# Testing for observation-dependent regime switching in mixture autoregressive models* 

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#### Abstract

Testing for regime switching when the regime switching probabilities are specified either as constants ('mixture models') or are governed by a finite-state Markov chain ('Markov switching models') are long-standing problems that have also attracted recent interest. This paper considers testing for regime switching when the regime switching probabilities are time-varying and depend on observed data ('observation-dependent regime switching'). Specifically, we consider the likelihood ratio test for observation-dependent regime switching in mixture autoregressive models. The testing problem is highly nonstandard, involving unidentified nuisance parameters under the null, parameters on the boundary, singular information matrices, and higher-order approximations of the log-likelihood. We derive the asymptotic null distribution of the likelihood ratio test statistic in a general mixture autoregressive setting using high-level conditions that allow for various forms of dependence of the regime switching probabilities on past observations, and we illustrate the theory using two particular mixture autoregressive models. The likelihood ratio test has a nonstandard asymptotic distribution that can easily be simulated, and Monte Carlo studies show the test to have good finite sample size and power properties.


JEL classification: C12, C22, C52.
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## 1 Introduction

Different regime switching models are in widespread use in economics, finance, and other fields. When the regime switching probabilities are constants, these models are often referred to as 'mixture models', and when these probabilities depend on past regimes and are governed by a finite-state Markov chain, the term (time-homogeneous) 'Markov switching models' is typically used. In this paper, we are interested in the case where the regime switching probabilities depend on observed data but not on past regimes, a case we refer to as 'observation-dependent regime switching'. Models of this kind can be viewed as special cases of time-inhomogeneous Markov switching models (in which regime switching probabilities depend on both past regimes and observed data). Overviews of regime switching models can be found, for example, in Frühwirth-Schnatter (2006) and Hamilton (2016). Of critical interest in all these models is whether the use of several regimes is warranted or if a single-regime model would suffice. Testing for regime switching in all these models is plagued by several irregular features such as unidentified parameters and parameters on the boundary, and is consequently notoriously difficult.

Tests for Markov switching have been considered by several authors in the econometrics literature. Hansen (1992) and Garcia (1998) both considered sup-type likelihood ratio (LR) tests in Markov switching models but they did not present complete solutions. Hansen derived a bound for the distribution of the LR statistic, leading to a conservative procedure, while Garcia did not treat all the non-standard features of the problem in detail. Cho and White (2007) analyzed the use of a LR statistic for a mixture model to test for Markov-switching type regime switching. They found their test based on a mixture model to have power against Markov switching alternatives even though it ignores the temporal dependence of the Markov chain. Carrasco, Hu, and Ploberger (2014) took a different approach and proposed an information matrix type test that they showed to be asymptotically optimal against Markov switching alternatives. Very recently, both Qu and Zhuo (2017) and Kasahara and Shimotsu (2018) have studied the LR statistic for regime switching in Markov switching models. None of these papers allow for the regime switching probabilities to depend on the observed data.

Regarding testing for mixture type regime switching, the existing literature is extensive, and several early references can be found, for instance, in McLachlan and Peel (2000, Sec 6.5.1). Most papers in this literature consider the case of independent observations without regressors. Notable exceptions allowing for regressors (but not dependent data) and having set-ups closer to the present paper are Zhu and Zhang $(2004,2006)$ and Kasahara and Shimotsu (2015) who consider (among other things) LR tests for regime switching. Further comparison to these works will be provided in later sections.

In contrast to testing for Markov switching or mixture type regime switching, the literature on testing for time-inhomogeneous Markov switching and observation-dependent regime switching seems to be almost nonexistent. To our knowledge, no tests have been obtained for time-inhomogeneous Markov switching models although these models have attracted some interest (see, e.g., Diebold, Lee, and Weinbach (1994), Filardo (1994), and Kim, Piger, and Startz (2008)). As to testing for observationdependent regime switching, the only two exceptions we are aware of are the unpublished PhD thesis of Jeffries (1998) and the recent paper of Shen and He (2015). Jeffries's thesis analyses the LR test in a specific (first-order) mixture autoregressive model; we will discuss his work further in later sections. Shen and He (2015) consider the case of independent observations with regressors and observationdependent regime switching, and propose an 'expectation maximization test' for regime switching.

In this paper we consider testing for observation-dependent regime switching in a time series context. Specifically, we analyze the asymptotic distribution of the LR statistic for testing a linear autoregressive model against a two-regime mixture autoregressive model with observation-dependent regime switching. Following Kasahara and Shimotsu (2015) it would also be possible to consider the more general testing problem that in a model with more than two regimes the number of regimes can be re-
duced. However, as even the case of two regimes is quite complex in our set-up, we leave this extension to future research.

To provide economic motivation, mixture autoregressive models with observation-dependent regime switching have been applied to economic data for instance in Wong and Li (2001), Lanne and Saikkonen (2003), Dueker, Sola, and Spagnolo (2007), Bec, Rahbek, and Shephard (2008), Dueker, Psaradakis, Sola, and Spagnolo (2011), and Kalliovirta, Meitz, and Saikkonen (2015, 2016) in classical likelihood framework and Villani, Kohn, and Giordani (2009) in Bayesian framework; further discussion of this previous work will be provided in Section 2. A major motivation for allowing the regime switching probabilities to depend on observed data stems from the desire to associate changes in regime to observable economic variables. In the words of Hamilton's (2016) Handbook of Macroeconomics chapter:

Many economic time series exhibit dramatic breaks associated with events such as economic recessions, financial panics, and currency crises. Such changes in regime may arise from tipping points or other nonlinear dynamics and are core to some of the most important questions in macroeconomics.

Associating changes in regime to observable economic variables has recently been discussed also by Chang, Choi, and Park (2017, pp. 127-128) who find it problematic that the regime switching mechanism of the conventional Markov switching model is completely independent of all other parts of the model and, in particular, of the realizations of the underlying time series. According to these authors,
[Markov switching models] assume that the Markov chain determining regimes is completely independent from all other parts of the model, which is extremely unrealistic in many cases. Note that exogenous regime switching in particular implies that future transitions between states are completely determined by the current state, and does not rely on the realizations of underlying time series. This is highly unlikely in many practical applications.

Time-inhomogeneous Markov switching models and their special cases, observation-dependent regime switching models, respond to this issue. Of these two alternatives the former are rather difficult to handle both theoretically and empirically. Observation-dependent regime switching models are much more feasible although rather challenging in our testing problem.

We consider mixture autoregressive (MAR) models in a rather general setting employing high-level conditions that allow for various forms of observation-dependent regime switching. As specific examples, we treat the so-called logistic MAR (LMAR) model of Wong and Li (2001) and (a version of the) Gaussian MAR (GMAR) model of Kalliovirta et al. (2015) in detail. The technical difficulties we face in analyzing the LR statistic are quite challenging. First, there are nuisance parameters that are unidentified under the null hypothesis. This is the classical Davies $(1977,1987)$ type problem. Second, under the null hypothesis, there are parameters on the boundary of the permissible parameter space. Such problems (also allowing for unidentified nuisance parameters under the null) are discussed in Andrews (1999, 2001). Third, the Fisher information matrix is (potentially) singular, preventing the use of conventional second-order expansions of the log-likelihood function to analyze the LR test statistic. Such problems are discussed by Rotnitzky, Cox, Bottai, and Robins (2000), and suitable reparameterizations and higher-order expansions are needed to analyze the LR statistic. A particular challenge in the present paper is to deal with these three problems simultaneously. Similar problems were also faced by Kasahara and Shimotsu $(2015,2018)$ when testing for mixture and Markov switching type regime switching. Inspired by the work of these authors we consider a suitably reparameterized model, write a higher-order expansion of the log-likelihood function as a quadratic function of the new parameters, and then derive the asymptotic distribution of the LR test statistic by slightly extending and adapting the arguments of Andrews $(1999,2001)$ and Zhu and Zhang $(2006)$ (who partially generalize results of Andrews). Our two examples demonstrate that, compared to the mixture type regime
switching considered by Kasahara and Shimotsu (2015), observation-dependent regime switching can either simplify or complicate the analysis of the LR statistic.

We contribute to the literature in several ways. (1) To the best of our knowledge, apart from the unpublished PhD thesis of Jeffries (1998), we are the first to study testing for observation-dependent regime switching using the LR statistic and among the rather few to allow for dependent observations. (2) We provide a general framework to cover various forms of observation-dependent regime switching, making our results potentially applicable to several models not explicitly discussed in the present paper. (3) From a methodological perspective, we slightly extend and adapt certain arguments of Andrews (1999, 2001) and Zhu and Zhang (2006), which may be of independent interest.

The rest of the paper is organized as follows. Section 2 reviews mixture autoregressive models. Section 3 analyzes the LR test statistic for testing a linear autoregressive model against a two-regime mixture autoregressive model. Simulating the asymptotic null distribution, parametric bootstrap, and a Monte Carlo study are discussed in Section 4, and Section 5 concludes. Appendices A-C contain technical details and proofs. Supplementary Appendices D-E, available from the authors upon request, contain further technical details omitted from the paper.

Finally, a few notational conventions are given. All vectors will be treated as column vectors and, for the sake of uncluttered notation, we shall write $x=\left(x_{1}, \ldots, x_{n}\right)$ for the (column) vector $x$ where the components $x_{i}$ may be either scalars or vectors (or both). For any vector or matrix $x$, the Euclidean norm is denoted by $\|x\|$. We let $X_{T \alpha}=o_{p \alpha}(1)$ and $X_{T \alpha}=O_{p \alpha}(1)$ stand for $\sup _{\alpha \in A}\left\|X_{T \alpha}\right\|=o_{p}(1)$ and $\sup _{\alpha \in A}\left\|X_{T \alpha}\right\|=O_{p}(1)$, respectively, and $\lambda_{\min }(\cdot)$ and $\lambda_{\max }(\cdot)$ for the smallest and largest eigenvalue of the indicated matrix.

## 2 Mixture autoregressive models

### 2.1 General formulation

Let $y_{t}(t=1,2, \ldots)$ be the real-valued time series of interest, $\mathcal{F}_{t-1}=\sigma\left(y_{s}, s<t\right)$ denote the $\sigma$-algebra generated by past $y_{t}$ 's, and $P_{t-1}(\cdot)$ signify the conditional probability of the indicated event given $\mathcal{F}_{t-1}$. In the general two component mixture autoregressive model we consider the $y_{t}$ 's are generated by

$$
\begin{equation*}
y_{t}=s_{t}\left(\tilde{\phi}_{0}+\sum_{i=1}^{p} \tilde{\phi}_{i} y_{t-i}+\tilde{\sigma}_{1} \varepsilon_{t}\right)+\left(1-s_{t}\right)\left(\tilde{\varphi}_{0}+\sum_{i=1}^{p} \tilde{\varphi}_{i} y_{t-i}+\tilde{\sigma}_{2} \varepsilon_{t}\right), \tag{1}
\end{equation*}
$$

where the parameters $\tilde{\sigma}_{1}$ and $\tilde{\sigma}_{2}$ are positive, and conditions required for the autoregressive parameters $\tilde{\phi}_{i}$ and $\tilde{\varphi}_{i}(i=1, \ldots, p)$ will be discussed later. Furthermore, $\varepsilon_{t}$ and $s_{t}$ are (unobserved) stochastic processes which satisfy the following conditions: (a) $\varepsilon_{t}$ is a sequence of independent standard normal random variables such that $\varepsilon_{t}$ is independent of $\left\{y_{t-j}, j>0\right\}$, (b) $s_{t}$ is a sequence of Bernoulli random variables such that, for each $t, P_{t-1}\left(s_{t}=1\right)=\alpha_{t}$ with $\alpha_{t}$ a function of $\boldsymbol{y}_{t-1}=\left(y_{t-1}, \ldots, y_{t-p}\right)$, and (c) conditional on $\mathcal{F}_{t-1}$, $s_{t}$ and $\varepsilon_{t}$ are independent. Extending model (1) to allow for multiple regimes is straightforward but not needed in the present paper; see Kalliovirta et al. (2015, Sec 2.1).

The conditional probabilities $\alpha_{t}$ and $1-\alpha_{t}\left(=P_{t-1}\left(s_{t}=0\right)\right)$ are referred to as mixing weights. They can be thought of as (conditional) probabilities that determine which one of the two autoregressive components of the mixture generates the next observation $y_{t}$. In condition (b) we assume the mixing weight $\alpha_{t}$ (and hence also the conditional distribution of $y_{t}$ given its past) to depend on only $p$ lags of $y_{t}$; allowing it to depend on more than $p$ lags is easy but would unnecessarily complicate the notation.

We assume that of the original parameters $\tilde{\phi}=\left(\tilde{\phi}_{0}, \tilde{\phi}_{1}, \ldots, \tilde{\phi}_{p}, \tilde{\sigma}_{1}^{2}\right)$ and $\tilde{\varphi}=\left(\tilde{\varphi}_{0}, \tilde{\varphi}_{1}, \ldots, \tilde{\varphi}_{p}, \tilde{\sigma}_{2}^{2}\right)$, $q_{1}$ parameters are a priori assumed the same in both regimes and the remaining $q_{2}$ parameters are potentially different in the two regimes (with $q_{1}+q_{2}=p+2$ ). For instance, one may assume that $\tilde{\phi}_{0}$
and $\tilde{\varphi}_{0}$ are equal, or alternatively that $\tilde{\sigma}_{1}^{2}$ and $\tilde{\sigma}_{2}^{2}$ are equal. To this end, let $\beta$ be a $q_{1} \times 1$ vector of common parameters, and let $\phi$ and $\varphi$ be $q_{2} \times 1$ vectors of (potentially) different parameters. Then, for some known $(p+2)$-dimensional permutation matrix $P,(\beta, \phi)=P \tilde{\phi}$ and $(\beta, \varphi)=P \tilde{\varphi}$. For simplicity, we assume that $\beta$ and $\phi$ are variation-free, requiring the autoregressive coefficients $\tilde{\phi}_{1}, \ldots, \tilde{\phi}_{p}$ to be contained in either $\beta$ or $\phi$ (the same variation-freeness is assumed of $\beta$ and $\varphi$ ). If there are no common coefficients in the two regimes, the parameter $\beta$ can be dropped and $\phi=\tilde{\phi}$ and $\varphi=\tilde{\varphi}$.

As for the mixing weight $\alpha_{t}$, in addition to past $y_{t}$ 's it depends on unknown parameters which may include components of the parameter vector $(\beta, \phi, \varphi)$ and an additional parameter $\alpha$ (scalar or vector). When this dependence needs to be emphasized we use the notation $\alpha_{t}(\alpha, \beta, \phi, \varphi)$.

Using equation (1) and the conditions following it, the conditional density function of $y_{t}$ given its past, $f\left(\cdot \mid \mathcal{F}_{t-1}\right)$, is obtained as

$$
\begin{equation*}
f\left(y_{t} \mid \mathcal{F}_{t-1}\right)=\alpha_{t} f_{t}(\beta, \phi)+\left(1-\alpha_{t}\right) f_{t}(\beta, \varphi) \tag{2}
\end{equation*}
$$

where the notation $f_{t}(\beta, \phi)$ signifies the density function of a (univariate) normal distribution with mean $\tilde{\phi}_{0}+\tilde{\phi}_{1} y_{t-1}+\cdots+\tilde{\phi}_{p} y_{t-p}$ and variance $\tilde{\sigma}_{1}^{2}$ evaluated at $y_{t}$, that is,

$$
\begin{equation*}
f_{t}(\beta, \phi)=\frac{1}{\tilde{\sigma}_{1}} \mathfrak{N}\left(\frac{y_{t}-\left(\tilde{\phi}_{0}+\tilde{\phi}_{1} y_{t-1}+\cdots+\tilde{\phi}_{p} y_{t-p}\right)}{\tilde{\sigma}_{1}}\right) \tag{3}
\end{equation*}
$$

with $\mathfrak{N}(u)=(2 \bar{\pi})^{-1 / 2} \exp \left(-u^{2} / 2\right)$ the density function of a standard normal random variable and $\bar{\pi}=3.14 \ldots$ the number pi. The notation $f_{t}(\beta, \varphi)$ is defined similarly by using the parameters $\tilde{\varphi}_{i}$ $(i=0, \ldots, p)$ and $\tilde{\sigma}_{2}^{2}$ instead of $\tilde{\phi}_{i}(i=0, \ldots, p)$ and $\tilde{\sigma}_{1}^{2}$. Thus, the distribution of $y_{t}$ given its past is specified as a mixture of two normal densities with time varying mixing weights $\alpha_{t}$ and $1-\alpha_{t}$.

To facilitate discussion, denote $\mu_{1, t}=\tilde{\phi}_{0}+\sum_{i=1}^{p} \tilde{\phi}_{i} y_{t-i}, \mu_{2, t}=\tilde{\varphi}_{0}+\sum_{i=1}^{p} \tilde{\varphi}_{i} y_{t-i}$, and $\mu_{t}=\alpha_{t} \mu_{1, t}+$ $\left(1-\alpha_{t}\right) \mu_{2, t}$. Now equation (1) can be written as

$$
\begin{equation*}
y_{t}=s_{t} \mu_{1, t}+\left(1-s_{t}\right) \mu_{2, t}+u_{t} \quad \text { with } \quad u_{t}=\left(s_{t} \tilde{\sigma}_{1}+\left(1-s_{t}\right) \tilde{\sigma}_{2}\right) \varepsilon_{t} \tag{4}
\end{equation*}
$$

and the conditional mean and variance of $y_{t}$ are given by (cf. Kalliovirta et al. (2015, p. 249))

$$
\begin{equation*}
E\left[y_{t} \mid \mathcal{F}_{t-1}\right]=\mu_{t} \quad \text { and } \quad \operatorname{Var}\left[y_{t} \mid \mathcal{F}_{t-1}\right]=\alpha_{t} \tilde{\sigma}_{1}^{2}+\left(1-\alpha_{t}\right) \tilde{\sigma}_{2}^{2}+\alpha_{t}\left(\mu_{1, t}-\mu_{t}\right)^{2}+\left(1-\alpha_{t}\right)\left(\mu_{2, t}-\mu_{t}\right)^{2} \tag{5}
\end{equation*}
$$

The expression of $u_{t}$ in (4) suggests that the normality assumption imposed on $\varepsilon_{t}$ in (1) is not quite as restrictive as it may seem at first sight: $u_{t}$ in (4) may be interpreted as an error term and is a martingale difference sequence with the conditional distribution being a mixture of two normal distributions. Such mixtures are fairly flexible and can allow for various departures from normality including skewness, excess kurtosis, and bimodality. One could consider non-normal distributions for $\varepsilon_{t}$ but this would significantly alter the derivations of the results that follow. The expressions in (5) facilitate comparison to smooth transition autoregressive (STAR) models (with two components and independent and identically distributed errors; see, e.g., van Dijk, Teräsvirta, and Franses (2002)). As $\alpha_{t}$ is a function of $p$ past $y_{t}$ 's, the conditional mean in (5) has a form similar to that in STAR models; in contrast, the conditional variance of model (1) can allow for some type of conditional heteroskedasticity whereas the conditional variance of the standard STAR model is constant.

Different mixture autoregressive models are obtained by different specifications of the mixing weights (or in our case the single mixing weight $\alpha_{t}$ ). In some of the proposed models more than two mixture components are allowed but for reasons to be discussed below we will not consider these extensions. If the mixing weights are assumed constant over time the general mixture autoregressive model introduced above reduces to a (two-component) MAR model studied by Wong and Li (2000).

In the LMAR model of Wong and Li (2001), a logistic transformation of the two mixing weights is assumed to be a linear function of past observed variables. Related two-regime mixture models with time-varying mixing weights have also been considered by Gouriéroux and Robert (2006), Dueker et al. (2007) and Bec et al. (2008) whereas Lanne and Saikkonen (2003), Villani et al. (2009), and Kalliovirta et al. (2015) have considered mixture autoregressive models in which multiple regimes are allowed. For a particular specification of the mixing weights and for a known number of regimes, standard likelihood based inference, information criteria, and quantile residual based diagnostics can be used for model selection, in particular for choosing the lag order $p$; for details, see Kalliovirta et al. (2015, Sec 3).

As discussed in the Introduction, a common and highly irregular problem with the application of mixture autoregressive models is determining the value of the (usually unknown) number of component models or regimes. To our knowledge the only solution presented for mixture autoregressive models is provided for a simple first order case with no intercept terms in the unpublished PhD thesis of Jeffries (1998). As discussed in the recent papers by Kasahara and Shimotsu (2012, 2015) and the references therein, some of the difficulties involved stem from properties of the normal distribution, and they are encountered even when the observations are a random sample from a mixture of (two) normal distributions. These difficulties also partly explain the complexity of our testing problem and why we only consider tests for the null hypothesis that a two component mixture autoregressive model reduces to a conventional linear autoregressive model. Following the ideas in Zhu and Zhang (2006) and Kasahara and Shimotsu (2012, 2015), we derive a LR test in the general set-up described above and apply it to two particular cases, the LMAR model of Wong and Li (2001) and the GMAR model of Kalliovirta et al. (2015). Next, we shall discuss these two models in more detail.

### 2.2 Two particular examples

LMAR Example. In Wong and Li's (2001) LMAR model the mixing weight $\alpha_{t}$ is defined as

$$
\alpha_{t}^{L}=\alpha_{t}^{L}(\alpha)=\frac{\exp \left(\alpha_{0}+\alpha_{1} y_{t-1}+\cdots+\alpha_{m} y_{t-m}\right)}{1+\exp \left(\alpha_{0}+\alpha_{1} y_{t-1}+\cdots+\alpha_{m} y_{t-m}\right)},
$$

where the vector $\alpha=\left(\alpha_{0}, \alpha_{1}, \ldots, \alpha_{m}\right)$ contains $m+1$ unknown parameters and the order $m(1 \leq m \leq p)$ is assumed known. Note that for this model, $E\left[y_{t} \mid \mathcal{F}_{t-1}\right]$ agrees with that of the logistic STAR model.

GMAR Example. In the GMAR model of Kalliovirta et al. (2015) the mixing weight is defined as

$$
\begin{equation*}
\alpha_{t}^{G}=\alpha_{t}^{G}(\alpha, \tilde{\phi}, \tilde{\varphi})=\frac{\alpha \mathrm{n}_{p}\left(\boldsymbol{y}_{t-1} ; \tilde{\phi}\right)}{\alpha \mathrm{n}_{p}\left(\boldsymbol{y}_{t-1} ; \tilde{\phi}\right)+(1-\alpha) \mathrm{n}_{p}\left(\boldsymbol{y}_{t-1} ; \tilde{\varphi}\right)}, \tag{6}
\end{equation*}
$$

where $\alpha \in(0,1)$ is an unknown parameter and $\mathrm{n}_{p}\left(\boldsymbol{y}_{t-1} ; \cdot\right)$ denotes the density function of a particular $p$-dimensional ( $p \geq 1$ ) normal distribution defined as follows. First, define the auxiliary Gaussian $\operatorname{AR}(p)$ processes (cf. equation (1))

$$
\nu_{1, t}=\tilde{\phi}_{0}+\sum_{i=1}^{p} \tilde{\phi}_{i} \nu_{1, t-i}+\tilde{\sigma}_{1} \varepsilon_{t} \quad \text { and } \quad \nu_{2, t}=\tilde{\varphi}_{0}+\sum_{i=1}^{p} \tilde{\varphi}_{i} \nu_{2, t-i}+\tilde{\sigma}_{2} \varepsilon_{t},
$$

where the autoregressive coefficients are assumed to satisfy

$$
\begin{equation*}
\tilde{\phi}(z):=1-\sum_{i=1}^{p} \tilde{\phi}_{i} z^{i} \neq 0 \text { for }|z| \leq 1 \quad \text { and } \quad \tilde{\varphi}(z):=1-\sum_{i=1}^{p} \tilde{\varphi}_{i} z^{i} \neq 0 \text { for }|z| \leq 1 \text {. } \tag{7}
\end{equation*}
$$

This condition implies that the processes $\nu_{1, t}$ and $\nu_{2, t}$ are stationary and that each of the two component models in (1) satisfies the usual stationarity condition of the conventional linear $\operatorname{AR}(p)$ model. Now
set $\boldsymbol{\nu}_{m, t}=\left(\nu_{m, t}, \ldots, \nu_{m, t-p+1}\right)$ and $\mathbf{1}_{p}=(1, \ldots, 1)(p \times 1)$, and let $\mu_{m} \mathbf{1}_{p}$ and $\boldsymbol{\Gamma}_{m, p}$ signify the mean vector and covariance matrix of $\boldsymbol{\nu}_{m, t}(m=1,2) .{ }^{1}$ The random vector $\boldsymbol{\nu}_{1, t}$ follows the $p$-dimensional multivariate normal distribution with density

$$
\begin{equation*}
\mathbf{n}_{p}\left(\boldsymbol{\nu}_{1, t} ; \tilde{\phi}\right)=(2 \bar{\pi})^{-p / 2} \operatorname{det}\left(\boldsymbol{\Gamma}_{1, p}\right)^{-1 / 2} \exp \left\{-\frac{1}{2}\left(\boldsymbol{\nu}_{1, t}-\mu_{1} \mathbf{1}_{p}\right)^{\prime} \boldsymbol{\Gamma}_{1, p}^{-1}\left(\boldsymbol{\nu}_{1, t}-\mu_{1} \mathbf{1}_{p}\right)\right\}, \tag{8}
\end{equation*}
$$

and the density of $\boldsymbol{\nu}_{2, t}$, denoted by $\mathrm{n}_{p}\left(\boldsymbol{\nu}_{2, t} ; \tilde{\varphi}\right)$, is defined similarly. Equation (1) and conditions (6)-(8) define the (two component) GMAR model; condition (7) is part of the definition of the model because it is used to define the mixing weights. ${ }^{2}$

## 3 Test procedure

We now consider a test procedure of the null hypothesis that a two component mixture autoregressive model reduces to a conventional linear autoregressive model.

### 3.1 The null hypothesis and the LR test statistic

We denote the conditional density function corresponding to the unrestricted model as (see (2))

$$
f_{2, t}(\alpha, \beta, \phi, \varphi):=f_{2}\left(y_{t} \mid \boldsymbol{y}_{t-1} ; \alpha, \beta, \phi, \varphi\right):=\alpha_{t}(\alpha, \beta, \phi, \varphi) f_{t}(\beta, \phi)+\left(1-\alpha_{t}(\alpha, \beta, \phi, \varphi)\right) f_{t}(\beta, \varphi),
$$

where we now make the dependence of the mixing weight on the parameters explicit. With this notation the $\log$-likelihood function of the model based on a sample $\left(y_{-p+1}, \ldots, y_{T}\right)$ (and conditional on the initial values $\left.\left(y_{-p+1}, \ldots, y_{0}\right)\right)$ is $L_{T}(\alpha, \beta, \phi, \varphi)=\sum_{t=1}^{T} l_{t}(\alpha, \beta, \phi, \varphi)$ where

$$
l_{t}(\alpha, \beta, \phi, \varphi)=\log \left[f_{2, t}(\alpha, \beta, \phi, \varphi)\right]=\log \left[\alpha_{t}(\alpha, \beta, \phi, \varphi) f_{t}(\beta, \phi)+\left(1-\alpha_{t}(\alpha, \beta, \phi, \varphi)\right) f_{t}(\beta, \varphi)\right] .
$$

The following assumption provides conditions on the data generation process, the parameter space of $(\alpha, \beta, \phi, \varphi)$, and the mixing weight $\alpha_{t}(\alpha, \beta, \phi, \varphi)$.

## Assumption 1.

(i) The $y_{t}$ 's are generated by a stationary linear Gaussian $A R(p)$ model with (the true but unknown) parameter value $\tilde{\phi}^{*}$ an interior point of $\tilde{\Phi}$, a compact subset of $\left\{\tilde{\phi}=\left(\tilde{\phi}_{0}, \tilde{\phi}_{1}, \ldots, \tilde{\phi}_{p}, \tilde{\sigma}^{2}\right) \in \mathbb{R}^{p+2}\right.$ : $\tilde{\phi}_{0} \in \mathbb{R} ; 1-\sum_{i=1}^{p} \tilde{\phi}_{i} z^{i} \neq 0$ for $\left.|z| \leq 1 ; \tilde{\sigma}^{2} \in(0, \infty)\right\}$.
(ii) The parameter space of $(\alpha, \beta, \phi, \varphi)$ is $A \times B \times \Phi \times \Phi$, where $A$ is a compact subset of $\mathbb{R}^{a}$ and $B$ and $\Phi$ are those compact subsets of $\mathbb{R}^{q_{1}}$ and $\mathbb{R}^{q_{2}}$, respectively, that satisfy $(\beta, \phi) \in B \times \Phi$ if and only if $P^{-1}(\beta, \phi) \in \tilde{\Phi}$ (here $P$ is as in the third paragraph of Section 2.1).
(iii) For all t and all $(\alpha, \beta, \phi, \varphi) \in A \times B \times \Phi \times \Phi$, the mixing weight $\alpha_{t}(\alpha, \beta, \phi, \varphi)$, is $\sigma\left(\boldsymbol{y}_{t-1}\right)$-measurable (with $\sigma\left(\boldsymbol{y}_{t-1}\right)$ denoting the $\sigma$-algebra generated by $\boldsymbol{y}_{t-1}$ ) and satisfies $\alpha_{t}(\alpha, \beta, \phi, \varphi) \in(0,1)$.

Corresponding to the null hypothesis, Assumption 1(i) requires the data to be generated by a stationary linear Gaussian $\operatorname{AR}(p)$ model. Assuming a compact parameter space in Assumptions 1(i)

[^1]and (ii) is a standard requirement which facilitates proofs. Assumption 1(ii) accommodates to the main cases of interest, namely $\beta=\tilde{\phi}_{0}, \beta=\left(\tilde{\phi}_{1}, \ldots, \tilde{\phi}_{p}\right), \beta=\tilde{\sigma}^{2}$, or any combination of these. ${ }^{3}$

Assumption 1(iii) implies that our two-component mixture autoregressive model reduces to a linear autoregression only when $\phi=\varphi$, regardless of the values of $\alpha \in A$ and $\beta \in B$. The null hypothesis to be tested is therefore $\phi=\varphi$ and the alternative is $\phi \neq \varphi$ or, more precisely,

$$
H_{0}:(\phi, \varphi) \in(\Phi \times \Phi)^{0}, \alpha \in A, \beta \in B \quad \text { vs. } \quad H_{1}:(\phi, \varphi) \in(\Phi \times \Phi) \backslash(\Phi \times \Phi)^{0}, \alpha \in A, \beta \in B
$$

where (here and in what follows, the superscript 0 refers to the model restricted by the null hypothesis)

$$
(\Phi \times \Phi)^{0}=\{(\phi, \varphi) \in \Phi \times \Phi: \phi=\varphi\}
$$

Note that under the null hypothesis the parameter $\alpha$ vanishes from the likelihood function and is therefore unidentified. Specifically, letting $f_{t}^{0}(\tilde{\phi}):=f^{0}\left(y_{t} \mid \boldsymbol{y}_{t-1} ; \tilde{\phi}\right)$ and $L_{T}^{0}(\tilde{\phi})$ denote the conditional density and log-likelihood functions of the restricted model, we have

$$
f_{t}^{0}(\tilde{\phi})=f_{2, t}(\alpha, \beta, \phi, \phi)=f_{t}(\tilde{\phi}) \quad \text { and } \quad L_{T}^{0}(\tilde{\phi})=\sum_{t=1}^{T} l_{t}^{0}(\tilde{\phi}) \quad \text { with } \quad l_{t}^{0}(\tilde{\phi})=\log \left[f_{t}(\tilde{\phi})\right]
$$

as in a linear Gaussian $\operatorname{AR}(p)$ model. As $f_{2, t}\left(\alpha, \beta^{*}, \phi^{*}, \phi^{*}\right)=f_{t}\left(\tilde{\phi}^{*}\right)$ for any $\alpha \in A$, in the unrestricted model the parameter vector $\left(\alpha, \beta^{*}, \phi^{*}, \phi^{*}\right)$ corresponds to the true model for any $\alpha \in A$.

