathbf{S}_{RK})\, \mathbf{m}_{4,t|T}(\boldsymbol{\theta}).
\end{aligned}$$

Similarly, in order to compute

$$\mathbf{s}_{\text{s}t|T}(\boldsymbol{\theta}) = c_3 E\left[\boldsymbol{\varepsilon}_t^{\text{GH}} | \mathbf{Y}_T, \boldsymbol{\theta}\right] + E\left[\boldsymbol{\varepsilon}_t^{\text{GH}} \varsigma_t^{\text{GH}} | \mathbf{Y}_T, \boldsymbol{\theta}\right],$$

we need the expectation of the first component, which is trivially $E\left[\boldsymbol{\varepsilon}_t^{\text{GH}} | \mathbf{Y}_T, \boldsymbol{\theta}\right] = \boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta})$. We also need

$$\begin{aligned}
E\left[\boldsymbol{\varepsilon}_t^{\text{GH}} \varsigma_t | \mathbf{Y}_T, \boldsymbol{\theta}\right] &= E[\boldsymbol{\varepsilon}_t^{\text{GH}} (\boldsymbol{\varepsilon}_t^{\text{GH}} \odot \boldsymbol{\varepsilon}_t^{\text{GH}})' | \mathbf{Y}_T, \boldsymbol{\theta}] \boldsymbol{\ell}_R \\
&= 2\boldsymbol{\Omega}_{t|T}^{\text{GH}}(\boldsymbol{\theta}) \boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta}) + tr[\boldsymbol{\Omega}_{t|T}^{\text{GH}}(\boldsymbol{\theta})] \boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta}) \\
&+ \boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta}) [\boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta}) \odot \boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta})]' \boldsymbol{\ell}_R,
\end{aligned}$$

where we have used the fact that $\varsigma_t^{\text{GH}} = \left[\boldsymbol{\varepsilon}_t^{\text{GH}} \odot \boldsymbol{\varepsilon}_t^{\text{GH}}\right]' \boldsymbol{\ell}_R$ in the first equality, and applied Lemma 1.ii. in the last one. Finally, we obtain the desired result by exploiting the fact that

$$vec\{\boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta}) [\boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta}) \odot \boldsymbol{\varepsilon}_{t|T}^{\text{GH}}(\boldsymbol{\theta})]'\} = (\mathbf{S}_{RK} \otimes \mathbf{S}_{RK})\, \mathbf{m}_{3,t|T}(\boldsymbol{\theta}),$$

after re-arranging terms. □

Lemma 6 *Let*

$$\boldsymbol{\kappa}_i(\boldsymbol{\theta}) = \sum_{j=-\infty}^{\infty} cov[\mathbf{m}_{it}(\boldsymbol{\theta}), \mathbf{m}_{it-j}(\boldsymbol{\theta})],$$

denote the autocovariance generating function of $\mathbf{m}_{it}(\boldsymbol{\theta})$ evaluated at one. Then,
i) The asymptotic variance of $\bar{s}_{\text{k}T}(\boldsymbol{\theta})$ is given by

$$\mathcal{C}_{\text{k}}(\boldsymbol{\theta}) = \mathbf{b}_4'(\boldsymbol{\theta}) \boldsymbol{\kappa}_4(\boldsymbol{\theta}) \mathbf{b}_4'(\boldsymbol{\theta}) - \mathbf{b}_2'(\boldsymbol{\theta}) \boldsymbol{\kappa}_2(\boldsymbol{\theta}) \mathbf{b}_2'(\boldsymbol{\theta}).$$

ii) The asymptotic variance of $\bar{s}_{\text{s}T}(\boldsymbol{\theta})$ is given by

$$\mathcal{C}_{\text{s}|\infty}(\boldsymbol{\theta}) = \mathbf{b}_3'(\boldsymbol{\theta}) \boldsymbol{\kappa}_3(\boldsymbol{\theta}) \mathbf{b}_3'(\boldsymbol{\theta}) - \mathbf{b}_1'(\boldsymbol{\theta}) \boldsymbol{\kappa}_1(\boldsymbol{\theta}) \mathbf{b}_1'(\boldsymbol{\theta}),$$

iii) $\sqrt{T}\bar{s}_{\text{k}T}(\boldsymbol{\theta})$ and $\sqrt{T}\bar{\mathbf{s}}_{\text{s}T}(\boldsymbol{\theta})$ are asymptotically independent.

Proof. Following the same steps as in Lemma 5, but conditioning on $\mathbf{Y}_\infty$ instead of $\mathbf{Y}_T$, we can obtain $\mathbf{s}_{\mathrm{k}t|\infty}(\boldsymbol{\theta}) = E\left[\mathbf{s}_{\mathrm{k}t}(\boldsymbol{\theta})\,|\mathbf{Y}_\infty, \boldsymbol{\theta}\right]$ and $\mathbf{s}_{\mathrm{s}t|\infty}(\boldsymbol{\theta}) = E\left[\mathbf{s}_{\mathrm{s}t}(\boldsymbol{\theta})\,|\mathbf{Y}_\infty, \boldsymbol{\theta}\right]$. Specifically, we can write

$$\begin{bmatrix} s_{\mathrm{k}t|\infty}(\boldsymbol{\theta}) - b_0(\boldsymbol{\theta}) \\ \mathbf{s}_{\mathrm{s}t|\infty}(\boldsymbol{\theta}) \end{bmatrix} = \mathbf{B}'(\boldsymbol{\theta})\mathbf{m}_t(\boldsymbol{\theta}) \text{ where } \mathbf{B}(\boldsymbol{\theta}) = \begin{bmatrix} \mathbf{0} & \mathbf{b}_1(\boldsymbol{\theta}) \\ \mathbf{b}_2(\boldsymbol{\theta}) & \mathbf{0} \\ \mathbf{0} & \mathbf{b}_3(\boldsymbol{\theta}) \\ \mathbf{b}_4(\boldsymbol{\theta}) & \mathbf{0} \end{bmatrix},$$

and $\mathbf{m}_t(\boldsymbol{\theta}) = [\mathbf{m}_{1t}(\boldsymbol{\theta}), \mathbf{m}_{2t|}(\boldsymbol{\theta}), \mathbf{m}_{3t}(\boldsymbol{\theta}), \mathbf{m}_{4t}(\boldsymbol{\theta})]'$, where

$$\begin{aligned}
b_0(\boldsymbol{\theta}) &= c_0 + \{c_1 + tr[\boldsymbol{\Omega}_\infty^{\mathrm{GH}}(\boldsymbol{\theta})]c_2\}tr[\boldsymbol{\Omega}_\infty^{\mathrm{GH}}(\boldsymbol{\theta})] + 2c_2 tr\{[\boldsymbol{\Omega}_\infty^{\mathrm{GH}}(\boldsymbol{\theta})]^2\}, \\
\mathbf{b}_1(\boldsymbol{\theta}) &= \{c_3 + tr[\boldsymbol{\Omega}_\infty^{\mathrm{GH}}(\boldsymbol{\theta})]\}\mathbf{S}'_{RK} + 2\mathbf{S}'_{RK}\boldsymbol{\Omega}_\infty^{\mathrm{GH}}(\boldsymbol{\theta}), \\
\mathbf{b}_2(\boldsymbol{\theta}) &= \{c_1 + 2tr[\boldsymbol{\Omega}_\infty^{\mathrm{GH}}(\boldsymbol{\theta})]c_2\}(\mathbf{S}'_{RK} \otimes \mathbf{S}'_{RK})vec(\mathbf{I}_R) + 4c_2(\mathbf{S}'_{RK} \otimes \mathbf{S}'_{RK})vec[\boldsymbol{\Omega}_\infty^{\mathrm{GH}}(\boldsymbol{\theta})], \\
\mathbf{b}_3(\boldsymbol{\theta}) &= \mathbf{S}'_{RK}\boldsymbol{\ell}_R \otimes \mathbf{S}'_{RK}, \\
\mathbf{b}_4(\boldsymbol{\theta}) &= c_2\left[\mathbf{S}'_{RK} \otimes \mathbf{S}'_{RK}\right]\boldsymbol{\ell}_{R^2},
\end{aligned}$$

with $\boldsymbol{\Omega}_\infty^{\mathrm{GH}}(\boldsymbol{\theta}) = \mathbf{S}_{RK}\boldsymbol{\Omega}_\infty(\boldsymbol{\theta})\mathbf{S}'_{RK}$ and

$$\begin{aligned}
\mathbf{m}_{1t}(\boldsymbol{\theta}) &= \boldsymbol{\varepsilon}_{t|\infty}^{\mathrm{GH}}(\boldsymbol{\theta}), \\
\mathbf{m}_{2t}(\boldsymbol{\theta}) &= vec[\boldsymbol{\varepsilon}_{t|\infty}^{\mathrm{GH}}(\boldsymbol{\theta})\boldsymbol{\varepsilon}_{t|\infty}^{\mathrm{GH}}(\boldsymbol{\theta})'], \\
\mathbf{m}_{3t}(\boldsymbol{\theta}) &= vec\{\boldsymbol{\varepsilon}_{t|\infty}^{\mathrm{GH}}(\boldsymbol{\theta})[\boldsymbol{\varepsilon}_{t|\infty}^{\mathrm{GH}}(\boldsymbol{\theta}) \odot \boldsymbol{\varepsilon}_{t|\infty}^{\mathrm{GH}}(\boldsymbol{\theta})]'\}, \\
\mathbf{m}_{4t}(\boldsymbol{\theta}) &= vec\{[\boldsymbol{\varepsilon}_{t|\infty}^{\mathrm{GH}}(\boldsymbol{\theta}) \odot \boldsymbol{\varepsilon}_{t|\infty}^{\mathrm{GH}}(\boldsymbol{\theta})][\boldsymbol{\varepsilon}_{t|\infty}^{\mathrm{GH}}(\boldsymbol{\theta}) \odot \boldsymbol{\varepsilon}_{t|\infty}^{\mathrm{GH}}(\boldsymbol{\theta})]'\}.
\end{aligned}$$

Next, we can use Lemma 4 to obtain $\boldsymbol{\Gamma}_j = E[\boldsymbol{\varepsilon}_{t|\infty}^{\mathrm{GH}}(\boldsymbol{\theta})\boldsymbol{\varepsilon}_{t-j|\infty}^{\mathrm{GH}}(\boldsymbol{\theta})']$, which corresponds to the j^{th} order autocovariance matrix of the Wiener-Kolmogorov filter for $\boldsymbol{\varepsilon}_t$ based on $\mathbf{Y}_\infty$ for any integer j. Further, we can apply Lemma 2 to obtain:

i) Covariance matrices with the first power:

$$cov\left[\mathbf{m}_{1t}(\boldsymbol{\theta}), \mathbf{m}_{2t-j}(\boldsymbol{\theta})\right] = \mathbf{0}, \tag{A1}$$

$$\begin{aligned}
cov\left[\mathbf{m}_{1t}(\boldsymbol{\theta}), \mathbf{m}_{3t-j}(\boldsymbol{\theta})\right] &= 2\left[\boldsymbol{\ell}_K \otimes vec'(\boldsymbol{\Gamma}_0)\right] \odot \left(\boldsymbol{\Gamma}_j \otimes \boldsymbol{\ell}'_K\right) \\
&\quad + \left[vecd'(\boldsymbol{\Gamma}_0) \otimes \mathbf{1}_{K\times K}\right] \odot \left(\boldsymbol{\ell}'_K \otimes \Gamma_j\right),
\end{aligned} \tag{A2}$$

$$cov\left[\mathbf{m}_{1t}(\boldsymbol{\theta}), \mathbf{m}_{4t-j}(\boldsymbol{\theta})\right] = \mathbf{0}, \tag{A3}$$

ii) Covariance matrices with the second power:

$$\begin{aligned}
cov\left[\mathbf{m}_{2t}(\boldsymbol{\theta}), \mathbf{m}_{2t-j}(\boldsymbol{\theta})\right] &= (\mathbf{1}_{K\times K} \otimes \boldsymbol{\Gamma}_j) \odot (\boldsymbol{\Gamma}_j \otimes \mathbf{1}_{K\times K}) \\
&\quad + \left(\boldsymbol{\ell}_K \otimes \boldsymbol{\Gamma}_j \otimes \boldsymbol{\ell}'_K\right) \odot \left(\boldsymbol{\ell}'_K \otimes \boldsymbol{\Gamma}_j \otimes \boldsymbol{\ell}_K\right),
\end{aligned}$$

$$cov\left[\mathbf{m}_{2t}(\boldsymbol{\theta}), \mathbf{m}_{3t-j}(\boldsymbol{\theta})\right] = \mathbf{0}, \tag{A4}$$

$$\begin{aligned}
cov\left[\mathbf{m}_{2t}(\boldsymbol{\theta}), \mathbf{m}_{4t-j}(\boldsymbol{\theta})\right] &= 4\left[\boldsymbol{\ell}_{K^2} \otimes vec'(\boldsymbol{\Gamma}_0)\right] \odot cov\left[\mathbf{m}_{2t}(\boldsymbol{\theta}), \mathbf{m}_{2t-j}(\boldsymbol{\theta})\right] \\
&\quad + 2\left[\boldsymbol{\ell}_{K^2} \otimes \boldsymbol{\ell}'_K \otimes vecd'(\boldsymbol{\Gamma}_0)\right] \odot \left(\boldsymbol{\ell}_K \otimes \boldsymbol{\Gamma}_j \otimes \boldsymbol{\ell}'_K\right) \odot (\boldsymbol{\Gamma}_j \otimes \mathbf{1}_{K\times K}) \\
&\quad + 2\left[\boldsymbol{\ell}_{K^2} \otimes vecd'(\boldsymbol{\Gamma}_0) \otimes \boldsymbol{\ell}'_K\right] \odot (\mathbf{1}_{K\times K} \otimes \boldsymbol{\Gamma}_j) \odot \left(\boldsymbol{\ell}'_K \otimes \boldsymbol{\Gamma}_j \otimes \boldsymbol{\ell}_K\right),
\end{aligned} \tag{A5}$$

iii) Covariance matrices with the third power:

$$
\begin{aligned}
cov\left[\mathbf{m}_{3t}(\boldsymbol{\theta}),\mathbf{m}_{3t-j}(\boldsymbol{\theta})\right] &= \left[vecd(\boldsymbol{\Gamma}_0)\otimes\boldsymbol{\ell}_K\otimes\boldsymbol{\ell}'_{K^2}\right]\odot\left\{\boldsymbol{\ell}_K\otimes cov\left[\mathbf{m}_{1t}(\boldsymbol{\theta}),\mathbf{m}_{3t-j}(\boldsymbol{\theta})\right]\right\}\\
&+2\left(\mathbf{1}_{K\times K}\otimes\boldsymbol{\Gamma}_j\right)\odot\left[\left(\boldsymbol{\Gamma}_j\odot\boldsymbol{\Gamma}_j\right)\otimes\mathbf{1}_{K\times K}\right]\\
&+2\left[vec(\boldsymbol{\Gamma}_0)\otimes\boldsymbol{\ell}'_{K^2}\right]\odot\left[\boldsymbol{\ell}_{K^2}\otimes vecd'(\boldsymbol{\Gamma}_0)\otimes\boldsymbol{\ell}'_K\right]\odot\left(\boldsymbol{\ell}'_K\otimes\boldsymbol{\Gamma}_j\otimes\boldsymbol{\ell}_K\right)\\
&+4\left[vec(\boldsymbol{\Gamma}_0)\otimes\boldsymbol{\ell}'_{K^2}\right]\odot\left[\boldsymbol{\ell}_{K^2}\otimes vec'(\boldsymbol{\Gamma}_0)\right]\odot\left(\boldsymbol{\Gamma}_j\otimes\mathbf{1}_{K\times K}\right)\\
&+4\left(\boldsymbol{\ell}_K\otimes\boldsymbol{\Gamma}_j\otimes\boldsymbol{\ell}'_K\right)\odot\left(\boldsymbol{\ell}'_K\otimes\boldsymbol{\Gamma}_j\otimes\boldsymbol{\ell}_K\right)\odot\left(\boldsymbol{\Gamma}_j\otimes\mathbf{1}_{K\times K}\right),
\end{aligned}
$$