As already indicated, Assumption 1(iii) implies that the restriction $\phi=\varphi$ is the only possibility to formulate the null hypothesis. However, this is not necessarily the case if (against Assumption 1(iii)) the mixing weight $\alpha_{t}(\alpha, \beta, \phi, \varphi)$ were allowed to take the boundary values zero and one. Of our two examples this would be possible for the GMAR (but not for the LMAR) model in which $\alpha_{t}(\alpha, \beta, \phi, \varphi)$ takes the boundary values zero and one when the parameter $\alpha$ takes these values (see (6)). In both cases a linear autoregression results and either the parameter $\phi$ or $\varphi$ is unidentified (see (2)). (The MAR model of Wong and Li (2000) provides a similar example.) It would be possible to obtain tests for the GMAR model by using a null hypothesis which specifies $\alpha=0$ or $\alpha=1$. However, as in Kasahara and Shimotsu $(2012,2015)^{4}$, this approach would require rather restrictive assumptions and would also lead to very complicated derivations. Therefore, we will not consider this option.

As the parameter $\alpha$ is unidentified under the null hypothesis, the appropriate likelihood ratio type test statistic is

$$
L R_{T}=\sup _{\alpha \in A} L R_{T}(\alpha)
$$

where, for each fixed $\alpha \in A$,

$$
L R_{T}(\alpha)=2\left[\sup _{(\beta, \phi, \varphi) \in B \times \Phi \times \Phi} L_{T}(\alpha, \beta, \phi, \varphi)-\sup _{\tilde{\phi} \in \tilde{\Phi}} L_{T}^{0}(\tilde{\phi})\right]
$$

Now let, for each fixed $\alpha \in A,\left(\hat{\beta}_{T \alpha}, \hat{\phi}_{T \alpha}, \hat{\varphi}_{T \alpha}\right)$ denote an (approximate) unrestricted maximum likelihood (ML) estimator of the parameter vector $(\beta, \phi, \varphi)$. We make the following assumption.

Assumption 2. The unrestricted ML estimator satisfies the following conditions:
(i) $L_{T}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\phi}_{T \alpha}, \hat{\varphi}_{T \alpha}\right)=\sup _{(\beta, \phi, \varphi) \in B \times \Phi \times \Phi} L_{T}(\alpha, \beta, \phi, \varphi)+o_{p \alpha}(1)$,
(ii) $\left(\hat{\beta}_{T \alpha}, \hat{\phi}_{T \alpha}, \hat{\varphi}_{T \alpha}\right)=\left(\beta^{*}, \phi^{*}, \phi^{*}\right)+o_{p \alpha}(1)$.

Assumption 2(i) means that $\left(\hat{\beta}_{T \alpha}, \hat{\phi}_{T \alpha}, \hat{\varphi}_{T \alpha}\right)$ is assumed to maximize the likelihood function only asymptotically. This assumption is technical and made for ease of exposition; see Andrews (1999) and

[^2]Zhu and Zhang (2006) for similar assumptions. Assumption 2(ii) is a high level condition on (uniform) consistency of the ML estimator and is analogous to Assumption 1 of Andrews (2001). It has to be verified on a case by case basis (typically relying on a uniform law of large numbers and a suitable identification condition); this is exemplified below for the LMAR and GMAR models.

As for the term $\sup _{\tilde{\phi} \in \tilde{\Phi}} L_{T}^{0}(\tilde{\phi})$ in the LR test statistic, note that $L_{T}^{0}(\tilde{\phi})$ is the (conditional) loglikelihood function of a linear Gaussian $\operatorname{AR}(p)$ model. Let $\hat{\tilde{\phi}}_{T}$ denote an (approximate) maximum likelihood estimator of the parameters of a linear $\operatorname{Gaussian} \operatorname{AR}(p)$ model, that is, $\hat{\tilde{\phi}}_{T}$ satisfies $^{5}$

$$
L_{T}^{0}\left(\hat{\tilde{\phi}}_{T}\right)=\sup _{\tilde{\phi} \in \tilde{\Phi}} L_{T}^{0}(\tilde{\phi})+o_{p}(1) \text { and } \hat{\tilde{\phi}}_{T}=\tilde{\phi}^{*}+o_{p}(1)
$$

Noting that $L_{T}\left(\alpha, \beta^{*}, \phi^{*}, \phi^{*}\right)=L_{T}^{0}\left(\tilde{\phi}^{*}\right)$ for any $\alpha$ now allows us to write $L R_{T}(\alpha)$ as

$$
\begin{equation*}
L R_{T}(\alpha)=2\left[L_{T}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\phi}_{T \alpha}, \hat{\varphi}_{T \alpha}\right)-L_{T}\left(\alpha, \beta^{*}, \phi^{*}, \phi^{*}\right)\right]-2\left[L_{T}^{0}\left(\hat{\tilde{\phi}}_{T}\right)-L_{T}^{0}\left(\tilde{\phi}^{*}\right)\right]+o_{p \alpha}(1) \tag{9}
\end{equation*}
$$

The analysis of the second term on the right hand side is standard while dealing with the first term is more demanding requiring a substantial amount of preparation.

### 3.1.1 Examples (continued)

LMAR Example. In the LMAR example, we assume there are no common parameters in the two regimes so that the parameter $\beta$ is omitted, $\phi=\tilde{\phi}, \varphi=\tilde{\varphi}, q_{1}=0$, and $q_{2}=p+2$. To satisfy conditions (ii) and (iii) in Assumption $1, A$ can be any compact subset of $\left\{\left(\alpha_{0}, \alpha_{1}, \ldots, \alpha_{m}\right) \in \mathbb{R}^{m+1}\right.$ : $\left.\left(\alpha_{1}, \ldots, \alpha_{m}\right) \neq(0, \ldots, 0)\right\}$ where $1 \leq m \leq p$. This ensures that the mixing weight $\alpha_{t}^{L}$ is not equal to a constant. For the verification of Assumption 2, see Appendix B.

Note that we only consider the 'real' LMAR case where $m \geq 1$ and $\left(\alpha_{1}, \ldots, \alpha_{m}\right) \neq(0, \ldots, 0)$. When $m=0$ or $\left(\alpha_{1}, \ldots, \alpha_{m}\right)=(0, \ldots, 0)$, the mixing weight $\alpha_{t}^{L}$ is constant and the model reduces to the MAR model of Wong and $\mathrm{Li}(2000)$ and then our testing problem requires different and more complicated analyses (we shall discuss this point more later). A similar restriction is made by Jeffries (1998) in his (first-order) logistic mixture autoregressive model to facilitate the derivation of the LR test (see the hypotheses at the end of p. 95 and the following discussion, as well as the end of p. 110).

GMAR Example. The GMAR model exemplifies the setting with common coefficients by assuming that the intercept terms in the two regimes are the same (note that this still allows for different means in the two regimes). As will be discussed in more detail in Section 3.3.1, this assumption is partly due to the fact that otherwise the derivation of the LR test would become extremely complicated. Hence, in this example $\beta=\tilde{\phi}_{0}\left(=\tilde{\varphi}_{0}\right), \phi=\left(\tilde{\phi}_{1}, \ldots, \tilde{\phi}_{p}, \tilde{\sigma}_{1}^{2}\right), \varphi=\left(\tilde{\varphi}_{1}, \ldots, \tilde{\varphi}_{p}, \tilde{\sigma}_{2}^{2}\right), q_{1}=1$, and $q_{2}=p+1$. To satisfy Assumptions 1(ii) and (iii) the parameter space $A$ of $\alpha$ can be any compact and convex subset of $(0,1)$ with $[\eta, 1-\eta]$ for some $\eta>0$ being the typical choice (this also rules out the possibility that $\alpha=0$ or $\alpha=1$ discussed after Assumption 1). For the verification of Assumption 2, see Appendix C.

We note that there are cases when the mixing weight $\alpha_{t}^{G}$ is constant and the model reduces to the MAR model of Wong and $\mathrm{Li}(2000) .{ }^{6}$ Unlike in the LMAR example, this fact does not complicate the derivation of our test. The reason seems to be that in the GMAR model the reduction occurs only for certain values of $\tilde{\phi}$ and $\tilde{\varphi}$ whereas in the LMAR model it occurs for all values of $\tilde{\phi}$ and $\tilde{\varphi}$.

[^3]
### 3.2 Reparameterized model

In standard testing problems the derivation of the asymptotic distribution of the LR test relies on a quadratic expansion of the log-likelihood function $L_{T}(\alpha, \beta, \phi, \varphi)$; when the parameter $\alpha$ is not identified under the null hypothesis, the relevant derivatives in this expansion are with respect to $(\beta, \phi, \varphi)$ for fixed values of $\alpha \in A$. In problems with a singular information matrix it turns out to be convenient to follow Rotnitzky et al. (2000) and Kasahara and Shimotsu (2012, 2015) and employ an appropriately reparameterized model. The employed reparameterization is model specific and aims to have two conveniences. First, it transforms the null hypothesis $\phi=\varphi$ into a point null hypothesis where some components of the parameter vector are restricted to zero and the rest are left unrestricted. Second, and more importantly, it simplifies derivations in cases where the conventional quadratic expansion of the $\log$-likelihood function breaks down because, under the null hypothesis, the scores of the parameters $(\beta, \phi, \varphi)$ are linearly dependent and, consequently, the (Fisher) information matrix is singular. As will be seen later, this is the case for the GMAR model but not for the LMAR model.

General requirements for the reparameterization are described in the following assumption. Only the parameters restricted by the null hypothesis, $\phi$ and $\varphi$, are reparameterized. The examples in this and the following subsection illustrate how the reparameterization could be chosen.

## Assumption 3.

(i) For every $\alpha \in A$, the mapping $(\pi, \varpi)=\boldsymbol{\pi}_{\alpha}(\phi, \varphi)$ from $\Phi \times \Phi$ to $\Pi_{\alpha}$ is one-to-one with $\boldsymbol{\pi}_{\alpha}$ and $\pi_{\alpha}^{-1}$ continuous.
(ii) For every $\alpha \in A, \boldsymbol{\pi}_{\alpha}\left((\Phi \times \Phi)^{0}\right)=\Phi \times\{0\}$ and $\boldsymbol{\pi}_{\alpha}\left(\phi^{*}, \phi^{*}\right)=\left(\pi^{*}, 0\right):=\left(\phi^{*}, 0\right)$.
(iii) $\left(\hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)=\left(\beta^{*}, \pi^{*}, 0\right)+o_{p \alpha}(1)$, where $\left(\hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right):=\pi_{\alpha}\left(\hat{\phi}_{T \alpha}, \hat{\varphi}_{T \alpha}\right)$.

We sometimes refer to the reparameterization described in Assumption 3 as the ' $\pi$-parameterization' and the original parameterization as the ' $\phi$-parameterization'. Note that the transformed parameters $\pi$ and $\varpi$ generally depend on $\alpha$ but, for brevity, we suppress this dependence from the notation. The parameter space of $(\pi, \varpi)$ also depends on $\alpha$ and is given, for any $\alpha \in A$, by

$$
\Pi_{\alpha}=\left\{(\pi, \varpi) \in \mathbb{R}^{2 q_{2}}:(\pi, \varpi)=\pi_{\alpha}(\phi, \varphi) \text { for some }(\phi, \varphi) \in \Phi \times \Phi\right\} .
$$

Assumption 3(i) contains some minimal regularity conditions for the reparameterization. By Assumption 3(ii), the null hypothesis $\phi=\varphi$ can be equivalently written as $\varpi=0$ or, more precisely, as

$$
H_{0}: \pi \in \Phi, \varpi=0, \alpha \in A, \beta \in B \quad \text { vs. } \quad H_{1}:(\pi, \varpi) \in \Pi_{\alpha} \backslash(\Phi \times\{0\}), \alpha \in A, \beta \in B .
$$

Under $H_{0}$ parameters $\beta$ and $\pi$ are identified but $\alpha$ is not. Assumption 3(iii) is a high level condition similar to Assumption 2(ii) from which it can be derived with appropriate additional assumptions. A simple Lipschitz condition similar to Andrews (1992, Assumption SE-1(b)), given in Lemma A. 1 in Appendix A, is one possibility.

To develop further notation, partition $\boldsymbol{\pi}_{\alpha}^{-1}(\pi, \varpi)$ into two $q_{2}$-dimensional components as $\boldsymbol{\pi}_{\alpha}^{-1}(\pi, \varpi)=$ $\left(\boldsymbol{\pi}_{\alpha, 1}^{-1}(\pi, \varpi), \boldsymbol{\pi}_{\alpha, 2}^{-1}(\pi, \varpi)\right)$, and define

$$
f_{2, t}^{\pi}(\alpha, \beta, \pi, \varpi):=f_{2}\left(y_{t} \mid \boldsymbol{y}_{t-1} ; \alpha, \beta, \phi, \varphi\right)=\alpha_{t}^{\pi}(\alpha, \beta, \pi, \varpi) f_{t}\left(\beta, \boldsymbol{\pi}_{\alpha, 1}^{-1}(\pi, \varpi)\right)+\left(1-\alpha_{t}^{\pi}(\alpha, \beta, \pi, \varpi)\right) f_{t}\left(\beta, \boldsymbol{\pi}_{\alpha, 2}^{-1}(\pi, \varpi)\right),
$$

where $\alpha_{t}^{\pi}(\alpha, \beta, \pi, \varpi)=\alpha_{t}\left(\alpha, \beta, \boldsymbol{\pi}_{\alpha, 1}^{-1}(\pi, \varpi), \boldsymbol{\pi}_{\alpha, 2}^{-1}(\pi, \varpi)\right)$ and the function $f_{t}(\cdot)$ is as in (2). The loglikelihood function of the reparameterized model can now be expressed as

$$
\begin{equation*}
L_{T}^{\pi}(\alpha, \beta, \pi, \varpi)=\sum_{t=1}^{T} l_{t}^{\pi}(\alpha, \beta, \pi, \varpi), \tag{10}
\end{equation*}
$$

where $l_{t}^{\pi}(\alpha, \beta, \pi, \varpi)=\log \left[f_{2, t}^{\pi}(\alpha, \beta, \pi, \varpi)\right]$, and in the $\pi$-parameterization equation (9) reads as

$$
\begin{equation*}
L R_{T}(\alpha)=2\left[L_{T}^{\pi}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)-L_{T}^{\pi}\left(\alpha, \beta^{*}, \phi^{*}, 0\right)\right]-2\left[L_{T}^{0}\left(\hat{\tilde{\phi}}_{T}\right)-L_{T}^{0}\left(\tilde{\phi}^{*}\right)\right]+o_{p \alpha}(1) \tag{11}
\end{equation*}
$$

### 3.2.1 Examples (continued)

LMAR Example. The reparameterization we employ in the LMAR model is

$$
(\pi, \varpi)=\boldsymbol{\pi}_{\alpha}(\phi, \varphi)=(\phi, \phi-\varphi) \text { so that }(\phi, \varphi)=\boldsymbol{\pi}_{\alpha}^{-1}(\pi, \varpi)=(\pi, \pi-\varpi)
$$

In this case the reparameterization (via $\left.\boldsymbol{\pi}_{\alpha}(\cdot, \cdot)\right)$ does not depend on $\alpha$, and the same is true for the parameter space of $(\pi, \varpi)$. Verification of Assumption 3 is straightforward using Lemma A. 1 (for details, see Appendix B). In the LMAR case, the only benefit of the reparameterization is to transform the null hypothesis into a point null hypothesis.

GMAR Example. We reparameterize the GMAR model by setting, for any fixed $\alpha \in A$,

$$
(\pi, \varpi)=\boldsymbol{\pi}_{\alpha}(\phi, \varphi)=(\alpha \phi+(1-\alpha) \varphi, \phi-\varphi) \text { so that }(\phi, \varphi)=\boldsymbol{\pi}_{\alpha}^{-1}(\pi, \varpi)=(\pi+(1-\alpha) \varpi, \pi-\alpha \varpi) .
$$

Verifying Assumption 3 is again straightforward using Lemma A. 1 (see Appendix C). In the GMAR model simplifying the null hypothesis is not the only benefit of the reparameterization. As discussed before Assumption 3, the relevant derivatives when expanding $L_{T}(\alpha, \beta, \phi, \varphi)$ are with respect to $(\beta, \phi, \varphi)$ and in the GMAR case these derivatives are linearly dependent under the null hypothesis. To see this and how the reparameterization affects this feature, note that straightforward differentiation yields

$$
\nabla_{(\beta, \phi, \varphi)} l_{t}(\alpha, \beta, \phi, \phi)=\left(\frac{\nabla_{\beta} f_{t}(\beta, \phi)}{f_{t}(\beta, \phi)}, \alpha \frac{\nabla_{\phi} f_{t}(\beta, \phi)}{f_{t}(\beta, \phi)},(1-\alpha) \frac{\nabla_{\phi} f_{t}(\beta, \phi)}{f_{t}(\beta, \phi)}\right)
$$

where the null hypothesis $\phi=\varphi$ is imposed and $\nabla$ denotes differentiation with respect to the indicated parameters. As $\left(f_{t}(\beta, \phi)\right)^{-1} \nabla_{(\beta, \phi)} f_{t}(\beta, \phi)$ is the score vector obtained from a linear Gaussian $\operatorname{AR}(p)$ model, it is clear that the covariance matrix of the $(2 p+3)$-dimensional vector $\nabla_{(\beta, \phi, \varphi)} l_{t}(\alpha, \beta, \phi, \phi)$, and hence the (Fisher) information matrix, is singular with rank $p+2$. In contrast to the above, in the $\pi$-parameterization the score vector is given by (see Supplementary Appendix F)

$$
\nabla_{(\beta, \pi, \varpi)} l_{t}^{\pi}(\alpha, \beta, \pi, 0)=\left(\frac{\nabla_{\beta} f_{t}(\beta, \pi)}{f_{t}(\beta, \pi)}, \frac{\nabla_{\pi} f_{t}(\beta, \pi)}{f_{t}(\beta, \pi)}, \mathbf{0}_{p+1}\right)
$$

when the null hypothesis $\varpi=0$ is imposed. Now the score of $\varpi$ is identically zero so that the reparameterization simplifies linear dependencies of the scores which turns out to be very useful in subsequent asymptotic analyses.

### 3.3 Quadratic expansion of the (reparameterized) log-likelihood function

As alluded to above, in standard testing problems the asymptotic analysis of a LR test statistic is based on a second-order Taylor expansion of the (average) log-likelihood function around the true parameter value. An essential assumption here is positive definiteness of the (limiting) information matrix but, as illustrated above, this assumption does not necessarily hold in our testing problem due to linear dependencies among the derivatives of the log-likelihood function. As in Rotnitzky et al. (2000), Zhu and Zhang (2006), and Kasahara and Shimotsu (2012, 2015), we thus consider a quadratic expansion of the log-likelihood function that is not based on a second-order Taylor expansion but (possibly) on
a higher-order Taylor expansion. The need for higher-order derivatives is illustrated by the GMAR example: as the score of $\varpi$ is identically zero, the second derivative (which turns out to be linearly independent of the score of $(\beta, \pi)$ ) now provides the first (nontrivial) local approximation for $\varpi$.

The following assumption ensures that the (reparameterized) log-likelihood function (10) is (at least) twice continuously differentiable.

Assumption 4. For some integer $k \geq 2$, and for every fixed $\alpha \in A$, the functions $\alpha_{t}(\alpha, \beta, \phi, \varphi)$ and $\boldsymbol{\pi}_{\alpha}^{-1}(\pi, \varpi)$ are $k$ times continuously differentiable (with respect to $(\beta, \phi, \varphi)$ and $(\pi, \varpi)$ in the interior of $B \times \Phi \times \Phi$ and $\Pi_{\alpha}$, respectively).

We assume the reparameterized log-likelihood function to have, for each $\alpha \in A$, a quadratic expansion in a transformed parameter vector $\boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)$ around $\left(\beta^{*}, \pi^{*}, 0\right)$ given by

$$
\begin{align*}
& L_{T}^{\pi}(\alpha, \beta, \pi, \varpi)-L_{T}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right) \\
& \quad=\left(T^{-1 / 2} S_{T \alpha}\right)^{\prime}\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)\right]-\frac{1}{2}\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)\right]^{\prime} \mathcal{I}_{\alpha}\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)\right]+R_{T}(\alpha, \beta, \pi, \varpi) \tag{12}
\end{align*}
$$

To illustrate this expansion, suppose the information matrix is positive definite so that the quantities on the right hand side are (typically) based on a second-order Taylor expansion with $S_{T \alpha}$ and $\mathcal{I}_{\alpha}$ functions of $\left(\alpha, \beta^{*}, \pi^{*}, 0\right)$. As already mentioned, this is the case for the LMAR model where (the following will be justified shortly) the parameter $\boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)$ is independent of $\alpha$ and given by $\left(\pi-\pi^{*}, \varpi\right)$ and, for each $\alpha \in A, S_{T \alpha}$ is the score vector, $\mathcal{I}_{\alpha}$ is the (positive definite) Fisher information matrix, and $R_{T}(\alpha, \beta, \pi, \varpi)$ is a remainder term. As the notation indicates, these three terms depend on $\alpha$, and in general they may involve partial derivatives of the log-likelihood function of order higher than two (this is the case for the GMAR model, as will be demonstrated shortly). Then it may also get more complicated to find the reparameterization of the previous subsection and the transformed parameter vector $\boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)$, as the examples of Kasahara and Shimotsu $(2015,2018)$ and the discussion on the GMAR model below show; one possibility is to consider the iterative procedure discussed by Rotnitzky et al. (2000, Secs 4.4, 4.5) (for an illuminating illustration of this approach, see Hallin and Ley (2014)).

Our next assumption provides further details on expansion (12). We use $\Rightarrow$ to signify weak convergence of a sequence of stochastic processes on a function space. In the assumption below, the weak convergence of interest is that of the process $S_{T \alpha}$ (indexed by $\alpha \in A$ ) to a process $S_{\alpha}$. The two function spaces relevant in this paper are $\mathcal{B}\left(A, \mathbb{R}^{k}\right)$ and $\mathcal{C}\left(A, \mathbb{R}^{k}\right)$, the former is the space of all $\mathbb{R}^{k}$-valued bounded functions defined on (the compact set) $A$ equipped with the uniform metric $\left(d(x, y)=\sup _{a \in A}\|x(a)-y(a)\|\right)$, and the latter is the same but with the continuity of the functions (with respect to $\alpha \in A$ ) also assumed.

Assumption 5. For each $\alpha \in A, L_{T}^{\pi}(\alpha, \beta, \pi, \varpi)$ has a quadratic expansion given in (12), where
(i) for each $\alpha \in A, \boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)$ is a mapping from $B \times \Pi_{\alpha}$ to $\Theta_{\alpha}=\left\{\boldsymbol{\theta} \in \mathbb{R}^{r}: \boldsymbol{\theta}=\boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)\right.$ for some $\left.(\beta, \pi, \varpi) \in B \times \Pi_{\alpha}\right\}$ such that (a) $\boldsymbol{\theta}\left(\alpha, \beta^{*}, \pi^{*}, 0\right)=0$ and (b) for all $\epsilon>0$, $\inf _{\alpha \in A} \inf _{(\beta, \pi, \varpi) \in B \times \Pi_{\alpha}:\left\|(\beta, \pi, \varpi)-\left(\beta^{*}, \pi^{*}, 0\right)\right\| \geq \epsilon}\|\boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)\| \geq \delta_{\epsilon}$ for some $\delta_{\epsilon}>0$.
(ii) $S_{T \alpha}=\sum_{t=1}^{T} s_{t \alpha}$ is a sequence of $\mathbb{R}^{r}$-valued $\mathcal{F}_{T}$-measurable stochastic processes indexed by $\alpha \in A$; $S_{T \alpha}$ does not depend on $(\beta, \pi, \varpi)$; $S_{T \alpha}$ has sample paths that are continuous as functions of $\alpha$; the process $T^{-1 / 2} S_{T \alpha}$ obeys $T^{-1 / 2} S_{T \bullet} \Rightarrow S \bullet$ for some mean zero $\mathbb{R}^{r}$-valued Gaussian process $\left\{S_{\alpha}: \alpha \in A\right\}$ that satisfies $E\left[S_{\alpha} S_{\alpha}^{\prime}\right]=E\left[s_{t \alpha} s_{t \alpha}^{\prime}\right]=\mathcal{I}_{\alpha}$ for all $\alpha \in A$ and has continuous sample paths (as functions of $\alpha$ ) with probability 1.
(iii) $\mathcal{I}_{\alpha}$ is, for each $\alpha \in A$, a non-random symmetric $r \times r$ matrix (independent of $(\beta, \pi, \varpi)$ ); $\mathcal{I}_{\alpha}$ is continuous as a function of $\alpha$ and such that $0<\inf _{\alpha \in A} \lambda_{\min }\left(\mathcal{I}_{\alpha}\right), \sup _{\alpha \in A} \lambda_{\max }\left(\mathcal{I}_{\alpha}\right)<\infty$.
(iv) $R_{T}(\alpha, \beta, \pi, \varpi)$ is a remainder term such that

$$
\sup _{(\beta, \pi, \varpi) \in B \times \Pi_{\alpha}:\left\|(\beta, \pi, \varpi)-\left(\beta^{*}, \pi^{*}, 0\right)\right\| \leq \gamma_{T}} \frac{\left|R_{T}(\alpha, \beta, \pi, \varpi)\right|}{\left(1+\left\|T^{1 / 2} \boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)\right\|\right)^{2}}=o_{p \alpha}(1)
$$

for all sequences of (non-random) positive scalars $\left\{\gamma_{T}, T \geq 1\right\}$ for which $\gamma_{T} \rightarrow 0$ as $T \rightarrow \infty$.
Assumption 5(i) describes the transformed parameter $\boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)$ with part (b) being an identification condition, 5 (ii) is the main ingredient needed to derive the limiting distribution of our LR test and would typically be verified relying on a functional central limit theorem, 5 (iii) imposes rather standard conditions on the counterpart of the information matrix, and 5 (iv) ensures that the remainder term $R_{T}(\alpha, \beta, \pi, \varpi)$ has no effect on the final result.

As in Andrews (1999, 2001), Zhu and Zhang (2006), and Kasahara and Shimotsu (2012, 2015), for further developments it will be convenient to write the expansion (12) in an alternative form as

$$
\begin{align*}
& L_{T}^{\pi}(\alpha, \beta, \pi, \varpi)-L_{T}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right) \\
& \quad=\frac{1}{2} Z_{T \alpha}^{\prime} \mathcal{I}_{\alpha} Z_{T \alpha}-\frac{1}{2}\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)-Z_{T \alpha}\right]^{\prime} \mathcal{I}_{\alpha}\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)-Z_{T \alpha}\right]+R_{T}(\alpha, \beta, \pi, \varpi), \tag{13}
\end{align*}
$$

where $Z_{T \alpha}=\mathcal{I}_{\alpha}^{-1} T^{-1 / 2} S_{T \alpha}$. Assumptions 5(ii) and (iii) imply the following facts (that will be justified in the proof of Lemma 1 in Appendix A): $Z_{T \alpha}$ is $\mathcal{F}_{T}$-measurable, independent of ( $\beta, \pi, \varpi$ ), continuous as a function of $\alpha$ with probability 1 , and $Z_{T \bullet} \Rightarrow Z_{\bullet}$ where the mean zero $\mathbb{R}^{r}$-valued Gaussian process $Z_{\alpha}=\mathcal{I}_{\alpha}^{-1} S_{\alpha}$ satisfies $E\left[Z_{\alpha} Z_{\alpha}^{\prime}\right]=\mathcal{I}_{\alpha}^{-1}$ for all $\alpha \in A$ and has continuous sample paths (as functions of $\alpha$ ) with probability 1 .

### 3.3.1 Examples (continued)

LMAR Example. For the LMAR model, expansion (12) (with the unnecessary $\beta$ being dropped everywhere) is obtained from a standard second-order Taylor expansion. Specifically, for an arbitrary fixed $\alpha \in A$, a standard second-order Taylor expansion of $L_{T}^{\pi}(\alpha, \pi, \varpi)=\sum_{t=1}^{T} l_{t}^{\pi}(\alpha, \pi, \varpi)$ around $\left(\pi^{*}, 0\right)$ with respect to the parameters $(\pi, \varpi)$ yields

$$
\begin{align*}
L_{T}^{\pi}(\alpha, \pi, \varpi)-L_{T}^{\pi}\left(\alpha, \pi^{*}, 0\right)= & ( \\
& \left.-\pi^{*}, \varpi\right)^{\prime} \nabla_{(\pi, \varpi)} L_{T}^{\pi}\left(\alpha, \pi^{*}, 0\right)  \tag{14}\\
& +\frac{1}{2}\left(\pi-\pi^{*}, \varpi\right)^{\prime} \nabla_{(\pi, \varpi)(\pi, \varpi)^{\prime}}^{2} L_{T}^{\pi}(\alpha, \dot{\pi}, \dot{\varpi})\left(\pi-\pi^{*}, \varpi\right),
\end{align*}
$$

where $(\dot{\pi}, \dot{\varpi})$ denotes a point between $(\pi, \varpi)$ and $\left.\left(\pi^{*}, 0\right), \nabla_{(\pi, \varpi)} L_{T}^{\pi}\left(\alpha, \pi^{*}, 0\right)=\sum_{t=1}^{T} \nabla_{(\pi, \varpi)}\right)_{t}^{\pi}\left(\alpha, \pi^{*}, 0\right)$ and $\nabla_{(\pi, \varpi)(\pi, \varpi)^{\prime}}^{2} L_{T}^{\pi}(\alpha, \pi, \varpi)=\sum_{t=1}^{T} \nabla_{(\pi, \varpi)(\pi, \varpi)^{\prime} l_{t}^{\pi}}^{2}(\alpha, \pi, \varpi)$ (explicit expressions for the required derivatives are provided in Appendix B), and $\nabla$ and $\nabla^{2}$ denote first- and second-order differentiation with respect to the indicated parameters. Set $\boldsymbol{\theta}(\alpha, \pi, \varpi)=\left(\pi-\pi^{*}, \varpi\right)=(\theta, \vartheta)$ and note that the parameter space $\Theta_{\alpha}=\Theta$ is independent of $\alpha$ and contains the origin, corresponding to the true model, as an interior point. Then define the vector $S_{T \alpha}$ and the matrix $\mathcal{I}_{\alpha} \mathrm{as}^{7}$

$$
\begin{align*}
S_{T \alpha} & =\nabla_{(\pi, w)} L_{T}^{\pi}\left(\alpha, \pi^{*}, 0\right)=\sum_{t=1}^{T}\left(\frac{\nabla f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)},-\left(1-\alpha_{1, t}^{L}(\alpha)\right) \frac{\nabla f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)}\right),  \tag{15}\\
\mathcal{I}_{\alpha} & =\left[\begin{array}{cc}
E\left[\frac{\nabla f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)} \frac{\nabla^{\prime} f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)}\right] & -E\left[\left(1-\alpha_{1, t}^{L}(\alpha)\right) \frac{\nabla f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)} \frac{\nabla^{\prime} f_{t}\left(\pi^{*}\right)}{f_{t} t\left(\pi^{*}\right)}\right] \\
-E\left[\left(1-\alpha_{1, t}^{L}(\alpha)\right) \frac{\nabla f_{t}\left(\left(\pi^{*}\right)\right.}{f_{t}\left(\pi^{*}\right)} \frac{\nabla^{\prime} f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)}\right] & E\left[\left(1-\alpha_{1, t}^{L}(\alpha)\right)^{2} \frac{\left.\nabla f_{t}(t)^{*}\right)}{f_{t}\left(\pi^{*}\right)} \frac{\nabla_{t} f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)}\right]
\end{array}\right] .
\end{align*}
$$

[^4]Adding and subtracting terms and reorganizing, expansion (14) can be written as

$$
\begin{align*}
L_{T}^{\pi}(\alpha, \pi, \varpi)-L_{T}^{\pi}\left(\alpha, \pi^{*}, 0\right)= & \left(T^{-1 / 2} S_{T \alpha}\right)^{\prime}\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \pi, \varpi)\right] \\
& -\frac{1}{2}\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \pi, \varpi)\right]^{\prime} \mathcal{I}_{\alpha}\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \pi, \varpi)\right]+R_{T}(\alpha, \pi, \varpi), \tag{16}
\end{align*}
$$

with the remainder term

$$
\begin{equation*}
R_{T}(\alpha, \pi, \varpi)=\frac{1}{2}\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \pi, \varpi)\right]^{\prime}\left[T^{-1} \nabla_{(\pi, \varpi)(\pi, \varpi)^{\prime}}^{2} L_{T}^{\pi}(\alpha, \dot{\pi}, \dot{\varpi})-\left(-\mathcal{I}_{\alpha}\right)\right]\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \pi, \varpi)\right] . \tag{17}
\end{equation*}
$$

These equations yield the expansion (12) in the case of the LMAR model. For details of verifying Assumptions 4 and 5, we refer to Appendix B.

As mentioned in the LMAR example of Section 3.1.1, the treatment of the special case where the mixing weight $\alpha_{t}^{L}$ is constant is more complicated than that of the 'real' LMAR case. Indeed, replacing the mixing weight $\alpha_{t}^{L}$ by a constant in the preceding expression of the score vector $S_{T \alpha}$ immediately shows that the second-order Taylor expansion (16) breaks down because, contrary to Assumption 5(iii), the components of $S_{T \alpha}$ are not linearly independent and, consequently, the Fisher information matrix $\mathcal{I}_{\alpha}$ is singular. Thus, a higher-order Taylor expansion is needed to analyze the LR test statistic.

To give an idea of how one could proceed, we first note that the partial derivatives of the loglikelihood function behave in the same way as their counterparts in Kasahara and Shimotsu (2015) where mixtures of normal regression models (with constant mixing weights) are considered (see particularly the discussion following their Proposition 1). This is due to the fact that in the special case of constant mixing weights the LMAR model is obtained from the model considered in Kasahara and Shimotsu (2015) by replacing the exogenous regressors therein by lagged values of $y_{t}$. Thus, the arguments employed in that paper could be used to obtain the asymptotic distribution of the LR test statistic. Instead of a conventional second-order Taylor expansion this would require a more complicated reparameterization and an expansion based on partial derivatives of the log-likelihood function up to order eight. As most of the details appear very similar to those in Kasahara and Shimotsu (2015) we have preferred not to pursue this matter in this paper.

The preceding discussion means that, in the case of the LMAR model, time varying mixing weights are beneficial when the purpose is to derive a LR test for the adequacy of a single-regime model. A similar observation was made already by Jeffries (1998, p. 80). However, this does not happen in all mixture autoregressive models with time varying mixing weights, as the following discussion on the GMAR model demonstrates.

GMAR Example. As alluded to earlier, in the case of the GMAR model the expansion (12) cannot be based on a second-order Taylor expansion of the log-likelihood function. A higher-order expansion is required, and similarly to Kasahara and Shimotsu (2012) the appropriate order turns out to be the fourth one with the elements of $\nabla_{\beta} l_{t}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right)$ and $\nabla_{\pi} l_{t}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right)$ and the distinctive elements of $\nabla_{\varpi \varpi} l_{t}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right)$ (suitably normalized) used to define the vector $S_{T \alpha}$. In Appendix C we present, for an arbitrary fixed $\alpha \in A$, the explicit form of a standard fourth-order Taylor expansion of $L_{T}^{\pi}(\alpha, \beta, \pi, \varpi)=\sum_{t=1}^{T} l_{t}^{\pi}(\alpha, \beta, \pi, \varpi)$ around $\left(\beta^{*}, \pi^{*}, 0\right)$ with respect to the parameters $(\beta, \pi, \varpi)$. Therein we also demonstrate that this fourth-order Taylor expansion can be written as a quadratic expansion of the form (12) (or (13)) with the different quantities appearing therein defined as follows.

Define the vector $\boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)$ in (12) as

$$
\boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)=\left[\begin{array}{c}
\theta(\alpha, \beta, \pi, \varpi) \\
\vartheta(\alpha, \beta, \pi, \varpi)
\end{array}\right]=\left[\begin{array}{c}
\beta-\beta^{*} \\
\pi-\pi^{*} \\
\alpha(1-\alpha) v(\varpi)
\end{array}\right],
$$

where $\theta(\alpha, \beta, \pi, \varpi)$ is $\left(q_{1}+q_{2}\right) \times 1$ and $\vartheta(\alpha, \beta, \pi, \varpi)$ is $q_{\vartheta} \times 1$ with $q_{\vartheta}=q_{2}\left(q_{2}+1\right) / 2$ (where $q_{1}=1$ and $\left.q_{2}=p+1\right)$ and where the vector $v(\varpi)$ contains the unique elements of $\varpi \varpi^{\prime}$, that is,

$$
v(\varpi)=\left(\varpi_{1}^{2}, \ldots, \varpi_{q_{2}}^{2}, \varpi_{1} \varpi_{2}, \ldots, \varpi_{1} \varpi_{q_{2}}, \varpi_{2} \varpi_{3}, \ldots, \varpi_{q_{2}-1} \varpi_{q_{2}}\right)
$$

(note that $v(\varpi)$ is just a re-ordering of $\operatorname{vech}\left(\varpi \varpi^{\prime}\right)$ ). The parameter space
$\Theta_{\alpha}=\left\{\boldsymbol{\theta}=(\theta, \vartheta) \in \mathbb{R}^{q_{1}+q_{2}+q_{\vartheta}}: \theta=\left(\beta-\beta^{*}, \pi-\pi^{*}\right), \vartheta=\alpha(1-\alpha) v(\varpi)\right.$ for some $\left.(\beta, \pi, \varpi) \in B \times \Pi_{\alpha}\right\}$
now depends on $\alpha$ and has the origin, corresponding to the true model, as a boundary point (due to the particular shape of the range of $\vartheta=\alpha(1-\alpha) v(\varpi))$; both of these features will complicate the subsequent analysis. Next set

$$
\begin{equation*}
S_{T}\left(=S_{T \alpha}\right)=\sum_{t=1}^{T} \tilde{\nabla}_{\boldsymbol{\theta}} l_{t}^{\pi *} \quad \text { where } \quad \tilde{\nabla}_{\boldsymbol{\theta}} l_{t}^{\pi *}=\left(\tilde{\nabla}_{\theta} l_{t}^{\pi *}, \tilde{\nabla}_{\vartheta} l_{t}^{\pi *}\right) \quad\left(\left(q_{1}+q_{2}+q_{\vartheta}\right) \times 1\right) \tag{18}
\end{equation*}
$$

with the component vectors $\tilde{\nabla}_{\theta} l_{t}^{\pi *}$ and $\tilde{\nabla}_{\vartheta} l_{t}^{\pi *}$ given by

$$
\begin{align*}
\tilde{\nabla}_{\theta} l_{t}^{\pi *}= & \left(\nabla_{\beta} l_{t}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right), \nabla_{\pi} l_{t}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right)\right), \\
\tilde{\nabla}_{\vartheta} l_{t}^{\pi *}= & \left(c_{11} \nabla_{\varpi_{1} \varpi_{1}}^{2} l_{t}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right), \ldots, c_{q_{2} q_{2}} \nabla_{\varpi_{q_{2}} \varpi_{q_{2}}}^{2} l_{t}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right),\right. \\
& \left.c_{12} \nabla_{\varpi_{1} \varpi_{2}}^{2} l_{t}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right), \ldots, c_{q_{2}-1, q_{2}} \nabla_{\varpi_{q_{2}-1} \varpi_{q_{2}}}^{2} l_{t}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right)\right) /(\alpha(1-\alpha)), \tag{19}
\end{align*}
$$

where $c_{i j}=1 / 2$ if $i=j$ and $c_{i j}=1$ if $i \neq j$. Explicit expressions for $\tilde{\nabla}_{\theta} l_{t}^{\pi *}$ and $\tilde{\nabla}_{\vartheta} l_{t}^{\pi *}$ can be found in Appendix C , and from them it can be seen that $S_{T}$ depends on $\left(\beta^{*}, \pi^{*}\right)$ only and not on $(\alpha, \beta, \pi, \varpi)$. The same is true for the matrix $\mathcal{I}\left(=\mathcal{I}_{\alpha}\right)\left(\left(q_{1}+q_{2}+q_{\vartheta}\right) \times\left(q_{1}+q_{2}+q_{\vartheta}\right)\right)$ whose expression is also given in Appendix C. Finally, an explicit expression of the remainder term $R_{T}(\alpha, \beta, \pi, \varpi)$ is given in Appendix C. For the verification of Assumptions 4 and 5, see Appendix C.

In the GMAR example we have assumed that the intercept terms $\tilde{\phi}_{0}$ and $\tilde{\varphi}_{0}$ in the two regimes are the same. We are now in a position to describe the difficulties that allowing for $\tilde{\phi}_{0} \neq \tilde{\varphi}_{0}$ (and, hence, dropping $\beta$ ) would entail. In this case, the additional parameter $\tilde{\varphi}_{0}$ would correspond to $\varpi_{1}$, the first component of $\varpi$. As in Section 3.2.1, it would again be the case that $\nabla_{\varpi} l_{t}^{\pi}\left(\alpha, \pi^{*}, 0\right)=0$, leading us to consider second derivatives. But now, due to the properties of the Gaussian distribution, it would be the case that $\nabla_{\varpi_{1} \varpi_{1}}^{2} l_{t}^{\pi}\left(\alpha, \pi^{*}, 0\right)$ is linearly dependent with the components of $\nabla_{\pi} l_{t}^{\pi}\left(\alpha, \pi^{*}, 0\right)$, making it unsuitable to serve as a component of $S_{T}$. A reparameterization more complicated than that used in Section 3.2.1 would be needed, with the aim of obtaining $\nabla_{\varpi_{1} \varpi_{1}}^{2} l_{t}^{\pi}\left(\alpha, \pi^{*}, 0\right)=0$ and, instead of $\nabla_{\varpi_{1} \varpi_{1}}^{2} l_{t}^{\pi}\left(\alpha, \pi^{*}, 0\right)$, using $\nabla_{\varpi_{1} \varpi_{1} \varpi_{1}}^{3} l_{t}^{\pi}\left(\alpha, \pi^{*}, 0\right)$ or perhaps $\nabla_{\varpi_{1} \varpi_{1} \varpi_{1} \varpi_{1}}^{4} l_{t}^{\pi}\left(\alpha, \pi^{*}, 0\right)$ as the counterpart of the score of the parameter $\varpi_{1}$. It turns out that (restricting the discussion to the case $p=1$ only) the third derivative is suitable when $\alpha \neq 1 / 2$ and $\tilde{\phi}_{1} \neq-1 / 2$, but that fourth (or higher) order derivatives are needed when $\alpha=1 / 2$ or $\tilde{\phi}_{1}=-1 / 2$. Similar difficulties (involving situations comparable to the cases $\alpha=1 / 2$ vs. $\alpha \neq 1 / 2$, but apparently not ones involving also a counterpart of $\tilde{\phi}_{1}$ ) were faced by Cho and White (2007, Sec 2.3.3) and Kasahara and Shimotsu (2018, Sec 6.2), whose analyses suggest that expanding the log-likelihood at least to the eighth order is required. As the required analysis gets excessively complicated, we have chosen to leave it for future research.

### 3.4 Asymptotic analysis of the quadratic expansion

We continue by analyzing expansion (13) evaluated at ( $\hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}$ ). Previously, a similar analysis was provided by Andrews (2001) but his approach is not directly applicable in our setting. The reason for this is that in the quadratic expansion in (13) the dependence of the parameter $\boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)$ and its
parameter space $\Theta_{\alpha}$ on the nuisance parameter $\alpha$ is not compatible with the formulation of Andrews (2001, eqn (3.3)). The results of Zhu and Zhang (2006) probably cover our case but instead of trying to verify their assumptions we prove the needed results directly by adapting the arguments used in Andrews (1999, 2001) and Zhu and Zhang (2006) to our setting. We proceed in several steps.

Asymptotic insignificance of the remainder term. We first establish that the remainder term $R_{T}(\alpha, \beta, \pi, \varpi)$, when evaluated at ( $\hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}$ ), has no effect on the asymptotic distribution of the quadratic expansion. A crucial ingredient in showing this is that the transformed parameter vector $\boldsymbol{\theta}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)$ is root- $T$ consistent in the sense that $\left\|T^{1 / 2} \boldsymbol{\theta}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)\right\|=O_{p \alpha}(1)$. This, together with Assumption 5(iv), allows us to obtain the result $R_{T}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)=o_{p \alpha}(1)$. The following lemma collects these results; assertion (iii) below is analogous to Andrews (1999, Thm 2b).

Lemma 1. If Assumptions 1-5 hold ${ }^{8}$, then (i) $\left\|T^{1 / 2} \boldsymbol{\theta}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)\right\|=O_{p \alpha}(1)$,
(ii) $R_{T}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)=o_{p \alpha}(1)$, and (iii)
$L_{T}^{\pi}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)-L_{T}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right)$
$=\frac{1}{2} Z_{T \alpha}^{\prime} \mathcal{I}_{\alpha} Z_{T \alpha}-\frac{1}{2}\left[T^{1 / 2} \boldsymbol{\theta}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)-Z_{T \alpha}\right]^{\prime} \mathcal{I}_{\alpha}\left[T^{1 / 2} \boldsymbol{\theta}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)-Z_{T \alpha}\right]+o_{p \alpha}(1)$.
Maximization of the likelihood vs. minimization of a related quadratic form. For later developments it would be convenient if the ML estimator ( $\hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}$ ) on the right hand side of (20) could be replaced by an (approximate) minimizer of the quadratic form $\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)-\right.$ $\left.Z_{T \alpha}\right]^{\prime} \mathcal{I}_{\alpha}\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)-Z_{T \alpha}\right]$. To this end, note that by definition (cf. Andrews (1999, eqn (3.6)))
$\inf _{(\beta, \pi, \varpi) \in B \times \Pi_{\alpha}}\left\{\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)-Z_{T \alpha}\right]^{\prime} \mathcal{I}_{\alpha}\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)-Z_{T \alpha}\right]\right\}=\inf _{\boldsymbol{\lambda} \in \Theta_{\alpha, T}}\left\{\left(\boldsymbol{\lambda}-Z_{T \alpha}\right)^{\prime} \mathcal{I}_{\alpha}\left(\boldsymbol{\lambda}-Z_{T \alpha}\right)\right\}$ where, for each $T$,

$$
\Theta_{\alpha, T}=\left\{\boldsymbol{\lambda} \in \mathbb{R}^{r}: \boldsymbol{\lambda}=T^{1 / 2} \boldsymbol{\theta} \text { for some } \boldsymbol{\theta} \in \Theta_{\alpha}\right\}
$$

and $\Theta_{\alpha}$ is as defined in Assumption 5(i). Now, for each $\alpha \in A$, let $\hat{\boldsymbol{\lambda}}_{T \alpha q}=T^{1 / 2} \boldsymbol{\theta}\left(\alpha, \hat{\beta}_{T \alpha q}, \hat{\pi}_{T \alpha q}, \hat{\varpi}_{T \alpha q}\right)$ (with the additional ' $q$ ' in the subscripts referring to quadratic form) denote an approximate minimizer of $\left(\boldsymbol{\lambda}-Z_{T \alpha}\right)^{\prime} \mathcal{I}_{\alpha}\left(\boldsymbol{\lambda}-Z_{T \alpha}\right)$ over $\Theta_{\alpha, T}$, that is, (cf. Zhu and Zhang (2006, eqn (12)))

$$
\begin{equation*}
\left(\hat{\boldsymbol{\lambda}}_{T \alpha q}-Z_{T \alpha}\right)^{\prime} \mathcal{I}_{\alpha}\left(\hat{\boldsymbol{\lambda}}_{T \alpha q}-Z_{T \alpha}\right)=\inf _{\boldsymbol{\lambda} \in \Theta_{\alpha, T}}\left\{\left(\boldsymbol{\lambda}-Z_{T \alpha}\right)^{\prime} \mathcal{I}_{\alpha}\left(\boldsymbol{\lambda}-Z_{T \alpha}\right)\right\}+o_{p \alpha}(1) . \tag{21}
\end{equation*}
$$

The following lemma justifies the discussed replacement of $\left(\hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)$ with $\hat{\boldsymbol{\lambda}}_{T \alpha q}$.
Lemma 2. If Assumptions 1-5 hold, then

$$
L_{T}^{\pi}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)-L_{T}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right)=\frac{1}{2} Z_{T \alpha}^{\prime} \mathcal{I}_{\alpha} Z_{T \alpha}-\frac{1}{2}\left(\hat{\boldsymbol{\lambda}}_{T \alpha q}-Z_{T \alpha}\right)^{\prime} \mathcal{I}_{\alpha}\left(\hat{\boldsymbol{\lambda}}_{T \alpha q}-Z_{T \alpha}\right)+o_{p \alpha}(1) .
$$

Approximating the parameter space with a cone. In (21) the quadratic form $\left(\boldsymbol{\lambda}-Z_{T \alpha}\right)^{\prime} \mathcal{I}_{\alpha}(\boldsymbol{\lambda}-$ $\left.Z_{T \alpha}\right)$ is minimized over the set $\Theta_{\alpha, T}$; this minimization can in fact be done over a simpler set. We first introduce some terminology and say that a collection of sets $\left\{\Gamma_{\alpha}, \alpha \in A\right\}$ (where for each $\alpha \in A$, $\Gamma_{\alpha} \subset \mathbb{R}^{r}$ ) is 'locally (at the origin) uniformly equal' to a set $\Lambda \subset \mathbb{R}^{r}$ if there exists a $\delta>0$ such that $\Gamma_{\alpha} \cap(-\delta, \delta)^{r}=\Lambda \cap(-\delta, \delta)^{r}$ for all $\alpha \in A$. Note that ' $\left\{\Gamma_{\alpha}, \alpha \in A\right\}$ is locally uniformly equal to $\Lambda$ ' implies that (i) 'for all $\alpha \in A, \Gamma_{\alpha}$ is locally equal to $\Lambda$ in the sense of Andrews (1999, p. 1359)', but the reverse does not hold; and also that (ii) $\left\{\Gamma_{\alpha}, \alpha \in A\right\}$ is uniformly approximated by the set $\Lambda$ in

[^5]the sense of Zhu and Zhang (2006, Defn 3). We also say that a set $\Lambda \subset \mathbb{R}^{r}$ is a 'cone' if $\lambda \in \Lambda$ implies that $a \lambda \in \Lambda$ for all positive real scalars $a$. We can now state the following assumption.

Assumption 6. The collection of sets $\left\{\Theta_{\alpha}, \alpha \in A\right\}$ is locally uniformly equal to a cone $\Lambda\left(\subset \mathbb{R}^{r}\right)$.
By Assumption $5(\mathrm{i})(\mathrm{a}), 0 \in \Theta_{\alpha}$ for all $\alpha \in A$, so that the cone $\Lambda$ in Assumption 6 necessarily contains $0\left(\in \mathbb{R}^{r}\right)$. The cone $\Lambda$ also does not depend on $\alpha$. Now we can establish the following result.

Lemma 3. If Assumptions 1-6 hold, then

$$
\inf _{\boldsymbol{\lambda} \in \Theta_{\alpha, T}}\left\{\left(\boldsymbol{\lambda}-Z_{T \alpha}\right)^{\prime} \mathcal{I}_{\alpha}\left(\boldsymbol{\lambda}-Z_{T \alpha}\right)\right\}=\inf _{\boldsymbol{\lambda} \in \Lambda}\left\{\left(\boldsymbol{\lambda}-Z_{T \alpha}\right)^{\prime} \mathcal{I}_{\alpha}\left(\boldsymbol{\lambda}-Z_{T \alpha}\right)\right\}+o_{p \alpha}(1) .
$$

Describing the limiting random process. From Lemmas 2-3 and the definition of $\hat{\boldsymbol{\lambda}}_{T \alpha q}$ we obtain

$$
\begin{equation*}
2\left[L_{T}^{\pi}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)-L_{T}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right)\right]=Z_{T \alpha}^{\prime} \mathcal{I}_{\alpha} Z_{T \alpha}-\inf _{\boldsymbol{\lambda} \in \Lambda}\left\{\left(\boldsymbol{\lambda}-Z_{T \alpha}\right)^{\prime} \mathcal{I}_{\alpha}\left(\boldsymbol{\lambda}-Z_{T \alpha}\right)\right\}+o_{p \alpha}(1) . \tag{22}
\end{equation*}
$$

The assumed weak convergence of $S_{T \alpha}$ (and hence that of $Z_{\alpha}=\mathcal{I}_{\alpha}^{-1} S_{\alpha}$ ) allows us to derive the weak limit of this random process described in the following lemma.

Lemma 4. If Assumptions 1-6 hold, then

$$
2\left[L_{T}^{\pi}\left(\bullet, \hat{\beta}_{T_{\bullet}}, \hat{\pi}_{T \bullet}, \hat{\varpi}_{T_{\bullet}}\right)-L_{T}^{\pi}\left(\bullet, \beta^{*}, \pi^{*}, 0\right)\right] \Rightarrow Z_{\bullet}^{\prime} \mathcal{I}_{\bullet} Z_{\bullet}-\inf _{\boldsymbol{\lambda} \in \Lambda}\left\{\left(\boldsymbol{\lambda}-Z_{\bullet}\right)^{\prime} \mathcal{I}_{\bullet}\left(\boldsymbol{\lambda}-Z_{\bullet}\right)\right\} .
$$

The limiting random process in Lemma 4 can be written in a somewhat simpler form (cf. Andrews (1999, Thm 4; 2001, Thm 2)). The motivation for this comes from the fact that in our applications $\boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)$ can be decomposed into two parts as $\boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)=(\theta(\alpha, \beta, \pi, \varpi), \vartheta(\alpha, \beta, \pi, \varpi))$ with $\theta(\alpha, \beta, \pi, \varpi) \in \mathbb{R}^{q_{\theta}}$ and $\vartheta(\alpha, \beta, \pi, \varpi) \in \mathbb{R}^{q_{\vartheta}}$ (with $q_{\theta}=q_{1}+q_{2}$ and $r=q_{\theta}+q_{\vartheta}$ ) such that (i) the values of $\theta(\alpha, \beta, \pi, \varpi)$ are not restricted by the null hypothesis and do not lie on the boundary of the parameter space and (ii) the values of $\vartheta(\alpha, \beta, \pi, \varpi)$ are restricted by the null hypothesis and potentially lie on the boundary of the parameter space. Specifically, we assume the following.

Assumption 7. The cone $\Lambda$ of Assumption 6 satisfies $\Lambda=\mathbb{R}^{q_{\theta}} \times \Lambda_{\vartheta}$ with $\Lambda_{\vartheta}$ a cone in $\mathbb{R}^{q_{\vartheta}}$.
Partition $S_{\alpha}, Z_{\alpha}, \boldsymbol{\lambda}$, and $\mathcal{I}_{\alpha}$ conformably with the partition $\boldsymbol{\theta}=(\theta, \vartheta)$ as

$$
S_{\alpha}=\left[\begin{array}{c}
S_{\theta \alpha} \\
S_{\vartheta \alpha}
\end{array}\right], \quad Z_{\alpha}=\left[\begin{array}{c}
Z_{\theta \alpha} \\
Z_{\vartheta \alpha}
\end{array}\right], \quad \boldsymbol{\lambda}=\left[\begin{array}{c}
\boldsymbol{\lambda}_{\theta} \\
\boldsymbol{\lambda}_{\vartheta}
\end{array}\right], \quad \mathcal{I}_{\alpha}=\left[\begin{array}{cc}
\mathcal{I}_{\theta \theta \alpha} & \mathcal{I}_{\theta \vartheta \alpha} \\
\mathcal{I}_{\vartheta \theta \alpha} & \mathcal{I}_{\vartheta \vartheta \alpha}
\end{array}\right]
$$

and let $\left(\mathcal{I}_{\alpha}^{-1}\right)_{\vartheta \vartheta}$ denote the $\left(q_{\vartheta} \times q_{\vartheta}\right)$ bottom right block of $\mathcal{I}_{\alpha}^{-1}$. Assumption 7 together with properties of partitioned matrices yields the following result.

Lemma 5. If Assumptions 1-7 hold, then

$$
\begin{aligned}
& Z_{\alpha}^{\prime} \mathcal{I}_{\alpha} Z_{\alpha}-\inf _{\boldsymbol{\lambda} \in \Lambda}\left\{\left(\boldsymbol{\lambda}-Z_{\alpha}\right)^{\prime} \mathcal{I}_{\alpha}\left(\boldsymbol{\lambda}-Z_{\alpha}\right)\right\} \\
& \quad=Z_{\vartheta \alpha}^{\prime}\left(\mathcal{I}_{\alpha}^{-1}\right)_{\vartheta \vartheta}^{-1} Z_{\vartheta \alpha}-\inf _{\boldsymbol{\lambda}_{\vartheta} \in \Lambda_{\vartheta}}\left\{\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)^{\prime}\left(\mathcal{I}_{\alpha}^{-1}\right)_{\vartheta \vartheta}^{-1}\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)\right\}+S_{\theta \alpha}^{\prime} \mathcal{I}_{\theta \theta \alpha}^{-1} S_{\theta \alpha} .
\end{aligned}
$$

Explicit expressions for $\left(\mathcal{I}_{\alpha}^{-1}\right)_{\vartheta \vartheta}$ and $Z_{\vartheta \alpha}$ in terms of $S_{\alpha}$ and $\mathcal{I}_{\alpha}$ are given in the proof of this lemma in Supplementary Appendix D.

### 3.5 The LR test statistic

### 3.5.1 Derivation of the test statistic

The previous subsection described the asymptotic behavior of $2\left[L_{T}^{\pi}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)-L_{T}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right)\right]$, the first term in the expression of $L R_{T}(\alpha)$ in (11). Now consider the second term, namely $2\left[L_{T}^{0}\left(\hat{\tilde{\phi}}_{T}\right)-\right.$ $\left.L_{T}^{0}\left(\tilde{\phi}^{*}\right)\right]$, corresponding to the model restricted by the null hypothesis. Recall that $L_{T}^{0}(\tilde{\phi})=\sum_{t=1}^{T} l_{t}^{0}(\tilde{\phi})$ with $l_{t}^{0}(\tilde{\phi})=\log \left[f_{t}(\tilde{\phi})\right]$ so that $\nabla_{\tilde{\phi}} l_{t}^{0}\left(\tilde{\phi}^{*}\right)=\left(f_{t}\left(\tilde{\phi}^{*}\right)\right)^{-1} \nabla f_{t}\left(\tilde{\phi}^{*}\right)$ with $\tilde{\phi}^{*}$ an interior point of $\tilde{\Phi}$. Denote the score vector and limiting information matrix by

$$
S_{T}^{0}=\sum_{t=1}^{T} \frac{\nabla f_{t}\left(\tilde{\phi}^{*}\right)}{f_{t}\left(\tilde{\phi}^{*}\right)}, \quad \mathcal{I}^{0}=E\left[\frac{\nabla f_{t}\left(\tilde{\phi}^{*}\right)}{f_{t}\left(\tilde{\phi}^{*}\right)} \frac{\nabla^{\prime} f_{t}\left(\tilde{\phi}^{*}\right)}{f_{t}\left(\tilde{\phi}^{*}\right)}\right],
$$

respectively. For the following assumption, partition the process $S_{T \alpha}$ of Assumption 5 as $S_{T \alpha}=$ $\left(S_{T \theta \alpha}, S_{T \vartheta \alpha}\right)$ (with $S_{T \theta \alpha} q_{\theta}$-dimensional and $S_{T \vartheta \alpha} q_{\vartheta}$-dimensional). The following simplifying assumption, which holds in our examples (see the expressions of $S_{T \alpha}$ in (15) and (18)), allows us to obtain a neat expression for the likelihood ratio test statistic in Theorem 1 below.

Assumption 8. $S_{T \theta \alpha}=S_{T}^{0}$.
Together with the earlier assumptions, Assumption 8 implies that $T^{-1 / 2} S_{T \theta \alpha}=T^{-1 / 2} S_{T}^{0} \xrightarrow{d} S^{0}$, a $q_{\theta}$-dimensional Gaussian random vector with mean zero and covariance matrix $\mathcal{I}_{\theta \theta \alpha}=E\left[S^{0} S^{0}\right]=\mathcal{I}^{0}$. Standard likelihood theory now implies the following result.
Lemma 6. If Assumptions 1-8 hold, then $2\left[L_{T}^{0}\left(\hat{\tilde{\phi}}_{T}\right)-L_{T}^{0}\left(\tilde{\phi}^{*}\right)\right] \xrightarrow{d} S^{0 \prime}\left(\mathcal{I}^{0}\right)^{-1} S^{0}$, and the convergence is joint with that in Lemma 4.

The preceding results, in particular Lemmas 4-6, now yield the distribution of the LR test statistic.
Theorem 1. If Assumptions 1-8 hold, then
(i) $L R_{T}(\bullet) \Rightarrow Z_{\vartheta \bullet}^{\prime}\left(\mathcal{I}_{\bullet}^{-1}\right)_{\vartheta \vartheta}^{-1} Z_{\vartheta \bullet}-\inf _{\boldsymbol{\lambda}_{\vartheta} \in \Lambda_{\vartheta}}\left\{\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \bullet}\right)^{\prime}\left(\mathcal{I}_{\bullet}^{-1}\right)_{\vartheta \vartheta}^{-1}\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \bullet}\right)\right\}$, and
(ii) $L R_{T}=\sup _{\alpha \in A} L R_{T}(\alpha) \xrightarrow{d} \sup _{\alpha \in A}\left\{Z_{\vartheta \alpha}^{\prime}\left(\mathcal{I}_{\alpha}^{-1}\right)_{\vartheta \vartheta}^{-1} Z_{\vartheta \alpha}-\inf _{\boldsymbol{\lambda}_{\vartheta} \in \Lambda_{\vartheta}}\left\{\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)^{\prime}\left(\mathcal{I}_{\alpha}^{-1}\right)_{\vartheta \vartheta}^{-1}\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)\right\}\right\}$.

The asymptotic distribution of the LR test statistic is similar to that in Andrews (2001, Thm 4). As we next discuss, this distribution simplifies in both the LMAR and the GMAR examples.

### 3.5.2 Examples (continued)

LMAR Example. As was noted in Section 3.3.1, the LMAR case is rather standard in the sense that a conventional second-order Taylor expansion with a nonsingular information matrix and with no parameters on the boundary was sufficient to study the LR test. The only nonstandard feature in this case is the presence of unidentified parameters under the null hypothesis. Validity of Assumptions 6-8 is easy to check (see Appendix B) with the cone $\Lambda$ of Assumption 6 equal to $\mathbb{R}^{r}$. Thus the infimum in the distribution of the $L R_{T}$ statistic in Theorem 1(ii) equals zero and the result therein simplifies to ${ }^{9}$

$$
L R_{T}=\sup _{\alpha \in A} L R_{T}(\alpha) \xrightarrow{d} \sup _{\alpha \in A}\left\{Z_{\vartheta \alpha}^{\prime}\left(\mathcal{I}_{\alpha}^{-1}\right)_{\vartheta \vartheta}^{-1} Z_{\vartheta \alpha}\right\} .
$$

For every fixed $\alpha \in A$, the quantity $Z_{\vartheta \alpha}^{\prime}\left(\mathcal{I}_{\alpha}^{-1}\right)_{\vartheta \vartheta}^{-1} Z_{\vartheta \alpha}$ is a chi-squared random variable, so that the limiting distribution is a supremum of a chi-squared process similarly as in, for example, Davies (1987), Hansen (1996, Thm 1), and Andrews (2001, eqn (5.7)).