$$
cov\left[\mathbf{m}_{3t}(\boldsymbol{\theta}),\mathbf{m}_{4t-j}(\boldsymbol{\theta})\right]=\mathbf{0}, \tag{A6}
$$

iv) Covariance matrix of the fourth power:

$$
\begin{aligned}
cov\left[\mathbf{m}_{4t}(\boldsymbol{\theta}),\mathbf{m}_{4t-j}(\boldsymbol{\theta})\right] &= 4cov\left[\mathbf{m}_{2t}(\boldsymbol{\theta}),\mathbf{m}_{2t-j}(\boldsymbol{\theta})\right]\odot cov\left[\mathbf{m}_{2t}(\boldsymbol{\theta}),\mathbf{m}_{2t-j}(\boldsymbol{\theta})\right]\\
&+4\left[vec(\boldsymbol{\Gamma}_0)\otimes\boldsymbol{\ell}'_{K^2}\right]\odot cov\left[\mathbf{m}_{2t}(\boldsymbol{\theta}),\mathbf{m}_{4t-j}(\boldsymbol{\theta})\right]\\
&+2\left[\boldsymbol{\ell}_K\otimes vecd(\boldsymbol{\Gamma}_0)\otimes\boldsymbol{\ell}'_{K^2}\right]\odot\left[\boldsymbol{\ell}_{K^2}\otimes\boldsymbol{\ell}_K\otimes vecd'(\boldsymbol{\Gamma}_0)\right]\odot\left(\boldsymbol{\Gamma}_j\otimes\mathbf{1}_{K\times K}\right)\odot\left(\boldsymbol{\Gamma}_j\otimes\mathbf{1}_{K\times K}\right)\\
&+2\left[\boldsymbol{\ell}_K\otimes vecd(\boldsymbol{\Gamma}_0)\otimes\boldsymbol{\ell}'_{K^2}\right]\odot\left[\boldsymbol{\ell}_{K^2}\otimes vecd'(\boldsymbol{\Gamma}_0\otimes\boldsymbol{\ell}_K)\right]\odot\left(\boldsymbol{\ell}'_K\otimes\boldsymbol{\Gamma}_j\otimes\boldsymbol{\ell}_K\right)\odot\left(\boldsymbol{\ell}'_K\otimes\boldsymbol{\Gamma}_j\otimes\boldsymbol{\ell}_K\right)\\
&+2\left[vecd(\boldsymbol{\Gamma}_0)\otimes\boldsymbol{\ell}_K\otimes\boldsymbol{\ell}'_{K^2}\right]\odot\left[\boldsymbol{\ell}_{K^2}\otimes\boldsymbol{\ell}_K\otimes vecd'(\boldsymbol{\Gamma}_0)\right]\odot\left(\boldsymbol{\ell}_K\otimes\boldsymbol{\Gamma}_j\otimes\boldsymbol{\ell}'_K\right)\odot\left(\boldsymbol{\ell}_K\otimes\boldsymbol{\Gamma}_j\otimes\boldsymbol{\ell}'_K\right)\\
&+2\left[vecd(\boldsymbol{\Gamma}_0)\otimes\boldsymbol{\ell}_K\otimes\boldsymbol{\ell}'_{K^2}\right]\odot\left[\boldsymbol{\ell}_{K^2}\otimes vecd'(\boldsymbol{\Gamma}_0\otimes\boldsymbol{\ell}_K)\right]\odot\left(\mathbf{1}_{K\times K}\otimes\boldsymbol{\Gamma}_j\right)\odot\left(\mathbf{1}_{K\times K}\otimes\boldsymbol{\Gamma}_j\right)\\
&+8\left[\boldsymbol{\ell}_K\otimes vecd(\boldsymbol{\Gamma}_0)\otimes\boldsymbol{\ell}'_{K^2}\right]\odot\left[\boldsymbol{\ell}_{K^2}\otimes vec'(\boldsymbol{\Gamma}_0)\right]\odot\left(\boldsymbol{\ell}'_K\otimes\boldsymbol{\Gamma}_j\otimes\boldsymbol{\ell}_K\right)\odot\left(\boldsymbol{\Gamma}_j\otimes\mathbf{1}_{K\times K}\right)\\
&+8\left[vecd(\boldsymbol{\Gamma}_0)\otimes\boldsymbol{\ell}_K\otimes\boldsymbol{\ell}'_{K^2}\right]\odot\left[\boldsymbol{\ell}_{K^2}\otimes vec'(\boldsymbol{\Gamma}_0)\right]\odot\left(\mathbf{1}_{K\times K}\otimes\boldsymbol{\Gamma}_j\right)\odot\left(\boldsymbol{\ell}_K\otimes\boldsymbol{\Gamma}_j\otimes\boldsymbol{\ell}'_K\right)\\
&+8\left(\boldsymbol{\Gamma}_j\otimes\mathbf{1}_{K\times K}\right)\odot\left(\mathbf{1}_{K\times K}\otimes\boldsymbol{\Gamma}_j\right)\odot\left(\boldsymbol{\ell}'_K\otimes\boldsymbol{\Gamma}_j\otimes\boldsymbol{\ell}_K\right)\odot\left(\boldsymbol{\ell}_K\otimes\boldsymbol{\Gamma}_j\otimes\boldsymbol{\ell}'_K\right).
\end{aligned}
$$

Then, we can show the asymptotic independence of the kurtosis and skewness components by noticing that

$$
\begin{aligned}
cov\left[\mathbf{s}_{\mathrm{s}t|\infty}(\boldsymbol{\theta}),s_{\mathrm{k}t-j|\infty}(\boldsymbol{\theta})\right] &= \mathbf{b}'_1 cov\left[\mathbf{m}_{1t}(\boldsymbol{\theta}),\mathbf{m}_{2t-j}(\boldsymbol{\theta})\right]\mathbf{b}_2\\
&+\mathbf{b}'_1 cov\left[\mathbf{m}_{1t}(\boldsymbol{\theta}),\mathbf{m}_{4t-j}(\boldsymbol{\theta})\right]\mathbf{b}_4\\
&+\mathbf{b}'_3 cov\left[\mathbf{m}_{3t}(\boldsymbol{\theta}),\mathbf{m}_{2t-j}(\boldsymbol{\theta})\right]\mathbf{b}_2\\
&+\mathbf{b}'_3 cov\left[\mathbf{m}_{3t}(\boldsymbol{\theta}),\mathbf{m}_{4t-j}(\boldsymbol{\theta})\right]\mathbf{b}_4\\
&=\mathbf{0},
\end{aligned}
$$

where the last equality follows from (A1), (A3), (A4) and (A6). Moreover, we can simplify even further the relevant expressions by exploiting the cancellation of cross-terms within the variance formulas,

$$
cov\left[\mathbf{m}_{1t}(\boldsymbol{\theta}),\mathbf{s}_{\mathrm{s}t-j}(\boldsymbol{\theta})\right]=\mathbf{0},\ \text{ and }\ cov\left[\mathbf{m}_{2t}(\boldsymbol{\theta}),s_{\mathrm{k}t-j|\infty}(\boldsymbol{\theta})\right]=\mathbf{0}. \tag{A7}
$$

For the sake of brevity, we prove the above equalities for the case when $R=K$; the proof for the case $R<K$ is similar, but more tedious.

To show the first equality in (A7), notice that for any j, we obtain

$$
cov\left[\mathbf{m}_{1t}(\boldsymbol{\theta}),\mathbf{m}_{1t-j}(\boldsymbol{\theta})\right]\mathbf{b}_1=-\left[2\boldsymbol{\Gamma}_j\boldsymbol{\Gamma}_0+tr(\boldsymbol{\Gamma}_0)\boldsymbol{\Gamma}_j\right]
$$

because $\boldsymbol{\Omega}_\infty = \mathbf{I}_K - \boldsymbol{\Gamma}_0$ and $\mathbf{b}_1 = -tr(\boldsymbol{\Gamma}_0)\mathbf{I}_K - \boldsymbol{\Gamma}_0$. The remaining part follows from exploiting the following equalities:

$$\boldsymbol{\Gamma}_j\boldsymbol{\Gamma}_0 = \left\{\left[\boldsymbol{\ell}_K \otimes vec'(\boldsymbol{\Gamma}_0)\right] \odot \left(\boldsymbol{\Gamma}_j \otimes \boldsymbol{\ell}'_K\right)\right\}(\boldsymbol{\ell}_K \otimes \mathbf{I}_K) \tag{A8}$$

and

$$tr(\boldsymbol{\Gamma}_0)\boldsymbol{\Gamma}_j = \left\{\left[vecd'(\boldsymbol{\Gamma}_0) \otimes \mathbf{1}_{K\times K}\right] \odot \left(\boldsymbol{\ell}'_K \otimes \Gamma_j\right)\right\}(\boldsymbol{\ell}_K \otimes \mathbf{I}_K). \tag{A9}$$

For instance, to show (A8), define

$$\mathbf{T}_K = \left[\mathbf{e}_1\mathbf{e}'_1 \quad \ldots \quad \mathbf{e}_K\mathbf{e}'_K\right],$$

with $(\mathbf{e}_1|\ldots|\mathbf{e}_K) = \mathbf{I}_K$, as the unique $K \times K^2$ "diagonalization" matrix that transforms $vec(\mathbf{A})$ into $vecd(\mathbf{A})$ as $vecd(\mathbf{A}) = \mathbf{T}'_K vec(\mathbf{A})$ (see Magnus (1988)). Similarly, let

$$\mathbf{T}_{K^2} = \left[(\mathbf{e}_1\mathbf{e}'_1 \otimes \mathbf{e}_1\mathbf{e}'_1) \quad (\mathbf{e}_1\mathbf{e}'_2 \otimes \mathbf{e}_1\mathbf{e}'_2) \quad \ldots \quad (\mathbf{e}_K\mathbf{e}'_{K-1} \otimes \mathbf{e}_K\mathbf{e}'_{K-1}) \quad (\mathbf{e}_K\mathbf{e}'_K \otimes \mathbf{e}_K\mathbf{e}'_K)\right],$$

which is $K^2 \times K^4$. Some straightforward algebra delivers the following key identities:

$$\begin{aligned}
\mathbf{e}'_i\mathbf{T}_K &= (\mathbf{e}_i \otimes \mathbf{e}_i)', \\
(\mathbf{e}_i \otimes \mathbf{e}_i)'\mathbf{T}_{K^2} &= (\mathbf{e}_i \otimes \mathbf{e}_i \otimes \mathbf{e}_i \otimes \mathbf{e}_i)', \\
\mathbf{T}'_{K^2}(\boldsymbol{\ell}_K \otimes \mathbf{I}_K)\mathbf{e}_i &= (\mathbf{I}_K \otimes \mathbf{e}_i \otimes \mathbf{I}_K \otimes \mathbf{e}_i)vec(\mathbf{I}_K), \\
\mathbf{T}'_{K^2}\boldsymbol{\ell}_{K^2} &= vec(\mathbf{I}_{K^2}),
\end{aligned}$$

for all $i = 1, \ldots, K$. Moreover, $\mathbf{T}_K$ and $\mathbf{T}_{K^2}$ have the important property that

$$(\mathbf{A} \odot \mathbf{B}) = \mathbf{T}_K(\mathbf{A} \otimes \mathbf{B})\mathbf{T}'_{K^2}$$

for any pair of $K \times K^2$ matrices $\mathbf{A}$ and $\mathbf{B}$. As a consequence, we have that for any pair of indices $i_1, i_2 = 1, \ldots, K$,

$$\begin{aligned}
\mathbf{e}'_{i_1}\{[\boldsymbol{\ell}_K \otimes vec'(\boldsymbol{\Gamma}_0)] \odot (\boldsymbol{\Gamma}_j \otimes \boldsymbol{\ell}'_K)\}(\boldsymbol{\ell}_K \otimes \mathbf{I}_K)\mathbf{e}_{i_2} &= \mathbf{e}'_{i_1}\mathbf{T}_K\{[\boldsymbol{\ell}_K \otimes vec'(\boldsymbol{\Gamma}_0)] \otimes (\boldsymbol{\Gamma}_j \otimes \boldsymbol{\ell}'_K)\} \\
&\quad \times \mathbf{T}'_{K^2}(\boldsymbol{\ell}_K \otimes \mathbf{I}_K)\mathbf{e}_{i_2} \\
&= (\mathbf{e}_{i_1} \otimes \mathbf{e}_{i_1})'\{[\boldsymbol{\ell}_K \otimes vec'(\boldsymbol{\Gamma}_0)] \otimes (\boldsymbol{\Gamma}_j \otimes \boldsymbol{\ell}'_K)\} \times (\mathbf{I}_K \otimes \mathbf{e}_{i_2} \otimes \mathbf{I}_K \otimes \mathbf{e}_{i_2})vec(\mathbf{I}_K) \\
&= \{\mathbf{e}'_{i_1}[\boldsymbol{\ell}_K \otimes vec'(\boldsymbol{\Gamma}_0)](\mathbf{I}_K \otimes \mathbf{e}_{i_2}) \otimes \mathbf{e}'_{i_1}(\boldsymbol{\Gamma}_j \otimes \boldsymbol{\ell}'_K)(\mathbf{I}_K \otimes \mathbf{e}_{i_2})\} \times vec(\mathbf{I}_K) \\
&= (\mathbf{e}'_{i_2}\boldsymbol{\Gamma}_0 \otimes \mathbf{e}'_{i_1}\boldsymbol{\Gamma}_j)vec(\mathbf{I}_K) = \mathbf{e}'_{i_1}\boldsymbol{\Gamma}_j\boldsymbol{\Gamma}_0\mathbf{e}_{i_2}.
\end{aligned}$$

But since i_1, i_2 are arbitrary, we can conclude that (A8) holds. Analogous calculations allow us to show (A9). Therefore (A8) and (A9), together with the fact that $\mathbf{b}_3 = \boldsymbol{\ell}_K \otimes \mathbf{I}_K$ and (A2), imply that

$$cov\left[\mathbf{m}_{1t}(\boldsymbol{\theta}), \mathbf{s}_{st-j|\infty}(\boldsymbol{\theta})\right] = cov[\mathbf{m}_{1t}(\boldsymbol{\theta}), \mathbf{m}'_{1t}(\boldsymbol{\theta})]\mathbf{b}'_1 + cov[\mathbf{m}_{1t}(\boldsymbol{\theta}), \mathbf{m}'_{3t}(\boldsymbol{\theta})]\mathbf{b}'_3 = \mathbf{0}.$$