[^6]GMAR Example. In Section 3.3.1 it was seen that in the GMAR example $Z_{\alpha}$ and $\mathcal{I}_{\alpha}$ do not depend on $\alpha$. As the cone $\Lambda$ of Assumption 6 does not depend on $\alpha$ either, the weak limit of $L R_{T}(\alpha)$ does not depend on $\alpha$. Therefore the result of Theorem 1 (validity of Assumptions 6-8 is checked in Appendix C) simplifies to

$$
L R_{T}=\sup _{\alpha \in A} L R_{T}(\alpha) \xrightarrow{d} Z_{\vartheta}^{\prime}\left(\mathcal{I}^{-1}\right)_{\vartheta \vartheta}^{-1} Z_{\vartheta}-\inf _{\boldsymbol{\lambda}_{\vartheta} \in \Lambda_{\vartheta}}\left\{\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta}\right)^{\prime}\left(\mathcal{I}^{-1}\right)_{\vartheta \vartheta}^{-1}\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta}\right)\right\}
$$

where the unnecessary $\alpha$ has been dropped from the notation. Here $Z_{\vartheta}$ follows an $q_{\vartheta}$-variate Gaussian distribution with covariance matrix $\left(\mathcal{I}^{-1}\right)_{\vartheta \vartheta}$, and the limiting distribution, which is sometimes referred to as the chi-bar-squared distribution, is similar to the one in Kasahara and Shimotsu (2012, Propn $3 \mathrm{c}, \mathrm{d}$ ). Note that the cone $\Lambda_{\vartheta}=v\left(\mathbb{R}^{q_{2}}\right)$ (see Appendix C) is not convex (in contrast to at least most of the examples in Andrews (2001), but similarly to Kasahara and Shimotsu (2012, Propn 3c,d)) and the dimension of this cone, $q_{\vartheta}=q_{2}\left(q_{2}+1\right) / 2$, may not be small either $\left(q_{\vartheta}=3,6,10, \ldots\right.$ for $\left.q_{2}=2,3,4, \ldots\right)$.

## 4 Simulating the asymptotic null distribution and the bootstrap

### 4.1 Simulating the asymptotic null distribution

Similarly to Hansen (1996) and Andrews (2001), the asymptotic null distribution of the LR statistic in Theorem 1 is typically application-specific and cannot be tabulated. Following these papers, one option is to use simulation methods to obtain critical values of the asymptotic null distribution. The following procedure is analogous to the ones used by Hansen (1996) in a rather general setting involving unidentified nuisance parameters (and also specifically in threshold models), Andrews (2001, Sec 5.6) in a random coefficient model, and Zhu and Zhang (2004, Sec 2.1) in a mixture regression model.

To set notation, assume the (observed or simulated) sample is ( $y_{1-p}, \ldots, y_{0}, y_{1}, \ldots, y_{T}$ ), obtain the estimator under the null hypothesis $\hat{\tilde{\phi}}_{T}=\left(\hat{\tilde{\phi}}_{0, T}, \hat{\tilde{\phi}}_{1, T}, \ldots, \hat{\tilde{\phi}}_{p, T}, \hat{\tilde{\sigma}}_{T}^{2}\right)$ (e.g. the OLS estimator), the estimator under the alternative hypothesis $\left(\hat{\alpha}_{T}, \hat{\beta}_{T}, \hat{\phi}_{T}, \hat{\varphi}_{T}\right)$, and the LR statistic $L R_{T}=2\left[L_{T}\left(\hat{\alpha}_{T}, \hat{\beta}_{T}, \hat{\phi}_{T}, \hat{\varphi}_{T}\right)-\right.$ $\left.L_{T}^{0}\left(\hat{\tilde{\phi}}_{T}\right)\right]$. Let $A_{G}$ be some finite grid of $\alpha$ values in $A$. For each fixed $\alpha \in A_{G}$, let $\hat{s}_{t \alpha}$ signify an empirical counterpart of $s_{t \alpha}$ (see Assumption 5) where the unknown parameter $\tilde{\phi}^{*}\left(\operatorname{or}\left(\beta^{*}, \pi^{*}\right)\right)$ is replaced by its consistent estimator under the null, $\hat{\tilde{\phi}}_{T}$. (The specific forms of $\hat{s}_{t \alpha}$ in the LMAR and GMAR examples are provided in Appendices B and C, respectively.) Set $\hat{\mathcal{I}}_{T \alpha}=T^{-1} \sum_{t=1}^{T} \hat{s}_{t \alpha} \hat{s}_{t \alpha}^{\prime}$. Now, for each $j=1, \ldots, J$ (where $J$ denotes the number of repetitions), do the following.
(i) Generate a sequence $\left\{v_{t j}\right\}_{t=1}^{T}$ of $T$ i.i.d. $N(0,1)$ random variables.
(ii) For each $\alpha \in A_{G}$, set $\hat{S}_{T \alpha}^{j}=\sum_{t=1}^{T} \hat{s}_{t \alpha} v_{t j}, \hat{Z}_{T \alpha}^{j}=\hat{\mathcal{I}}_{T \alpha}^{-1} T^{-1 / 2} \hat{S}_{T \alpha}^{j}$, and (using similar partitioning notation as before)

$$
\widehat{L R}_{T}^{j}(\alpha)=\hat{Z}_{T \vartheta \alpha}^{j \prime}\left(\hat{\mathcal{I}}_{T \alpha}^{-1}\right)_{\vartheta \vartheta}^{-1} \hat{Z}_{T \vartheta \alpha}^{j}-\inf _{\boldsymbol{\lambda}_{\vartheta} \in \Lambda_{\vartheta}}\left\{\left(\boldsymbol{\lambda}_{\vartheta}-\hat{Z}_{T \vartheta \alpha}^{j}\right)^{\prime}\left(\hat{\mathcal{I}}_{T \alpha}^{-1}\right)_{\vartheta \vartheta}^{-1}\left(\boldsymbol{\lambda}_{\vartheta}-\hat{Z}_{T \vartheta \alpha}^{j}\right)\right\} ;
$$

here the minimization of the quadratic form over the cone $\Lambda_{\vartheta}$ has to be performed numerically.
(iii) Set $\widehat{L R}_{T, A_{G}}^{j}=\max _{\alpha \in A_{G}} \widehat{L R}_{T}^{j}(\alpha)$.

This yields a sample $\left\{\widehat{L R}_{T, A_{G}}^{1}, \ldots, \widehat{L R}_{T, A_{G}}^{J}\right\}$ of $J$ realizations. An approximate $p$-value corresponding to an observed LR test statistic $L R_{T}$ is computed as $J^{-1} \sum_{j=1}^{J} \mathbf{1}\left(\widehat{L R}_{T, A_{G}}^{j}>L R_{T}\right)$ (here $\mathbf{1}(\cdot)$ denotes the indicator function). The precision of this approximation can be controlled by choosing $J$ large enough, see Hansen (1996) (in the illustration below we use $J=2500$ ).

### 4.2 Parametric bootstrap

An alternative to simulating the asymptotic null distribution is to use bootstrap. The following restricted parametric bootstrap procedure was found to work well in simulations. Related bootstrap procedures for LR tests are discussed by Qu and Zhuo (2017) and Kasahara and Shimotsu (2018) when testing for regime switching in Markov switching models.

For each $b=1, \ldots, B$ (where $B$ denotes the number of bootstrap repetitions), do the following.
(i) Generate a sequence $\left\{\varepsilon_{t}^{\star b}\right\}_{t=1}^{T}$ of $T$ i.i.d. $N(0,1)$ bootstrap errors.
(ii) To generate the $b$ th bootstrap sample, set $\left(y_{1-p}^{\star b}, \ldots, y_{0}^{\star b}\right)=\left(y_{1-p}, \ldots, y_{0}\right)$ as initial values and obtain $\left(y_{1}^{\star b}, \ldots, y_{T}^{\star b}\right)$ using the recursion $y_{t}^{\star b}=\hat{\tilde{\phi}}_{0, T}+\sum_{i=1}^{p} \hat{\tilde{\phi}}_{i, T} y_{t-i}^{\star b}+\hat{\tilde{\sigma}}_{T} \varepsilon_{t}^{\star b}$.
(iii) Using the bootstrap sample $\left(y_{1-p}^{\star b}, \ldots, y_{T}^{\star b}\right)$, perform estimation under the null and alternative hypotheses and obtain $\hat{\tilde{\phi}}_{T}^{\star b}=\left(\hat{\tilde{\phi}}_{0, T}^{\star b}, \hat{\dot{\phi}}_{1, T}^{\star b}, \ldots, \hat{\tilde{\phi}}_{p, T}^{\star b}, \hat{\tilde{\sigma}}_{T}^{\star 22}\right)$ and $\left(\hat{\alpha}_{T}^{\star b}, \hat{\beta}_{T}^{\star b}, \hat{\phi}_{T}^{\star b}, \hat{\varphi}_{T}^{\star b}\right)$. The bth bootstrap LR statistic is $L R_{T}^{\star b}=2\left[L_{T}\left(\hat{\alpha}_{T}^{\star b}, \hat{\beta}_{T}^{\star b}, \hat{\phi}_{T}^{\star b}, \hat{\varphi}_{T}^{\star b}\right)-L_{T}^{0}\left(\hat{\tilde{\phi}}_{T}^{\star b}\right)\right]$.
This yields $B$ (conditionally) independent bootstrap statistics $L R_{T}^{\star 1}, \ldots, L R_{T}^{\star B}$ and an approximate bootstrap $p$-value is computed as $B^{-1} \sum_{b=1}^{B} \mathbf{1}\left(L R_{T}^{\star b}>L R_{T}\right)$.

Our (somewhat limited) simulation results below indicate that this procedure works well in the LMAR and GMAR examples. Nevertheless, one should be cautious when using bootstrap in mixture models as the validity of bootstrap in the presence of parameters on the boundary and singular information matrices is not clear (see, e.g., Andrews (2000)). We also note that additional (unreported) simulations indicated that restricted residual bootstrap in which the bootstrap errors $\left\{\varepsilon_{t}^{\star b}\right\}_{t=1}^{T}$ are obtained by sampling with replacement from $\left\{\hat{\varepsilon}_{t}\right\}_{t=1}^{T}$, the residuals from restricted estimation (i.e., $\left.\hat{\varepsilon}_{t}=\left(y_{t}-\hat{\tilde{\phi}}_{0, T}-\sum_{i=1}^{p} \hat{\tilde{\phi}}_{i, T} y_{t-i}\right) / \hat{\tilde{\sigma}}_{T}, t=1, \ldots, T\right)$, may not work well in the present setting.

### 4.3 A small Monte Carlo study

We now study how well simulating the asymptotic null distribution and using the parametric bootstrap procedure work in finite samples. The results are presented in Tables 1 and 2. We consider two LR test statistics, one based on an estimated LMAR model, and another based on an estimated GMAR model (as in our two examples in the preceding sections). In all simulations, we use an autoregressive order $p=1, J=2500$ repetitions, $B=1000$ bootstrap repetitions, and three different sample sizes: $T=250,500$, and 1000 . For the bootstrap, fewer parameter combinations and a smaller number of repetitions are used to keep the computational load manageable.

The top parts of Tables 1 and 2 present results for size simulations. Data is generated from an $\operatorname{AR}(1)$ model (for a range of different parameter values shown in Tables 1 and 2 ) and $\operatorname{AR}(1), \operatorname{LMAR}(1)$, and $\operatorname{GMAR}(1)$ models are estimated (LMAR with $m=1 ;$ GMAR with the restriction $\tilde{\phi}_{0}=\tilde{\varphi}_{0}$; in estimation of the mixture models we use a genetic algorithm as singularity of the information matrix may render gradient based methods unreliable). Two LR test statistics are calculated based on the estimated LMAR and GMAR models, respectively, and labelled 'LMAR $L R_{T}$ ' and 'GMAR $L R_{T}$ '. For Table 1, simulation-based $p$-values are computed based on the asymptotic distributions in Section 3.5.2 and using the simulation procedure in Section 4.1; for Table 2, bootstrap $p$-values are computed as described in Section 4.2. Using nominal levels $10 \%, 5 \%$, and $1 \%$, a reject/not-reject decision is recorded. This exercise is repeated 2500 times for Table 1 and 1000 times for Table 2, and the six rightmost columns in Tables 1 and 2 present the empirical rejection frequencies (for the LMAR $L R_{T}$ and GMAR $L R_{T}$ tests and the three nominal levels used).

As can be seen from the results in Table 1 (top part), based on the simulated asymptotic null distribution the LMAR $L R_{T}$ test's size is satisfactory overall, typically being somewhat oversized for the


Table 1: Results of a Monte Carlo study, simulating the asymptotic null distribution. Empirical rejection frequencies (six rightmost columns) for LMAR $L R_{T}$ and GMAR $L R_{T}$ tests for nominal sizes $10 \%, 5 \%$, and $1 \%$. Different rows correspond to results with data generated from AR, GMAR, or LMAR models (with parameter values used shown in the table) with different sample sizes ( $T=250,500$, or 1000).

| DGP | Parameter values |  |  |  |  |  |  |  | $T$ | LMAR $L R_{T}$ |  |  | GMAR $L R_{T}$ |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  |  |  |  |  |  |  |  |  | 10\% | $5 \%$ | 1\% | 10\% | $5 \%$ | 1\% |
| AR | $\tilde{\phi}_{0}$ | $\tilde{\phi}_{1}$ |  |  |  |  |  |  |  |  |  |  |  |  |  |
| GMAR | 0 | 0.40 | 1 |  |  |  |  |  | 250 | 0.11 | 0.05 | 0.015 | 0.10 | 0.05 | 0.008 |
|  |  |  |  |  |  |  |  |  | 500 | 0.10 | 0.05 | 0.007 | 0.10 | 0.05 | 0.009 |
|  |  |  |  |  |  |  |  |  | 1000 | 0.11 | 0.06 | 0.007 | 0.11 | 0.06 | 0.008 |
|  | 0 | 0.60 | 1 |  |  |  |  |  | 250 | 0.10 | 0.05 | 0.010 | 0.11 | 0.05 | 0.008 |
|  |  |  |  |  |  |  |  |  | 500 | 0.08 | 0.04 | 0.012 | 0.10 | 0.05 | 0.012 |
|  |  |  |  |  |  |  |  |  | 1000 | 0.12 | 0.06 | 0.006 | 0.10 | 0.05 | 0.011 |
|  | 0 | 0.80 | 1 |  |  |  |  |  | 250 | 0.10 | 0.05 | 0.009 | 0.09 | 0.03 | 0.005 |
|  |  |  |  |  |  |  |  |  | 500 | 0.09 | 0.05 | 0.011 | 0.09 | 0.05 | 0.007 |
|  |  |  |  |  |  |  |  |  | 1000 | 0.11 | 0.06 | 0.013 | 0.10 | 0.05 | 0.012 |
|  | 0 | 0.90 | 1 |  |  |  |  |  | 250 | 0.11 | 0.05 | 0.006 | 0.10 | 0.05 | 0.010 |
|  |  |  |  |  |  |  |  |  | 500 | 0.11 | 0.05 | 0.015 | 0.11 | 0.06 | 0.007 |
|  |  |  |  |  |  |  |  |  | 1000 | 0.10 | 0.04 | 0.007 | 0.09 | 0.03 | 0.005 |
|  | $\tilde{\phi}_{0}$ | $\tilde{\phi}_{1}$ | $\tilde{\sigma}_{1}^{2}$ | $\tilde{\varphi}_{1}$ | $\tilde{\sigma}_{2}^{2}$ | $\alpha$ |  |  |  |  |  |  |  |  |  |
| LMAR | 0 | 0.2 | 1 | 0.8 | 1 | 0.33 |  |  | 250 | 0.22 | 0.13 | 0.03 | 0.57 | 0.42 | 0.20 |
|  |  |  |  |  |  |  |  |  | 500 | 0.32 | 0.21 | 0.09 | 0.82 | 0.71 | 0.47 |
|  |  |  |  |  |  |  |  |  | 1000 | 0.48 | 0.38 | 0.22 | 0.97 | 0.94 | 0.84 |
|  | 0 | 0.5 | 1 | 0.5 | 3 | 0.33 |  |  | 250 | 0.35 | 0.25 | 0.09 | 0.67 | 0.51 | 0.28 |
|  |  |  |  |  |  |  |  |  | 500 | 0.57 | 0.45 | 0.24 | 0.89 | 0.81 | 0.59 |
|  |  |  |  |  |  |  |  |  | 1000 | 0.78 | 0.69 | 0.48 | 0.98 | 0.97 | 0.90 |
|  | $\tilde{\phi}_{0}$ | $\tilde{\phi}_{1}$ | $\tilde{\sigma}_{1}^{2}$ | $\tilde{\varphi}_{0}$ | $\tilde{\varphi}_{1}$ | $\tilde{\sigma}_{2}^{2}$ | $\alpha_{0}$ | $\alpha_{1}$ |  |  |  |  |  |  |  |
|  | 0 | 0.2 | 1 | 0 | 0.8 | 1 | 0 | 1 | 250 | 0.78 | 0.69 | 0.42 | 0.69 | 0.55 | 0.31 |
|  |  |  |  |  |  |  |  |  | 500 | 0.98 | 0.95 | 0.88 | 0.90 | 0.83 | 0.66 |
|  |  |  |  |  |  |  |  |  | 1000 | 1.00 | 1.00 | 1.00 | 0.99 | 0.99 | 0.95 |
|  | 0 | 0.5 | 1 | 0 | 0.5 | 3 | 0 | 1 | 250 | 0.88 | 0.79 | 0.59 | 0.58 | 0.46 | 0.22 |
|  |  |  |  |  |  |  |  |  | 500 | 0.99 | 0.98 | 0.94 | 0.83 | 0.72 | 0.49 |
|  |  |  |  |  |  |  |  |  | 1000 | 1.00 | 1.00 | 1.00 | 0.97 | 0.94 | 0.84 |

Table 2: Results of a Monte Carlo study, parametric bootstrap. Empirical rejection frequencies (six rightmost columns) for LMAR $L R_{T}$ and GMAR $L R_{T}$ tests for nominal sizes $10 \%, 5 \%$, and $1 \%$. Different rows correspond to results with data generated from AR, GMAR, or LMAR models (with parameter values used shown in the table) with different sample sizes $(T=250,500$, or 1000$)$.
smallest sample size $T=250$, and somewhat conservative for the larger sample sizes $(T=500,1000)$. The parameter values used in simulation do not seem to have a large effect on the size. The GMAR $L R_{T}$ test, on the other hand, appears to be slightly oversized across all sample sizes and parameter values used. The top part of Table 2 indicates that parametric bootstrap works very well already for small sample sizes.

The lower parts of Tables 1 and 2 present results for power simulations. Data is generated either from a GMAR model or from an LMAR model (for a range of different parameter values shown in Tables 1 and 2), and empirical rejection frequencies are calculated as above. Both the LMAR $L R_{T}$ test and the GMAR $L R_{T}$ test appear to have good overall power, and the powers based on asymptotic and bootstrap $p$-values in Tables 1 and 2 are quite close to each other. As expected, when the two regimes differ more from each other, the tests have higher power, and the same happens when sample size is increased. Besides having good power against the 'right' alternatives, the tests also turn out to have decent power against 'wrong' alternatives: When data is generated from the GMAR (resp., LMAR) model, the LMAR $L R_{T}$ (resp., GMAR $L R_{T}$ ) test rejects reasonably often (the GMAR $L R_{T}$ test in particular seems capable of picking up LMAR type regime switching).

As a computational remark we note that the LMAR $L R_{T}$ and GMAR $L R_{T}$ tests and their simulation-based asymptotic $p$-values are reasonably straightforward to compute in a matter of sec-
onds using a standard, modern desktop computer. In contrast, a drawback of computing bootstrap $p$-values is that it requires repeated estimation of the mixture model under the alternative, leading to a notable increase, roughly $B$-fold, in computation time. The value of having an asymptotic distribution available would be particularly pronounced if one were to consider testing for the need of mixture models in a larger collection of different data sets in an automated manner.

## 5 Conclusions

This paper has studied the asymptotic distribution of the LR test statistic for testing a linear autoregressive model against a two-regime mixture autoregressive model. A distinguishing feature of the paper is that the regime switching probabilities are observation-dependent. Technical challenges resulting from unidentified parameters under the null, parameters on the boundary, and singularity of the information matrix were dealt with by considering an appropriately reparameterized model and higher-order expansions of the log-likelihood function. The resulting asymptotic distribution of the LR test statistic is non-standard and application-specific. One can either simulate the asymptotic null distribution or use parametric bootstrap, and a Monte Carlo study indicated the proposed tests to have good size and power properties.

The general theory of the paper was illustrated using two concrete examples, the LMAR model of Wong and Li (2001) and (a version of the) GMAR model of Kalliovirta et al. (2015). Considering other mixture AR models, as well as the general GMAR model, is left for future research. This paper was concerned with testing linearity against a two-regime model, and considering tests of $M \geq 2$ regimes versus $M+1$ regimes, similarly as in Kasahara and Shimotsu (2015) in a related setting, forms another interesting research topic. Extensions of the testing theory to mixture autoregressive models with exogeneous variables or to mixtures of vector autoregressions could also be considered.

## Appendix

## A Details for the general results

Lemma A.1. If Assumptions 2(ii) and 3(i,ii) hold, a sufficient condition for Assumption 3(iii) is that $\left\|\boldsymbol{\pi}_{\alpha}(\phi, \varphi)-\boldsymbol{\pi}_{\alpha}\left(\phi^{*}, \phi^{*}\right)\right\| \leq C h\left(\left\|(\phi, \varphi)-\left(\phi^{*}, \phi^{*}\right)\right\|_{*}\right)$ for all $(\phi, \varphi) \in \Phi \times \Phi$ and some $C<\infty$, where $h:[0, \infty) \rightarrow[0, \infty)$ is a strictly increasing function such that $h(x) \downarrow 0$ as $x \downarrow 0$ and $\|\cdot\|_{*}$ is any vector norm on $\mathbb{R}^{2 q_{2}}$.

Proof of Lemma A.1. See Supplementary Appendix D.
Proof of Lemma 1. Set $\boldsymbol{\theta}_{T \alpha}=\mathcal{I}_{\alpha}^{1 / 2} T^{1 / 2} \boldsymbol{\theta}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)$, and rewrite equation (12) evaluated at $\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)$ as
$o_{p \alpha}(1) \leq L_{T}^{\pi}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)-L_{T}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right)=\left(\mathcal{I}_{\alpha}^{1 / 2} Z_{T \alpha}\right)^{\prime} \boldsymbol{\theta}_{T \alpha}-\frac{1}{2}\left\|\boldsymbol{\theta}_{T \alpha}\right\|^{2}+R_{T}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)$
(the lower bound is due to Assumption 2(i) and the definitions of ( $\hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}$ ) and $L_{T}^{\pi}(\alpha, \beta, \pi, \varpi)$ ). As to the first term on the right hand side, note that

$$
\begin{equation*}
\sup _{\alpha \in A}\left\|\mathcal{I}_{\alpha}^{1 / 2} Z_{T \alpha}\right\|=\sup _{\alpha \in A}\left(Z_{T \alpha}^{\prime} \mathcal{I}_{\alpha} Z_{T \alpha}\right)^{1 / 2} \leq \sup _{\alpha \in A}\left(\lambda_{\max }\left(\mathcal{I}_{\alpha}\right)\right)^{1 / 2} \sup _{\alpha \in A}\left\|Z_{T \alpha}\right\| \leq C \sup _{\alpha \in A}\left\|Z_{T \alpha}\right\|=O_{p}(1) \tag{24}
\end{equation*}
$$

where the latter inequality holds with some finite $C$ in view of Assumption 5(iii), and the last equality will be justified below. Thus

$$
\begin{equation*}
\mathcal{I}_{\alpha}^{1 / 2} Z_{T \alpha}=O_{p \alpha}(1) \tag{25}
\end{equation*}
$$

a result which also implies $\left\|\left(\mathcal{I}_{\alpha}^{1 / 2} Z_{T \alpha}\right)^{\prime} \boldsymbol{\theta}_{T \alpha}\right\| \leq\left\|\boldsymbol{\theta}_{T \alpha}\right\|\left\|\mathcal{I}_{\alpha}^{1 / 2} Z_{T \alpha}\right\|=\left\|\boldsymbol{\theta}_{T \alpha}\right\| O_{p \alpha}(1)$.
Next consider the third term on the right hand side of (23), where the assumption $\left(\hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)=$ $\left(\beta^{*}, \pi^{*}, 0\right)+o_{p \alpha}(1)$ (see Assumption $3(\mathrm{iii})$ ) allows us to choose a sequence $\left\{\gamma_{T}, T \geq 1\right\}$ of (nonrandom) positive scalars converging to zero slowly enough to ensure that $P\left(\sup _{\alpha \in A} \|\left(\hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)-\right.$ $\left.\left(\beta^{*}, \pi^{*}, 0\right) \| \leq \gamma_{T}\right) \rightarrow 1$, and with this sequence $\left\{\gamma_{T}, T \geq 1\right\}$, Assumption 5(iv) implies that

$$
\begin{equation*}
\left|R_{T}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)\right|=\left(1+\left\|T^{1 / 2} \boldsymbol{\theta}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)\right\|\right)^{2} o_{p \alpha}(1) \tag{26}
\end{equation*}
$$

(cf. Pakes and Pollard (1989, proof of Thm 3.3)). Here $\left\|T^{1 / 2} \boldsymbol{\theta}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)\right\|=\left\|\mathcal{I}_{\alpha}^{-1 / 2} \boldsymbol{\theta}_{T \alpha}\right\|$ so that, as $0<\inf _{\alpha \in A} \lambda_{\min }\left(\mathcal{I}_{\alpha}\right)$ by Assumption 5(iii),

$$
\begin{equation*}
\left|R_{T}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)\right| \leq o_{p \alpha}(1)+\left\|\boldsymbol{\theta}_{T \alpha}\right\| o_{p \alpha}(1)+\left\|\boldsymbol{\theta}_{T \alpha}\right\|^{2} o_{p \alpha}(1) \tag{27}
\end{equation*}
$$

Combining the results above (see (23), the inequality below (25), and (27)), organizing terms, and absorbing $\left\|\boldsymbol{\theta}_{T \alpha}\right\| o_{p \alpha}(1)$ into $\left\|\boldsymbol{\theta}_{T \alpha}\right\| O_{p \alpha}(1)$ leads to

$$
\begin{equation*}
0 \leq\left\|\boldsymbol{\theta}_{T \alpha}\right\| O_{p \alpha}(1)+o_{p \alpha}(1)+\left(o_{p \alpha}(1)-\frac{1}{2}\right)\left\|\boldsymbol{\theta}_{T \alpha}\right\|^{2} \tag{28}
\end{equation*}
$$

In Supplementary Appendix D we show that the last term on the right hand side of (28) is dominated by $-\frac{1}{4}\left\|\boldsymbol{\theta}_{T \alpha}\right\|^{2}+o_{p \alpha}(1)$ so that (absorbing constants into the $O_{p \alpha}(1)$ and $o_{p \alpha}(1)$ terms) $\left\|\boldsymbol{\theta}_{T \alpha}\right\|^{2} \leq$ $2\left\|\boldsymbol{\theta}_{T \alpha}\right\| O_{p \alpha}(1)+o_{p \alpha}(1)$. Denoting the $O_{p \alpha}(1)$ term on the majorant side with $\xi_{T \alpha}$ and reorganizing,

$$
\left(\left\|\boldsymbol{\theta}_{T \alpha}\right\|-\xi_{T \alpha}\right)^{2}=\left\|\boldsymbol{\theta}_{T \alpha}\right\|^{2}-2\left\|\boldsymbol{\theta}_{T \alpha}\right\| \xi_{T \alpha}+\xi_{T \alpha}^{2} \leq \xi_{T \alpha}^{2}+o_{p \alpha}(1)=O_{p \alpha}(1)
$$

Taking square roots yields $\left\|\boldsymbol{\theta}_{T \alpha}\right\|=O_{p \alpha}(1)$ so that $\left\|T^{1 / 2} \boldsymbol{\theta}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)\right\|=O_{p \alpha}(1)$, establishing part (i) of the lemma. Part (ii) follows from $\left\|\boldsymbol{\theta}_{T \alpha}\right\|=O_{p \alpha}(1)$ and (27). Part (iii) follows from part (ii).

To complete the proof of Lemma 1, we now justify the last equality in (24). By Assumption 5(ii), $T^{-1 / 2} S_{T \bullet} \Rightarrow S \bullet$ in $\mathcal{C}\left(A, \mathbb{R}^{r}\right)$, and by the continuous mapping theorem (justification in Supplementary Appendix D), $Z_{T \alpha}=\mathcal{I}_{\alpha}^{-1} T^{-1 / 2} S_{T \alpha}$ converges weakly in $\mathcal{C}\left(A, \mathbb{R}^{r}\right)$ to a mean zero $\mathbb{R}^{r}$-valued Gaussian process $Z_{\alpha}=\mathcal{I}_{\alpha}^{-1} S_{\alpha}$ whose sample paths are continuous in $\alpha$ with probability one and that has $E\left[Z_{\alpha} Z_{\alpha}^{\prime}\right]=\mathcal{I}_{\alpha}^{-1}$ for all $\alpha \in A$. A further application of the continuous mapping theorem (justification in Supplementary Appendix D) implies that $\sup _{\alpha \in A}\left\|Z_{T \alpha}\right\|$ converges in distribution in $\mathbb{R}$ and, as all probability measures on $\mathbb{R}$ are tight, the limit must be tight. This justifies the last equality in (24).

Proof of Lemma 2. By the definition of $\hat{\boldsymbol{\lambda}}_{T \alpha q}$, the fact that $\mathbf{0} \in \Theta_{\alpha, T}$, and (25),

$$
\left\|\mathcal{I}_{\alpha}^{1 / 2}\left(\hat{\boldsymbol{\lambda}}_{T \alpha q}-Z_{T \alpha}\right)\right\|^{2}=\left(\hat{\boldsymbol{\lambda}}_{T \alpha q}-Z_{T \alpha}\right)^{\prime} \mathcal{I}_{\alpha}\left(\hat{\boldsymbol{\lambda}}_{T \alpha q}-Z_{T \alpha}\right) \leq Z_{T \alpha}^{\prime} \mathcal{I}_{\alpha} Z_{T \alpha}+o_{p \alpha}(1)=O_{p \alpha}(1)
$$

implying that $\mathcal{I}_{\alpha}^{1 / 2}\left(\hat{\boldsymbol{\lambda}}_{T \alpha q}-Z_{T \alpha}\right)=O_{p \alpha}(1)$. A further use of $(25)$ and the condition $0<\inf _{\alpha \in A} \lambda_{\min }\left(\mathcal{I}_{\alpha}\right)$ in Assumption $5(\mathrm{iii})$ yields $\hat{\boldsymbol{\lambda}}_{T \alpha q}=O_{p \alpha}(1)$. We next show that $R_{T}\left(\alpha, \hat{\beta}_{T \alpha q}, \hat{\pi}_{T \alpha q}, \hat{\varpi}_{T \alpha q}\right)=o_{p \alpha}(1)$. First, $\hat{\boldsymbol{\lambda}}_{T \alpha q}=T^{1 / 2} \boldsymbol{\theta}\left(\alpha, \hat{\beta}_{T \alpha q}, \hat{\pi}_{T \alpha q}, \hat{\varpi}_{T \alpha q}\right)$ and $\hat{\boldsymbol{\lambda}}_{T \alpha q}=O_{p \alpha}(1)$ imply that $\boldsymbol{\theta}\left(\alpha, \hat{\beta}_{T \alpha q}, \hat{\pi}_{T \alpha q}, \hat{\varpi}_{T \alpha q}\right)=$ $o_{p \alpha}(1)$. Second, to show that $\left(\hat{\beta}_{T \alpha q}, \hat{\pi}_{T \alpha q}, \hat{\varpi}_{T \alpha q}\right)=\left(\beta^{*}, \pi^{*}, 0\right)+o_{p \alpha}(1)$, pick arbitrary $\epsilon, \delta>0$, and conclude from Assumption $5(\mathrm{i})(\mathrm{b})$ that $\inf _{\alpha \in A} \inf _{(\beta, \pi, \varpi) \in B \times \Pi_{\alpha}:\left\|(\beta, \pi, \varpi)-\left(\beta^{*}, \pi^{*}, 0\right)\right\| \geq \epsilon}\|\boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)\| \geq$ $\delta_{\epsilon}>0$ for some $\delta_{\epsilon}$. Now, as $\boldsymbol{\theta}\left(\alpha, \hat{\beta}_{T \alpha q}, \hat{\pi}_{T \alpha q}, \hat{\varpi}_{T \alpha q}\right)=o_{p \alpha}(1)$, we can find a $T_{\delta, \delta_{\epsilon}}$ such that for all $T \geq T_{\delta, \delta_{\epsilon}}, P\left(\sup _{\alpha \in A}\left\|\boldsymbol{\theta}\left(\alpha, \hat{\beta}_{T \alpha q}, \hat{\pi}_{T \alpha q}, \hat{\varpi}_{T \alpha q}\right)\right\|<\delta_{\epsilon}\right)>1-\delta$. Note that whenever the event
$\left\{\sup _{\alpha \in A}\left\|\boldsymbol{\theta}\left(\alpha, \hat{\beta}_{T \alpha q}, \hat{\pi}_{T \alpha q}, \hat{\varpi}_{T \alpha q}\right)\right\|<\delta_{\epsilon}\right\}$ occurs, the event $\left\{\sup _{\alpha \in A}\left\|\left(\hat{\beta}_{T \alpha q}, \hat{\pi}_{T \alpha q}, \hat{\varpi}_{T \alpha q}\right)-\left(\beta^{*}, \pi^{*}, 0\right)\right\|<\right.$ $\epsilon\}$ must also occur (if, on the contrary, $\left\|\left(\hat{\beta}_{T \alpha q}, \hat{\pi}_{T \alpha q}, \hat{\varpi}_{T \alpha q}\right)-\left(\beta^{*}, \pi^{*}, 0\right)\right\| \geq \epsilon$ for some $\alpha \in A$, then necessarily $\left.\left\|\boldsymbol{\theta}\left(\alpha, \hat{\beta}_{T \alpha q}, \hat{\pi}_{T \alpha q}, \hat{\varpi}_{T \alpha q}\right)\right\| \geq \delta_{\epsilon}\right)$ ). We therefore conclude that for all $T \geq T_{\delta, \delta_{\epsilon}}$, $1-\delta<P\left(\sup _{\alpha \in A}\left\|\boldsymbol{\theta}\left(\alpha, \hat{\beta}_{T \alpha q}, \hat{\pi}_{T \alpha q}, \hat{\varpi}_{T \alpha q}\right)\right\|<\delta_{\epsilon}\right) \leq P\left(\sup _{\alpha \in A}\left\|\left(\hat{\beta}_{T \alpha q}, \hat{\pi}_{T \alpha q}, \hat{\varpi}_{T \alpha q}\right)-\left(\beta^{*}, \pi^{*}, 0\right)\right\|<\epsilon\right)$, so that $\sup _{\alpha \in A}\left\|\left(\hat{\beta}_{T \alpha q}, \hat{\pi}_{T \alpha q}, \hat{\varpi}_{T \alpha q}\right)-\left(\beta^{*}, \pi^{*}, 0\right)\right\|=o_{p}(1)$, as desired. Third, as $\left(\hat{\beta}_{T \alpha q}, \hat{\pi}_{T \alpha q}, \hat{\varpi}_{T \alpha q}\right)=$ $\left(\beta^{*}, \pi^{*}, 0\right)+o_{p \alpha}(1)$, using the same argument as in the derivation of (26) now leads to

$$
\left|R_{T}\left(\alpha, \hat{\beta}_{T \alpha q}, \hat{\pi}_{T \alpha q}, \hat{\varpi}_{T \alpha q}\right)\right|=\left(1+\left\|T^{1 / 2} \boldsymbol{\theta}\left(\alpha, \hat{\beta}_{T \alpha q}, \hat{\pi}_{T \alpha q}, \hat{\varpi}_{T \alpha q}\right)\right\|\right)^{2} o_{p \alpha}(1)=\left(1+\left\|\hat{\boldsymbol{\lambda}}_{T \alpha q}\right\|\right)^{2} o_{p \alpha}(1)
$$

Together with $\hat{\boldsymbol{\lambda}}_{T \alpha q}=O_{p \alpha}(1)$ established above, this yields $R_{T}\left(\alpha, \hat{\beta}_{T \alpha q}, \hat{\pi}_{T \alpha q}, \hat{\varpi}_{T \alpha q}\right)=o_{p \alpha}(1)$.
Now, by expansion (13), the definitions of ( $\hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}$ ) and $\hat{\boldsymbol{\lambda}}_{T \alpha q}$, and making use of results $R_{T}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)=o_{p \alpha}(1)$ and $R_{T}\left(\alpha, \hat{\beta}_{T \alpha q}, \hat{\pi}_{T \alpha q}, \hat{\varpi}_{T \alpha q}\right)=o_{p \alpha}(1)$,

$$
\begin{aligned}
o_{p \alpha}(1) \leq & L_{T}^{\pi}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)-L_{T}^{\pi}\left(\alpha, \hat{\beta}_{T \alpha q}, \hat{\pi}_{T \alpha q}, \hat{\varpi}_{T \alpha q}\right) \\
= & \frac{1}{2}\left(\hat{\boldsymbol{\lambda}}_{T \alpha q}-Z_{T \alpha}\right)^{\prime} \mathcal{I}_{\alpha}\left(\hat{\boldsymbol{\lambda}}_{T \alpha q}-Z_{T \alpha}\right) \\
& -\frac{1}{2}\left[T^{1 / 2} \boldsymbol{\theta}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)-Z_{T \alpha}\right]^{\prime} \mathcal{I}_{\alpha}\left[T^{1 / 2} \boldsymbol{\theta}\left(\alpha, \hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)-Z_{T \alpha}\right]+o_{p \alpha}(1) \\
\leq & o_{p \alpha}(1),
\end{aligned}
$$

implying, by (20), the desired result.
Proof of Lemma 3. For any vectors $a, b \in \mathbb{R}^{r}$, we denote $\|a-b\|_{\mathcal{I}_{\alpha}^{-1}}=\left[(a-b)^{\prime} \mathcal{I}_{\alpha}(a-b)\right]^{1 / 2}$, and for any point $p \in \mathbb{R}^{r}$ and a set $S \subset \mathbb{R}^{r}$, we define $\|p-S\|_{\mathcal{I}_{\alpha}^{-1}}$ via $\|p-S\|_{\mathcal{I}_{\alpha}^{-1}}^{2}=\inf _{s \in S}\|s-p\|_{\mathcal{I}_{\alpha}^{-1}}^{2}=$ $\inf _{s \in S}\left[(s-p)^{\prime} \mathcal{I}_{\alpha}(s-p)\right]$. With this notation, we need to prove that $\left\|Z_{T \alpha}-\Theta_{\alpha, T}\right\|_{\mathcal{I}_{\alpha}^{-1}}^{2}=\left\|Z_{T \alpha}-\Lambda\right\|_{\mathcal{I}_{\alpha}^{-1}}^{2}+$ $o_{p \alpha}(1)$. First note that, because $\Lambda$ is a cone, we have, for any $T$,

$$
\left\|Z_{T \alpha}-\Lambda\right\|_{\mathcal{I}_{\alpha}^{-1}}^{2}=T \inf _{\boldsymbol{\lambda} \in \Lambda}\left\{\left(T^{-1 / 2} \boldsymbol{\lambda}-T^{-1 / 2} Z_{T \alpha}\right)^{\prime} \mathcal{I}_{\alpha}\left(T^{-1 / 2} \boldsymbol{\lambda}-T^{-1 / 2} Z_{T \alpha}\right)\right\}=T\left\|T^{-1 / 2} Z_{T \alpha}-\Lambda\right\|_{\mathcal{I}_{\alpha}^{-1}}^{2}
$$

Similarly, by the definitions of $\Theta_{\alpha, T}$ and $\Theta_{\alpha}$,

$$
\left\|Z_{T \alpha}-\Theta_{\alpha, T}\right\|_{\mathcal{I}_{\alpha}^{-1}}^{2}=\inf _{\boldsymbol{\lambda} \in \Theta_{\alpha}}\left\{\left(T^{1 / 2} \boldsymbol{\lambda}-Z_{T \alpha}\right)^{\prime} \mathcal{I}_{\alpha}\left(T^{1 / 2} \boldsymbol{\lambda}-Z_{T \alpha}\right)\right\}=T\left\|T^{-1 / 2} Z_{T \alpha}-\Theta_{\alpha}\right\|_{\mathcal{I}_{\alpha}^{-1}}^{2}
$$

Now let $G_{T}(\alpha, \boldsymbol{x})=T\left\|\boldsymbol{x}-\Theta_{\alpha}\right\|_{\mathcal{I}_{\alpha}^{-1}}^{2}-T\|\boldsymbol{x}-\Lambda\|_{\mathcal{I}_{\alpha}^{-1}}^{2}$ define a (non-random) function on $A \times \mathbb{R}^{r}$. Because $\left\{\Theta_{\alpha}, \alpha \in A\right\}$ is locally uniformly equal to the cone $\Lambda$, we can find a $\delta>0$ such that $\Theta_{\alpha} \cap(-\delta, \delta)^{r}=\Lambda \cap(-\delta, \delta)^{r}$ for all $\alpha \in A$. Furthermore, $\mathbf{0} \in \Theta_{\alpha}$ and $\mathbf{0} \in \Lambda$ (here $\mathbf{0} \in \mathbb{R}^{r}$ ). Therefore, we can find a neighborhood $N_{\mathbf{0}}$ of $\mathbf{0}$ such that for all $(\alpha, \boldsymbol{x}) \in A \times N_{\mathbf{0}}$,

$$
G_{T}(\alpha, \boldsymbol{x})=T\left\|\boldsymbol{x}-\Theta_{\alpha} \cap(-\delta, \delta)^{r}\right\|_{\mathcal{I}_{\alpha}^{-1}}^{2}-T\left\|\boldsymbol{x}-\Lambda \cap(-\delta, \delta)^{r}\right\|_{\mathcal{I}_{\alpha}^{-1}}^{2}=0
$$

Define $G_{T}(\alpha)$, a random function of $\alpha$, as $G_{T}(\alpha)=G_{T}\left(\alpha, T^{-1 / 2} Z_{T \alpha}\right)$. In the proof of Lemma 1 it was shown that $\sup _{\alpha \in A}\left\|Z_{T \alpha}\right\|=O_{p}(1)($ see $(24))$ so that $T^{-1 / 2}\left\|Z_{T \alpha}\right\|=o_{p \alpha}(1)$. Therefore, for all $\epsilon>0$,

$$
\begin{aligned}
P\left(\sup _{\alpha \in A}\left|G_{T}(\alpha)\right|>\epsilon\right) & \leq P\left(\sup _{\alpha \in A}\left|G_{T}(\alpha)\right|>\epsilon ; T^{-1 / 2} \sup _{\alpha \in A}\left\|Z_{T \alpha}\right\| \in N_{0}\right) \\
& +P\left(\sup _{\alpha \in A}\left|G_{T}(\alpha)\right|>\epsilon ; T^{-1 / 2} \sup _{\alpha \in A}\left\|Z_{T \alpha}\right\| \notin N_{0}\right)
\end{aligned}
$$

The first term on the right hand side equals zero as $G_{T}(\alpha, \boldsymbol{x})=0$ for all $(\alpha, \boldsymbol{x}) \in A \times N_{\mathbf{0}}$ and the second one is dominated by $P\left(T^{-1 / 2} \sup _{\alpha \in A}\left\|Z_{T \alpha}\right\| \notin N_{0}\right)$ which converges to zero because $T^{-1 / 2}\left\|Z_{T \alpha}\right\|=$ $o_{p \alpha}(1)$. Thus $\sup _{\alpha \in A}\left|G_{T}(\alpha)\right|=o_{p}(1)$, implying $\left\|Z_{T \alpha}-\Theta_{\alpha, T}\right\|_{\mathcal{I}_{\alpha}^{-1}}^{2}=\left\|Z_{T \alpha}-\Lambda\right\|_{\mathcal{I}_{\alpha}^{-1}}^{2}+o_{p \alpha}(1)$.

Proof of Lemma 4. It was shown in the proof of Lemma 1 that $Z_{T_{\bullet}} \Rightarrow Z_{\bullet}$ in $\mathcal{C}\left(A, \mathbb{R}^{r}\right)$. Therefore also $\left(Z_{T_{\bullet}}, \mathcal{I}_{\bullet}\right) \Rightarrow\left(Z_{\bullet}, \mathcal{I}_{\bullet}\right)$ in $\mathcal{C}\left(A, \mathbb{R}^{r}\right) \times\left\{\mathcal{I}_{\alpha}\right\}$ (Billingsley (1999, Thm 3.9)). As the function $g: \mathcal{C}\left(A, \mathbb{R}^{r}\right) \times$ $\left\{\mathcal{I}_{\alpha}\right\} \rightarrow \mathcal{B}(A, \mathbb{R})$ mapping $\left(x_{\bullet}, \mathcal{I}_{\bullet}\right)\left(\in \mathcal{C}\left(A, \mathbb{R}^{r}\right) \times\left\{\mathcal{I}_{\alpha}\right\}\right)$ to $x_{\bullet}^{\prime} \mathcal{I}_{\bullet} x_{\bullet}-\inf _{\boldsymbol{\lambda} \in \Lambda}\left\{\left(\boldsymbol{\lambda}-x_{\bullet}\right)^{\prime} \mathcal{I}_{\bullet}\left(\boldsymbol{\lambda}-x_{\bullet}\right)\right\}$ is continuous (justification in Supplementary Appendix D), the continuous mapping theorem is applicable. This, together with Billingsley (1999, Thm 3.1) (for which it is necessary that the remainder term in (22) is $o_{p \alpha}(1)$ and not only $\left.o_{p}(1)\right)$, implies that $2\left[L_{T}^{\pi}\left(\bullet, \hat{\beta}_{T_{\bullet}}, \hat{\pi}_{T_{\bullet}}, \hat{\varpi}_{T_{\bullet}}\right)-L_{T}^{\pi}\left(\bullet, \beta^{*}, \pi^{*}, 0\right)\right] \Rightarrow$ $Z_{\bullet}^{\prime} \mathcal{I}_{\bullet} Z_{\bullet}-\inf _{\boldsymbol{\lambda} \in \Lambda}\left\{\left(\boldsymbol{\lambda}-Z_{\bullet}\right)^{\prime} \mathcal{I}_{\bullet}\left(\boldsymbol{\lambda}-Z_{\bullet}\right)\right\}$.

Proofs of Lemmas 5 \& 6. Proof of Lemma 5 consists of straightforward matrix algebra. The arguments required for proving Lemma 6 are standard. For details, see Supplementary Appendix D.

Proof of Theorem 1. Under Assumption 8, the random process $S_{\theta \alpha}^{\prime} \mathcal{I}_{\theta \theta \alpha}^{-1} S_{\theta \alpha}$ in Lemma 5 coincides with the random variable $S^{0 \prime}\left(\mathcal{I}^{0}\right)^{-1} S^{0}$ in Lemma 6. Therefore the expression of $L R_{T}(\alpha)$ in (9), Lemmas 4,5 , and 6 , and Billingsley (1999, Thm 3.1) (for which it is necessary that the remainder term in (9) is $o_{p \alpha}(1)$ and not only $\left.o_{p}(1)\right)$ imply the weak convergence result for $L R_{T}(\alpha)$. The result for $L R_{T}$ follows from the continuous mapping theorem.

## B Details for the LMAR example

In this appendix it appears convenient to denote $\alpha_{1, t}^{L}$ instead of $\alpha_{t}^{L}$ and to set $\alpha_{2, t}^{L}=1-\alpha_{1, t}^{L}$. In some cases we also include the argument $\alpha$ and denote $\alpha_{1, t}^{L}(\alpha)$ and $\alpha_{2, t}^{L}(\alpha)$. The same notation is employed in the Supplementary Appendix and a similar modification is used in the case of the GMAR model.

Assumptions 1-4. Assumption 1(i) is assumed to hold, 1(ii) holds as $A$ is compact, and 1(iii) holds by the definition of the mixing weight. For the verification of Assumption 2, see the GMAR example in Appendix C; the LMAR case is treated there as well. To verify Assumption 3, note first that conditions (i) and (ii) clearly hold, and for condition (iii), we have $\boldsymbol{\pi}_{\alpha}(\phi, \varphi)-\boldsymbol{\pi}_{\alpha}\left(\phi^{*}, \phi^{*}\right)=$ $(\phi, \phi-\varphi)-\left(\phi^{*}, 0\right)=\left(\phi-\phi^{*},\left(\phi-\phi^{*}\right)-\left(\varphi-\phi^{*}\right)\right)$. Choosing $\|x\|_{*}=\|x\|_{1}=\sum_{i=1}^{2 q}\left|x_{i}\right|$ and using the triangle inequality it is straightforward to check that $\left\|\left(\phi-\phi^{*},\left(\phi-\phi^{*}\right)-\left(\varphi-\phi^{*}\right)\right)\right\| \leq 2\left\|\phi-\phi^{*}\right\|+$ $\left\|\varphi-\phi^{*}\right\| \leq 2\left\|(\phi, \varphi)-\left(\phi^{*}, \phi^{*}\right)\right\|_{1}$. Thus condition (iii) holds by Lemma A.1. Regarding Assumption 4, as $\alpha_{1, t}^{L}$ does not depend on $(\phi, \varphi)$ and $\boldsymbol{\pi}_{\alpha}^{-1}(\pi, \varpi)=(\pi, \pi-\varpi)$, the required differentiability conditions hold for all positive integers $k$.

Assumption 5: Computation of the required derivatives. As $\alpha_{1, t}^{L}$ does not depend on $(\phi, \varphi)$ and $\boldsymbol{\pi}_{\alpha}^{-1}(\pi, \varpi)=(\pi, \pi-\varpi)$, the quantities $f_{2, t}^{\pi}(\alpha, \pi, \varpi)$ and $l_{t}^{\pi}(\alpha, \pi, \varpi)$ take the form

$$
\begin{aligned}
f_{2, t}^{\pi}(\alpha, \pi, \varpi) & =\alpha_{1, t}^{L}(\alpha) f_{t}(\pi)+\left(1-\alpha_{1, t}^{L}(\alpha)\right) f_{t}(\pi-\varpi), \\
l_{t}^{\pi}(\alpha, \pi, \varpi) & =\log \left[\alpha_{1, t}^{L}(\alpha) f_{t}(\pi)+\left(1-\alpha_{1, t}^{L}(\alpha)\right) f_{t}(\pi-\varpi)\right] .
\end{aligned}
$$

Straightforward differentiation yields the following expressions for the first and second partial derivatives with respect to $\pi$ and $\varpi$ (recall that $\nabla$ and $\nabla^{2}$ denote first and second order differentiation with respect to the indicated parameters; $\nabla f_{t}(\cdot)$ denotes differentiation of $f_{t}(\cdot)$ in (3) with respect to $\tilde{\phi}$ ):

$$
\begin{aligned}
\nabla_{\pi} l_{t}^{\pi}(\alpha, \pi, \varpi) & =\left[\alpha_{1, t}^{L}(\alpha) \nabla f_{t}(\pi)+\alpha_{2, t}^{L}(\alpha) \nabla f_{t}(\pi-\varpi)\right] / f_{2, t}^{\pi}(\alpha, \pi, \varpi), \\
\nabla_{\varpi} l_{t}^{\pi}(\alpha, \pi, \varpi) & =\left[-\alpha_{2, t}^{L}(\alpha) \nabla f_{t}(\pi-\varpi)\right] / f_{2, t}^{\pi}(\alpha, \pi, \varpi),
\end{aligned}
$$

$$
\begin{aligned}
\nabla_{\pi \pi^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi) & =\alpha_{1, t}^{L}(\alpha)\left(\frac{\nabla^{2} f_{t}(\pi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)}-\frac{\nabla f_{t}(\pi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)} \frac{\alpha_{1, t}^{L}(\alpha) \nabla^{\prime} f_{t}(\pi)+\alpha_{2, t}^{L}(\alpha) \nabla^{\prime} f_{t}(\pi-\varpi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)}\right) \\
& +\alpha_{2, t}^{L}(\alpha)\left(\frac{\nabla^{2} f_{t}(\pi-\varpi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)}-\frac{\nabla f_{t}(\pi-\varpi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)} \frac{\alpha_{1, t}^{L}(\alpha) \nabla^{\prime} f_{t}(\pi)+\alpha_{2, t}^{L}(\alpha) \nabla^{\prime} f_{t}(\pi-\varpi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)}\right), \\
\nabla_{\pi \varpi^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi) & =\alpha_{1, t}^{L}(\alpha)\left(-\frac{\nabla f_{t}(\pi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)}-\frac{-\alpha_{2, t}^{L}(\alpha) \nabla^{\prime} f_{t}(\pi-\varpi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)}\right) \\
& +\alpha_{2, t}^{L}(\alpha)\left(-\frac{\nabla^{2} f_{t}(\pi-\varpi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)}+\frac{\nabla f_{t}(\pi-\varpi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)} \frac{\alpha_{2, t}^{L}(\alpha) \nabla^{\prime} f_{t}(\pi-\varpi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)}\right), \\
\nabla_{\varpi \varpi^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi) & =-\alpha_{2, t}^{L}(\alpha)\left(-\frac{\nabla^{2} f_{t}(\pi-\varpi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)}+\frac{\nabla f_{t}(\pi-\varpi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)} \frac{\alpha_{2, t}^{L}(\alpha) \nabla^{\prime} f_{t}(\pi-\varpi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)}\right) .
\end{aligned}
$$

The corresponding expressions evaluated at $(\alpha, \pi, \varpi)=\left(\alpha, \pi^{*}, 0\right)$ take the form

$$
\begin{align*}
\nabla_{(\pi, \varpi)_{t}}^{\pi}\left(\alpha, \pi^{*}, 0\right) & =\left(1,-\left(1-\alpha_{1, t}^{L}(\alpha)\right)\right) \otimes \frac{\nabla f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)},  \tag{29}\\
\nabla_{\pi \pi}^{2} l_{t}^{\pi}\left(\alpha, \pi^{*}, 0\right) & =\frac{\nabla^{2} f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)}-\frac{\nabla f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)} \frac{\nabla^{\prime} f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)},  \tag{30}\\
\nabla_{\pi \varpi \prime}^{2} l_{t}^{\pi}\left(\alpha, \pi^{*}, 0\right) & =-\left(1-\alpha_{1, t}^{L}(\alpha)\right) \frac{\nabla^{2} f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)}+\left(1-\alpha_{1, t}^{L}(\alpha)\right) \frac{\nabla f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)} \frac{\nabla^{\prime} f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)},  \tag{31}\\
\nabla_{\varpi \varpi \prime}^{2} l_{t}^{\pi}\left(\alpha, \pi^{*}, 0\right) & =\left(1-\alpha_{1, t}^{L}(\alpha)\right) \frac{\nabla^{2} f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)}-\left(1-\alpha_{1, t}^{L}(\alpha)\right)^{2} \frac{\nabla f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)} \frac{\nabla^{\prime} f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)} . \tag{32}
\end{align*}
$$

Assumption 5: Verifying the assumption. Omitting the unnecessary $\beta$ from $\boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)$, we have $\boldsymbol{\theta}(\alpha, \pi, \varpi)=\left(\pi-\pi^{*}, \varpi\right)=(\theta, \vartheta)$ so that part (i) is clearly satisfied with the parameter space

$$
\begin{aligned}
\Theta_{\alpha}=\Theta & =\left\{\boldsymbol{\theta}=(\theta, \vartheta) \in \mathbb{R}^{2 q_{2}}: \theta=\pi-\pi^{*}, \vartheta=\varpi \text { for some }(\pi, \varpi) \in \Pi\right\} \\
& =\left\{\boldsymbol{\theta}=(\theta, \vartheta) \in \mathbb{R}^{2 q_{2}}: \theta=\phi-\phi^{*}, \vartheta=\phi-\varphi \text { for some }(\phi, \varphi) \in \Phi \times \Phi\right\}
\end{aligned}
$$

independent of $\alpha$ and with $0\left(\in \mathbb{R}^{2 q_{2}}\right)$ an interior point of $\Theta$. The first two requirements in part (ii) are similarly clear, whereas the third requirement follows from the continuity of $\alpha_{1, t}^{L}(\alpha)$ in $\alpha$. The weak convergence requirement in part (ii) is verified in Supplementary Appendix E.1.

For part (iii), first consider the positive definiteness of $\mathcal{I}_{\alpha}$ for a fixed $\alpha \in A$. It suffices to show that

$$
\begin{equation*}
a^{\prime} \nabla f_{t}\left(\pi^{*}\right) / f_{t}\left(\pi^{*}\right)=\alpha_{2, t}^{L}(\alpha) b^{\prime} \nabla f_{t}\left(\pi^{*}\right) / f_{t}\left(\pi^{*}\right) \quad \text { a.s. } \tag{33}
\end{equation*}
$$

only if $a=\left(a_{1}, \ldots, a_{q_{2}}\right)=0$ and $b=\left(b_{1}, \ldots, b_{q_{2}}\right)=0$. For brevity, denote $g_{t}(\pi)=\left[y_{t}-\left(\pi_{1}+\pi_{2} y_{t-1}+\right.\right.$ $\left.\left.\cdots+\pi_{p+1} y_{t-p}\right)\right] / \pi_{p+2}^{1 / 2}=\left[y_{t}-\left(\tilde{\phi}_{0}+\tilde{\phi}_{1} y_{t-1}+\cdots+\tilde{\phi}_{p} y_{t-p}\right)\right] / \tilde{\sigma}_{1}$ and $\boldsymbol{z}_{t-1}=\left(1, \boldsymbol{y}_{t-1}\right)$ so that $g_{t}\left(\pi^{*}\right)=\varepsilon_{t}$. Straightforward differentiation yields

$$
\frac{\nabla f_{t}(\pi)}{f_{t}(\pi)}=\left[\begin{array}{c}
\boldsymbol{z}_{t-1} g_{t}(\pi) / \tilde{\sigma}_{1}  \tag{34}\\
\left(g_{t}^{2}(\pi)-1\right) /\left(2 \tilde{\sigma}_{1}^{2}\right)
\end{array}\right] \text { so that } \frac{\nabla f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)}=\left[\begin{array}{c}
\boldsymbol{z}_{t-1} \varepsilon_{t} / \tilde{\sigma}_{1}^{*} \\
\left(\varepsilon_{t}^{2}-1\right) /\left(2 \tilde{\sigma}_{1}^{* 2}\right)
\end{array}\right]
$$

(where $\tilde{\sigma}_{1}^{2}=\pi_{p+2}$ ). Multiplying both sides of equation (33) by $\left(\varepsilon_{t}^{2}-1\right) 2 \tilde{\sigma}_{1}^{* 2}$, taking expectations conditional on $\mathcal{F}_{t-1}$, and making use of the fact that odd moments of the normal distribution are zero, yields $a_{q_{2}} E\left[\left(\varepsilon_{t}^{2}-1\right)^{2}\right]=\alpha_{2, t}^{L}(\alpha) b_{q_{2}} E\left[\left(\varepsilon_{t}^{2}-1\right)^{2}\right]$ a.s. Because $\alpha_{2, t}^{L}(\alpha) \neq 0$ and not equal to a constant (see Section 3.1.1), it follows that $a_{q_{2}}=b_{q_{2}}=0$. Therefore, equation (33) (multiplied by $\sigma_{1}^{*}$ ) now reduces to $\left(a_{1}, \ldots, a_{q_{2}-1}\right)^{\prime} \boldsymbol{z}_{t-1} \varepsilon_{t}=\alpha_{2, t}^{L}(\alpha)\left(b_{1}, \ldots, b_{q_{2}-1}\right)^{\prime} \boldsymbol{z}_{t-1} \varepsilon_{t}$ a.s. Multiplying this equation by $\varepsilon_{t}$,
taking expectations conditional on $\mathcal{F}_{t-1}$, and dividing by $E\left[\varepsilon_{t}^{2}\right]=\sigma_{1}^{* 2}$ yields $\left(a_{1}, \ldots, a_{q_{2}-1}\right)^{\prime} \boldsymbol{z}_{t-1}=$ $\alpha_{2, t}^{L}(\alpha)\left(b_{1}, \ldots, b_{q_{2}-1}\right)^{\prime} \boldsymbol{z}_{t-1}$ a.s. This is clearly impossible unless $a_{1}=\ldots=a_{q_{2}-1}=b_{1}=\ldots=$ $b_{q_{2}-1}=0$, because $\alpha_{2, t}^{L}(\alpha)$ is a positive and strictly decreasing function of $\alpha^{\prime} \boldsymbol{z}_{t-1}\left(\neq \alpha_{0}\right)$ and because $\left(a_{1}, \ldots, a_{q_{2}-1}\right)^{\prime} \boldsymbol{z}_{t-1}$ and $\left(b_{1}, \ldots, b_{q_{2}-1}\right)^{\prime} \boldsymbol{z}_{t-1}$ are normally distributed or constants (if only $a_{1}$ and $b_{1}$ are nonzero). Therefore, $a=b=0$, so that $\mathcal{I}_{\alpha}$ is positive definite (for any fixed $\alpha \in A$ ).

To complete the verification of part (iii), we show that $\mathcal{I}_{\alpha}$ is a continuous function of $\alpha$ and such that $0<\inf _{\alpha \in A} \lambda_{\min }\left(\mathcal{I}_{\alpha}\right)$ and $\sup _{\alpha \in A} \lambda_{\max }\left(\mathcal{I}_{\alpha}\right)<\infty$. For continuity, let $\alpha_{n}$ be a sequence of points in $A$ converging to $\alpha_{\bullet} \in A$. It suffices to demonstrate that $\lim _{n \rightarrow \infty} E\left[\alpha_{2, t}^{L}\left(\alpha_{n}\right) \frac{\nabla f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)} \frac{\nabla^{\prime} f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)}\right]=$ $E\left[\alpha_{2, t}^{L}\left(\alpha_{\bullet}\right) \frac{\nabla f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)} \frac{\nabla^{\prime} f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)}\right]$ and similarly with $\alpha_{2, t}^{L}(\cdot)$ replaced by its square. This, however, is an immediate consequence of the dominated convergence theorem because $\alpha_{2, t}^{L}(\alpha)$ is a continuous positive function of $\alpha$ and smaller than 1, and because $E\left[\left\|\frac{\nabla f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)} \frac{\nabla^{\prime} f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)}\right\|\right]<\infty$ due to Lemma F. 1 (in Supplementary Appendix F.5). The statements on the eigenvalues follow from the continuity of $\mathcal{I}_{\alpha}$, the compactness of its domain $A$, and the positive definiteness of $\mathcal{I}_{\alpha}$ for all fixed $\alpha \in A$ shown above.