As for the second equality in (A7), again given that $\boldsymbol{\Omega}_\infty = \mathbf{I}_K - \boldsymbol{\Gamma}_0$ and

$$\mathbf{b}_2 = -\frac{1}{2} tr(\boldsymbol{\Gamma}_0) vec(\mathbf{I}_K) - vec(\boldsymbol{\Gamma}_0),$$

we can then use the same tedious but straightforward arguments as before to show that

$$cov\left[\mathbf{m}_{2t}(\boldsymbol{\theta}), \mathbf{m}_{2t-j}(\boldsymbol{\theta})\right] vec(\boldsymbol{\Gamma}_0) = \{[\boldsymbol{\ell}_{K^2} \otimes vec'(\boldsymbol{\Gamma}_0)] \odot cov\left[\mathbf{m}_{2t}(\boldsymbol{\theta}), \mathbf{m}_{2t-j}(\boldsymbol{\theta})\right]\}\boldsymbol{\ell}_{K^2},$$

$$tr(\boldsymbol{\Gamma}_0) cov\left[\mathbf{m}_{2t}(\boldsymbol{\theta}), \mathbf{m}_{2t-j}(\boldsymbol{\theta})\right] vec(\mathbf{I}_K) = \{\left[\boldsymbol{\ell}_{K^2} \otimes \boldsymbol{\ell}'_K \otimes vecd'(\boldsymbol{\Gamma}_0)\right] \odot \left(\boldsymbol{\ell}_K \otimes \boldsymbol{\Gamma}_j \otimes \boldsymbol{\ell}'_K\right)$$
$$\odot (\boldsymbol{\Gamma}_j \otimes \mathbf{1}_{K\times K})\}\boldsymbol{\ell}_{K^2} + \{[\boldsymbol{\ell}_{K^2} \otimes vecd'(\boldsymbol{\Gamma}_0) \otimes \boldsymbol{\ell}'_K] \odot (\mathbf{1}_{K\times K} \otimes \boldsymbol{\Gamma}_j) \odot \boldsymbol{\ell}'_K \otimes \boldsymbol{\Gamma}_j \otimes \boldsymbol{\ell}_K)\}\boldsymbol{\ell}_{K^2},$$

which, together with the fact that $\mathbf{b}_4 = \boldsymbol{\ell}_{K^2}/4$ and (A5), imply that

$$cov\left[\mathbf{m}_{2t}(\boldsymbol{\theta}), s_{\mathrm{k}t-j|\infty}(\boldsymbol{\theta})\right] = cov\left[\mathbf{m}_{2t}(\boldsymbol{\theta}), \mathbf{m}'_{2t}(\boldsymbol{\theta})\right] \mathbf{b}'_2 + cov\left[\mathbf{m}_{2t}(\boldsymbol{\theta}), \mathbf{m}'_{4t}(\boldsymbol{\theta})\right] \mathbf{b}'_4 = \mathbf{0},$$

as desired. This allows us to write

$$\lim_{T\to\infty} V \begin{bmatrix} \sqrt{T}\bar{s}_{\mathrm{k}T}(\boldsymbol{\theta}) \\ \sqrt{T}\bar{\mathbf{s}}_{\mathrm{s}T}(\boldsymbol{\theta}) \end{bmatrix} = \begin{bmatrix} \mathcal{C}_{\mathrm{k}}(\boldsymbol{\theta}) & \mathbf{0} \\ \mathbf{0} & \mathcal{C}_{\mathrm{s}}(\boldsymbol{\theta}) \end{bmatrix}$$

where the expressions for $\mathcal{C}_{\mathrm{k}}(\boldsymbol{\theta})$ and $\mathcal{C}_{\mathrm{s}}(\boldsymbol{\theta})$ can be found in the statement of the Lemma. □

Lemma 7 *Let $\bar{\mathbf{s}}_{\mathrm{MV}T}(\boldsymbol{\theta})$ denote the Gaussian ML score with respect to the conditional mean and variance parameters $\boldsymbol{\theta}$. Then,*

$$\begin{aligned} i)\ \lim_{T\to\infty} cov[\sqrt{T}\bar{\mathbf{s}}_{\mathrm{MV}T}(\boldsymbol{\theta}), \sqrt{T}\bar{s}_{\mathrm{k}T}(\boldsymbol{\theta})|\boldsymbol{\theta}] &= \mathbf{0}, \\ ii)\ \lim_{T\to\infty} cov[\sqrt{T}\bar{\mathbf{s}}_{\mathrm{MV}T}(\boldsymbol{\theta}), \sqrt{T}\bar{\mathbf{s}}_{\mathrm{s}T}(\boldsymbol{\theta})|\boldsymbol{\theta}] &= \mathbf{0}. \end{aligned}$$

Proof. As shown in Mencía and Sentana (2012), the score of the latent model with respect to the mean-variance parameter vector $\boldsymbol{\theta}$ converges to the Gaussian score as we approach the null hypothesis along any of the possible directions through which the GH distribution approaches Gaussianity. This observation combined with the EM principle provides a very convenient way of studying explicitly the score with respect to $\boldsymbol{\theta}$. For ease of exposition assume $\boldsymbol{\xi}_0 = \mathbf{0}$. Then,

$$\begin{aligned} \mathbf{Y}_T &= [\ell_T \otimes \boldsymbol{\pi}(\boldsymbol{\theta})] + [\mathbf{I}_T \otimes \mathbf{H}(\boldsymbol{\theta})]\{\mathbf{I}_{MT} - [\mathbf{C}_T \otimes \mathbf{F}(\boldsymbol{\theta})]\}^{-1}[\mathbf{I}_T \otimes \mathbf{M}(\boldsymbol{\theta})]\mathbf{E}_T \\ &\equiv \boldsymbol{\Pi}_T(\boldsymbol{\theta}) + \mathbf{D}_T(\boldsymbol{\theta})\mathbf{E}_T. \end{aligned}$$

where we have defined

$$\mathbf{C}_T \equiv \begin{bmatrix} \mathbf{0} & \mathbf{I}_{T-1} \\ \mathbf{0} & \mathbf{0} \end{bmatrix},$$
$$\boldsymbol{\Pi}_T(\boldsymbol{\theta}) \equiv \ell_T \otimes \boldsymbol{\pi}(\boldsymbol{\theta}),$$

and

$$\mathbf{D}_T(\boldsymbol{\theta}) \equiv [\mathbf{I}_T \otimes \mathbf{H}(\boldsymbol{\theta})]\{\mathbf{I}_{MT} - [\mathbf{C}_T \otimes \mathbf{F}(\boldsymbol{\theta})]\}^{-1}[\mathbf{I}_T \otimes \mathbf{M}(\boldsymbol{\theta})]. \tag{A10}$$

Under our assumption that no linear combination of $\mathbf{Y}_T$ has zero variance, the matrix $\mathbf{D}_T(\boldsymbol{\theta})$ has full row-rank. Let $\mathbf{D}^*_T(\boldsymbol{\theta})$ be a $(K-N)T \times KT$ matrix of differentiable functions such that

$$\tilde{\mathbf{D}}_T(\boldsymbol{\theta}) = [\mathbf{D}'_T(\boldsymbol{\theta}), (\mathbf{D}^*_T(\boldsymbol{\theta}))']' \tag{A11}$$

is non-singular. Let $\tilde{\mathbf{\Pi}}_T(\boldsymbol{\theta}) = [\mathbf{\Pi}_T'(\boldsymbol{\theta}), \mathbf{0}']'$ and define

$$\tilde{\mathbf{Y}}_T \equiv \tilde{\mathbf{\Pi}}_T(\boldsymbol{\theta}) + \tilde{\mathbf{D}}_T(\boldsymbol{\theta})\mathbf{E}_T.$$

This reasoning delivers the following alternative state space representation under the null,

$$\mathbf{Y}_T = \mathbf{S}_{NT,KT}\tilde{\mathbf{Y}}_T, \quad \text{with} \quad \tilde{\mathbf{Y}}_T \sim N[\tilde{\mathbf{\Pi}}_T(\boldsymbol{\theta}), \tilde{\mathbf{D}}_T'(\boldsymbol{\theta})\tilde{\mathbf{D}}_T(\boldsymbol{\theta})].$$

We now apply the EM principle to the previous representation noting that the measurement equation contains no unknown parameters. The score of the latent model is

$$\begin{aligned}\frac{\partial \ln f_{\tilde{\mathbf{Y}}}(\tilde{\mathbf{Y}}_T|\boldsymbol{\theta})}{\partial \boldsymbol{\theta}} &= \frac{\partial \tilde{\mathbf{\Pi}}_T'(\boldsymbol{\theta})}{\partial \boldsymbol{\theta}}\tilde{\mathbf{D}}_T(\boldsymbol{\theta})\mathbf{E}_T + \frac{1}{2}\frac{\partial vec'[\tilde{\mathbf{D}}_T'(\boldsymbol{\theta})\tilde{\mathbf{D}}_T(\boldsymbol{\theta})]}{\partial \boldsymbol{\theta}}[\tilde{\mathbf{D}}_T'(\boldsymbol{\theta}) \otimes \tilde{\mathbf{D}}_T'(\boldsymbol{\theta})]vec\left(\mathbf{E}_T\mathbf{E}_T' - \mathbf{I}_{KT}\right), \\ &\equiv \mathbf{b}_{\text{MV1},T}(\boldsymbol{\theta})\mathbf{E}_T + \mathbf{b}_{\text{MV2},T}(\boldsymbol{\theta})vec\left(\mathbf{E}_T\mathbf{E}_T' - \mathbf{I}_{KT}\right).\end{aligned}$$

where we have used $\mathbf{E}_T = \tilde{\mathbf{D}}_T^{-1}(\boldsymbol{\theta})[\tilde{\mathbf{Y}}_T - \tilde{\mathbf{\Pi}}_T(\boldsymbol{\theta})]$, with $\mathbf{b}_{\text{MV1},T}(\boldsymbol{\theta})$ and $\mathbf{b}_{\text{MV2},T}(\boldsymbol{\theta})$ defined in the obvious way. Smoothing the score above we obtain

$$\begin{aligned}\bar{\mathbf{s}}_{\text{MV}T}(\boldsymbol{\theta}) &\equiv \frac{1}{T}\frac{\partial \ln f_{\mathbf{Y}}(\mathbf{Y}_T|\boldsymbol{\theta})}{\partial \boldsymbol{\theta}} \\ &= \frac{1}{T}\mathbf{b}_{\text{MV1},T}(\boldsymbol{\theta})E[\mathbf{E}_T|\mathbf{Y}_T, \boldsymbol{\theta}] + \frac{1}{T}\mathbf{b}_{\text{MV2},T}(\boldsymbol{\theta})vec\{E[\mathbf{E}_T\mathbf{E}_T'|\mathbf{Y}_T, \boldsymbol{\theta}] - \mathbf{I}_{KT}\}.\end{aligned}$$

But since

$$E[\mathbf{E}_T\mathbf{E}_T'|\mathbf{Y}_T, \boldsymbol{\theta}] = V[\mathbf{E}_T|\mathbf{Y}_T, \boldsymbol{\theta}] + E[\mathbf{E}_T|\mathbf{Y}_T, \boldsymbol{\theta}]E[\mathbf{E}_T'|\mathbf{Y}_T, \boldsymbol{\theta}]$$

and $V[\mathbf{E}_T|\mathbf{Y}_T, \boldsymbol{\theta}]$ does not depend on $\mathbf{Y}_T$, it is clear that $\bar{\mathbf{s}}_{\text{MV}T}(\boldsymbol{\theta})$ is a linear combination of $\mathbf{E}_{T|T}(\boldsymbol{\theta}) \equiv E[\mathbf{E}_T|\mathbf{Y}_T, \boldsymbol{\theta}]$ and $vec[\mathbf{E}_{T|T}(\boldsymbol{\theta})\mathbf{E}_{T|T}'(\boldsymbol{\theta})]$ (with coefficients possibly varying with T). If we then replace Kalman smoothed variables by their Wiener-Kolmogorov counterparts and the coefficients of the linear combination by their limits as $T \to \infty$, we obtain

$$\begin{aligned}\bar{\mathbf{s}}_{\text{MV}T}(\boldsymbol{\theta}) &= \mathbf{b}_{\text{MV0}}(\boldsymbol{\theta}) + \mathbf{b}_{\text{MV1}}(\boldsymbol{\theta})'\frac{1}{T}\sum_{t=1}^{T}\boldsymbol{\varepsilon}_{t|\infty}(\boldsymbol{\theta}) \\ &\quad + \sum_{\ell=0}^{T-1}\left\{\mathbf{b}_{\text{MV2}}(\boldsymbol{\theta})'\frac{1}{T}\sum_{t=\ell+1}^{T}vec\left[\boldsymbol{\varepsilon}_{t|\infty}(\boldsymbol{\theta})\boldsymbol{\varepsilon}_{t-\ell|\infty}'(\boldsymbol{\theta})\right]\right\}.\end{aligned}$$

The rest of the proof follows from (A7) and from the fact that

$$cov\{vec[\boldsymbol{\varepsilon}_{t|\infty}(\boldsymbol{\theta})\boldsymbol{\varepsilon}_{t-\ell|\infty}'(\boldsymbol{\theta})], \mathbf{s}_{st-j}(\boldsymbol{\theta})\} = \mathbf{0} \text{ and } cov\{vec[\boldsymbol{\varepsilon}_{t|\infty}(\boldsymbol{\theta})\boldsymbol{\varepsilon}_{t-\ell|\infty}'(\boldsymbol{\theta})], s_{kt-j}(\boldsymbol{\theta})\} = \mathbf{0},$$

which can be established by an argument analogous to that of (A7). □

Proposition 1

To simplify the exposition, we focus on the case where $\boldsymbol{\theta}$ is fixed and known, so that and the task is to derive the scores with respect to the shape parameters $\boldsymbol{\varphi}$ only. We further assume

that $\boldsymbol{\xi}_0 = \mathbf{0}$ and $\pi(\boldsymbol{\theta}) = \mathbf{0}$. These assumptions are not essential to the argument and may be removed at the cost of more notation. We can then write

$$\mathbf{Y}_T = \mathbf{D}_T(\boldsymbol{\theta})\mathbf{E}_T,$$

where $\mathbf{D}_T(\boldsymbol{\theta})$ is given in (A10). Note the $(NT \times KT)$ matrix $\mathbf{D}_T(\boldsymbol{\theta})$ does not depend on $\boldsymbol{\varphi}$, although it depends on $\boldsymbol{\theta}$. Notice also that $f_{\mathbf{E}}(\mathbf{E}_T|\varphi)$ is continuous in $\mathbf{E}_T$ and differentiable in $\boldsymbol{\varphi}$ by construction because of the properties of the GH distribution $D(\mathbf{0}, \mathbf{I}_K, \boldsymbol{\varphi})$. Given that we are assuming $N \leq K$, we require the additional assumption that $\mathbf{D}_T(\boldsymbol{\theta})$ has full row rank. As in Lemma 7, we define $\mathbf{D}_T^*(\boldsymbol{\theta})$ such that **(**A11**)** is invertible. Similarly, we define the random vector $\mathbf{Y}_T^* = \mathbf{D}_T^*(\boldsymbol{\theta})\mathbf{E}_T$. Hence,

$$\tilde{\mathbf{Y}}_T \equiv \begin{bmatrix} \mathbf{Y}_T \\ \mathbf{Y}_T^* \end{bmatrix} = \tilde{\mathbf{D}}_T(\boldsymbol{\theta})\mathbf{E}_T$$

will be a KT-dimensional random vector with density $f_{\tilde{\mathbf{Y}}}(\tilde{\mathbf{Y}}_T|\boldsymbol{\varphi})$ with respect to Lebesgue measure on $\mathbb{R}^{KT}$ given by the usual change-of-variable formula,

$$f_{\tilde{\mathbf{Y}}}(\tilde{\mathbf{Y}}_T|\boldsymbol{\varphi}) = \frac{f_{\mathbf{E}}(\tilde{\mathbf{D}}_T^{-1}(\boldsymbol{\theta})\tilde{\mathbf{Y}}_T|\boldsymbol{\varphi})}{|\det[\tilde{\mathbf{D}}_T(\boldsymbol{\theta})]|}.$$

Moreover, the density $f_{\tilde{\mathbf{Y}}}$ is continuous in $\tilde{\mathbf{Y}}_T$ and differentiable in $\boldsymbol{\varphi}$, and the marginal density

$$f_{\mathbf{Y}}(\mathbf{Y}_T|\boldsymbol{\varphi}) = \int_{\mathbb{R}^{(K-N)T}} f_{\tilde{\mathbf{Y}}}(\tilde{\mathbf{Y}}_T|\boldsymbol{\varphi}) d\mathbf{Y}_T^*$$

is continuous in $\mathbf{Y}_T$ and differentiable in $\boldsymbol{\varphi}$ as well. Taking logs, differentiating with respect to $\boldsymbol{\varphi}$ on both sides of the foregoing equation, and exchanging the orders of the differentiation and integration operators on the right-hand side by virtue of theorem 16.8 in Billingsley (1995), we conclude that

$$\frac{\partial \ln f_{\mathbf{Y}}(\mathbf{Y}_T|\boldsymbol{\varphi})}{\partial \boldsymbol{\varphi}} = \int_{\mathbb{R}^{(K-N)T}} \frac{\partial \ln f_{\mathbf{E}}(\tilde{\mathbf{D}}_T^{-1}(\boldsymbol{\theta})\tilde{\mathbf{Y}}_T|\boldsymbol{\varphi})}{\partial \boldsymbol{\varphi}} \frac{f_{\tilde{\mathbf{Y}}}(\tilde{\mathbf{Y}}_T|\boldsymbol{\varphi})}{f_{\mathbf{Y}}(\mathbf{Y}_T|\boldsymbol{\varphi})} d\mathbf{Y}_T^*, \tag{A12}$$

for all $\mathbf{Y}_T$ and $\boldsymbol{\varphi}$ for which $f_{\mathbf{Y}}(\mathbf{Y}_T|\boldsymbol{\varphi}) > 0$ (and this holds for almost all $\boldsymbol{\varphi}$). The function $f_{\mathbf{Y}^*|\mathbf{Y}}(\mathbf{Y}_T^*|\mathbf{Y}_T, \boldsymbol{\varphi}) \equiv f_{\tilde{\mathbf{Y}}}(\tilde{\mathbf{Y}}_T|\boldsymbol{\varphi})/f_{\mathbf{Y}}(\mathbf{Y}_T|\boldsymbol{\varphi})$ is the conditional density of $\mathbf{Y}_T^*$ given $\mathbf{Y}_T$, which is a continuous density with respect to Lebesgue measure on $\mathbb{R}^{(K-N)T}$. In that precise sense, we write

$$E\left[\left.\frac{\partial \ln f_{\mathbf{E}}(\mathbf{E}_T|\boldsymbol{\varphi})}{\partial \boldsymbol{\varphi}}\right| \mathbf{Y}_T, \boldsymbol{\varphi}\right] = \int_{\mathbb{R}^{(K-N)T}} \frac{\partial \ln f_{\mathbf{E}}(\tilde{\mathbf{D}}_T^{-1}(\boldsymbol{\theta})\tilde{\mathbf{Y}}_T|\boldsymbol{\varphi})}{\partial \boldsymbol{\varphi}} \frac{f_{\tilde{\mathbf{Y}}}(\tilde{\mathbf{Y}}_T|\boldsymbol{\varphi})}{f_{\mathbf{Y}}(\mathbf{Y}_T|\boldsymbol{\varphi})} d\mathbf{Y}_T^*.$$

Importantly, the value of the integral in (A12) is independent of the choice of $\mathbf{D}_T^*$. To see this, multiply both sides by $f_{\mathbf{Y}}(\mathbf{Y}_T|\boldsymbol{\varphi})$ and integrate with respect to $\mathbf{Y}_T$,

$$\begin{aligned} &\int_{\mathbb{R}^{NT}} \int_{\mathbb{R}^{(K-N)T}} \frac{\partial \ln f_{\mathbf{E}}(\tilde{\mathbf{D}}_T^{-1}(\boldsymbol{\theta})\tilde{\mathbf{Y}}_T|\boldsymbol{\varphi})}{\partial \boldsymbol{\varphi}} f_{\tilde{\mathbf{Y}}}(\tilde{\mathbf{Y}}_T|\boldsymbol{\varphi}) d\mathbf{Y}_T^* d\mathbf{Y}_T \\ &= \int_{\mathbb{R}^{KT}} \frac{\partial \ln f_{\mathbf{E}}(\tilde{\mathbf{D}}_T^{-1}(\boldsymbol{\theta})\tilde{\mathbf{Y}}_T|\boldsymbol{\varphi})}{\partial \boldsymbol{\varphi}} f_{\tilde{\mathbf{Y}}}(\tilde{\mathbf{Y}}_T|\boldsymbol{\varphi}) d\tilde{\mathbf{Y}}_T = E\left[\left.\frac{\partial \ln f_{\mathbf{E}}(\mathbf{E}_T|\boldsymbol{\varphi})}{\partial \boldsymbol{\varphi}}\right| \boldsymbol{\varphi}\right] \end{aligned}$$

by Fubini's theorem. For all possible choices of $\mathbf{D}_T^*(\boldsymbol{\theta})$ we obtain a version of

$$E\left[\frac{\partial \ln f_{\mathbf{E}}(\mathbf{E}_T|\boldsymbol{\varphi})}{\partial \boldsymbol{\varphi}}\middle|\mathbf{Y}_T,\boldsymbol{\varphi}\right],$$

whose uniqueess follows from the a.s. equality of conditional expectations. Therefore, equation (5) holds. □

Proposition 2

It follows from Lemma 5 when $\boldsymbol{\beta}=\mathbf{0}$. □

Proposition 3

It follows from Lemma 6.i and 7.i. □

Proposition 4

It follows from Propositions 2 and 3. □

Proposition 5

It is a rewriting of Lemma 5. □

Proposition 6

It follows from Lemma 6.i, , 6.ii, 6.iii and 7.ii. □

Proposition 7

It follows by combining the arguments in the proof of Proposition 5 in Mencía and Sentana (2012) with the results in Propositions 5 and 6. □

B Asymptotic equivalence of smoothed scores sample moments

Consider the model (1)-(2) with $\boldsymbol{\varepsilon}_t \sim N(\mathbf{0},\mathbf{I}_K)$ and, to save notation, assume (i) $\boldsymbol{\pi}=\mathbf{0}$. To facilitate exposition we further assume that (ii) $\det(\mathbf{I}_M-\mathbf{F}z)=0$ implies $|z|>1$. This condition can be removed at the cost of considerably complicating the analysis.

Under these assumptions, the MA(∞) representation of $\{\mathbf{y}_t\}$ is

$$\mathbf{y}_t=\sum_{s=-\infty}^{\infty}\mathbf{D}(s)\boldsymbol{\varepsilon}_{t-s} \text{ for all } t,$$

where $\mathbf{D}(s)=\mathbf{H}\mathbf{F}^s\mathbf{M}$ for all $s\geq 0$, and $\mathbf{D}(s)=\mathbf{0}$ whenever $s<0$.

Let $\mathcal{F}_T=\sigma(\{\mathbf{y}_t\}_{|t|\leq T})$ denote the σ-field generated by $\{\mathbf{y}_t\}_{|t|\leq T}$. Also, let $\mathcal{F}_\infty=\sigma(\cup_{T=0}^{\infty}\mathcal{F}_T)$. It is well known that the assumption of Gaussianity implies existence of sequences of $K\times N$

matrices $\{\mathbf{A}_{t|T}(\tau)\}$ for all t and T, and $\{\mathbf{A}(\tau)\}$ with $\mathbf{A}_{t|T}(\tau) = \mathbf{0}$ whenever $|t| > \tau$, such that

$$\boldsymbol{\varepsilon}_{t|T} = E(\boldsymbol{\varepsilon}_t|\mathcal{F}_T) = \sum_{\tau=-T}^{T} \mathbf{A}_{t|T}(\tau)\mathbf{y}_{t-\tau}, \text{ for all } t \text{ and } T,$$

$$\boldsymbol{\varepsilon}_{t|\infty} = E(\boldsymbol{\varepsilon}_t|\mathcal{F}_\infty) = \sum_{\tau=-\infty}^{\infty} \mathbf{A}(\tau)\mathbf{y}_{t-\tau}, \text{ for all } t.$$

For any real matrix $\mathbf{A}$, let $\|\mathbf{A}\| = \sqrt{tr(\mathbf{A}'\mathbf{A})}$ be its Frobenius norm.

The purpose of this appendix is to show that:

Proposition 8 *As* $T^* \equiv 2T + 1 \to \infty$,

$$\frac{1}{\sqrt{T^*}} \sum_{|t|\leq T} (\boldsymbol{\varepsilon}_{t|\infty} - \boldsymbol{\varepsilon}_{t|T}) = o_{\mathrm{P}}(1).$$

In the proof of Proposition 8, we will make use of the following:

Lemma 8 *The following three properties hold:*
i) (L_2*-optimality) Any* $\mathcal{F}_T$*-measurable function* $\tilde{\boldsymbol{\varepsilon}}_T$ *satisfies*

$$E\left(\left\|\sum_{|t|\leq T} (\boldsymbol{\varepsilon}_{t|\infty} - \boldsymbol{\varepsilon}_{t|T})\right\|^2\right) \leq E\left(\left\|\sum_{|t|\leq T} \boldsymbol{\varepsilon}_{t|\infty} - \tilde{\boldsymbol{\varepsilon}}_T\right\|^2\right) \text{ for all } T.$$

ii) (Geometric decay of $\mathbf{A}$*) For some* $\rho_\alpha \in (0,1)$, $C_\alpha > 0$ *and all* τ, $\|\mathbf{A}(\tau)\| \leq C_\alpha \rho_\alpha^{|\tau|}$. *Hence,*

$$\sum_{\tau=-\infty}^{\infty} \|\mathbf{A}(\tau)\| < \infty.$$

iii)(Geometric decay of $\mathbf{D}$*) For some* $\rho_\delta \in (0,1)$, $C_\delta > 0$ *and all* s, $\|\mathbf{D}(s)\| \leq C_\delta \rho_\delta^{|s|}$. *Hence,*

$$\sum_{s=-\infty}^{\infty} \|\mathbf{D}(s)\| < \infty.$$

Proof of Lemma 8.

Property (i) is a consequence of the fact that

$$\boldsymbol{\varepsilon}_{t|T} = E(\boldsymbol{\varepsilon}_{t|\infty}|\mathcal{F}_T) \text{ for all } t \text{ and } T$$

by virtue of the law of iterated expectations, and the standard result that an expectation conditional on $\mathcal{F}_T$ minimizes the L_2-distance to the set of $\mathcal{F}_T$-measurable functions.

In turn, Property (ii) follows from the fact that $\boldsymbol{\varepsilon}_{t|\infty}$ is a VARMA process. Hence,

$$\sum_{|\tau|>T} \|\mathbf{A}(\tau)\| \leq 2C_\alpha \sum_{\tau=T+1}^{\infty} \rho_\alpha^\tau = \frac{2C_\alpha}{1-\rho_\alpha} \rho_\alpha^{T+1} \to 0 \text{ as } T \to \infty,$$

implying $\sum_{\tau=-\infty}^{\infty} \|\mathbf{A}(\tau)\| < \infty$.

To establish property (iii), note that $\|\mathbf{D}_s\| \leq \|\mathbf{H}\|\|\mathbf{F}\|^s\|\mathbf{M}\| \leq \sqrt{M}\|\mathbf{H}\|\|\mathbf{M}\||\lambda_{\mathbf{F}}|^s$, where we have denoted by $\lambda_{\mathbf{F}}$ the largest eigenvalue of $\mathbf{F}$. By assumption, $|\lambda_{\mathbf{F}}| < 1$, so $C_\delta = \sqrt{M}\|\mathbf{H}\|\|\mathbf{M}\|$ and $\rho_\delta = |\lambda_{\mathbf{F}}|$. Finally,

$$\sum_{|s|>S} \|\mathbf{D}(s)\| \leq 2C_\delta \sum_{s=S+1}^{\infty} \rho_\delta^s = \frac{2C_\delta}{1-\rho_\delta}\rho_\delta^{S+1} \to 0 \text{ as } S \to \infty,$$

implying $\sum_{s=-\infty}^{\infty} \|\mathbf{D}(s)\| < \infty$. □

Proof of Proposition 8.

Fix some $\epsilon > 0$ and $k = 1, \ldots, K$ and define the event

$$\mathcal{E}_{k,T} \equiv \left\{ \left| \sum_{|t|\leq T} (\varepsilon_{k,t|\infty} - \varepsilon_{k,t|T}) \right| > \sqrt{T^*}\epsilon \right\}$$

By Chebyshev-Bienaymé's inequality,

$$\Pr(\mathcal{E}_{k,T}) \leq \frac{1}{T^*\epsilon^2} V\left[\sum_{|t|\leq T} (\varepsilon_{k,t|\infty} - \varepsilon_{k,t|T}) \right] \leq \frac{1}{T^*\epsilon^2} E\left(\left\| \sum_{|t|\leq T} (\boldsymbol{\varepsilon}_{t|\infty} - \boldsymbol{\varepsilon}_{t|T}) \right\|^2 \right).$$

Further, Lemma 8.i implies that for any $\mathcal{F}_T$-measurable function $\tilde{\boldsymbol{\varepsilon}}_T$,

$$\Pr(\mathcal{E}_{k,T}) \leq \frac{1}{T^*\epsilon^2} E\left(\left\| \sum_{|t|\leq T} \boldsymbol{\varepsilon}_{t|\infty} - \tilde{\boldsymbol{\varepsilon}}_T \right\|^2 \right).$$

Therefore, the proof will be completed if we establish that, for some suitable choice of $\tilde{\boldsymbol{\varepsilon}}_T$,

$$E\left(\left\| \sum_{|t|\leq T} \boldsymbol{\varepsilon}_{t|\infty} - \tilde{\boldsymbol{\varepsilon}}_T \right\|^2 \right) = o(T).$$

To do so, consider the linear $\mathcal{F}_T$-measurable variable

$$\tilde{\boldsymbol{\varepsilon}}_T = \sum_{|t|\leq T} \sum_{|\tau|\leq T} \mathbf{A}(\tau)\mathbf{y}_{t-\tau}.$$

We have

$$\Delta_T = \sum_{|t|\leq T} \boldsymbol{\varepsilon}_{t|\infty} - \tilde{\boldsymbol{\varepsilon}}_T = \sum_{|t|\leq T} \sum_{|\tau|>T} \mathbf{A}(\tau)\mathbf{y}_{t-\tau} = \sum_{r=-\infty}^{\infty} \boldsymbol{\Phi}_T(r)\mathbf{y}_r,$$

where $\boldsymbol{\Phi}_T(0) = \mathbf{0}$,

$$\boldsymbol{\Phi}_T(r) = \sum_{j=\max\{1,r-2T\}}^{r} \mathbf{A}[-(T+j)], \text{ for } r > 0,$$
$$\boldsymbol{\Phi}_T(r) = \sum_{j=\max\{1,r-2T\}}^{r} \mathbf{A}(T+j), \text{ for } r < 0,$$

implying that

$$\|\boldsymbol{\Phi}_T(r)\| \leq C_\alpha \rho_\alpha^{T+1}/(1-\rho_\alpha), \text{ for } |r| \leq 2T+1, \text{ and}$$
$$\|\boldsymbol{\Phi}_T(r)\| \leq C_\alpha \rho_\alpha^{r-T}/(1-\rho_\alpha), \text{ for } |r| > 2T+1,$$

whence it follows immediately that $\sum_{r=-\infty}^{\infty} \|\mathbf{\Phi}_T(r)\| < C_\phi T \rho_\alpha^T$ for some constant $C_\phi > 0$. Finally,

$$\begin{aligned}\sqrt{E(\|\Delta_T\|^2)} &= \sqrt{\sum_{r=-\infty}^{\infty} \left\| \sum_{s=-\infty}^{\infty} \mathbf{\Phi}_T(r)\mathbf{\Psi}(s-r) \right\|^2} \\ &\leq \sum_{r=-\infty}^{\infty} \sum_{s=-\infty}^{\infty} \|\mathbf{\Phi}_T(r)\| \|\mathbf{D}(s-r)\| \\ &\leq \left(\sum_{r=-\infty}^{\infty} \|\mathbf{\Phi}_T(r)\| \right) \left(\sum_{s=-\infty}^{\infty} \|\mathbf{D}(s)\| \right) < \infty,\end{aligned}$$

where the last inequality follows from Lemma 8.ii and 8.iii. As a consequence of this, $E(\|\Delta_T\|^2) = o(T)$. □