As for Assumption 5(iv), based on the expression of the remainder term in (17) it suffices to show that for all sequences of (non-random) positive scalars $\left\{\gamma_{T}, T \geq 1\right\}$ for which $\gamma_{T} \rightarrow 0$ as $T \rightarrow \infty$,

$$
\begin{equation*}
\sup _{(\pi, \varpi) \in \Pi:\left\|(\pi, \varpi)-\left(\pi^{*}, 0\right)\right\| \leq \gamma_{T}}\left\|T^{-1} \nabla_{(\pi, \varpi)(\pi, \varpi)^{\prime}}^{2} L_{T}^{\pi}(\alpha, \pi, \varpi)-\left(-\mathcal{I}_{\alpha}\right)\right\|=o_{p \alpha}(1) \tag{35}
\end{equation*}
$$

as noted above, the parameter space of $(\pi, \varpi)$, denoted by $\Pi$, is independent of $\alpha$. First we show that a uniform law of large numbers applies to the matrix $T^{-1} \nabla_{(\pi, \varpi)(\pi, \varpi)^{\prime}}^{2} L_{T}^{\pi}(\alpha, \pi, \varpi)$ on $A \times \Pi$, that is,

$$
\begin{equation*}
\sup _{\alpha \in A} \sup _{(\pi, \varpi) \in \Pi}\left\|T^{-1} \nabla_{(\pi, \varpi)(\pi, \varpi)^{\prime}}^{2} L_{T}^{\pi}(\alpha, \pi, \varpi)-E\left[\nabla_{(\pi, \varpi)(\pi, \varpi)^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)\right]\right\|=o_{p}(1) \tag{36}
\end{equation*}
$$

As $T^{-1} \nabla_{(\pi, \varpi)(\pi, \varpi)^{\prime}}^{2} L_{T}^{\pi}(\alpha, \pi, \varpi)=T^{-1} \sum_{t=1}^{T} \nabla_{(\pi, \varpi)(\pi, \varpi)^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)$ with $\nabla_{(\pi, \varpi)(\pi, \varpi)^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)$ a function of the stationary and ergodic process $\left(y_{t}, \boldsymbol{y}_{t-1}\right)$ (by Assumption 1(i)), we only need to establish that $E\left[\sup _{\alpha \in A} \sup _{(\pi, \varpi) \in \Pi}\left\|\nabla_{(\pi, \varpi)(\pi, \varpi)}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)\right\|\right]<\infty$ (see Ranga Rao (1962)); we verify this in Supplementary Appendix E.2. Furthermore, using the dominated convergence theorem and arguments similar to those used in part (iii), it can be shown that $E\left[\nabla_{(\pi, \varpi)(\pi, \varpi)^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)\right]$ is a (uniformly) continuous function of $(\alpha, \pi, \varpi)$ on $A \times \Pi$. Now note that the left hand side of (35) is dominated by

$$
\begin{aligned}
& \sup _{(\pi, \varpi) \in \Pi}\left\|T^{-1} \nabla_{(\pi, \varpi)(\pi, \varpi)^{\prime}}^{2} L_{T}^{\pi}(\alpha, \pi, \varpi)-E\left[\nabla_{(\pi, \varpi)(\pi, \varpi)^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)\right]\right\| \\
& \quad+\sup _{(\pi, \varpi) \in \Pi:\left\|(\pi, \varpi)-\left(\pi^{*}, 0\right)\right\| \leq \gamma_{T}}\left\|E\left[\nabla_{(\pi, \varpi)(\pi, \varpi)^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)\right]-\left(-\mathcal{I}_{\alpha}\right)\right\| .
\end{aligned}
$$

The former term is, due to (36), of order $o_{p \alpha}(1)$. Regarding the latter term, the supremum of it over $\alpha \in A$ converges to zero due to the uniform continuity of $E\left[\nabla_{(\pi, \varpi)(\pi, \varpi)^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)\right]$ and the fact that $E\left[\nabla_{(\pi, \varpi)(\pi, \varpi)^{\prime}}^{2} l_{t}^{\pi}\left(\alpha, \pi^{*}, 0\right)\right]=-\mathcal{I}_{\alpha}$ (this fact follows from the expression of $\nabla_{(\pi, \varpi)(\pi, \varpi)^{\prime}}^{2} l_{t}^{\pi}\left(\alpha, \pi^{*}, 0\right)$ given in (30)-(32) and Lemma F.3). Thus, we have verified (35), and hence Assumption 5(iv).

Assumptions 6-8. That $\Theta$ is locally (uniformly) equal to the cone $\Lambda=\mathbb{R}^{2 q}$ follows from the expression of the set $\Theta$ given in the verification of Assumption 5 above and the fact that $0\left(\in \mathbb{R}^{2 q}\right)$ is an interior point of $\Theta$. Assumption 7 is clear, as Assumption 6 holds with the cone $\Lambda=\mathbb{R}^{2 q}$. Assumption 8 is clear from the verification of Assumption 5.

Expression of $\hat{s}_{t \alpha}$ in Section 4.1. Let $\hat{\varepsilon}_{t}$ denote the OLS residuals rescaled by the estimated standard deviation, i.e., $\hat{\varepsilon}_{t}=\left(y_{t}-\hat{\tilde{\phi}}_{T, 0}-\hat{\tilde{\phi}}_{T, 1} y_{t-1}-\cdots-\hat{\tilde{\phi}}_{T, p} y_{t-p}\right) / \hat{\tilde{\sigma}}_{T}$, and set $\hat{s}_{t \alpha}=\left(\nabla f_{t}\left(\hat{\tilde{\phi}}_{T}\right) / f_{t}\left(\hat{\tilde{\phi}}_{T}\right),-(1-\right.$ $\left.\left.\alpha_{1, t}^{L}(\alpha)\right) \nabla f_{t}\left(\hat{\tilde{\phi}}_{T}\right) / f_{t}\left(\hat{\tilde{\phi}}_{T}\right)\right)$ with $\nabla f_{t}\left(\hat{\tilde{\phi}}_{T}\right) / f_{t}\left(\hat{\tilde{\phi}}_{T}\right)=\left(\boldsymbol{z}_{t-1} \hat{\varepsilon}_{t} / \hat{\tilde{\sigma}}_{T},\left(\hat{\varepsilon}_{t}^{2}-1\right) /\left(2 \hat{\tilde{\sigma}}_{T}^{2}\right)\right)$ (see (15) and (34)).

## C Details for the GMAR example

Assumption 1. Assumption 1(i) is assumed to hold. Assumption 1(ii) holds as $A$ is a compact subset of $(0,1)$. Assumption 1 (iii) holds by the definition of the mixing weight.

Assumption 2. For each fixed $\alpha \in A$, compactness of $B$ and $\Phi$ together with the continuity of $L_{T}(\alpha, \beta, \phi, \varphi)$ ensures the existence of a measurable maximizer $\left(\hat{\beta}_{T \alpha}, \hat{\phi}_{T \alpha}, \hat{\varphi}_{T \alpha}\right)$. Hence part (i) holds (with the $o_{p \alpha}(1)$ term equal to zero). (In the GMAR case, this maximizer is not unique when $\alpha=1 / 2$, but this does not matter for Assumption 2.) To prove that $\sup _{\alpha \in A}\left\|\left(\hat{\beta}_{T \alpha}, \hat{\phi}_{T \alpha}, \hat{\varphi}_{T \alpha}\right)-\left(\beta^{*}, \phi^{*}, \phi^{*}\right)\right\| \xrightarrow{p} 0$, by Andrews (1993, Lemma A1) it suffices to show that (a) $\sup _{(\alpha, \beta, \phi, \varphi) \in A \times B \times \Phi \times \Phi} \mid T^{-1} L_{T}(\alpha, \beta, \phi, \varphi)-$ $E\left[l_{t}(\alpha, \beta, \phi, \varphi)\right] \mid \xrightarrow{p} 0$ as $T \rightarrow \infty$ and that (b) for every neighborhood $N\left(\beta^{*}, \phi^{*}, \phi^{*}\right)$ of $\left(\beta^{*}, \phi^{*}, \phi^{*}\right)$, $\sup _{\alpha \in A^{\prime}} \sup _{(\beta, \phi, \varphi) \in B \times \Phi \times \Phi \backslash N\left(\beta^{*}, \phi^{*}, \phi^{*}\right)}\left(E\left[l_{t}(\alpha, \beta, \phi, \varphi)\right]-E\left[l_{t}\left(\alpha, \beta^{*}, \phi^{*}, \phi^{*}\right)\right]\right)<0$.

Property (a) can be verified by using the uniform law of large numbers given in Ranga Rao (1962). As $T^{-1} L_{T}(\alpha, \beta, \phi, \varphi)=T^{-1} \sum_{t=1}^{T} l_{t}(\alpha, \beta, \phi, \varphi)$ with $l_{t}(\alpha, \beta, \phi, \varphi)$ a function of the stationary and ergodic process $\left(y_{t}, \boldsymbol{y}_{t-1}\right)$, we only need to show that $E\left[\sup _{(\alpha, \beta, \phi, \varphi) \in A \times B \times \Phi \times \Phi}\left|l_{t}(\alpha, \beta, \phi, \varphi)\right|\right]<\infty$. Making use of Assumption 1, it is easy to show that $\left.C_{1} \exp \left\{-C_{2}\left(1+y_{t}^{2}+\cdots+y_{t-p}^{2}\right)\right\}\right) \leq f_{t}(\beta, \phi) \leq C_{2}$ for some $0<C_{1}, C_{2}<\infty$ and for all $(\beta, \phi) \in B \times \Phi$, so that $\log \left(C_{1}\right)-C_{2}\left(1+y_{t}^{2}+\cdots+y_{t-p}^{2}\right) \leq$ $l_{t}(\alpha, \beta, \phi, \varphi) \leq \log \left(C_{2}\right)$ for all $(\alpha, \beta, \phi, \varphi) \in A \times B \times \Phi \times \Phi$ (cf. Kalliovirta et al. (2015, pp. 264-265)); this holds in both the LMAR and GMAR cases. The required moment condition follows from this.

As for property (b), the uniform law of large numbers used above also delivers the continuity of the limit function $E\left[l_{t}(\alpha, \beta, \phi, \varphi)\right]$ on the compact set $A \times B \times \Phi \times \Phi$. Therefore it suffices to show that, for each fixed $\alpha \in A, E\left[l_{t}(\alpha, \beta, \phi, \varphi)\right]-E\left[l_{t}\left(\alpha, \beta^{*}, \phi^{*}, \phi^{*}\right)\right] \leq 0$ with equality if and only if $(\beta, \phi, \varphi)=\left(\beta^{*}, \phi^{*}, \phi^{*}\right)$. For the GMAR model, this can be straightforwardly shown with arguments used in the proof of Theorem 2 in Kalliovirta et al. (2015). To this end, define $\mathrm{n}_{1}\left(\nu_{1, t} \mid \boldsymbol{\nu}_{1, t-1} ;(\beta, \phi)\right)=$ $\left(2 \bar{\pi} \tilde{\sigma}_{1}^{2}\right)^{-1 / 2} \exp \left(-\left(\nu_{1, t}-\tilde{\phi}_{0}-\sum_{i=1}^{p} \tilde{\phi}_{i} \nu_{1, t-i}\right)^{2} /\left(2 \tilde{\sigma}_{1}^{2}\right)\right)$, where $\nu_{1, t}$ is the auxiliary $\operatorname{Gaussian} \operatorname{AR}(p)$ process introduced in the GMAR example of Section 2.2. Clearly, $\mathrm{n}_{1}\left(\nu_{1, t} \mid \boldsymbol{\nu}_{1, t-1} ;(\beta, \phi)\right)$ is the conditional density of $\nu_{1, t}$ given $\boldsymbol{\nu}_{1, t-1}=\left(\nu_{1, t-1}, \ldots, \nu_{1, t-p}\right)$. The notation $\mathrm{n}_{1}\left(\nu_{2, t} \mid \boldsymbol{\nu}_{2, t-1} ;(\beta, \varphi)\right)$ is defined similarly by using $\tilde{\varphi}_{i}$ and $\tilde{\sigma}_{2}^{2}$ instead of $\tilde{\phi}_{i}$ and $\tilde{\sigma}_{1}^{2}$.

Now, in the same way as in the above-mentioned proof of Kalliovirta et al. (2015) we can use arguments based on the Kullback-Leibler divergence and conclude that, for each fixed $\alpha \in A, E\left[l_{t}(\alpha, \beta, \phi, \varphi)\right]-$ $E\left[l_{t}\left(\alpha, \beta^{*}, \phi^{*}, \phi^{*}\right)\right] \leq 0$ with equality if and only if for almost all $(y, \boldsymbol{y}) \in \mathbb{R} \times \mathbb{R}^{p}$

$$
\begin{equation*}
\alpha_{1}^{G} \mathrm{n}_{1}(y \mid \boldsymbol{y} ;(\beta, \phi))+\alpha_{2}^{G} \mathrm{n}_{1}(y \mid \boldsymbol{y} ;(\beta, \varphi))=\mathrm{n}_{1}\left(y \mid \boldsymbol{y} ;\left(\beta^{*}, \phi^{*}\right)\right), \tag{37}
\end{equation*}
$$

where we use $\alpha_{m}^{G}$ to stand for $\alpha_{m, t}^{G}$ but with $\boldsymbol{y}_{t-1}$ therein replaced by $\boldsymbol{y}(m=1,2)$. Using well-known results on identification of finite mixtures of Gaussian distributions we find that, for each fixed $\alpha \in A$, and for each fixed $\boldsymbol{y} \in \mathbb{R}^{p}$ at a time, $\mathrm{n}_{1}(y \mid \boldsymbol{y} ;(\beta, \phi))=\mathrm{n}_{1}(y \mid \boldsymbol{y} ;(\beta, \varphi))=\mathrm{n}_{1}\left(y \mid \boldsymbol{y} ;\left(\beta^{*}, \phi^{*}\right)\right)$ for almost all $y$. Using the arguments following equation (A.4) in Kalliovirta et al. (2015) we can now establish the desired result $(\beta, \phi, \varphi)=\left(\beta^{*}, \phi^{*}, \phi^{*}\right)$.

The arguments used for the GMAR model above can also be used for the LMAR model, but two things are worth noting. First, the proof given for the GMAR model above goes through even when there are no common parameters so that $\phi$ and $\varphi$ could be used in place of $(\beta, \phi)$ and $(\beta, \varphi)$. Second, equation (37) can be obtained in the same way as in the GMAR case even though the derivation of the related equation (A.4) in Kalliovirta et al. (2015) made use of the explicit expression of the stationary density of $\left(y_{t}, \boldsymbol{y}_{t-1}\right)$ which is known for the GMAR model but, in general, unknown for the LMAR model. The reason for this is that the null hypothesis is here assumed to hold so that $y_{t}$ is a linear Gaussian $\operatorname{AR}(p)$ process, implying that $\left(y_{t}, \boldsymbol{y}_{t-1}\right)$ is normally distributed with density function a $p+1$ dimensional counterpart of the $p$ dimensional normal density function $\mathrm{n}_{p}\left(\boldsymbol{\nu}_{1, t} ; \tilde{\phi}\right)$ defined in equation
(8) (see the GMAR example of Section 2.2). After observing these two facts we can proceed in the same way as in the GMAR case and conclude that equation (37) holds also for the LMAR model as long as we replace the mixing weights of the GMAR model with those of the LMAR model. As the arguments employed in the proof of the GMAR case after equation (37) made no use of the mixing weights they apply also to the LMAR model and can be used to complete the proof.

Assumptions 3 and 4. Conditions (i) and (ii) of Assumption 3 clearly hold and condition (iii) can be verified in the same way as in the case of the LMAR model. Specifically, we have
$\boldsymbol{\pi}_{\alpha}(\phi, \varphi)-\boldsymbol{\pi}_{\alpha}\left(\phi^{*}, \phi^{*}\right)=(\alpha \phi+(1-\alpha) \varphi, \phi-\varphi)-\left(\phi^{*}, 0\right)=\left(\alpha\left(\phi-\phi^{*}\right)+(1-\alpha)\left(\varphi-\phi^{*}\right),\left(\phi-\phi^{*}\right)-\left(\varphi-\phi^{*}\right)\right)$,
and choosing $\|x\|_{*}=\|x\|_{1}=\sum_{i=1}^{2 q}\left|x_{i}\right|$ it can straightforwardly be seen that condition (iii) holds by Lemma A.1. Regarding Assumption 4, based on the expression of $\alpha_{t}^{G}$ and the definition of $\boldsymbol{\pi}_{\alpha}^{-1}(\pi, \varpi)=$ $(\pi+(1-\alpha) \varpi, \pi-\alpha \varpi)$, the required differentiability holds for all positive integers $k$.

Assumption 5: Derivation of expansion (12). First note that, as $\boldsymbol{\pi}_{\alpha}^{-1}(\pi, \varpi)=(\pi+(1-\alpha) \varpi, \pi-$ $\alpha \varpi)$, the reparameterized mixing weight in the GMAR model is given by

$$
\alpha_{1, t}^{G \pi}(\alpha, \beta, \pi, \varpi)=\frac{\alpha \mathrm{n}_{p}\left(\boldsymbol{y}_{t-1} ;(\beta, \pi+(1-\alpha) \varpi)\right)}{\alpha \mathrm{n}_{p}\left(\boldsymbol{y}_{t-1} ;(\beta, \pi+(1-\alpha) \varpi)\right)+(1-\alpha) \mathrm{n}_{p}\left(\boldsymbol{y}_{t-1} ;(\beta, \pi-\alpha \varpi)\right)}
$$

and the quantities $f_{2, t}^{\pi}(\alpha, \beta, \pi, \varpi)$ and $l_{t}^{\pi}(\alpha, \beta, \pi, \varpi)$ take the form

$$
\begin{aligned}
f_{2, t}^{\pi}(\alpha, \beta, \pi, \varpi) & =\alpha_{1, t}^{G \pi}(\alpha, \beta, \pi, \varpi) f_{t}(\beta, \pi+(1-\alpha) \varpi)+\left(1-\alpha_{1, t}^{G \pi}(\alpha, \beta, \pi, \varpi)\right) f_{t}(\beta, \pi-\alpha \varpi) \\
l_{t}^{\pi}(\alpha, \beta, \pi, \varpi) & =\log \left[\alpha_{1, t}^{G \pi}(\alpha, \beta, \pi, \varpi) f_{t}(\beta, \pi+(1-\alpha) \varpi)+\left(1-\alpha_{1, t}^{G \pi}(\alpha, \beta, \pi, \varpi)\right) f_{t}(\beta, \pi-\alpha \varpi)\right]
\end{aligned}
$$

Partial derivatives of $l_{t}^{\pi}(\alpha, \beta, \pi, \varpi)$ can be obtained with straightforward differentation; as the calculations are rather lengthy, they are relegated to Supplementary Appendix F.1.

Now consider, for an arbitrary fixed $\alpha \in A$, a standard fourth-order Taylor expansion of $L_{T}^{\pi}(\alpha, \beta, \pi, \varpi)=$ $\sum_{t=1}^{T} l_{t}^{\pi}(\alpha, \beta, \pi, \varpi)$ around $\left(\beta^{*}, \pi^{*}, 0\right)$ with respect to $(\beta, \pi, \varpi)$. For brevity, write $\tilde{\pi}=(\beta, \pi)$ and $\tilde{\pi}^{*}=\left(\beta^{*}, \pi^{*}\right)$. Collecting terms that turn out to be asymptotically negligible into a remainder term,

$$
\begin{align*}
L_{T}^{\pi}(\alpha, \tilde{\pi}, \varpi)-L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right) & =\left(\tilde{\pi}-\tilde{\pi}^{*}\right)^{\prime} \nabla_{\tilde{\pi}} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)+\frac{1}{2!}\left(\tilde{\pi}-\tilde{\pi}^{*}\right)^{\prime} \nabla_{\tilde{\pi} \tilde{\pi}^{\prime}}^{2} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)\left(\tilde{\pi}-\tilde{\pi}^{*}\right) \\
& +\frac{1}{2!} \varpi^{\prime} \nabla_{\varpi \varpi^{\prime}}^{2} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right) \varpi+\frac{3}{3!} \sum_{i=1}^{q_{1}+q_{2}} \sum_{j=1}^{q_{2}} \sum_{k=1}^{q_{2}} \nabla_{\tilde{\pi}_{i} \varpi_{j} \varpi_{k}}^{3} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)\left(\tilde{\pi}_{i}-\tilde{\pi}_{i}^{*}\right) \varpi_{j} \varpi_{k} \\
& +\frac{1}{4!} \sum_{i=1}^{q_{2}} \sum_{j=1}^{q_{2}} \sum_{k=1}^{q_{2}} \sum_{l=1}^{q_{2}} \nabla_{\varpi_{i} \varpi_{j} \varpi_{k} \varpi_{l}}^{4} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right) \varpi_{i} \varpi_{j} \varpi_{k} \varpi_{l}+R_{T}^{(1)}(\alpha, \tilde{\pi}, \varpi) \tag{38}
\end{align*}
$$

with an explicit expression of the remainder term $R_{T}^{(1)}(\alpha, \tilde{\pi}, \varpi)$ available in Supplementary Appendix F.2. Therein we also demonstrate that this fourth-order Taylor expansion can be written as a quadratic expansion of the form (12) given by

$$
\begin{align*}
& L_{T}^{\pi}(\alpha, \beta, \pi, \varpi)-L_{T}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right) \\
& =S_{T}^{\prime} \boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)-\frac{1}{2}\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)\right]^{\prime} \mathcal{I}\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)\right]+R_{T}(\alpha, \beta, \pi, \varpi) \tag{39}
\end{align*}
$$

or, setting $Z_{T}=\mathcal{I}^{-1} T^{-1 / 2} S_{T}$, in an alternative form corresponding to (13), with an explicit expression of the remainder term $R_{T}(\alpha, \beta, \pi, \varpi)$ available in Supplementary Appendix F.2.

The required derivatives are available in Supplementary Appendix F.1. Here we only present the derivatives that appear in the expression of $S_{T}$ in (18), that is, the components of $\tilde{\nabla}_{\boldsymbol{\theta}} l_{t}^{\pi *}=$
( $\tilde{\nabla}_{\theta} l_{t}^{\pi *}, \tilde{\nabla}_{\vartheta} l_{t}^{\pi *}$ ). From Supplementary Appendix F. 1 we obtain (here $\nabla_{i}$ denotes the $i$ th component of a derivative, and $\nabla \mathrm{n}_{p}(\cdot)$ denotes differentiation of $\mathrm{n}_{p}(\cdot)$ in (8) with respect to $\left.\tilde{\phi}=\left(\tilde{\phi}_{0}, \tilde{\phi}_{1}, \ldots, \tilde{\phi}_{p}, \tilde{\sigma}_{1}^{2}\right)\right)$

$$
\begin{align*}
\nabla_{\beta} l_{t}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right) & =\frac{\nabla_{1} f_{t}^{*}}{f_{t}^{*}}  \tag{40}\\
\nabla_{\pi_{i}} l_{t}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right) & =\frac{\nabla_{i+1} f_{t}^{*}}{f_{t}^{*}}  \tag{41}\\
\nabla_{w_{i} \varpi_{j}}^{2} l_{t}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right) & =\alpha(1-\alpha)\left[\frac{\nabla_{i+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{j+1} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{i+1} f_{t}^{*}}{f_{t}^{*}} \frac{\nabla_{j+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}+\frac{\nabla_{i+1, j+1}^{2} f_{t}^{*}}{f_{t}^{*}}\right] \tag{42}
\end{align*}
$$

where $i, j=1, \ldots, p+1$; for brevity, we denote $f_{t}^{*}=f_{t}\left(\beta^{*}, \pi^{*}\right), \mathrm{n}_{p}^{*}=\mathrm{n}_{p}\left(\beta^{*}, \pi^{*}\right)$, and similarly for their derivatives. Explicit expressions for the derivatives of $f_{t}$ and $\mathrm{n}_{p}$ are in Supplementary Appendix F.3.

Assumption 5: Verifying the assumption. As $\boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)=\left(\beta-\beta^{*}, \pi-\pi^{*}, \alpha(1-\alpha) v(\varpi)\right)$, part (i.a) is clearly satisfied. Assumption 1 requires $\alpha$ to be bounded away from zero and one, so that also part (i.b) is satisfied. For part (ii), notice from (18), (19), and (40)-(42) that $S_{T}=\sum_{t=1}^{T} s_{t}$ with

$$
\begin{equation*}
s_{t}=\left(\frac{\nabla f_{t}^{*}}{f_{t}^{*}}, c_{11} X_{t, 1,1}^{*}, \ldots, c_{q_{2}-1, q_{2}} X_{t, q_{2}-1, q_{2}}^{*}\right) \tag{43}
\end{equation*}
$$

where, for $i, j \in\left\{1, \ldots, q_{2}\right\}$, the $c_{i j}$ 's are as in Section 3.3.1 ( $c_{i j}=1 / 2$ if $i=j$ and $c_{i j}=1$ if $\left.i \neq j\right)$ and

$$
\begin{equation*}
X_{t, i, j}^{*}=\frac{\nabla_{i+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{j+1} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{i+1} f_{t}^{*}}{f_{t}^{*}} \frac{\nabla_{j+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}+\frac{\nabla_{i+1, j+1}^{2} f_{t}^{*}}{f_{t}^{*}} \tag{44}
\end{equation*}
$$

(note that $s_{t}$, and hence $S_{T}$, does not involve $\alpha$ as it cancels out from the expressions in (19) and (42)). Therefore the first three requirements in part (ii) are clearly satisfied. For the weak convergence requirement in part (ii) it now suffices to show that $T^{-1 / 2} S_{T} \xrightarrow{d} S$ in $\mathbb{R}^{r}$ for some multivariate Gaussian random vector $S$ with mean zero and $E\left[S S^{\prime}\right]=\mathcal{I}$. To this end, $s_{t}$ clearly forms a stationary and ergodic process. Moreover, due to Lemma F. 3 in Supplementary Appendix F.5, $E\left[\nabla_{i} f_{t}^{*} / f_{t}^{*} \mid \boldsymbol{y}_{t-1}\right]=$ $E\left[\nabla_{i j}^{2} f_{t}^{*} / f_{t}^{*} \mid \boldsymbol{y}_{t-1}\right]=0$ for any $i, j \in\left\{1, \ldots, q_{2}\right\}$ so that $s_{t}$ is a martingale difference sequence. From the expression of $\mathcal{I}$ in (51) in Supplementary Appendix F. 2 it is clear that $E\left[s_{t} s_{t}^{\prime}\right]=\mathcal{I}$. Positive definiteness of $\mathcal{I}$ is proven in Supplementary Appendix F.3. The stated convergence result now follows from the central limit theorem of Billingsley (1961) in conjunction with the Cramér-Wold device.

For part (iii), it suffices to show the finiteness and positive definiteness of $\mathcal{I}$; these are proven in Supplementary Appendices F. 2 and F.3. Part (iv) is proven in Supplementary Appendix F.4.

Assumption 6. By the definition of the set $\Theta_{\alpha}$ and the transformation $(\phi, \varphi) \rightarrow(\pi, \varpi)$, the set $\Theta_{\alpha}$ (see Sec 3.3.1) can equivalently be expressed as

$$
\begin{aligned}
& \Theta_{\alpha}=\left\{\boldsymbol{\theta}=(\theta, \vartheta) \in \mathbb{R}^{q_{1}+q_{2}+q_{\vartheta}}: \theta=\left(\beta-\beta^{*}, \alpha\left(\phi-\phi^{*}\right)+(1-\alpha)\left(\varphi-\phi^{*}\right)\right), \vartheta=\alpha(1-\alpha) v(\phi-\varphi)\right. \\
&\left.\quad \text { for some }(\beta, \phi, \varphi) \in B \times \Phi^{2}\right\} .
\end{aligned}
$$

We aim to show that the collection of sets $\left\{\Theta_{\alpha}, \alpha \in A\right\}$ is locally uniformly equal to the cone $\Lambda=\mathbb{R}^{q_{1}+q_{2}} \times v\left(\mathbb{R}^{q_{2}}\right)$ where $v\left(\mathbb{R}^{q_{2}}\right)=\left\{v(\varpi): \varpi \in \mathbb{R}^{q_{2}}\right\}$. Let $\bar{S}\left(\left(\beta^{*}, \phi^{*}\right), \delta\right)$ denote a closed $\left(q_{1}+q_{2}\right)$ sphere centered at $\left(\beta^{*}, \phi^{*}\right)$ with radius $\delta$, and $\bar{S}\left(\phi^{*}, \delta\right)$ a similar $q_{2}$-sphere. As $\left(\beta^{*}, \phi^{*}\right)$ is an interior point of $B \times \Phi$, we can find a $\delta_{1}>0$ such that $\bar{S}\left(\left(\beta^{*}, \phi^{*}\right), \delta_{1}\right) \times \bar{S}\left(\phi^{*}, \delta_{1}\right) \subset B \times \Phi^{2}$. By the definitions of the transformations $(\alpha, \beta, \phi, \varphi) \rightarrow \theta(\alpha, \beta, \phi, \varphi)$ and $(\alpha, \beta, \phi, \varphi) \rightarrow \vartheta(\alpha, \beta, \phi, \varphi)$ (defined implicitly by the definition of $\Theta_{\alpha}$ above), we can find a $\delta_{2} \in(0,1)$ such that $\left(-\delta_{2}, \delta_{2}\right)^{q_{1}+q_{2}} \times v\left(\left(-\delta_{2}, \delta_{2}\right)^{q_{2}}\right) \subset$
$\bigcap_{\alpha \in A} \Theta_{\alpha}\left(\delta_{1}\right) \subset \bigcap_{\alpha \in A} \Theta_{\alpha}$ where

$$
\begin{aligned}
& \Theta_{\alpha}\left(\delta_{1}\right)=\left\{\boldsymbol{\theta}=(\theta, \vartheta) \in \mathbb{R}^{q_{1}+q_{2}+q_{\vartheta}}: \theta=\left(\beta-\beta^{*}, \alpha\left(\phi-\phi^{*}\right)+(1-\alpha)\left(\varphi-\phi^{*}\right)\right), \vartheta=\alpha(1-\alpha) v(\phi-\varphi)\right. \\
& \\
& \text { for some } \left.(\beta, \phi, \varphi) \in \bar{S}\left(\left(\beta^{*}, \phi^{*}\right), \delta_{1}\right) \times \bar{S}\left(\phi^{*}, \delta_{1}\right)\right\} .
\end{aligned}
$$

Thus $\left(-\delta_{2}, \delta_{2}\right)^{q_{1}+q_{2}} \times v\left(\left(-\delta_{2}, \delta_{2}\right)^{q_{2}}\right) \subset \Theta_{\alpha}$ for all $\alpha \in A$ so that (as $\delta_{2}<1$ )

$$
\left(-\delta_{2}^{2}, \delta_{2}^{2}\right)^{q_{1}+q_{2}+q_{v}} \cap \Theta_{\alpha}=\left(-\delta_{2}^{2}, \delta_{2}^{2}\right)^{q_{1}+q_{2}} \times v\left(\left(-\delta_{2}, \delta_{2}\right)^{q_{2}}\right)=\left(-\delta_{2}^{2}, \delta_{2}^{2}\right)^{q_{1}+q_{2}+q_{\vartheta}} \cap\left[\mathbb{R}^{q_{1}+q_{2}} \times v\left(\mathbb{R}^{q_{2}}\right)\right]
$$

for all $\alpha \in A$. Hence, $\left\{\Theta_{\alpha}, \alpha \in A\right\}$ is locally uniformly equal to the cone $\Lambda=\mathbb{R}^{q_{1}+q_{2}} \times v\left(\mathbb{R}^{q_{2}}\right)$.
Assumptions 7 and 8. Assumption 7 is satisfied as Assumption 6 holds with the cone $\Lambda=\mathbb{R}^{q_{1}+q_{2}} \times$ $v\left(\mathbb{R}^{q_{2}}\right)$; note that now $q_{\vartheta}=q_{2}\left(q_{2}+1\right) / 2$. Assumption 8 is clear from the verification of Assumption 5 .