C An algorithm for computing the asymptotic variance

Consider a VARMA process with scalar VAR part for the K_x-dimensional process $\mathbf{x}_t$,

$$\phi(L)\mathbf{x}_t = \mathbf{\Theta}(L)\mathbf{u}_t$$

where $\phi(z) = 1 - \phi_1 z - \cdots - \phi_p z^p$ and $\mathbf{\Theta}(z) = \mathbf{\Theta}_0 + \mathbf{\Theta}_1 z + \ldots \mathbf{\Theta}_q z^q$. The error process $\mathbf{u}_t$ is assumed to be K-dimensional white noise, i.e. $E(\mathbf{u}_t) = \mathbf{0}$, $E(\mathbf{u}_t\mathbf{u}_t') = \mathbf{\Sigma}$, $E(\mathbf{u}_t\mathbf{u}_{t-j}') = \mathbf{0}$ for $j \neq 0$. Next, write the VARMA process in companion VAR(1) form as

$$\mathbf{X}_t = \mathbf{A}\mathbf{X}_{t-1} + \mathbf{Q}\mathbf{u}_t,$$

where $\mathbf{X}_t = (\mathbf{x}_t, ..., \mathbf{x}_{t-p+1}, \mathbf{u}_t, ..., \mathbf{u}_{t-p+1})'$,

$$\mathbf{A} = \begin{pmatrix} \bar{\mathbf{\Phi}} \otimes \mathbf{I}_{K_x} & \mathbf{e}_1 \otimes \bar{\mathbf{\Theta}} \\ \mathbf{0} & \mathbf{J}_q \otimes \mathbf{I}_K \end{pmatrix}, \quad \mathbf{Q} = \begin{pmatrix} \mathbf{\Theta}_0 \\ \mathbf{0} \\ \mathbf{I}_K \\ \mathbf{0} \end{pmatrix},$$

with $\mathbf{e}_1$ being the first vector of the canonical basis in $\mathbb{R}^p$,

$$\bar{\mathbf{\Phi}} = \begin{pmatrix} \phi_1 & \cdots & \phi_{p-1} & \phi_p \\ 1 & 0 & \cdots & 0 \\ & \ddots & & \vdots \\ 0 & & 1 & 0 \end{pmatrix}, \quad \bar{\mathbf{\Theta}} = \begin{pmatrix} \mathbf{\Theta}_1 & \cdots & \mathbf{\Theta}_q \end{pmatrix}, \text{ and } \mathbf{J}_q = \begin{pmatrix} \mathbf{0} & \mathbf{0} \\ \mathbf{I}_{q-1} & \mathbf{0} \end{pmatrix}.$$

Suppose we can find an invertible matrix $\mathbf{C}$ and a block diagonal matrix $\mathbf{\Lambda}$ (with Jordan blocks) such that $\mathbf{A} = \mathbf{C}\mathbf{\Lambda}\mathbf{C}^{-1}$. Then, we can transform the original system by defining $\mathbf{Z}_t = \mathbf{C}^{-1}\mathbf{X}_t$, a possibly complex-valued stochastic process that satisfies

$$\mathbf{Z}_t = \mathbf{\Lambda}\mathbf{Z}_{t-1} + \boldsymbol{\eta}_t,$$

with $\boldsymbol{\eta}_t = \mathbf{C}^{-1}\mathbf{Q}\mathbf{u}_t$ being white-noise (and possibly complex-valued). Then, it can be shown that a computationally convenient decomposition of $\mathbf{A}$ is given by

$$\mathbf{A} = \mathbf{C}\mathbf{\Lambda}\mathbf{C}^{-1}, \tag{C13}$$

where

$$\mathbf{C} = \begin{pmatrix} \bar{\mathbf{C}} \otimes \mathbf{I}_{K_x} & -(\bar{\mathbf{\Phi}}^{-q} \otimes \mathbf{I}_{K_x})\mathbf{\Theta}^* \\ \mathbf{0} & \mathbf{I}_{Kq} \end{pmatrix}, \mathbf{\Lambda} = \begin{pmatrix} \bar{\mathbf{\Lambda}} \otimes \mathbf{I}_{K_x} & \mathbf{0} \\ \mathbf{0} & \mathbf{J}_q \otimes \mathbf{I}_K \end{pmatrix}$$

and

$$\mathbf{C}^{-1} = \begin{pmatrix} \bar{\mathbf{C}}^{-1} \otimes \mathbf{I}_{K_x} & (\bar{\mathbf{C}}^{-1}\bar{\mathbf{\Phi}}^{-q} \otimes \mathbf{I}_{K_x})\mathbf{\Theta}^* \\ \mathbf{0} & \mathbf{I}_{Kq} \end{pmatrix},$$

with $\bar{\mathbf{\Phi}} = \bar{\mathbf{C}}\bar{\mathbf{\Lambda}}\bar{\mathbf{C}}^{-1}$ providing the Jordan decomposition of $\bar{\mathbf{\Phi}}$, and

$$\mathbf{\Theta}^* = \sum_{h=1}^{q} (\bar{\mathbf{\Phi}}^{q-h}\mathbf{e}_1 \otimes \bar{\mathbf{\Theta}})(\mathbf{J}_q^{h-1} \otimes \mathbf{I}_K).$$

Notice that the decomposition outlined above is convenient to handle large systems because it reduces substantially the size of the matrices for which the Jordan decomposition needs to be performed.

We can also show that the autocovariance function of the Wiener-Kolmogorov filter derived in Lemma 4 is the autocovariance function of the stable solution to the difference equation embodied in its VARMA representation. For that reason, we decompose $\mathbf{A}$ as in (C13), with the absolute values of the eigenvalues in decreasing order. But since we have assumed no unit roots, we will have that $K_{\mathrm{S}} = K_x p + Kq - K_{\mathrm{U}}$, where K_{U} is the number of roots outside the unit circle and K_{S} the number of roots inside the unit circle.

Let $\mathbf{R} = \mathbf{C}\mathbf{Q}\overline{\mathbf{Q}}'\overline{\mathbf{C}}'$ denote the variance-covariance matrix of $\boldsymbol{\eta}_t$. We can partition the system into its unstable and stable parts as follows:

$$\mathbf{Z}_t = \begin{pmatrix} \mathbf{Z}_{\mathrm{U}t} \\ \mathbf{Z}_{\mathrm{S}t} \end{pmatrix}, \ \boldsymbol{\eta}_t = \begin{pmatrix} \boldsymbol{\eta}_{\mathrm{U}t} \\ \boldsymbol{\eta}_{\mathrm{S}t} \end{pmatrix}, \ \mathbf{\Lambda} = \begin{pmatrix} \mathbf{\Lambda}_{\mathrm{UU}} & \mathbf{0} \\ \mathbf{0} & \mathbf{\Lambda}_{\mathrm{SS}} \end{pmatrix}, \text{ and } \mathbf{R} = \begin{pmatrix} \mathbf{R}_{\mathrm{UU}} & \mathbf{R}_{\mathrm{US}} \\ \mathbf{R}_{\mathrm{SU}} & \mathbf{R}_{\mathrm{SS}} \end{pmatrix}.$$

Next, if we write

$$\mathbf{Z}_{\mathrm{U}t} = \mathbf{\Lambda}_{\mathrm{UU}}^{-1}(\mathbf{Z}_{\mathrm{U}t+1} - \boldsymbol{\eta}_{\mathrm{U}t+1}) \quad \text{and} \quad \mathbf{Z}_{\mathrm{S}t} = \mathbf{\Lambda}_{\mathrm{SS}}\mathbf{Z}_{\mathrm{S}t-1} + \boldsymbol{\eta}_{\mathrm{S}t},$$

and partition

$$\mathbf{\Gamma}_{\mathbf{Z}}(j) = \begin{bmatrix} \mathbf{\Gamma}_{\mathrm{UU}}(j) & \mathbf{\Gamma}_{\mathrm{US}}(j) \\ \mathbf{\Gamma}_{\mathrm{SU}}(j) & \mathbf{\Gamma}_{\mathrm{SS}}(j) \end{bmatrix} = \begin{bmatrix} E(\mathbf{Z}_{\mathrm{U}t}\bar{\mathbf{Z}}'_{\mathrm{U}t-j}) & E(\mathbf{Z}_{\mathrm{U}t}\bar{\mathbf{Z}}'_{\mathrm{S}t-j}) \\ E(\mathbf{Z}_{\mathrm{S}t}\bar{\mathbf{Z}}'_{\mathrm{U}t-j}) & E(\mathbf{Z}_{\mathrm{S}t}\bar{\mathbf{Z}}'_{\mathrm{S}t-j}) \end{bmatrix},$$

we can show that the autocovariance function of $\mathbf{Z}_t$ can be computed from

$$\begin{aligned}
vec[\mathbf{\Gamma}_{\mathrm{UU}}(0)] &= [\mathbf{I}_{K_{\mathrm{U}}^2} - (\mathbf{\Lambda}_{\mathrm{UU}}^{-1} \otimes \mathbf{\Lambda}_{\mathrm{UU}}^{-1})]^{-1} vec[\mathbf{\Lambda}_{\mathrm{UU}}^{-1}\mathbf{R}_{\mathrm{UU}}(\overline{\mathbf{\Lambda}}_{\mathrm{UU}}^{-1})'], \\
\mathbf{\Gamma}_{\mathrm{UU}}(j) &= \mathbf{\Gamma}_{\mathrm{UU}}(0)(\overline{\mathbf{\Lambda}}_{\mathrm{UU}}^{-j})', \quad \text{for } j > 0 \\
\mathbf{\Gamma}_{\mathrm{UU}}(j) &= \overline{\mathbf{\Gamma}}'_{\mathrm{UU}}(-j), \quad \text{for } j < 0 \\
vec[\mathbf{\Gamma}_{\mathrm{SS}}(0)] &= [\mathbf{I}_{K_{\mathrm{S}}^2} - (\mathbf{\Lambda}_{\mathrm{SS}} \otimes \mathbf{\Lambda}_{\mathrm{SS}})]^{-1} vec\left(\mathbf{R}_{\mathrm{SS}}\right), \\
\mathbf{\Gamma}_{\mathrm{SS}}(j) &= \mathbf{\Lambda}_{\mathrm{SS}}^{j}\mathbf{\Gamma}_{\mathrm{SS}}(0), \quad \text{for } j > 0 \\
\mathbf{\Gamma}_{\mathrm{SS}}(j) &= \overline{\mathbf{\Gamma}}'_{\mathrm{SS}}(-j), \quad \text{for } j < 0 \\
\mathbf{\Gamma}_{\mathrm{SU}}(j) &= -\sum_{h=1}^{j} (\overline{\mathbf{\Lambda}}_{\mathrm{SS}}^{j-h})'\mathbf{R}_{\mathrm{SU}}(\overline{\mathbf{\Lambda}}_{\mathrm{UU}}^{-h})', \quad \text{for } j > 0 \\
\mathbf{\Gamma}_{\mathrm{SU}}(j) &= \mathbf{0}, \quad \text{for } j \leq 0, \quad \text{and} \\
\mathbf{\Gamma}_{\mathrm{US}}(j) &= \overline{\mathbf{\Gamma}}'_{\mathrm{SU}}(-j).
\end{aligned}$$

Finally, we can recover the autocovariance function of $\mathbf{X}_t$ from

$$\boldsymbol{\Gamma}_{\mathbf{X}}(j) = E[\mathbf{X}_t \mathbf{X}'_{t-j}] = E[(\mathbf{C}\mathbf{Z}_t)(\overline{\mathbf{C}\mathbf{Z}}_{t-j})'] = \mathbf{C}\Gamma_{\mathbf{Z}}(j)\overline{\mathbf{C}}'.$$

Obviously, the autocovariance function of $\mathbf{x}_t$ is the first block of $\boldsymbol{\Gamma}_{\mathbf{X}}$.

D A Gibbs sampler algorithm for the common trend model with asymmetric Student t innovations

In this section, we develop a Gibbs sampler for the model we use in the empirical application in section 7.3, namely

$$\begin{aligned}
\mathbf{y}_t &= \mathbf{H}\boldsymbol{\xi}_t, \\
\boldsymbol{\xi}_t &= \mathbf{c}(\boldsymbol{\theta}) + \mathbf{F}(\boldsymbol{\theta})\boldsymbol{\xi}_{t-1} + \mathbf{M}(\boldsymbol{\theta})\boldsymbol{\varepsilon}_t, \\
\boldsymbol{\varepsilon}_t &= \boldsymbol{\alpha}(\boldsymbol{\varphi}) + \zeta_t^{-1}\boldsymbol{\Upsilon}(\boldsymbol{\varphi})\boldsymbol{\beta} + \zeta_t^{-1/2}\boldsymbol{\Upsilon}^{1/2}(\boldsymbol{\varphi})\mathbf{z}_t, \\
\zeta_t|\boldsymbol{\theta}, \boldsymbol{\varphi} &\sim iid\ \Gamma(\nu/2, 1/2), \\
\mathbf{z}_t|\boldsymbol{\theta}, \boldsymbol{\varphi} &\sim iid\ N(\mathbf{0}, \mathbf{I}_K),
\end{aligned}$$

where $\boldsymbol{\theta}$ are mean-variance parameters and $\boldsymbol{\varphi} = (\nu, \boldsymbol{\beta}')'$ are shape parameters describing the asymmetric Student t distribution (a member of the *GH* family of distributions). More specifically, $\boldsymbol{\theta} = (\mu, \delta, \rho_x, \rho_{\epsilon_E}, \rho_{\epsilon_I}, \sigma_x^2, \sigma_{v_E}^2, \sigma_{v_I}^2)'$,

$$\mathbf{y}_t = \begin{pmatrix} y_{Et} \\ y_{It} \end{pmatrix}, \boldsymbol{\xi}_t = \begin{pmatrix} x_t \\ x_{t-1} \\ \epsilon_{Et} \\ \epsilon_{It} \end{pmatrix}, \boldsymbol{\varepsilon}_t = \begin{pmatrix} f_t \\ v_{Et} \\ v_{It} \end{pmatrix}, \mathbf{H} = \begin{pmatrix} 1 & 0 & 1 & 0 \\ 1 & 0 & 0 & 1 \end{pmatrix},$$

$$\mathbf{c}(\boldsymbol{\theta}) = \begin{pmatrix} (1-\rho_x)\mu \\ 0 \\ (1-\rho_{\epsilon_E})\delta/2 \\ -(1-\rho_{\epsilon_I})\delta/2 \end{pmatrix}, \mathbf{F}(\boldsymbol{\theta}) = \begin{pmatrix} 1+\rho_x & -\rho_x & 0 & 0 \\ 1 & 0 & 0 & 0 \\ 0 & 0 & \rho_{\epsilon_E} & 0 \\ 0 & 0 & 0 & \rho_{\epsilon_E} \end{pmatrix}, \mathbf{M}(\boldsymbol{\theta}) = \begin{pmatrix} \sigma_x & 0 & 0 \\ 0 & 0 & 0 \\ 0 & \sigma_{v_E} & 0 \\ 0 & 0 & \sigma_{v_I} \end{pmatrix},$$

$$\boldsymbol{\alpha}(\boldsymbol{\varphi}) = -a(\boldsymbol{\varphi})\boldsymbol{\beta}, \text{ and } \boldsymbol{\Upsilon}(\boldsymbol{\varphi}) = (\nu-2)\left\{\mathbf{I}_K + \frac{[a(\boldsymbol{\varphi})-1]}{\boldsymbol{\beta}'\boldsymbol{\beta}}\boldsymbol{\beta}\boldsymbol{\beta}'\right\}, \text{ with}$$

$$a(\boldsymbol{\varphi}) = \frac{-(\nu-4) + \sqrt{(\nu-4)^2 + 8(\nu-4)\boldsymbol{\beta}'\boldsymbol{\beta}}}{4\boldsymbol{\beta}'\boldsymbol{\beta}}.$$

We produce draws from the posterior distribution by means of a Gibbs sampler in which we augment the original parameter space, consisting of $\boldsymbol{\theta}$ and $\boldsymbol{\varphi}$, with the state variables $\boldsymbol{\xi}_{0:T} = \{\boldsymbol{\xi}_t\}_{t=0}^T$ and the mixing variables $\zeta_{1:T} = \{\zeta_t\}_{t=1}^T$. Throughout, we implicitly assume prior independence between $\boldsymbol{\theta}$ and $\boldsymbol{\varphi}$.

Given $\mathbf{y}_{1:T}$ and initial values $(\boldsymbol{\theta}^0, \boldsymbol{\varphi}^0, \boldsymbol{\xi}_{0:T}^0)$, we draw, for $s = 1, \ldots, S$, in the following way:

Block I: $\zeta_{1:T}^s \sim p(\zeta_{1:T}|\boldsymbol{\theta}^{s-1}, \boldsymbol{\xi}_{0:T}^{s-1}, \boldsymbol{\varphi}^{s-1}, \mathbf{y}_{1:T})$, which is given by

$$\begin{aligned}\zeta_t|\boldsymbol{\theta}, \boldsymbol{\varphi}, \boldsymbol{\xi}_{0:T}, \mathbf{y}_{1:T} &\sim GIG\left(\frac{K+\nu}{2}, \sqrt{(\nu-2)a(\boldsymbol{\varphi})\boldsymbol{\beta}'\boldsymbol{\beta}}, \sqrt{\mathrm{q}_t+1}\right),\\ \mathrm{q}_t &= \mathbf{p}_t'\boldsymbol{\Upsilon}^{-1}(\boldsymbol{\varphi})\mathbf{p}_t,\\ \mathbf{p}_t &= [\mathbf{M}'(\boldsymbol{\theta})\mathbf{M}(\boldsymbol{\theta})]^{-1}\mathbf{M}'(\boldsymbol{\theta})[\boldsymbol{\xi}_t - \mathbf{c}(\boldsymbol{\theta}) - \mathbf{F}(\boldsymbol{\theta})\boldsymbol{\xi}_{t-1}] + a(\boldsymbol{\varphi})\boldsymbol{\beta}.\end{aligned}$$

Dapugnar (1989) developed a generator of GIG pseudo-random numbers based on the ratio-of-uniforms method. In our practical implementation, we switch to a generator of gamma pseudo-random numbers whenever the norm of $\boldsymbol{\beta}$ is below the square root of $\beta_{\text{tolerance}} = 10^{-3}$ as the generator may become inefficient and unstable when the GIG distribution approaches the gamma. We also set $a(\boldsymbol{\varphi}) = 1$ and $\boldsymbol{\Upsilon}(\boldsymbol{\varphi}) = (\nu-2)\mathbf{I}_K$ for small values of the norm of $\boldsymbol{\beta}$.

Block II: $\boldsymbol{\xi}_{0:T}^s \sim p(\boldsymbol{\xi}_{0:T}|\boldsymbol{\theta}^{s-1}, \boldsymbol{\varphi}^{s-1}, \zeta_{1:T}^s, \mathbf{y}_{1:T})$, which is obtained from a modified version of the simulation smoother in Durbin and Koopman (2002) (see also Koopman and Durbin (1998)). Specifically, we proceed as follows. First of all, we note that, conditional on $\boldsymbol{\theta}$, $\boldsymbol{\varphi}$ and $\zeta_{1:T}$, the system above admits the following representation as a Gaussian linear state space model:

$$\begin{aligned}\mathbf{y}_t &= \mathbf{H}\boldsymbol{\xi}_t,\\ \boldsymbol{\xi}_t &= \mathbf{c}_t(\boldsymbol{\theta}, \boldsymbol{\varphi}) + \mathbf{F}(\boldsymbol{\theta})\boldsymbol{\xi}_{t-1} + \mathbf{M}_t(\boldsymbol{\theta}, \boldsymbol{\varphi})\mathbf{z}_t,\end{aligned}$$

where

$$\begin{aligned}\mathbf{c}_t(\boldsymbol{\theta}, \boldsymbol{\varphi}) &= \mathbf{c}(\boldsymbol{\theta}) + \mathbf{M}(\boldsymbol{\theta})[\boldsymbol{\alpha}(\boldsymbol{\varphi}) + \zeta_t^{-1}\boldsymbol{\Upsilon}(\boldsymbol{\varphi})\boldsymbol{\beta}],\\ \mathbf{M}_t(\boldsymbol{\theta}, \boldsymbol{\varphi}) &= \zeta_t^{-1/2}\mathbf{M}(\boldsymbol{\theta})\boldsymbol{\Upsilon}^{1/2}(\boldsymbol{\varphi}).\end{aligned}$$

The algorithm has three parts:

1. We draw $\{\mathbf{z}_t^+\}_{t=1}^T$ from $\mathbf{z}_t^+ \sim iid\ N(\mathbf{0}, \mathbf{I}_K)$ and $\boldsymbol{\xi}_0^+ \sim N(\boldsymbol{\xi}_{0|0}, \mathbf{P}_{0|0})$. We compute $\{\mathbf{y}_t^+\}_{t=1}^T$ and $\{\boldsymbol{\xi}_t^+\}_{t=1}^T$ by means of the recursion

$$\begin{aligned}\boldsymbol{\xi}_t^+ &= \mathbf{c}_t(\boldsymbol{\theta}, \boldsymbol{\varphi}) + \mathbf{F}(\boldsymbol{\theta})\boldsymbol{\xi}_{t-1}^+ + \mathbf{M}_t(\boldsymbol{\theta}, \boldsymbol{\varphi})\mathbf{z}_t^+,\\ \mathbf{y}_t^+ &= \mathbf{y}_t - \mathbf{H}\boldsymbol{\xi}_t^+.\end{aligned}$$

2. We run the Kalman filter followed by the Kalman smoother, storing the sequence of smoothed states $\{\hat{\boldsymbol{\xi}}_t\}_{t=0}^T$, where we denote $\hat{\boldsymbol{\xi}}_t = \boldsymbol{\xi}_{t|T}$. Specifically, for $t = 1, \ldots, T$ we first compute

$$\begin{aligned}\mathbf{K}_t &= \mathbf{P}_{t|t-1}\mathbf{H}'(\mathbf{H}\mathbf{P}_{t|t-1}\mathbf{H}')^{-1},\\ \mathbf{P}_{t|t} &= (\mathbf{I}_M - \mathbf{K}_t\mathbf{H})\mathbf{P}_{t|t-1},\\ \mathbf{P}_{t+1|t} &= \mathbf{F}(\boldsymbol{\theta})\mathbf{P}_{t|t}\mathbf{F}(\boldsymbol{\theta})' + \mathbf{M}_{t+1}(\boldsymbol{\theta}, \boldsymbol{\varphi})\mathbf{M}_{t+1}'(\boldsymbol{\theta}, \boldsymbol{\varphi}),\\ \boldsymbol{\xi}_{t|t} &= \boldsymbol{\xi}_{t|t-1} + \mathbf{K}_t(\mathbf{y}_t^+ - \mathbf{H}\boldsymbol{\xi}_{t|t-1}),\\ \boldsymbol{\xi}_{t+1|t} &= \mathbf{F}(\boldsymbol{\theta})\boldsymbol{\xi}_{t|t}.\end{aligned}$$

Then, for $\tau = 1, \ldots, T-1$ we compute

$$\begin{aligned}\mathbf{J}_{T-\tau} &= \mathbf{P}_{T-\tau|T-\tau}\mathbf{F}(\boldsymbol{\theta})'\mathbf{P}_{T-\tau+1|T-\tau}^{-1},\\ \hat{\boldsymbol{\xi}}_{T-\tau} &= \boldsymbol{\xi}_{T-\tau|T-\tau} + \mathbf{J}_{T-\tau}(\hat{\boldsymbol{\xi}}_{T-\tau+1} - \boldsymbol{\xi}_{T-\tau+1|T-\tau+1}).\end{aligned}$$

Notice that we have neglected the time-varying constants in the state-transition equation (see Jarocinski (2015) for details).

3. We compute $\{\boldsymbol{\xi}_t^*\}_{t=0}^T$ as $\boldsymbol{\xi}_t^* = \boldsymbol{\xi}_t^+ + \hat{\boldsymbol{\xi}}_t$ for $t = 0, \ldots, T$.

It turns out $\boldsymbol{\xi}_{0:T}^*$ is a draw from $p(\boldsymbol{\xi}_{0:T}|\boldsymbol{\theta}, \boldsymbol{\varphi}, \zeta_{1:T}, \mathbf{y}_{1:T})$ as desired.

Block III: $\boldsymbol{\varphi}^s \sim p(\boldsymbol{\varphi}|\boldsymbol{\theta}^{s-1}, \boldsymbol{\xi}_{0:T}^s, \zeta_{1:T}^s, \mathbf{y}_{1:T})$, which we obtain by implementing an Adaptive Rejection Metropolis Sampler (ARMS, see Gilks and Wild (1992) and Gilks, Best, and Tan (1995)). We note that $\boldsymbol{\varepsilon}_{1:T}^{s-1} = \{\boldsymbol{\varepsilon}_t^{s-1}\}_{t=1}^T$, where

$$\boldsymbol{\varepsilon}_t^{s-1} = [\mathbf{M}'(\boldsymbol{\theta}^{s-1})\mathbf{M}(\boldsymbol{\theta}^{s-1})]^{-1}\mathbf{M}'(\boldsymbol{\theta}^{s-1})[\boldsymbol{\xi}_t^{s-1} - \mathbf{c}(\boldsymbol{\theta}^{s-1}) - \mathbf{F}(\boldsymbol{\theta}^{s-1})\boldsymbol{\xi}_{t-1}^{s-1}],$$

has the sufficiency property $\boldsymbol{\varphi}|\boldsymbol{\theta}^{s-1}, \boldsymbol{\xi}_{0:T}^{s-1}, \zeta_{1:T}^s, \mathbf{y}_{1:T} \sim \boldsymbol{\varphi}|\boldsymbol{\varepsilon}_{1:T}^{s-1}, \zeta_{1:T}^s$. In addition,

$$\begin{aligned}p(\boldsymbol{\varphi}|\boldsymbol{\varepsilon}_{1:T}, \zeta_{1:T}) &\propto \left[\prod_{t=1}^T p(\boldsymbol{\varepsilon}_t|\boldsymbol{\varepsilon}_{1:t-1}, \boldsymbol{\varphi}, \zeta_{1:t})p(\zeta_t|\boldsymbol{\varepsilon}_{1:t-1}, \boldsymbol{\varphi}, \zeta_{1:t-1})\right] p(\boldsymbol{\varphi}),\\ \boldsymbol{\varepsilon}_t|\boldsymbol{\varepsilon}_{1:t-1}, \boldsymbol{\varphi}, \zeta_{1:t} &\sim N[\boldsymbol{\alpha}(\boldsymbol{\varphi}) + \zeta_t^{-1}\boldsymbol{\Upsilon}(\boldsymbol{\varphi})\boldsymbol{\beta}, \zeta_t^{-1}\boldsymbol{\Upsilon}(\boldsymbol{\varphi})],\\ \zeta_t|\boldsymbol{\varepsilon}_{1:t-1}, \boldsymbol{\varphi}, \zeta_{1:t-1} &\sim \Gamma\left(\nu/2, 1/2\right).\end{aligned}$$

Thus, the log-likelihood we employ (up to an additive term constant in $\boldsymbol{\varphi}$) is

$$\mathcal{L}(\boldsymbol{\varphi}) = -\frac{T}{2}\log\{\det[\boldsymbol{\Upsilon}(\boldsymbol{\varphi})]\} - \frac{1}{2}\sum_{t=1}^T \tilde{\boldsymbol{\varepsilon}}_t'\tilde{\boldsymbol{\varepsilon}}_t - T\left[\frac{\nu}{2}\log(2) + \log\Gamma\left(\frac{\nu}{2}\right)\right] + \frac{\nu}{2}\sum_{t=1}^T \log(\zeta_t),$$

where

$$\tilde{\boldsymbol{\varepsilon}}_t \equiv \zeta_t^{1/2}\boldsymbol{\Upsilon}^{-1/2}(\boldsymbol{\varphi})[\boldsymbol{\varepsilon}_t - \boldsymbol{\alpha}(\boldsymbol{\varphi}) - \zeta_t^{-1}\boldsymbol{\Upsilon}(\boldsymbol{\varphi})\boldsymbol{\beta}].$$

We apply ARMS to each parameter in turn. Let ϑ be the result of applying a certain transformation to the specific entry of $\boldsymbol{\varphi}$ being updated. In particular, for the parameter ν we let $\vartheta = \nu_{\min}/\nu$ (we take $\nu_{\min} = 4$) while for β_j we use $\vartheta = [1 + \exp(-\beta_j)]^{-1}$, $j = x, 1, 2$. The transformation is chosen in all cases to ensure $\vartheta \in [0, 1]$.

Let ϑ^0 be the starting value and $\mathcal{L}^0$ its log-posterior. ARMS updates ϑ^0 to ϑ^1 as follows:

1. Construct a grid $\vartheta_1, \ldots, \vartheta_{n_{\text{ARMS}}}$ and compute their log-posteriors $\mathcal{L}_1, \ldots, \mathcal{L}_{n_{\text{ARMS}}}$.

2. Form the piecewise-linear function h given by

$$\begin{aligned}h(\vartheta) &= \max\left\{\mathrm{L}_j(\vartheta), \min\left[\mathrm{L}_{j-1}(\vartheta), \mathrm{L}_{j+1}(\vartheta)\right]\right\}, \qquad \vartheta_j < \vartheta \leq \vartheta_{j+1},\\ \mathrm{L}_j(\vartheta) &= \mathcal{L}_j + \mathcal{L}_{j+1}\frac{(\vartheta - \vartheta_j)}{(\vartheta_{j+1} - \vartheta_j)}.\end{aligned}$$

Next, draw ϑ^* from the piecewise exponencial distribution with density proportional to $\exp[h(\vartheta)]$. In other words, draw first a sub-interval and, conditioning on it, from a scaled truncated exponential distribution. Compute the associated log-posterior $\mathcal{L}^*$.

3. Draw $u_{\text{ARS}} \sim \text{U}[0,1]$. If $\log(u_{\text{ARS}}) > \mathcal{L}^* - h(\vartheta^*)$, augment the grid of ϑ by ϑ^* and that of $\mathcal{L}$ by $\mathcal{L}^*$ and go back to 2. Otherwise, move on to 4.

4. Draw $u_{\text{MH}} \sim \text{U}[0,1]$. If $\log(u_{\text{MH}}) > \mathcal{L}^* - \mathcal{L}^0$, set $\vartheta^1 = \vartheta^0$. Otherwise, set $\vartheta^1 = \vartheta^*$.

In the implementation, each draw $\boldsymbol{\varphi}^s$ is obtained by repeating the algorithm above n_{MH} times before proceeding with the Gibbs sampler.

We have also considered Slice Sampling (SS, see Neal (2003)) as an alternative method to update ϑ^0 to ϑ^1. The alternative sampling is done as follows:

1. Draw $e \sim \exp(1)$ (so that $y = \exp(\mathcal{L}^0 - e) \sim \text{U}[0, \exp(\mathcal{L}^0)]$).

2. Given a positive real number w, draw $u \sim \text{U}[0,1]$ and let $\vartheta_{\text{L}} = \max\{\vartheta^0 - uw, 0\}$ and $\vartheta_{\text{R}} = \min\{\text{L} + w, 1\}$. Let $\mathcal{L}_{\text{L}}$ and $\mathcal{L}_{\text{R}}$ be their respective log-posteriors.
   Given an integer m_{SS}, draw $v \sim \text{U}[0,1]$ and form $m_{\text{L}} = v m_{\text{SS}}$ and $m_{\text{R}} = m_{\text{SS}} - 1 - m_{\text{L}}$. While $\mathcal{L}_{\text{L}} > \mathcal{L}^0 - e$ and $m_{\text{L}} > 0$ update ϑ_{L} to $\max\{\vartheta_{\text{L}} - w, 0\}$ (recomputing $\mathcal{L}_{\text{L}}$) and m_{L} to $m_{\text{L}} - 1$. Likewise, update ϑ_{R} to $\min\{\vartheta_{\text{R}} + w, 1\}$ (recomputing $\mathcal{L}_{\text{R}}$) and m_{R} to $m_{\text{R}} - 1$ while $\mathcal{L}_{\text{R}} > \mathcal{L}^0 - e$ and $m_{\text{R}} > 0$.

3. Draw $\vartheta^* \sim \text{U}[\vartheta_{\text{L}}, \vartheta_{\text{R}}]$ and let $\mathcal{L}^*$ be its log-posterior. While $\mathcal{L}^* < \mathcal{L}^0 - e$, either $\vartheta^* < \vartheta^0$, in which case update ϑ_{L} to ϑ^*, or $\vartheta^* \geq \vartheta^0$, in which case update ϑ_{R} to ϑ^*. Re-draw $\vartheta^* \sim \text{U}[\vartheta_{\text{L}}, \vartheta_{\text{R}}]$ and re-compute $\mathcal{L}^*$. When this process terminates, set $\vartheta^1 = \vartheta^*$.

We report the output of the algorithm based on ARMS but our results are robust to the sampling method.