Expression of $\hat{s}_{t \alpha}$ in Section 4.1. Let $\hat{\varepsilon}_{t}$ and $\nabla f_{t}\left(\hat{\tilde{\phi}}_{T}\right) / f_{t}\left(\hat{\tilde{\phi}}_{T}\right)$ be as in the LMAR example (see Appendix B) and set (see (43) and (44)) $\hat{s}_{t}=\left(\nabla f_{t}\left(\hat{\tilde{\phi}}_{T}\right) / f_{t}\left(\hat{\tilde{\phi}}_{T}\right), c_{11} X_{t, 1,1}\left(\hat{\tilde{\phi}}_{T}\right), \ldots, c_{q_{2}-1, q_{2}} X_{t, q_{2}-1, q_{2}}\left(\hat{\tilde{\phi}}_{T}\right)\right)$ where, for $i, j \in\left\{1, \ldots, q_{2}\right\}$, the $c_{i j}$ 's are as in Section 3.3.1 ( $c_{i j}=1 / 2$ if $i=j$ and $c_{i j}=1$ if $i \neq j$ ) and

$$
X_{t, i, j}\left(\hat{\tilde{\phi}}_{T}\right)=\frac{\nabla_{i+1} \mathrm{n}_{p}\left(\hat{\tilde{\phi}}_{T}\right)}{\mathrm{n}_{p}\left(\hat{\tilde{\phi}}_{T}\right)} \frac{\nabla_{j+1} f_{t}\left(\hat{\tilde{\phi}}_{T}\right)}{f_{t}\left(\hat{\tilde{\phi}}_{T}\right)}+\frac{\nabla_{i+1} f_{t}\left(\hat{\tilde{\phi}}_{T}\right)}{f_{t}\left(\hat{\tilde{\phi}}_{T}\right)} \frac{\nabla_{j+1} \mathrm{n}_{p}\left(\hat{\tilde{\phi}}_{T}\right)}{\mathrm{n}_{p}\left(\hat{\tilde{\phi}}_{T}\right)}+\frac{\nabla_{i+1, j+1}^{2} f_{t}\left(\hat{\tilde{\phi}}_{T}\right)}{f_{t}\left(\hat{\tilde{\phi}}_{T}\right)} .
$$

Explicit expressions for the elements of $\nabla^{2} f_{t}\left(\hat{\tilde{\phi}}_{T}\right) / f_{t}\left(\hat{\tilde{\phi}}_{T}\right)$ can be obtained from (52) in Supplementary Appendix F. 3 by replacing $\varepsilon_{t}$ and $\tilde{\sigma}_{1}^{*}$ therein with $\hat{\varepsilon}_{t}$ and $\hat{\tilde{\sigma}}_{T}$, respectively. Expressions for the elements of $\nabla \mathrm{n}_{p}\left(\hat{\tilde{\phi}}_{T}\right) / \mathrm{n}_{p}\left(\hat{\tilde{\phi}}_{T}\right)$ can be obtained by evaluating (53) in Supplementary Appendix F. 3 at $\tilde{\phi}=\hat{\tilde{\phi}}_{T}$.

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## Supplementary Appendix to <br> 'Testing for observation-dependent regime switching in mixture autoregressive models' by Meitz and Saikkonen (not meant for publication)

## D Further details for the general results

Proof of Lemma A.1. Note that $\sup _{\alpha \in A}\left\|\left(\hat{\beta}_{T \alpha}, \hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)-\left(\beta^{*}, \pi^{*}, 0\right)\right\| \leq \sup _{\alpha \in A}\left\|\hat{\beta}_{T \alpha}-\beta^{*}\right\|+$ $\sup _{\alpha \in A}\left\|\left(\hat{\pi}_{T \alpha}, \hat{\varpi}_{T \alpha}\right)-\left(\pi^{*}, 0\right)\right\|$. The former term on the majorant side is $o_{p}(1)$ by Assumption 2(ii). The latter term equals $\sup _{\alpha \in A}\left\|\boldsymbol{\pi}_{\alpha}\left(\hat{\phi}_{T \alpha}, \hat{\varphi}_{T \alpha}\right)-\boldsymbol{\pi}_{\alpha}\left(\phi^{*}, \phi^{*}\right)\right\|$ which, due to the assumptions made, can be bounded by $C \sup _{\alpha \in A} h\left(\left\|\left(\hat{\phi}_{T \alpha}, \hat{\varphi}_{T \alpha}\right)-\left(\phi^{*}, \phi^{*}\right)\right\|_{*}\right) \leq C h\left(\sup _{\alpha \in A}\left\|\left(\hat{\phi}_{T \alpha}, \hat{\varphi}_{T \alpha}\right)-\left(\phi^{*}, \phi^{*}\right)\right\|_{*}\right)$. By Assumption 2(ii) and the equivalence of all vector norms on $\mathbb{R}^{2 q_{2}}$, the majorant side is $o_{p}(1)$.

Proof of Lemma 1, further details. To justify that the last term on the right hand side of (28) is dominated by $-\frac{1}{4}\left\|\boldsymbol{\theta}_{T \alpha}\right\|^{2}+o_{p \alpha}(1)$, note first that

$$
\left(o_{p \alpha}(1)-\frac{1}{2}\right)\left\|\boldsymbol{\theta}_{T \alpha}\right\|^{2}=\left(o_{p \alpha}(1)-\frac{1}{4}\right)\left\|\boldsymbol{\theta}_{T \alpha}\right\|^{2}-\frac{1}{4}\left\|\boldsymbol{\theta}_{T \alpha}\right\|^{2}:=W_{T \alpha}\left\|\boldsymbol{\theta}_{T \alpha}\right\|^{2}-\frac{1}{4}\left\|\boldsymbol{\theta}_{T \alpha}\right\|^{2},
$$

where $W_{T \alpha}=-\frac{1}{4}+o_{p \alpha}(1)$. Thus, $P\left(\sup _{\alpha \in A} W_{T \alpha} \leq 0\right) \rightarrow 1$ and (here $\mathbf{1}(\cdot)$ denotes the indicator function)

$$
\begin{aligned}
\sup _{\alpha \in A} W_{T \alpha}\left\|\boldsymbol{\theta}_{T \alpha}\right\|^{2} & =\sup _{\alpha \in A} W_{T \alpha}\left\|\boldsymbol{\theta}_{T \alpha}\right\|^{2} \mathbf{1}\left(\sup _{\alpha \in A} W_{T \alpha} \leq 0\right)+\sup _{\alpha \in A} W_{T \alpha}\left\|\boldsymbol{\theta}_{T \alpha}\right\|^{2} \mathbf{1}\left(\sup _{\alpha \in A} W_{T \alpha}>0\right) \\
& \leq \sup _{\alpha \in A} W_{T \alpha}\left\|\boldsymbol{\theta}_{T \alpha}\right\|^{2} \mathbf{1}\left(\sup _{\alpha \in A} W_{T \alpha}>0\right),
\end{aligned}
$$

where the last term is non-negative and positive with probability that is at most $P\left(\sup _{\alpha \in A} W_{T \alpha}>0\right) \rightarrow$ 0 . Thus, combining the above derivations yields the desired result $\left(o_{p \alpha}(1)-\frac{1}{2}\right)\left\|\boldsymbol{\theta}_{T \alpha}\right\|^{2} \leq-\frac{1}{4}\left\|\boldsymbol{\theta}_{T \alpha}\right\|^{2}+$ $o_{p \alpha}(1)$.

To justify the use of the continuous mapping theorem, note that in the first instance it is applied with the function $g: \mathcal{C}\left(A, \mathbb{R}^{r}\right) \times\left\{\mathcal{I}_{\alpha}\right\} \rightarrow \mathcal{C}\left(A, \mathbb{R}^{r}\right)$ mapping $\left(x_{\bullet}, \mathcal{I}_{\bullet}\right)$ to $\mathcal{I}_{\bullet}^{-1} x_{\bullet}$. Here $\mathcal{I}_{\alpha}^{-1} x_{\alpha}$ is continuous in $\alpha$ by Assumption 5(iii). Also, the latter set in the product $\mathcal{C}\left(A, \mathbb{R}^{r}\right) \times\left\{\mathcal{I}_{\alpha}\right\}$ contains only the non-random function $\mathcal{I}_{\alpha}$; this product space can be equipped with essentially the same metric as $\mathcal{C}\left(A, \mathbb{R}^{r}\right)$; cf. Andrews and Ploberger (1994, p. 1392 and 1407) and Zhu and Zhang (2006, proof of Theorem 5). In the second instance, the continuous mapping theorem is applied with the function $g: \mathcal{B}\left(A, \mathbb{R}^{r}\right) \rightarrow \mathbb{R}$ mapping $x_{\bullet}\left(\in \mathcal{B}\left(A, \mathbb{R}^{r}\right)\right)$ to $\sup _{\alpha \in A}\left\|x_{\alpha}\right\|$. For continuity, we need to establish that if a sequence $x_{n} \bullet$ converges to $x_{\bullet}$ in $\mathcal{B}\left(A, \mathbb{R}^{r}\right)$, then $g\left(x_{n \bullet}\right) \rightarrow g\left(x_{\bullet}\right)$ in $\mathbb{R}$ (i.e., if $\sup _{\alpha \in A}\left\|x_{n \alpha}-x_{\alpha}\right\| \rightarrow$ 0 , then $\left|\sup _{\alpha \in A}\left\|x_{n \alpha}\right\|-\sup _{\alpha \in A}\left\|x_{\alpha}\right\|\right| \rightarrow 0$ ). The triangle inequality implies that $\sup _{\alpha \in A}\left\|x_{n \alpha}\right\| \leq$ $\sup _{\alpha \in A}\left\|x_{n \alpha}-x_{\alpha}\right\|+\sup _{\alpha \in A}\left\|x_{\alpha}\right\|$, as well as the same result with $x_{n \alpha}$ and $x_{\alpha}$ interchanged, and the desired result follows from these inequalities.

Proof of Lemma 4, further details. It remains to verify the continuity mentioned in the proof. For simplicity, consider the continuity of the functions $g_{1}: \mathcal{C}\left(A, \mathbb{R}^{r}\right) \times\left\{\mathcal{I}_{\alpha}\right\} \rightarrow \mathcal{B}(A, \mathbb{R})$ mapping $\left(x_{\bullet}, \mathcal{I}_{\bullet}\right)$ to $x_{\bullet}^{\prime} \mathcal{I}_{\bullet} x_{\bullet}$ and $g_{2}: \mathcal{C}\left(A, \mathbb{R}^{r}\right) \times\left\{\mathcal{I}_{\alpha}\right\} \rightarrow \mathcal{B}(A, \mathbb{R})$ mapping $\left(x_{\bullet}, \mathcal{I}_{\bullet}\right)$ to $\inf _{\boldsymbol{\lambda} \in \Lambda}\left\{\left(\boldsymbol{\lambda}-x_{\bullet}\right)^{\prime} \mathcal{I}_{\bullet}\left(\boldsymbol{\lambda}-x_{\bullet}\right)\right\}$ separately. For $g_{1}$, continuity is rather clear, for if a sequence $\left(x_{n_{\bullet}}, \mathcal{I}_{\boldsymbol{\bullet}}\right)$ converges to $\left(x_{\bullet}, \mathcal{I}_{\boldsymbol{\bullet}}\right)$ in $\mathcal{C}\left(A, \mathbb{R}^{r}\right) \times\left\{\mathcal{I}_{\alpha}\right\}$, then $g_{1}\left(\left(x_{n}, \mathcal{I}_{\bullet}\right)\right) \rightarrow g_{1}\left(\left(x_{\bullet}, \mathcal{I}_{\bullet}\right)\right)$ in $\mathcal{B}(A, \mathbb{R})$ (i.e., if $\sup _{\alpha \in A}\left\|x_{n \alpha}-x_{\alpha}\right\| \rightarrow 0$, then
$\left.\sup _{\alpha \in A}\left|x_{n \alpha}^{\prime} \mathcal{I}_{\alpha} x_{n \alpha}-x_{\alpha}^{\prime} \mathcal{I}_{\alpha} x_{\alpha}\right| \rightarrow 0\right)$. For the continuity of $g_{2}$, suppose that $\sup _{\alpha \in A}\left\|x_{n \alpha}-x_{\alpha}\right\| \rightarrow 0$, and consider $\sup _{\alpha \in A}\left|\inf _{\boldsymbol{\lambda} \in \Lambda}\left\{\left(\boldsymbol{\lambda}-x_{n \alpha}\right)^{\prime} \mathcal{I}_{\alpha}\left(\boldsymbol{\lambda}-x_{n \alpha}\right)\right\}-\inf _{\boldsymbol{\lambda} \in \Lambda}\left\{\left(\boldsymbol{\lambda}-x_{\alpha}\right)^{\prime} \mathcal{I}_{\alpha}\left(\boldsymbol{\lambda}-x_{\alpha}\right)\right\}\right|$. Noting that

$$
\inf _{\boldsymbol{\lambda} \in \Lambda}\left\{\left(\boldsymbol{\lambda}-x_{n \alpha}\right)^{\prime} \mathcal{I}_{\alpha}\left(\boldsymbol{\lambda}-x_{n \alpha}\right)\right\}=\left\{\inf _{\boldsymbol{\lambda} \in \Lambda}\left\|\mathcal{I}_{\alpha}^{1 / 2}\left(\boldsymbol{\lambda}-x_{n \alpha}\right)\right\|\right\}^{2}
$$

and similarly for the other infimum, we need to consider

$$
\begin{equation*}
\sup _{\alpha \in A}\left\{\left|\inf _{\boldsymbol{\lambda} \in \Lambda}\left\|\mathcal{I}_{\alpha}^{1 / 2}\left(\boldsymbol{\lambda}-x_{n \alpha}\right)\right\|-\inf _{\boldsymbol{\lambda} \in \Lambda}\left\|\mathcal{I}_{\alpha}^{1 / 2}\left(\boldsymbol{\lambda}-x_{\alpha}\right)\right\|\right|\left(\inf _{\boldsymbol{\lambda} \in \Lambda}\left\|\mathcal{I}_{\alpha}^{1 / 2}\left(\boldsymbol{\lambda}-x_{n \alpha}\right)\right\|+\inf _{\boldsymbol{\lambda} \in \Lambda}\left\|\mathcal{I}_{\alpha}^{1 / 2}\left(\boldsymbol{\lambda}-x_{\alpha}\right)\right\|\right)\right\} . \tag{45}
\end{equation*}
$$

Using the triangle inequality and properties of the Euclidean vector norm,
$\left\|\mathcal{I}_{\alpha}^{1 / 2}\left(\boldsymbol{\lambda}-x_{n \alpha}\right)\right\| \leq\left\|\mathcal{I}_{\alpha}^{1 / 2}\left(\boldsymbol{\lambda}-x_{\alpha}\right)\right\|+\left\|\mathcal{I}_{\alpha}^{1 / 2}\left(x_{n \alpha}-x_{\alpha}\right)\right\| \leq\left\|\mathcal{I}_{\alpha}^{1 / 2}\left(\boldsymbol{\lambda}-x_{\alpha}\right)\right\|+\left(\lambda_{\max }\left(\mathcal{I}_{\alpha}\right)\right)^{1 / 2}\left\|x_{n \alpha}-x_{\alpha}\right\|$,
and similarly with $x_{n \alpha}$ and $x_{\alpha}$ exchanged, so that

$$
\left|\inf _{\boldsymbol{\lambda} \in \Lambda}\left\|\mathcal{I}_{\alpha}^{1 / 2}\left(\boldsymbol{\lambda}-x_{n \alpha}\right)\right\|-\inf _{\boldsymbol{\lambda} \in \Lambda}\left\|\mathcal{I}_{\alpha}^{1 / 2}\left(\boldsymbol{\lambda}-x_{\alpha}\right)\right\|\right| \leq\left(\lambda_{\max }\left(\mathcal{I}_{\alpha}\right)\right)^{1 / 2}\left\|x_{n \alpha}-x_{\alpha}\right\|
$$

As was noted after Assumption 6, the cone $\Lambda$ contains the origin, so that the term in (45) in parentheses is dominated by $\left(\lambda_{\max }\left(\mathcal{I}_{\alpha}\right)\right)^{1 / 2}\left(\left\|x_{n \alpha}\right\|+\left\|x_{\alpha}\right\|\right)$. Now, due to Assumption 5(iii), the fact that $x_{n \bullet}, x_{\bullet}$ are bounded, and the assumed $\sup _{\alpha \in A}\left\|x_{n \alpha}-x_{\alpha}\right\| \rightarrow 0$, the quantity in (45) converges to zero.

Proof of Lemma 5. For brevity and clarity, within this proof we use somewhat simplified notation and let

$$
\mathcal{I}_{\alpha}^{-1}=\left[\begin{array}{cc}
A & B \\
B^{\prime} & C
\end{array}\right]
$$

denote the partition of $\mathcal{I}_{\alpha}^{-1}$ (so that, e.g., $C$ is shorthand for $\left.\left(\mathcal{I}_{\alpha}^{-1}\right)_{\vartheta \vartheta}\right)$. This implies that $\mathcal{I}_{\alpha}$ can be expressed as

$$
\mathcal{I}_{\alpha}=\left[\begin{array}{cc}
D^{-1} & -D^{-1} B C^{-1} \\
-C^{-1} B^{\prime} D^{-1} & C^{-1}+C^{-1} B^{\prime} D^{-1} B C^{-1}
\end{array}\right]
$$

where $D=A-B C^{-1} B^{\prime}$ (thus, e.g., $D^{-1}=\mathcal{I}_{\theta \theta \alpha}$ ). Note also that $A, C$, and $D$ are symmetric (as $\mathcal{I}_{\alpha}$ is symmetric).

First note that $S_{\alpha}=\mathcal{I}_{\alpha} Z_{\alpha}$ can be expressed as

$$
S_{\alpha}=\left[\begin{array}{cc}
D^{-1} & -D^{-1} B C^{-1} \\
-C^{-1} B^{\prime} D^{-1} & C^{-1}+C^{-1} B^{\prime} D^{-1} B C^{-1}
\end{array}\right]\left[\begin{array}{l}
Z_{\theta \alpha} \\
Z_{\vartheta \alpha}
\end{array}\right]=\left[\begin{array}{c}
D^{-1} Z_{\theta \alpha}-D^{-1} B C^{-1} Z_{\vartheta \alpha} \\
-C^{-1} B^{\prime} D^{-1} Z_{\theta \alpha}+C^{-1} Z_{\vartheta \alpha}+C^{-1} B^{\prime} D^{-1} B C^{-1} Z_{\vartheta \alpha}
\end{array}\right]
$$

so that $S_{\theta \alpha}^{\prime} D S_{\theta \alpha}$ equals

$$
\begin{aligned}
S_{\theta \alpha}^{\prime} D S_{\theta \alpha} & =\left(D^{-1} Z_{\theta \alpha}-D^{-1} B C^{-1} Z_{\vartheta \alpha}\right)^{\prime} D\left(D^{-1} Z_{\theta \alpha}-D^{-1} B C^{-1} Z_{\vartheta \alpha}\right) \\
& =Z_{\theta \alpha}^{\prime} D^{-1} Z_{\theta \alpha}-Z_{\theta \alpha}^{\prime} D^{-1} B C^{-1} Z_{\vartheta \alpha}-Z_{\vartheta \alpha}^{\prime} C^{-1} B^{\prime} D^{-1} Z_{\theta \alpha}+Z_{\vartheta \alpha}^{\prime} C^{-1} B^{\prime} D^{-1} B C^{-1} Z_{\vartheta \alpha} .
\end{aligned}
$$

Now, since $Z_{\alpha}^{\prime} \mathcal{I}_{\alpha} Z_{\alpha}$ can be written as

$$
\begin{aligned}
Z_{\alpha}^{\prime} \mathcal{I}_{\alpha} Z_{\alpha} & =\left[\begin{array}{l}
Z_{\theta \alpha} \\
Z_{\vartheta \alpha}
\end{array}\right]^{\prime}\left[\begin{array}{cc}
D^{-1} & -D^{-1} B C^{-1} \\
-C^{-1} B^{\prime} D^{-1} & C^{-1}+C^{-1} B^{\prime} D^{-1} B C^{-1}
\end{array}\right]\left[\begin{array}{l}
Z_{\theta \alpha} \\
Z_{\vartheta \alpha}
\end{array}\right] \\
& =Z_{\theta \alpha}^{\prime} D^{-1} Z_{\theta \alpha}-Z_{\theta \alpha}^{\prime} D^{-1} B C^{-1} Z_{\vartheta \alpha}-Z_{\vartheta \alpha}^{\prime} C^{-1} B^{\prime} D^{-1} Z_{\theta \alpha} \\
& +Z_{\vartheta \alpha}^{\prime} C^{-1} Z_{\vartheta \alpha}+Z_{\vartheta \alpha}^{\prime} C^{-1} B^{\prime} D^{-1} B C^{-1} Z_{\vartheta \alpha},
\end{aligned}
$$

we obtain

$$
\begin{equation*}
Z_{\alpha}^{\prime} \mathcal{I}_{\alpha} Z_{\alpha}=Z_{\vartheta \alpha}^{\prime} C^{-1} Z_{\vartheta \alpha}+S_{\theta \alpha}^{\prime} D S_{\theta \alpha} . \tag{46}
\end{equation*}
$$

Now consider $\inf _{\boldsymbol{\lambda} \in \Lambda}\left\{\left(\boldsymbol{\lambda}-Z_{\alpha}\right)^{\prime} \mathcal{I}_{\alpha}\left(\boldsymbol{\lambda}-Z_{\alpha}\right)\right\}$. Similarly as above,

$$
\begin{aligned}
& \left(\boldsymbol{\lambda}-Z_{\alpha}\right)^{\prime} \mathcal{I}_{\alpha}\left(\boldsymbol{\lambda}-Z_{\alpha}\right) \\
& =\left(\boldsymbol{\lambda}_{\theta}-Z_{\theta \alpha}\right)^{\prime} D^{-1}\left(\boldsymbol{\lambda}_{\theta}-Z_{\theta \alpha}\right)-\left(\boldsymbol{\lambda}_{\theta}-Z_{\theta \alpha}\right)^{\prime} D^{-1} B C^{-1}\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)-\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)^{\prime} C^{-1} B^{\prime} D^{-1}\left(\boldsymbol{\lambda}_{\theta}-Z_{\theta \alpha}\right) \\
& +\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)^{\prime} C^{-1}\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)+\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)^{\prime} C^{-1} B^{\prime} D^{-1} B C^{-1}\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right) \\
& =\left(\boldsymbol{\lambda}_{\theta}-Z_{\theta \alpha}\right)^{\prime} D^{-1}\left(\boldsymbol{\lambda}_{\theta}-Z_{\theta \alpha}\right)-\left(\boldsymbol{\lambda}_{\theta}-Z_{\theta \alpha}\right)^{\prime} D^{-1}\left[B C^{-1}\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)\right]-\left[B C^{-1}\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)\right]^{\prime} D^{-1}\left(\boldsymbol{\lambda}_{\theta}-Z_{\theta \alpha}\right) \\
& +\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)^{\prime} C^{-1}\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)+\left[B C^{-1}\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)\right]^{\prime} D^{-1}\left[B C^{-1}\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)\right] \\
& =\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)^{\prime} C^{-1}\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)+\left\{\left(\boldsymbol{\lambda}_{\theta}-Z_{\theta \alpha}\right)-\left[B C^{-1}\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)\right]\right\}^{\prime} D^{-1}\left\{\left(\boldsymbol{\lambda}_{\theta}-Z_{\theta \alpha}\right)-\left[B C^{-1}\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)\right]\right\} .
\end{aligned}
$$

Now, for any fixed $\boldsymbol{\lambda}_{\vartheta} \in \mathbb{R}^{q_{\vartheta}}$, Assumption 7 implies that

$$
\inf _{\boldsymbol{\lambda}_{\theta} \in \mathbb{R}^{q_{\theta}}}\left\{\left(\boldsymbol{\lambda}_{\theta}-Z_{\theta \alpha}\right)-\left[B C^{-1}\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)\right]\right\}^{\prime} D^{-1}\left\{\left(\boldsymbol{\lambda}_{\theta}-Z_{\theta \alpha}\right)-\left[B C^{-1}\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)\right]\right\}=0
$$

(cf. Andrews (1999, eqn. (7.35))) so that

$$
\begin{equation*}
\inf _{\boldsymbol{\lambda} \in \Lambda}\left\{\left(\boldsymbol{\lambda}-Z_{\alpha}\right)^{\prime} \mathcal{I}_{\alpha}\left(\boldsymbol{\lambda}-Z_{\alpha}\right)\right\}=\inf _{\boldsymbol{\lambda}_{\vartheta} \in \Lambda_{\vartheta}}\left\{\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)^{\prime} C^{-1}\left(\boldsymbol{\lambda}_{\vartheta}-Z_{\vartheta \alpha}\right)\right\} . \tag{47}
\end{equation*}
$$

Combining (46) and (47) and recalling that $C^{-1}=\left(\mathcal{I}_{\alpha}^{-1}\right)_{\vartheta \vartheta}^{-1}$ and $D=\mathcal{I}_{\theta \theta \alpha}^{-1}$ yields the equality stated in the lemma.

Finally, $\left(\mathcal{I}_{\alpha}^{-1}\right)_{\vartheta \vartheta}$ and $Z_{\vartheta \alpha}$ can be expressed as

$$
\begin{aligned}
\left(\mathcal{I}_{\alpha}^{-1}\right)_{\vartheta \vartheta} & =\left(\mathcal{I}_{\vartheta \vartheta \alpha}-\mathcal{I}_{\vartheta \theta \alpha} \mathcal{I}_{\theta \theta}^{-1} \mathcal{I}_{\theta \vartheta \alpha}\right)^{-1}\left[=\mathcal{I}_{\vartheta \vartheta \alpha}^{-1}+\mathcal{I}_{\vartheta \vartheta \alpha}^{-1} \mathcal{I}_{\vartheta \theta \alpha}\left(\mathcal{I}_{\theta \theta \alpha}-\mathcal{I}_{\theta \vartheta \alpha} \mathcal{I}_{\vartheta \vartheta \alpha}^{-1} \mathcal{I}_{\vartheta \theta \alpha}\right)^{-1} \mathcal{I}_{\theta \vartheta \alpha} \mathcal{I}_{\vartheta \vartheta \alpha}^{-1}\right], \\
Z_{\vartheta \alpha} & =\left(\mathcal{I}_{\alpha}^{-1}\right)_{\vartheta \vartheta}\left(S_{\vartheta \alpha}-\mathcal{I}_{\vartheta \theta \alpha} \mathcal{I}_{\theta \theta \alpha}^{-1} S_{\theta \alpha}\right)\left[=\mathcal{I}_{\vartheta \vartheta \alpha}^{-1} S_{\vartheta \alpha}+\mathcal{I}_{\vartheta \vartheta \alpha}^{-1} \mathcal{I}_{\vartheta \theta \alpha}\left(\mathcal{I}_{\theta \theta \alpha}-\mathcal{I}_{\theta \vartheta \alpha} \mathcal{I}_{\vartheta \vartheta \alpha}^{-1} \mathcal{I}_{\vartheta \theta \alpha}\right)^{-1}\left(\mathcal{I}_{\theta \vartheta \alpha} \mathcal{I}_{\vartheta \vartheta \alpha}^{-1} S_{\vartheta \alpha}-S_{\theta \alpha}\right)\right],
\end{aligned}
$$

where the two different expressions result from two different ways of writing the inverse of a partitioned matrix.

Proof of Lemma 6. The required arguments are standard but presented for completeness and to contrast them with arguments that lead to Lemma 4. The reparameterization described in Assumption 3 is unnecessary and the original $\tilde{\phi}$-parameterization may be used (alternatively, consider the identity mapping $\pi=\boldsymbol{\pi}(\tilde{\phi})=\tilde{\phi})$. As for the quadratic expansion of the log-likelihood function, let $\boldsymbol{\theta}(\tilde{\phi})=$ $\left(\tilde{\phi}-\tilde{\phi}^{*}\right)$ take the role of $\boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)$, and note that straightforward derivations (similar to those used in the LMAR example in Section 3.3.1) yield

$$
\begin{aligned}
L_{T}^{0}(\tilde{\phi})-L_{T}^{0}\left(\tilde{\phi}^{*}\right) & =\left(T^{-1 / 2} S_{T}^{0}\right)^{\prime}\left[T^{1 / 2} \boldsymbol{\theta}(\tilde{\phi})\right]-\frac{1}{2}\left[T^{1 / 2} \boldsymbol{\theta}(\tilde{\phi})\right]^{\prime} \mathcal{I}^{0}\left[T^{1 / 2} \boldsymbol{\theta}(\tilde{\phi})\right]+R_{T}(\tilde{\phi}), \\
R_{T}(\tilde{\phi}) & =\frac{1}{2}\left[T^{1 / 2} \boldsymbol{\theta}(\tilde{\phi})\right]^{\prime}\left[T^{-1} \nabla_{\phi \phi^{\prime}} L_{T}^{0}(\tilde{\phi})-\left(-\mathcal{I}^{0}\right)\right]\left[T^{1 / 2} \boldsymbol{\theta}(\tilde{\phi})\right],
\end{aligned}
$$

with $\dot{\phi}$ denoting a point between $\tilde{\phi}$ and $\tilde{\phi}^{*}$. Validity of Assumption 5 follows from the arguments used in connection with the LMAR example together with Assumption 8. Assumption 6 holds with $\Lambda=\mathbb{R}^{p+2}$. Arguments analogous to those that lead to Lemma 4 now yield the stated convergence result, and the convergence is joint as in both cases it follows from the weak convergence result $T^{-1 / 2} S_{\boldsymbol{\bullet}} \Rightarrow S_{\bullet}$.

## E Further details for the LMAR example

## E. 1 Verification of Assumption 5(ii), further details

As for the weak convergence requirement in part (ii), we rely on Theorem 2 (and the remarks that follow it) in Andrews and Ploberger (1995). As can be seen from the proof of their Theorem 2, it suffices to verify their conditions EP1(a), EP1(e), and EP4 (omitting the weakly exogeneous $X_{t}$ variables therein). Under the null hypothesis, $y_{t}$ is a linear Gaussian $\operatorname{AR}(p)$ process so that condition EP1(a) is satisfied with geometrically declining mixing numbers. To check condition EP1(e), we show that $E\left[\sup _{\alpha \in A} \sup _{(\pi, \varpi) \in \Pi}\left|l_{t}^{\pi}(\alpha, \pi, \varpi)\right|\right]<\infty, E\left[\sup _{\alpha \in A}\left\|\nabla_{(\pi, \varpi)} l_{t}^{\pi}\left(\alpha, \pi^{*}, 0\right)\right\|^{r}\right]<\infty$ for any positive $r$, and $E\left[\sup _{\alpha \in A} \sup _{(\pi, \varpi) \in \Pi}\left\|\nabla_{(\pi, \varpi)(\pi, \varpi)^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)\right\|\right]<\infty$. The first of these moment conditions follows from the arguments used to verify our Assumption 2 (see the verification of this Assumption for the GMAR model; the details for the LMAR model are presented there). The second holds due to the expression of $\nabla_{(\pi, \varpi)} l_{t}^{\pi}\left(\alpha, \pi^{*}, 0\right)$ in (29), the fact that $0<\alpha_{1, t}^{L}(\alpha)<1$, and Lemma F.1. The third is verified below in Supplementary Appedix E.2. As the compactness requirement of condition EP4(a) holds by our Assumption 1(ii), it remains to verify EP4(b). To this end, note that for arbitrary $a, b \in A$,

$$
\left\|\nabla_{(\pi, \varpi)} l_{t}^{\pi}\left(a, \pi^{*}, 0\right)-\nabla_{(\pi, \varpi)} l_{t}^{\pi}\left(b, \pi^{*}, 0\right)\right\|^{r}=\left|\alpha_{1, t}^{L}(a)-\alpha_{1, t}^{L}(b)\right|^{r}\left\|\frac{\nabla f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)}\right\|^{r}
$$

By straightforward differentiation, $\nabla_{\alpha} \alpha_{1, t}^{L}(\alpha)=\alpha_{1, t}^{L}(\alpha)\left(1-\alpha_{1, t}^{L}(\alpha)\right)\left(1, y_{t-1}, \ldots, y_{t-m}\right)$, so that by the mean value theorem

$$
\alpha_{1, t}^{L}(a)-\alpha_{1, t}^{L}(b)=\alpha_{1, t}^{L}\left(c_{a, b}\right)\left(1-\alpha_{1, t}^{L}\left(c_{a, b}\right)\right)\left(1, y_{t-1}, \ldots, y_{t-m}\right)^{\prime}(a-b)
$$

for some $c_{a, b} \in \mathbb{R}^{m+1}$ between $a$ and $b$ (as $A$ is not necessarily convex, $c_{a, b}$ does not necessarily belong to $A$, but this has no effect in what follows as the expression $\alpha_{1, t}^{L}\left(c_{a, b}\right)$ is nevertheless well defined for all $c_{a, b} \in \mathbb{R}^{m+1}$ ). Setting $B_{t}=1+\left|y_{t-1}\right|+\ldots+\left|y_{t-m}\right|$ and noting that $0<\alpha_{1, t}^{L}\left(c_{a, b}\right)<1$ this implies that

$$
\left|\alpha_{1, t}^{L}(a)-\alpha_{1, t}^{L}(b)\right| \leq\left|\left(1, y_{t-1}, \ldots, y_{t-m}\right)^{\prime}(a-b)\right| \leq\left(1+\left|y_{t-1}\right|+\ldots+\left|y_{t-m}\right|\right)\|a-b\|=B_{t}\|a-b\|
$$

Hence

$$
E\left[\sup _{a, b \in A,\|a-b\|<\delta}\left\|\nabla_{(\pi, \varpi)} l_{t}^{\pi}\left(a, \pi^{*}, 0\right)-\nabla_{(\pi, \varpi)} l_{t}^{\pi}\left(b, \pi^{*}, 0\right)\right\|^{r}\right]<\delta^{r} E\left[B_{t}^{r}\left\|\frac{\nabla f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)}\right\|^{r}\right]
$$

where on the majorant side the expectation is finite (due to the fact that the $y_{t}$ 's possess moments of all orders, see also the proof of Lemma F.1). Hence condition EP4(b) holds, and the desired weak convergence follows.