Block IV: $\boldsymbol{\theta}^s \sim p(\boldsymbol{\theta}|\boldsymbol{\xi}^s_{0:T}, \boldsymbol{\varphi}^s, \zeta^s_{1:T}, \mathbf{y}_{1:T})$, which is obtained in blocks. First, we note the sufficiency property $\boldsymbol{\theta}|\boldsymbol{\xi}^s_{0:T}, \boldsymbol{\varphi}^s, \zeta^s_{1:T}, \mathbf{y}_{1:T} \sim \boldsymbol{\theta}|\boldsymbol{\xi}^s_{0:T}, \boldsymbol{\varphi}^s, \zeta^s_{1:T}$. Next, we partition $\boldsymbol{\theta} = (\boldsymbol{\theta}'_c, \boldsymbol{\theta}'_\rho, \boldsymbol{\theta}'_\gamma)'$, with $\boldsymbol{\theta}_c = (\mu, \delta)'$, $\boldsymbol{\theta}_\rho = (\rho_x, \rho_{\epsilon_E}, \rho_{\epsilon_I})'$ and $\boldsymbol{\theta}_\sigma = (\sigma^2_x, \sigma^2_{v_E}, \sigma^2_{v_I})'$. We proceed as follows:

1. We set a Gaussian prior on $\boldsymbol{\theta}_c$ given by $\boldsymbol{\theta}_c \sim N(\underline{\mathbf{c}}, \underline{\mathbf{S}}_c)$ and we draw from the posterior $\boldsymbol{\theta}_c|\boldsymbol{\theta}^{s-1}_\rho, \boldsymbol{\theta}^{s-1}_\sigma, \boldsymbol{\xi}^s_{0:T}, \boldsymbol{\varphi}^s, \zeta^s_{1:T}$, which is

   $$\boldsymbol{\theta}_c|\boldsymbol{\theta}_\rho, \boldsymbol{\theta}_\sigma, \boldsymbol{\xi}_{0:T}, \boldsymbol{\varphi}, \zeta_{1:T} \sim N\left(\bar{\mathbf{c}}, \bar{\mathbf{S}}_c\right), \text{ with } \bar{\mathbf{c}} = \bar{\mathbf{S}}_c(\underline{\mathbf{S}}_c^{-1}\underline{\mathbf{c}} + \hat{\mathbf{S}}_c^{-1}\hat{\mathbf{c}}) \text{ and } \bar{\mathbf{S}}_c = (\underline{\mathbf{S}}_c^{-1} + \hat{\mathbf{S}}_c^{-1})^{-1},$$

   where

   $$\hat{\mathbf{S}}_c = \left[\sum_{t=1}^T \zeta_t\right]^{-1} [\mathbf{D}_c \mathbf{M}(\boldsymbol{\theta}) \boldsymbol{\Upsilon}(\boldsymbol{\varphi}) \mathbf{M}(\boldsymbol{\theta}) \mathbf{D}_c], \text{ and } \hat{\mathbf{c}} = \left[\sum_{t=1}^T \zeta_t\right]^{-1} \left[\sum_{t=1}^T \sqrt{\zeta_t} \mathbf{Y}_{ct}\right],$$

with

$$\mathbf{D}_c = \begin{bmatrix} 1-\rho_x & 0 & 0 & 0 \\ 0 & 0 & 1-\rho_{\epsilon_E} & -(1-\rho_{\epsilon_I}) \end{bmatrix}$$

and

$$\mathbf{Y}_{ct} = \mathbf{D}_c\{\boldsymbol{\xi}_t - \mathbf{F}(\boldsymbol{\theta})\boldsymbol{\xi}_{t-1} - \mathbf{M}(\boldsymbol{\theta})[\boldsymbol{\alpha}(\boldsymbol{\varphi}) + \zeta_t^{-1}\boldsymbol{\Upsilon}(\boldsymbol{\varphi})\boldsymbol{\beta}]\}.$$

2. We set a Gaussian prior on $\boldsymbol{\theta}_\rho$ given by $\boldsymbol{\theta}_\rho \sim N(\underline{\boldsymbol{\rho}}, \underline{\mathbf{S}}_\rho)$ and we draw from the posterior $\boldsymbol{\theta}_\rho|\boldsymbol{\theta}_c^s, \boldsymbol{\theta}_\sigma^{s-1}, \boldsymbol{\xi}_{0:T}^s, \boldsymbol{\varphi}^s, \zeta_{1:T}^s$, which is

   $$\boldsymbol{\theta}_\rho|\boldsymbol{\theta}_c, \boldsymbol{\theta}_\sigma, \boldsymbol{\xi}_{0:T}, \boldsymbol{\varphi}, \zeta_{1:T} \sim N\left(\bar{\boldsymbol{\rho}}, \bar{\mathbf{S}}_\rho\right), \text{ with } \bar{\boldsymbol{\rho}} = \bar{\mathbf{S}}_\rho(\underline{\mathbf{S}}_\rho^{-1}\underline{\boldsymbol{\rho}} + \hat{\mathbf{S}}_\rho^{-1}\hat{\boldsymbol{\rho}}) \text{ and } \bar{\mathbf{S}}_\rho = (\underline{\mathbf{S}}_\rho^{-1} + \hat{\mathbf{S}}_\rho^{-1})^{-1},$$

   where

   $$\hat{\mathbf{S}}_\rho = \left[\sum_{t=1}^T \mathbf{X}_{\rho t}^2\right]^{-1} \text{ and } \hat{\boldsymbol{\rho}} = \left[\sum_{t=1}^T \mathbf{X}_{\rho t}^2\right]^{-1}\left[\sum_{t=1}^T \mathbf{X}_{\rho t}\mathbf{Y}_{\rho t}\right],$$

   with

   $$\begin{aligned} \mathbf{X}_{\rho t} &= \zeta_t^{1/2}\text{diag}[\mathbf{W}_\rho \mathbf{D}_\rho\{\boldsymbol{\xi}_{t-1} - [\mathbf{I}_K - \mathbf{D}_\rho\mathbf{F}(\boldsymbol{\theta})]^{-1}\mathbf{C}(\boldsymbol{\theta})\}], \\ \mathbf{Y}_{\rho t} &= \zeta_t^{1/2}\mathbf{W}_\rho\mathbf{D}_\rho\{\boldsymbol{\xi}_t - [\mathbf{I}_K - \mathbf{D}_\rho\mathbf{F}(\boldsymbol{\theta})]^{-1}\mathbf{C}(\boldsymbol{\theta}) - \mathbf{M}(\boldsymbol{\theta})[\boldsymbol{\alpha}(\boldsymbol{\varphi}) + \zeta_t^{-1}\boldsymbol{\Upsilon}(\boldsymbol{\varphi})\boldsymbol{\beta}]\}, \end{aligned}$$

   $$\mathbf{W}_\rho = [\mathbf{D}_\rho\mathbf{M}(\boldsymbol{\theta})\boldsymbol{\Upsilon}^{1/2}(\boldsymbol{\varphi})]^{-1} \text{ and } \mathbf{D}_\rho \equiv \begin{bmatrix} 1 & -1 & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \end{bmatrix}.$$

3. We set an inverse gamma prior on $\boldsymbol{\theta}_\sigma$ given by $\sigma_j^{-2} \sim \Gamma(\nu_j/2, \varsigma_j/2)$, for $j = x, v_E, v_I$, with these parameters being prior-independent across j. However, for the purposes of generating draws from the posterior distribution $\boldsymbol{\theta}_\sigma|\boldsymbol{\theta}_c^s, \boldsymbol{\theta}_\rho^s, \boldsymbol{\xi}_{0:T}^s, \boldsymbol{\varphi}^s, \zeta_{1:T}^s$, we need to consider two separate cases.

   If $\boldsymbol{\beta} = \mathbf{0}$, the three parameters are posterior-independent and direct sampling can be implemented because the prior conjugates with the likelihood. More formally,

   $$\sigma_j^{-2}|\boldsymbol{\theta}_c, \boldsymbol{\theta}_\rho, \boldsymbol{\xi}_{0:T}, \boldsymbol{\varphi}, \zeta_{1:T} \sim \Gamma\left[\frac{T+\nu_j}{2}, \frac{1}{2}\left(\sum_{t=1}^T \eta_{jt}^2 + \varsigma_j\right)\right], \text{ for } j = x, v_E, v_I,$$

   where

   $$\boldsymbol{\eta}_t = \begin{bmatrix} \eta_{xt} \\ \eta_{v_E t} \\ \eta_{v_I t} \end{bmatrix} = \zeta_t^{1/2}\boldsymbol{\Upsilon}^{-1/2}(\boldsymbol{\theta})\mathbf{D}_\rho[\boldsymbol{\xi}_t - \mathbf{C}(\boldsymbol{\theta}) - \mathbf{F}(\boldsymbol{\theta})\boldsymbol{\xi}_{t-1}].$$

   On the other hand, if $\boldsymbol{\beta} \neq \mathbf{0}$, direct sampling is not available. In this case, we generate draws from the posterior distribution by componentwise application of ARMS. The log-likelihood that we employ (up to an additive term constant in $\boldsymbol{\theta}_\sigma$) is

   $$\mathcal{L}(\boldsymbol{\theta}_\sigma) = -\frac{T}{2}[\log(\sigma_x^2) + \log(\sigma_{v_E}^2) + \log(\sigma_{v_I}^2)] - \frac{1}{2}\sum_{t=1}^T \tilde{\boldsymbol{\eta}}_t'\tilde{\boldsymbol{\eta}}_t,$$

   with

   $$\begin{aligned} \tilde{\boldsymbol{\eta}}_t &= \zeta_t^{1/2}\boldsymbol{\Upsilon}^{-1/2}(\boldsymbol{\varphi})\text{diag}(\sigma_x^{-1}, \sigma_{v_E}^{-1}, \sigma_{v_I}^{-1})\mathbf{D}_\rho \\ &\quad \times[\boldsymbol{\xi}_t - \mathbf{C}(\boldsymbol{\theta}) - \mathbf{F}(\boldsymbol{\theta})\boldsymbol{\xi}_{t-1} - \boldsymbol{\alpha}(\boldsymbol{\varphi}) - \zeta_t^{-1}\boldsymbol{\Upsilon}(\boldsymbol{\varphi})\boldsymbol{\beta}]. \end{aligned}$$

The procedure is exactly as explained above. Again, we also employed slice sampling, with our results being robust to this variation.

E Additional Monte Carlo results (HAC)

Table E1: Monte Carlo rejection rates (in %) under null and alternative hypotheses for the bivariate cointegrated, dynamic single factor model ($T = 100$)

		Panel A: Null hypothesis			Panel B: Alternative hypotheses (5%) Student t			asymmetric Student t		
		1%	5%	10%	J	S_f	$S_{\mathbf{v}}$	J	S_f	$S_{\mathbf{v}}$
H_J	Kt	0.17	1.71	4.56	24.76	2.29	15.76	30.25	2.94	19.88
	Sk	2.75	9.55	16.99	8.60	9.67	8.55	20.60	12.29	15.96
	GH	2.20	7.98	14.20	15.89	8.31	12.47	34.06	10.92	23.73
$H_{S_\mathbf{f}}$	Kt	0.17	1.67	4.62	5.14	3.75	2.13	9.13	4.50	2.94
	Sk	1.35	6.37	12.77	6.16	6.34	6.13	11.00	12.46	6.39
	GH	0.75	4.06	9.02	6.07	5.31	4.34	13.72	11.78	5.19
$H_{S_\mathbf{v}}$	Kt	0.17	1.64	4.77	18.89	1.73	17.50	24.15	2.00	20.95
	Sk	1.67	7.55	13.97	6.65	7.33	7.06	14.73	7.91	17.73
	GH	1.17	5.80	11.44	12.19	5.74	11.72	25.37	5.83	27.46
Red	Kt	0.25	1.89	5.13	24.17	2.51	13.68	28.05	3.40	17.73
	Sk	1.52	6.57	13.12	6.80	7.07	6.41	15.39	8.46	6.81
	GH	1.00	5.21	10.87	14.35	5.94	9.68	28.67	7.84	12.08

Notes: Results based on 10,000 samples of size $T = 100$ from model (16) with $\rho_x = .5$, $\rho_{\epsilon_E} = .2$, $\rho_{\epsilon_I} = .8$, $\sigma_f^2 = 1$ and $\sigma_{v_i}^2$ chosen such that $q_E = 2$ and $q_I = .5$, where $q_i = \sigma_x^2/\sigma_{\epsilon_i}^2$ represents the signal-to-noise ratio for y_{it} for $i = E, I$. The column labels J, S_f, $S_{\mathbf{v}}$ refer to the alternative $\boldsymbol{\varepsilon}_t \sim GH(\eta, \psi, \boldsymbol{\beta})$ (i.e. $R = 3$), $f_t \sim GH(\eta, \psi, \beta)$, $\mathbf{v}_t \sim N(\mathbf{0}, \mathbf{I}_N)$ ($R = 1$) and $\mathbf{v}_t \sim GH(\eta, \psi, \boldsymbol{\beta})$, $f_t \sim N(0, 1)$ ($R = 2$), respectively. The row labels H_J, H_{S_f}, and $H_{S_\mathbf{v}}$ refer to the score tests in Propositions 4 and 7 corresponding to the J, S_f, and $S_{\mathbf{v}}$ alternative hypotheses, while Red denotes the reduced form tests discussed in section 5.4.2. In Panel B, Student t refers to the DGP for the GH being symmetric Student t with 8 degrees of freedom and, analogously, asymmetric Student t to the asymmetric Student t with 8 degrees of freedom and skewness vector $\boldsymbol{\beta} = -\boldsymbol{\ell}_R$. For each of those labels, Kt and Sk refer to the kurtosis and skewness components of the corresponding test statistics, while GH indicates the sum of the two.

Table E2: Monte Carlo rejection rates (in %) under null and alternative hypotheses for the bivariate cointegrated, dynamic single factor model ($T = 250$)

		Panel A: Null hypothesis			Panel B: Alternative hypotheses (5%) Student t			asymmetric Student t		
		1%	5%	10%	J	S_f	$S_{\mathbf{v}}$	J	S_f	$S_{\mathbf{v}}$
H_J	Kt	0.12	1.79	5.15	61.57	3.57	42.37	65.21	5.22	49.08
	Sk	1.61	7.05	12.87	6.15	7.45	6.23	43.14	14.99	29.48
	GH	1.17	5.77	10.80	27.25	6.48	18.65	68.37	14.98	52.18
$H_{S_\mathbf{f}}$	Kt	0.13	1.73	4.97	11.68	6.51	2.35	20.10	8.69	4.21
	Sk	1.13	6.06	12.15	5.64	5.42	5.90	18.32	27.79	7.48
	GH	0.56	3.94	8.70	8.23	5.77	4.21	27.72	28.67	6.26
$H_{S_\mathbf{v}}$	Kt	0.12	1.48	4.86	50.02	1.99	46.26	58.78	2.45	52.70
	Sk	1.31	5.94	11.84	4.88	6.36	4.98	32.34	6.23	39.44
	GH	0.95	4.38	9.29	22.91	4.98	21.18	59.80	5.03	63.81
Red	Kt	0.15	1.85	5.65	59.11	3.96	34.63	61.54	5.98	43.07
	Sk	1.29	6.14	11.61	5.63	6.12	5.25	35.12	12.29	6.82
	GH	0.85	4.70	9.70	28.61	5.47	16.30	62.18	12.54	24.03

Notes: Results based on 10,000 samples of size $T = 250$ from model (16) with $\rho_x = .5$, $\rho_{\epsilon_E} = .2$, $\rho_{\epsilon_I} = .8$, $\sigma_f^2 = 1$ and $\sigma_{v_i}^2$ chosen such that $q_E = 2$ and $q_I = .5$, where $q_i = \sigma_x^2/\sigma_{\epsilon_i}^2$ represents the signal-to-noise ratio for y_{it} for $i = E, I$. The column labels J, S_f, $S_{\mathbf{v}}$ refer to the alternative $\boldsymbol{\varepsilon}_t \sim GH(\eta, \psi, \boldsymbol{\beta})$ (i.e. $R = 3$), $f_t \sim GH(\eta, \psi, \beta)$, $\mathbf{v}_t \sim N(\mathbf{0}, \mathbf{I}_N)$ ($R = 1$) and $\mathbf{v}_t \sim GH(\eta, \psi, \boldsymbol{\beta})$, $f_t \sim N(0, 1)$ ($R = 2$), respectively. The row labels H_J, H_{S_f}, and $H_{S_\mathbf{v}}$ refer to the score tests in Propositions 4 and 7 corresponding to the J, S_f, and $S_{\mathbf{v}}$ alternative hypotheses, while Red denotes the reduced form tests discussed in section 5.4.2. In Panel B, Student t refers to the DGP for the GH being symmetric Student t with 8 degrees of freedom and, analogously, asymmetric Student t to the asymmetric Student t with 8 degrees of freedom and skewness vector $\boldsymbol{\beta} = -\boldsymbol{\ell}_R$. For each of those labels, Kt and Sk refer to the kurtosis and skewness components of the corresponding test statistics, while GH indicates the sum of the two.

Table E3: Monte Carlo rejection rates (in %) under the null and alternative hypotheses for the local-level model

		Panel A: Null hypothesis			Panel B: Alternative hypotheses (5%) Student t			asymmetric Student t		
		1%	5%	10%	J	S_f	S_v	J	S_f	S_v
H_J	Kt	0.08	1.52	4.77	23.84	7.06	3.41	35.64	11.68	6.77
	Sk	1.33	6.15	12.01	4.90	5.70	5.38	41.94	24.15	11.18
	GH	0.73	4.65	9.57	11.60	6.60	4.98	57.08	29.14	12.57
H_{S_f}	Kt	0.10	1.60	5.08	17.91	8.72	2.24	27.32	12.57	3.53
	Sk	0.96	6.01	11.76	5.29	5.49	5.44	47.23	33.42	5.09
	GH	0.52	3.64	8.25	10.89	6.60	3.79	57.18	36.83	4.49
H_{S_v}	Kt	0.17	1.67	4.78	14.25	2.95	4.71	31.21	5.41	8.68
	Sk	1.03	5.41	10.78	3.94	5.18	4.61	24.47	4.22	15.97
	GH	0.51	3.44	7.72	8.27	3.95	4.41	41.34	4.47	18.39
Red	Kt	0.05	1.46	5.20	22.68	6.98	2.85	33.28	11.69	5.45
	Sk	1.07	5.75	11.45	4.65	5.50	5.01	55.45	28.16	6.84
	GH	0.53	3.48	8.12	12.35	6.38	4.02	64.79	31.23	6.66

Notes: Results based on 10,000 samples of size $T = 250$ from the local-level model discussed in section 5.3 in which the signal-to-noise ratio $q = \sigma_f^2/\sigma_v^2$ is set to 2. The column labels J, S_f, $S_{\mathbf{v}}$ refer to the alternative $\boldsymbol{\varepsilon}_t \sim GH(\eta, \psi, \boldsymbol{\beta})$ $(R = 2)$, $f_t \sim GH(\eta, \psi, \beta)$, $v_t \sim N(0, 1)$ $(R = 1)$ and $v_t \sim GH(\eta, \psi, \beta)$, $f_t \sim N(0, 1)$ $(R = 1)$, respectively. The row labels H_J, H_{S_f}, and H_{S_v} refer to the score tests in Propositions 4 and 7 corresponding to the J, S_f, and S_v alternative hypotheses, Red denotes the reduced form tests discussed in section 5.4.2, while HK denotes the original Harvey and Koopman (1992) tests discussed in section 5.4.1. In Panel B, Student t refers to the DGP for the GH being symmetric Student t with 8 degrees of freedom and, analogously, asymmetric Student t to the asymmetric Student t with 8 degrees of freedom and skewness vector $\boldsymbol{\beta} = -\boldsymbol{\ell}_R$. For each of those labels, Kt and Sk refer to the kurtosis and skewness components of the corresponding test statistics, while GH indicates the sum of the two.

Table E4: Monte Carlo rejection rates (in %) under null and alternative hypotheses for the multivariate local-level model

		Panel A: Null hypothesis			Panel B: Alternative hypotheses (5%) Student t			asymmetric Student t		
		1%	5%	10%	J	S_f	$S_{\mathbf{v}}$	J	S_f	$S_{\mathbf{v}}$
H_J	Kt	0.23	3.15	7.37	96.21	6.69	95.78	93.94	10.63	95.13
	Sk	8.95	21.52	31.50	18.46	20.76	18.73	69.96	33.29	40.30
	GH	8.63	20.71	30.31	75.83	21.25	74.13	92.99	35.27	85.33
$H_{S_\mathbf{f}}$	Kt	0.06	1.78	5.56	32.25	30.09	1.68	35.83	34.12	2.19
	Sk	1.18	5.73	11.57	4.49	4.44	5.59	44.24	60.98	5.06
	GH	0.62	3.68	8.33	17.26	16.19	3.77	58.05	68.64	3.65
$H_{S_\mathbf{v}}$	Kt	0.29	2.82	7.06	95.86	3.07	95.80	95.52	2.90	95.60
	Sk	7.73	18.94	28.70	16.00	18.45	16.05	63.85	19.16	40.45
	GH	7.42	18.34	27.33	74.66	17.18	74.22	93.08	18.06	86.57
Red	Kt	0.25	3.13	7.36	96.71	6.96	95.77	93.87	11.11	95.14
	Sk	7.74	18.94	28.44	16.14	18.56	16.26	62.83	31.52	26.57
	GH	7.11	17.91	27.21	75.61	19.07	73.54	91.69	33.83	79.37

Notes: Results based on 10,000 samples of size $T = 250$ from a 10-variate version of the local-level model with $\boldsymbol{\pi} = \mathbf{0}$, $\mathbf{c} = \boldsymbol{\ell}_{10}$ and $\boldsymbol{\gamma} = q^{-1}\boldsymbol{\ell}_{10}$, where q reflects the signal-to-noise ratio, which we set to 2. The column labels J, S_f, $S_{\mathbf{v}}$ refer to the alternative $\boldsymbol{\varepsilon}_t \sim GH(\eta, \psi, \boldsymbol{\beta})$ (i.e. $R = 11$), $f_t \sim GH(\eta, \psi, \beta)$, $\mathbf{v}_t \sim N(\mathbf{0}, \mathbf{I}_N)$ ($R = 1$) and $\mathbf{v}_t \sim GH(\eta, \psi, \boldsymbol{\beta})$, $f_t \sim N(0, 1)$ ($R = 10$), respectively. The row labels H_J, H_{S_f}, and $H_{S_\mathbf{v}}$ refer to the score tests in Propositions 4 and 7 corresponding to the J, S_f, and $S_{\mathbf{v}}$ alternative hypotheses. In Panel B, Student t refers to the DGP for the GH being symmetric Student t with 8 degrees of freedom and, analogously, asymmetric Student t to the asymmetric Student t with 8 degrees of freedom and skewness vector $\boldsymbol{\beta} = -\boldsymbol{\ell}_R$. For each of those labels, Kt and Sk refer to the kurtosis and skewness components of the corresponding test statistics, while GH indicates the sum of the two.

F Inferring real output from GDP and GDI over a long sample

Table F1: Parameter estimates and normality tests over the postwar period

Panel A: ML estimates

Param.	estimate	std. err.
μ	0.755	0.110
δ	0.304	0.031
α_x	0.493	0.059
α_{ϵ_E}	0.265	0.196
α_{ϵ_I}	0.939	0.024
σ_f^2	0.526	0.054
$\sigma_{v_E}^2$	0.076	0.021
$\sigma_{v_I}^2$	0.093	0.019

Panel B: Normality tests

		statistic	p-value
$H_{S_\mathbf{f}}$	Kt	19.061	0.000
	Sk	1.161	0.281
	GH	20.221	0.000
$H_{S_\mathbf{v}}$	Kt	6.537	0.005
	Sk	3.859	0.145
	GH	10.396	0.011
H_R	Kt	13.266	0.000
	Sk	1.232	0.540
	GH	14.498	0.002

Notes: Data: Quarterly real GDP and GDI from 1952Q1 to 2015Q2. Model: Bivariate cointegrated, dynamic single factor model (16); see section 7 for parameter definitions. In Panel A, estimates are Gaussian ML of the bivariate Gaussian likelihood of the stationary transformation $\Delta y_{Et} + \Delta y_{It}$ and $y_{Et} - y_{It}$ in the time domain. Standard errors are obtained from the asymptotic information matrix, which is computed using its frequency domain closed-form expression. In Panel B, the row labels H_{S_f} and $H_{S_\mathbf{v}}$ refer to the score tests in Propositions 4 and 7 corresponding to the S_f and $S_\mathbf{v}$ alternative hypotheses, respectively, while Red denotes the reduced form tests discussed in section 5.4.2. For each of those labels, Kt and Sk refer to the kurtosis and skewness components of the corresponding test statistics, while GH indicates the sum of the two.

Figure F1: Smoothed innovations and influence functions for the kurtosis and skewness tests: Sample 1952Q1 to 2015Q2.

Figure F1a: Smoothed innovations for the underlying factor

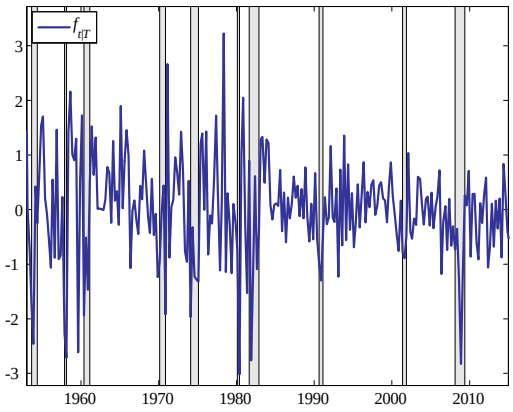


Figure F1b: Smoothed innovations for the measurement errors

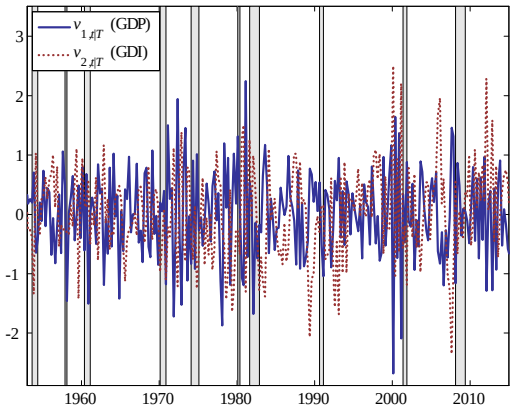


Figure F1c: Influence functions for the underlying factor (kurtosis)

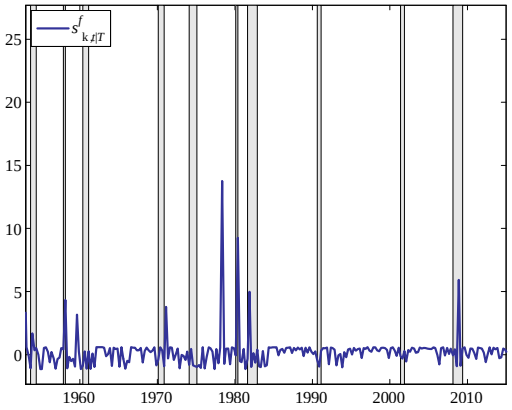


Figure F1d: Influence functions for the measurement errors (kurtosis)

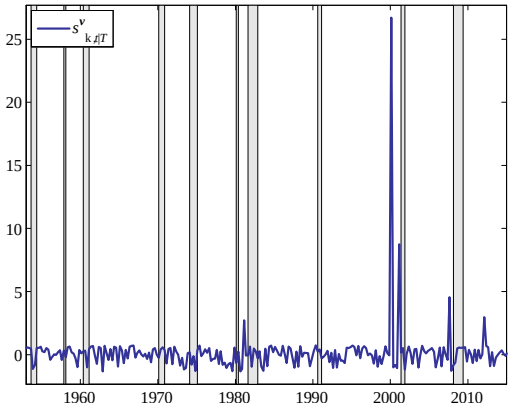


Figure F1e: Influence functions for the underlying factor (skewness)

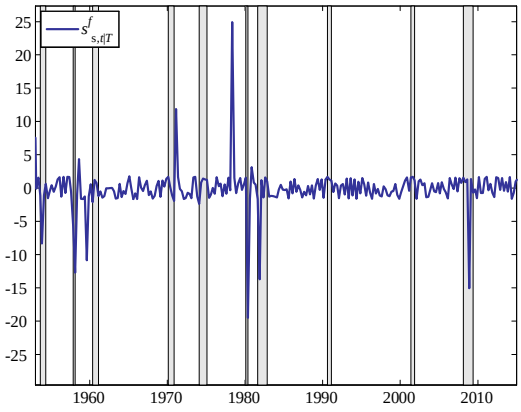
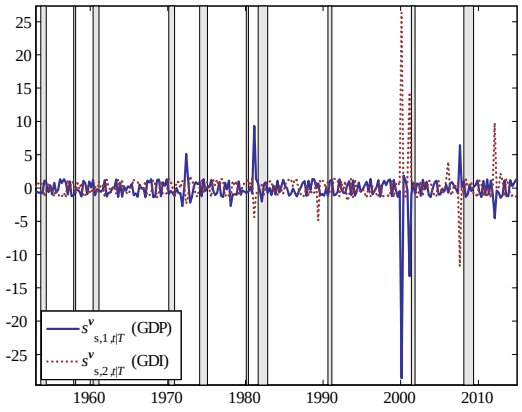


Figure F1f: Influence functions for the measurement errors (skewness)



Notes: Smoothed innovations and influence functions were obtained from fitting the bivariate cointegrated, dynamic single factor model (16) to the quarterly real GDP and GDI from 1952Q1 to 2015Q2; see Table F1 for parameter estimates. Shaded areas represent NBER recessions.

Figure F2: Posterior densities of shape parameters under the asymmetric Student t alternative: Sample 1952Q1 to 2015Q2

Figure F2a: η

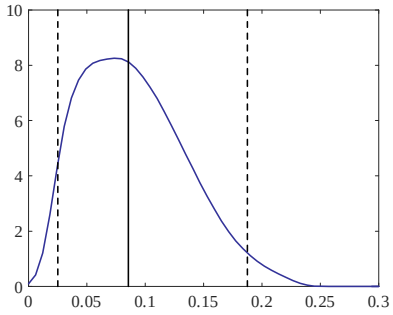


Figure F2b: β_x

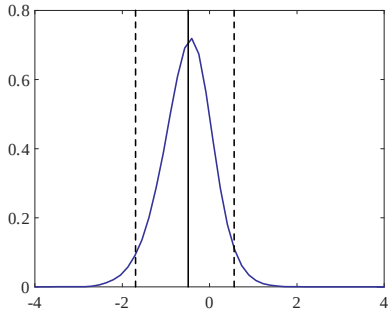


Figure F2c: β_{v_E}

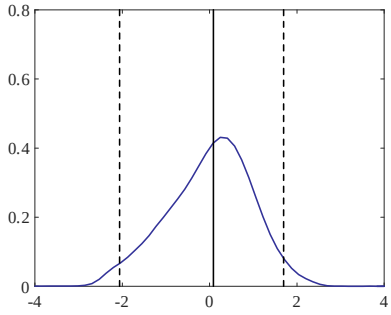
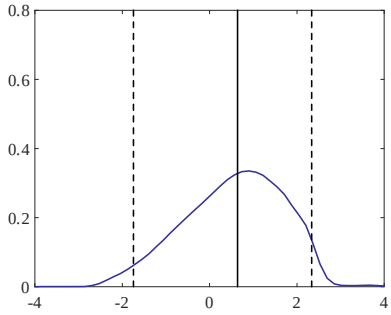


Figure F2d: β_{v_I}



Notes: Model: Bivariate cointegrated, dynamic single factor model (16) with multivariate asymmetric Student t innovations; see Section 7 for parameter definitions. η refers to the reciprocal of degrees of freedom while β_x (β_{v_E}) [β_{v_I}] refers to the skewness parameter of the "true GDP" (expenditure) [income] measure. Solid vertical lines refer to the median values while dashed lines report the 2.5% and 97.5% quantiles.

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