## E. 2 Verification of Assumption 5(iv), further details

It remains to show that $E\left[\sup _{\alpha \in A} \sup _{(\pi, \varpi) \in \Pi}\left\|\nabla_{(\pi, \varpi)(\pi, \varpi)^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)\right\|\right]<\infty$. This in turn follows if we show the same with $\nabla_{(\pi, \varpi)(\pi, \varpi)}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)$ replaced by $\nabla_{\pi \pi^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi), \nabla_{\varpi \varpi \prime}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)$, and $\nabla_{\pi \varpi^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)$. Consider the expression of $\nabla_{\pi \pi^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)$ given in Appendix B and recall that $0<\alpha_{1, t}^{L}(\alpha), \alpha_{2, t}^{L}(\alpha)<1$ and $\nabla f_{t}(\pi)=f_{t}(\pi) \nabla_{\pi} l_{t}^{0}(\pi)$ with $l_{t}^{0}(\pi)=\log \left[f_{t}(\pi)\right]$. Then we can, for instance, write

$$
\left\|\frac{\alpha_{1, t}^{L}(\alpha) \nabla f_{t}(\pi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)}\right\|=\left\|\frac{\alpha_{1, t}^{L}(\alpha) f_{t}(\pi) \nabla_{\pi} l_{t}^{0}(\pi)}{\alpha_{1, t}^{L}(\alpha) f_{t}(\pi)+\alpha_{2, t}^{L}(\alpha) f_{t}(\pi-\varpi)}\right\| \leq\left\|\nabla_{\pi} l_{t}^{0}(\pi)\right\|
$$

and

$$
\begin{aligned}
\left\|\frac{\alpha_{1, t}^{L}(\alpha) \nabla^{\prime} f_{t}(\pi)+\alpha_{2, t}^{L}(\alpha) \nabla^{\prime} f_{t}(\pi-\varpi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)}\right\| & =\left\|\frac{\alpha_{1, t}^{L}(\alpha) f_{t}(\pi) \nabla_{\pi} l_{t}^{0}(\pi)+\alpha_{2, t}^{L}(\alpha) f_{t}(\pi-\varpi) \nabla_{\pi} l_{t}^{0}(\pi-\varpi)}{\alpha_{1, t}^{L}(\alpha) f_{t}(\pi)+\alpha_{2, t}^{L}(\alpha) f_{t}(\pi-\varpi)}\right\| \\
& \leq\left\|\nabla_{\pi} l_{t}^{0}(\pi)\right\|+\left\|\nabla_{\pi} l_{t}^{0}(\pi-\varpi)\right\|
\end{aligned}
$$

As similar inequalities can be obtained for the second term of $\nabla_{\pi \pi^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)$, we get

$$
\begin{aligned}
\left\|\nabla_{\pi \pi^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)\right\| & \leq\left\|\alpha_{1, t}^{L}(\alpha) \frac{\nabla^{2} f_{t}(\pi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)}\right\|+\left\|\nabla_{\pi} l_{t}^{0}(\pi)\right\|^{2}+\left\|\nabla_{\pi} l_{t}^{0}(\pi)\right\|\left\|\nabla_{\pi} l_{t}^{0}(\pi-\varpi)\right\| \\
& +\left\|\alpha_{2, t}^{L}(\alpha) \frac{\nabla^{2} f_{t}(\pi-\varpi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)}\right\|+\left\|\nabla_{\pi} l_{t}^{0}(\pi-\varpi)\right\|^{2}+\left\|\nabla_{\pi} l_{t}^{0}(\pi)\right\|\left\|\nabla_{\pi} l_{t}^{0}(\pi-\varpi)\right\| .
\end{aligned}
$$

Next note that

$$
\frac{\nabla^{2} f_{t}(\pi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)}=\frac{\nabla\left(f_{t}(\pi) \nabla_{\pi} l_{t}^{0}(\pi)\right)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)}=\frac{f_{t}(\pi) \nabla_{\pi} l_{t}^{0}(\pi) \nabla_{\pi^{\prime}} l_{t}^{0}(\pi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)}+\frac{f_{t}(\pi) \nabla_{\pi \pi^{\prime}}^{2} l_{t}^{0}(\pi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)},
$$

so that arguments similar to those already used above give

$$
\left\|\alpha_{1, t}^{L}(\alpha) \frac{\nabla^{2} f_{t}(\pi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)}\right\| \leq\left\|\nabla_{\pi} l_{t}^{0}(\pi)\right\|^{2}+\left\|\nabla_{\pi \pi^{\prime}}^{2} l_{t}^{0}(\pi)\right\|
$$

and

$$
\left\|\alpha_{2, t}^{L}(\alpha) \frac{\nabla^{2} f_{t}(\pi-\varpi)}{f_{2, t}^{\pi}(\alpha, \pi, \varpi)}\right\| \leq\left\|\nabla_{\pi} l_{t}^{0}(\pi-\varpi)\right\|^{2}+\left\|\nabla_{\pi \pi^{\prime}}^{2} l_{t}^{0}(\pi-\varpi)\right\| .
$$

Hence, we can conclude that

$$
\begin{gathered}
\left\|\nabla_{\pi \pi^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)\right\| \leq 2\left\|\nabla_{\pi} l_{t}^{0}(\pi)\right\|^{2}+2\left\|\nabla_{\pi} l_{t}^{0}(\pi-\varpi)\right\|^{2}+2\left\|\nabla_{\pi} l_{t}^{0}(\pi)\right\|\left\|\nabla_{\pi} l_{t}^{0}(\pi-\varpi)\right\| \\
+\left\|\nabla_{\pi \pi^{\prime}}^{2} l_{t}^{0}(\pi)\right\|+\left\|\nabla_{\pi \pi^{\prime}}^{2} l_{t}^{0}(\pi-\varpi)\right\|
\end{gathered}
$$

To bound the expression on the dominant side, note that $\nabla_{\pi} l_{t}^{0}(\pi)=\frac{\nabla f_{t}(\pi)}{f_{t}(\pi)}$ and $\nabla_{\pi \pi^{\prime}}^{2} l_{t}^{0}(\pi)=\frac{\nabla^{2} f_{t}(\pi)}{f_{t}(\pi)}-$ $\frac{\nabla f_{t}(\pi)}{f_{t}(\pi)} \frac{\nabla^{\prime} f_{t}(\pi)}{f_{t}(\pi)}$ so that Lemma F. 1 ensures that $E\left[\sup _{\alpha \in A} \sup _{(\pi, \varpi) \in \Pi}\left\|\nabla_{\pi \pi^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)\right\|\right]<\infty$. An inspection of the expressions of $\nabla_{\pi \varpi^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)$ and $\nabla_{\varpi \varpi^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)$ in Appendix B shows that a similar result can be obtained with $\nabla_{\pi \pi^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)$ replaced by $\nabla_{\pi \varpi m^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)$ and $\nabla_{\varpi \varpi^{\prime}}^{2} l_{t}^{\pi}(\alpha, \pi, \varpi)$, yielding the desired result.

## F Further details for the GMAR example

## F. 1 Partial derivatives of the reparameterized log-likelihood function

Here we present certain partial derivatives of $l_{t}^{\pi}(\alpha, \beta, \pi, \varpi)$ with respect to $(\beta, \pi, \varpi)$. For brevity, set $\tilde{\pi}=(\beta, \pi)$ (and similarly $\tilde{\pi}^{*}=\left(\beta^{*}, \pi^{*}\right)$ ), so that the desired derivatives are with respect to $\tilde{\pi}$ and $\varpi$ or, elementwise, with respect to $\tilde{\pi}_{i}$ and $\varpi_{j}$ for $i=1, \ldots, p+2$ and $j=1, \ldots, p+1$. In the derivative expressions below, the subindices in $\tilde{\pi}$ and $\varpi$ are tacitly assumed to be withing these ranges. For brevity, denote $l_{t}^{\pi *}=l_{t}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right), f_{t}^{*}=f_{t}\left(\tilde{\pi}^{*}\right), \mathrm{n}_{p}^{*}=\mathrm{n}_{p}\left(\tilde{\pi}^{*}\right)$, and similarly for their partial derivatives.

The following derivatives are obtained with straightforward (but tedious and lengthy) differentiation. The necessary calculations for the first- and second-order derivatives are presented in Supplementary Appendix F.7, but for brevity we omit the detailed calculations for the third- and fourth-order derivatives.

First- and second-order derivatives:

$$
\begin{aligned}
\nabla_{\tilde{\pi}_{i}} l_{t}^{\pi *} & =\frac{\nabla_{i} f_{t}^{*}}{f_{t}^{*}} \\
\nabla_{\varpi_{j}} l_{t}^{\pi *} & =0 \\
\nabla_{\tilde{\pi}_{i} \tilde{\pi}_{j}}^{2} l_{t}^{\pi *} & =\frac{\nabla_{i j}^{2} f_{t}^{*}}{f_{t}^{*}}-\frac{\nabla_{i} f_{t}^{*}}{f_{t}^{*}} \frac{\nabla_{j} f_{t}^{*}}{f_{t}^{*}} . \\
\nabla_{\tilde{\pi}_{i} \varpi_{j}}^{2} l_{t}^{\pi *} & =0 \\
\nabla_{w_{i} \varpi_{j}}^{2} l_{t}^{\pi *} & =\alpha_{1} \alpha_{2}\left[\frac{\nabla_{i+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{j+1} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{i+1} f_{t}^{*}}{f_{t}^{*}} \frac{\nabla_{j+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}+\frac{\nabla_{i+1, j+1}^{2} f_{t}^{*}}{f_{t}^{*}}\right]
\end{aligned}
$$

Third-order derivatives:

$$
\begin{aligned}
& \nabla_{\tilde{\pi}_{i} \tilde{\pi}_{j} \tilde{\pi}_{k}}^{3} l_{t}^{\pi *}=\frac{\nabla_{i j k}^{3} f_{t}^{*}}{f_{t}^{*}}-\frac{\nabla_{i j}^{2} f_{t}^{*}}{f_{t}^{*}} \frac{\nabla_{k} f_{t}^{*}}{f_{t}^{*}}-\frac{\nabla_{i k}^{2} f_{t}^{*}}{f_{t}^{*}} \frac{\nabla_{j} f_{t}^{*}}{f_{t}^{*}}-\frac{\nabla_{j k}^{2} f_{t}^{*}}{f_{t}^{*}} \frac{\nabla_{i} f_{t}^{*}}{f_{t}^{*}}+2 \frac{\nabla_{i} f_{t}^{*}}{f_{t}^{*}} \frac{\nabla_{j} f_{t}^{*}}{f_{t}^{*}} \frac{\nabla_{k} f_{t}^{*}}{f_{t}^{*}} \\
& \nabla_{\tilde{\pi}_{i} \tilde{\pi}_{j} \omega_{k}}^{3} l_{t}^{\pi *}=-\alpha_{1} \alpha_{2} \frac{\nabla_{i} f_{t}^{*}}{f_{t}^{*}}\left(\frac{\nabla_{j, k+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}-\frac{\nabla_{j} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{k+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}\right) \\
& \nabla_{\tilde{\pi}_{i} \omega_{j} \omega_{k}}^{3} l_{t}^{\pi *}=\alpha_{1} \alpha_{2}\left[\frac{\nabla_{k+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}\left(\frac{\nabla_{i, j+1} f_{t}^{*}}{f_{t}^{*}}-\frac{\nabla_{i} f_{t}^{*}}{f_{t}^{*}} \frac{\nabla_{j+1} f_{t}^{*}}{f_{t}^{*}}\right)+\frac{\nabla_{j+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}\left(\frac{\nabla_{i, k+1} f_{t}^{*}}{f_{t}^{*}}-\frac{\nabla_{i} f_{t}^{*}}{f_{t}^{*}} \frac{\nabla_{k+1} f_{t}^{*}}{f_{t}^{*}}\right)\right. \\
& +\left(\frac{\nabla_{i, j+1, k+1}^{3} f_{t}^{*}}{f_{t}^{*}}-\frac{\nabla_{i} f_{t}^{*}}{f_{t}^{*}} \frac{\nabla_{j+1, k+1}^{2} f_{t}^{*}}{f_{t}^{*}}\right)+\left(\frac{\nabla_{i, k+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}-\frac{\nabla_{i} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{k+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}\right) \frac{\nabla_{j+1} f_{t}^{*}}{f_{t}^{*}} \\
& \left.+\left(\frac{\nabla_{i, j+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}-\frac{\nabla_{i} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{j+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}\right) \frac{\nabla_{k+1} f_{t}^{*}}{f_{t}^{*}}\right] \\
& \nabla_{w_{i} \varpi_{j} \varpi_{k}}^{3} l_{t}^{\pi *}=\alpha_{1} \alpha_{2}\left(\alpha_{2}-\alpha_{1}\right)\left[\frac{\nabla_{i+1, j+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{k+1} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{i+1, k+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{j+1} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{j+1, k+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{i+1} f_{t}^{*}}{f_{t}^{*}}\right. \\
& +\frac{\nabla_{i+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{j+1, k+1}^{2} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{j+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{i+1, k+1}^{2} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{k+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{i+1, j+1}^{2} f_{t}^{*}}{f_{t}^{*}} \\
& \left.+\frac{\nabla_{i+1, j+1, k+1}^{3} f_{t}^{*}}{f_{t}^{*}}\right]
\end{aligned}
$$

Fourth-order derivative (fourth-order derivatives with respect to $\tilde{\pi}$ will not be explicitly needed, and thus we omit their expressions):

$$
\begin{aligned}
& \nabla_{\varpi_{i} \varpi_{j} \varpi_{k} \varpi_{l}}^{4} l_{t}^{\pi *} \\
& =-\alpha_{1}^{2} \alpha_{2}^{2}\left[\left(\frac{\nabla_{i+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{j+1, k+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}+\frac{\nabla_{j+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{i+1, k+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}+\frac{\nabla_{k+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{i+1, j+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}\right) \frac{\nabla_{l+1} f_{t}^{*}}{f_{t}^{*}}\right. \\
& +\left(\frac{\nabla_{i+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{j+1, l+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}+\frac{\nabla_{j+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{i+1, l+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}+\frac{\nabla_{l+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{i+1, j+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}\right) \frac{\nabla_{k+1} f_{t}^{*}}{f_{t}^{*}} \\
& +\left(\frac{\nabla_{i+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{k+1, l+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}+\frac{\nabla_{k+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{i+1, l+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}+\frac{\nabla_{l+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{i+1, k+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}\right) \frac{\nabla_{j+1} f_{t}^{*}}{f_{t}^{*}} \\
& +\left(\frac{\nabla_{j+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{k+1, l+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}+\frac{\nabla_{k+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{j+1, l+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}+\frac{\nabla_{l+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{j+1, k+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}}\right) \frac{\nabla_{i+1} f_{t}^{*}}{f_{t}^{*}} \\
& +\left(\frac{\nabla_{i+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{j+1} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{j+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{i+1} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{i+1, j+1} f_{t}^{*}}{f_{t}^{*}}\right) \\
& \times\left(\frac{\nabla_{k+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{l+1} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{l+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{k+1} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{k+1, l+1}^{2} f_{t}^{*}}{f_{t}^{*}}\right) \\
& +\left(\frac{\nabla_{i+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{k+1} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{k+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{i+1} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{i+1, k+1} f_{t}^{*}}{f_{t}^{*}}\right) \\
& \times\left(\frac{\nabla_{j+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{l+1} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{l+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{j+1} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{j+1, l+1} f_{t}^{*}}{f_{t}^{*}}\right) \\
& +\left(\frac{\nabla_{i+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{l+1} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{l+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{i+1} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{i+1, l+1}^{2} f_{t}^{*}}{f_{t}^{*}}\right) \\
& \left.\times\left(\frac{\nabla_{j+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{k+1} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{k+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{j+1} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{j+1, k+1} f_{t}^{*}}{f_{t}^{*}}\right)\right] \\
& +\alpha_{1} \alpha_{2}\left(1-3 \alpha_{1} \alpha_{2}\right)\left[\frac{\nabla_{i+1, j+1, k+1}^{3} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{l+1} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{i+1, j+1, l+1}^{3} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{k+1} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{i+1, k+1, l+1}^{3} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{j+1} f_{t}^{*}}{f_{t}^{*}}\right. \\
& +\frac{\nabla_{j+1, k+1, l+1}^{3} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{i+1} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{i+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{j+1, k+1, l+1}^{3} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{j+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{i+1, k+1, l+1}^{3} f_{t}^{*}}{f_{t}^{*}} \\
& \left.+\frac{\nabla_{k+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{i+1, j+1, l+1}^{3} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{l+1} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{i+1, j+1, k+1}^{3} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{i+1, j+1, k+1, l+1}^{4} f_{t}^{*}}{f_{t}^{*}}\right] \\
& +\alpha_{1} \alpha_{2}\left(\alpha_{2}-\alpha_{1}\right)^{2}\left[\frac{\nabla_{i+1, j+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{k+1, l+1}^{2} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{i+1, k+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{j+1, l+1}^{2} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{i+1, l+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{j+1, k+1}^{2} f_{t}^{*}}{f_{t}^{*}}\right. \\
& \left.+\frac{\nabla_{j+1, k+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{i+1, l+1}^{2} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{j+1, l+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{i+1, k+1}^{2} f_{t}^{*}}{f_{t}^{*}}+\frac{\nabla_{k+1, l+1}^{2} \mathrm{n}_{p}^{*}}{\mathrm{n}_{p}^{*}} \frac{\nabla_{i+1, j+1}^{2} f_{t}^{*}}{f_{t}^{*}}\right]
\end{aligned}
$$

## F. 2 Fourth-order expansion of the log-likelihood function

Justification of (38) and expression of $R_{T}^{(1)}(\alpha, \tilde{\pi}, \varpi)$. Straightforward calculation yields the fourth-order Taylor expansion (38) with the remainder term (for brevity, we again write $\tilde{\pi}=(\beta, \pi)$ and $\left.\tilde{\pi}^{*}=\left(\beta^{*}, \pi^{*}\right)\right)$

$$
\begin{align*}
& R_{T}^{(1)}(\alpha, \tilde{\pi}, \varpi)=\varpi^{\prime} \nabla_{\varpi} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)+\frac{2}{2!}\left(\tilde{\pi}-\tilde{\pi}^{*}\right)^{\prime} \nabla_{\tilde{\pi} \varpi^{\prime}}^{2} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right) \varpi \\
&+\frac{1}{3!} \sum_{i=1}^{q_{1}+q_{2}} \sum_{j=1}^{q_{1}+q_{2}} \sum_{k=1}^{q_{1}+q_{2}} \nabla_{\tilde{\pi}_{i}}^{3} \tilde{\pi}_{j} \tilde{\pi}_{k} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)\left(\tilde{\pi}_{i}-\tilde{\pi}_{i}^{*}\right)\left(\tilde{\pi}_{j}-\tilde{\pi}_{j}^{*}\right)\left(\tilde{\pi}_{k}-\tilde{\pi}_{k}^{*}\right) \\
&+\frac{3}{3!} \sum_{i=1}^{q_{1}+q_{2}} \sum_{j=1}^{q_{1}+q_{2}} \sum_{k=1}^{q_{2}} \nabla_{\tilde{\pi}_{i} \tilde{\pi}_{j} \varpi_{k}}^{3} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)\left(\tilde{\pi}_{i}-\tilde{\pi}_{i}^{*}\right)\left(\tilde{\pi}_{j}-\tilde{\pi}_{j}^{*}\right) \varpi_{k} \\
&+\frac{1}{3!} \sum_{i=1}^{q_{2}} \sum_{j=1}^{q_{2}} \sum_{k=1}^{q_{2}} \nabla_{\varpi_{i}}^{3} \varpi_{j} \varpi_{k} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right) \varpi_{i} \varpi_{j} \varpi_{k} \\
&+\frac{1}{4!} \sum_{i=1}^{q_{1}+q_{2}} \sum_{j=1}^{q_{1}+q_{2}} \sum_{k=1}^{q_{1}+q_{2}} \sum_{l=1}^{q_{1}+q_{2}} \nabla_{\pi_{i}}^{4} \tilde{\pi}_{j} \tilde{\pi}_{k} \tilde{\pi}_{l} L_{T}^{\pi}(\alpha, \dot{\tilde{\pi}}, \dot{\varpi})\left(\tilde{\pi}_{i}-\tilde{\pi}_{i}^{*}\right)\left(\tilde{\pi}_{j}-\tilde{\pi}_{j}^{*}\right)\left(\tilde{\pi}_{k}-\tilde{\pi}_{k}^{*}\right)\left(\tilde{\pi}_{l}-\tilde{\pi}_{l}^{*}\right) \\
&+\frac{4}{4!} \sum_{i=1}^{q_{1}+q_{2}} \sum_{j=1}^{q_{1}+q_{2}} \sum_{k=1}^{q_{1}+q_{2}} \sum_{l=1}^{q_{2}} \nabla_{\tilde{\pi}_{i}}^{4} \tilde{\pi}_{j} \tilde{\pi}_{k} \varpi_{l} \\
& \sum_{i=1}^{\pi}(\alpha, \dot{\tilde{\pi}}, \dot{\varpi})\left(\tilde{\pi}_{i}-\tilde{\pi}_{i}^{*}\right)\left(\tilde{\pi}_{j}-\tilde{\pi}_{j}^{*}\right)\left(\tilde{\pi}_{k}-\tilde{\pi}_{k}^{*}\right) \varpi_{l} \\
&+\frac{4}{4!} \sum_{i=1}^{q_{1}+q_{2}} \sum_{k=1}^{q_{1}+q_{2}} \sum_{j=1}^{q_{2}} \sum_{k=1}^{q_{2}} \nabla_{\tilde{\pi}_{i}}^{4} \tilde{\pi}_{j} \varpi_{k} \varpi_{l} L_{T}^{\pi}(\alpha, \dot{\tilde{\pi}}, \dot{\varpi})\left(\tilde{\pi}_{i}-\tilde{\pi}_{i}^{*}\right)\left(\tilde{\pi}_{j}-\tilde{\pi}_{j}^{*}\right) \varpi_{k} \varpi_{l} \\
& q_{2} \varpi_{j} \varpi_{k} \varpi_{l} L_{T}^{\pi}(\alpha, \dot{\tilde{\pi}}, \dot{\varpi})\left(\tilde{\pi}_{i}-\tilde{\pi}_{i}^{*}\right) \varpi_{j} \varpi_{k} \varpi_{l}  \tag{48}\\
& q_{2} \\
&+\frac{1}{4!} \sum_{i=1}^{q_{2}} \sum_{j=1}^{q_{2}} \sum_{k=1}^{q_{2}} \sum_{l=1}^{q_{2}}\left(\nabla_{\varpi_{i} \varpi_{j} \varpi_{k} \varpi_{l}}^{4} L_{T}^{\pi}(\alpha, \dot{\tilde{\pi}}, \dot{\varpi})-\nabla_{\varpi_{i} \varpi_{j} \varpi_{k} \varpi_{l}} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)\right) \varpi_{i} \varpi_{j} \varpi_{k} \varpi_{l},
\end{align*}
$$

where $(\dot{\tilde{\pi}}, \dot{\varpi})$ denotes a point between $(\tilde{\pi}, \varpi)$ and $\left(\tilde{\pi}^{*}, 0\right)$.
Justification of (39) and expression of $R_{T}(\alpha, \tilde{\pi}, \varpi)$. We begin with some useful notation. Let $\mathfrak{I}$ denote the index set

$$
\mathfrak{I}=\left((1,1),(2,2), \ldots,\left(q_{2}, q_{2}\right),(1,2),(1,3), \ldots,\left(1, q_{2}\right),(2,3), \ldots,\left(q_{2}-1, q_{2}\right)\right) .
$$

For any scalars (or $d \times 1$ column vectors) $A_{i j}$ indexed by $i$ and $j$ (here and elsewhere it is tacitly assumed these indices belong to $\left.\left\{1, \ldots, q_{2}\right\}\right)$, let $\left[A_{i j}\right]_{(i, j) \in \mathcal{I}}$ denote the following $1 \times q_{2}\left(q_{2}+1\right) / 2$ row vector (or $d \times q_{2}\left(q_{2}+1\right) / 2$ matrix):

$$
\left[A_{i j}\right]_{(i, j) \in \mathcal{I}}=\left[A_{11}: \cdots: A_{q_{2} q_{2}}: A_{12}: \cdots: A_{q_{2}-1, q_{2}}\right] .
$$

For instance, $v(\varpi)=\left[\varpi_{i} \varpi_{j}\right]_{(i, j) \in \mathfrak{I}}^{\prime}$. Similarly, for any scalars $A_{i j k l}$ indexed by $i, j, k, l$, let $\left[A_{i j k l}\right]_{(i, j, k, l) \in \mathfrak{J} \times \mathfrak{J}}$
denote the following $q_{2}\left(q_{2}+1\right) / 2 \times q_{2}\left(q_{2}+1\right) / 2$ matrix

$$
\left[A_{i j k l}\right]_{(i, j, k, l) \in \mathfrak{I} \times \mathfrak{I}}=\left[\begin{array}{cccccc}
A_{1111} & \cdots & A_{11 q_{2} q_{2}} & A_{1112} & \cdots & A_{1,1, q_{2}-1, q_{2}} \\
\vdots & \ddots & \vdots & \vdots & \ddots & \vdots \\
A_{q_{2} q_{2} 11} & \cdots & A_{q_{2} q_{2} q_{2} q_{2}} & A_{q_{2} q_{2} 12} & \cdots & A_{q_{2}, q_{2}, q_{2}-1, q_{2}} \\
A_{1211} & \cdots & A_{12 q_{2} q_{2}} & A_{1212} & \cdots & A_{1,2, q_{2}-1, q_{2}} \\
\vdots & \ddots & \vdots & \vdots & \ddots & \vdots \\
A_{q_{2}-1, q_{2}, 1,1} & \cdots & A_{q_{2}-1, q_{2}, q_{2}, q_{2}} & A_{q_{2}-1, q_{2}, 1,2} & \cdots & A_{q_{2}-1, q_{2}, q_{2}-1, q_{2}}
\end{array}\right] .
$$

With this notation, and for any scalars $A_{i j k l}$ and $B_{i j}$ such that $A_{i j k l}=A_{j i k l}, A_{i j k l}=A_{i j l k}$, and $B_{i j}=B_{j i}$ for all $i, j, k, l$, it holds that ${ }^{10}$

$$
\begin{equation*}
\left[B_{i j}\right]_{(i, j) \in \mathfrak{J}}\left[c_{i j} c_{k l} A_{i j k l}\right]_{(i, j, k, l) \in \mathfrak{J} \times \mathfrak{I}}\left[B_{k l}\right]_{(k, l) \in \mathcal{I}}^{\prime}=\frac{1}{4} \sum_{i=1}^{q_{2}} \sum_{j=1}^{q_{2}} \sum_{k=1}^{q_{2}} \sum_{l=1}^{q_{2}} A_{i j k l} B_{i j} B_{k l}, \tag{49}
\end{equation*}
$$

where the $c_{i j}$ 's are as in Section 3.3.1 $\left(c_{i j}=1 / 2\right.$ if $i=j$ and $c_{i j}=1$ if $\left.i \neq j\right)$.
Now, to obtain (39), introduce the matrix

$$
\mathcal{J}_{T}=\left[\begin{array}{cc}
\mathcal{J}_{T, \tilde{\pi} \tilde{\pi}} & \mathcal{J}_{T, \tilde{\pi} \omega \varpi}^{\prime} \\
\mathcal{J}_{T, \tilde{\pi} \omega \omega} & \mathcal{J}_{T, \omega \omega \omega \mathrm{~m}}
\end{array}\right]
$$

where the matrices $\mathcal{J}_{T, \tilde{\pi} \tilde{\pi}}\left(\left(q_{1}+q_{2}\right) \times\left(q_{1}+q_{2}\right)\right), \mathcal{J}_{T, \tilde{\pi} \omega \omega}^{\prime}\left(\left(q_{1}+q_{2}\right) \times q_{\vartheta}\right)$, and $\mathcal{J}_{T, \omega \varpi \omega \varpi}\left(q_{\vartheta} \times q_{\vartheta}\right)$ are defined as follows (here $\nabla_{\tilde{\pi} \tilde{\pi}}^{2} L_{T}^{\pi *}$ stands for $\nabla_{\tilde{\tilde{\pi}} \tilde{\pi}^{2}}^{2} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)$ etc.)

$$
\begin{aligned}
& \mathcal{J}_{T, \tilde{\pi} \tilde{\pi}}=-T^{-1} \nabla_{\tilde{\pi} \tilde{\pi}^{\prime}}^{2} L_{T}^{\pi *}, \\
& \mathcal{J}_{T, \tilde{\pi} \varpi \varpi}^{\prime}=-T^{-1} \frac{1}{\alpha_{1} \alpha_{2}}\left[c_{i j} \nabla_{\tilde{\pi} \varpi_{i} \varpi_{j}}^{3} L_{T}^{\pi *}\right]_{(i, j) \in \mathcal{I}} \\
& =-T^{-1} \frac{1}{\alpha_{1} \alpha_{2}}\left[c_{11} \nabla_{\tilde{\pi} \omega_{1} \omega_{1}}^{3} L_{T}^{\pi *}: \cdots: c_{q_{2}-1, q_{2}} \nabla_{\tilde{\pi} \varpi_{q_{2}-1 \omega_{q_{2}}}^{3}}^{\left.L_{T}^{\pi *}\right],}\right. \\
& \mathcal{J}_{T, \varpi \varpi \varpi \varpi}=-T^{-1} \frac{8}{4!} \frac{1}{\alpha_{1}^{2} \alpha_{2}^{2}}\left[c_{i j} c_{k l} \nabla_{\varpi_{i} \varpi_{j} \omega_{k} \omega_{l}}^{4} L_{T}^{\pi *}\right]_{(i, j, k, l) \in \mathfrak{I} \times \mathcal{I}} \\
& =-T^{-1} \frac{8}{4!} \frac{1}{\alpha_{1}^{2} \alpha_{2}^{2}}\left[\begin{array}{ccc}
c_{11} c_{11} \nabla_{\varpi_{1} \varpi_{1} \varpi_{1} \varpi_{1}}^{4} L_{T}^{\pi *} & \cdots & c_{q_{2}-1, q_{2}} c_{11} \nabla_{\varpi_{q_{2}-1} \varpi_{q_{2}} \varpi_{1} \varpi_{1}}^{4} L_{T}^{\pi *} \\
\vdots & \ddots & \vdots \\
c_{11} c_{q_{2}-1, q_{2}} \nabla_{\varpi_{1} \varpi_{1} \varpi_{q_{2}-1} \varpi_{q_{2}}}^{4} L_{T}^{\pi *} & \cdots & c_{q_{2}-1, q_{2}} c_{q_{2}-1, q_{2}} \nabla_{\varpi_{q_{2}-1 \varpi_{q_{2}} \varpi_{q_{2}-1} \varpi_{q_{2}}}^{4}} L_{T}^{\pi *}
\end{array}\right] .
\end{aligned}
$$

[^7]where the properties $A_{i j k l}=A_{j i k l}, A_{i j k l}=A_{i j k}$, and $B_{i j}=B_{j i}$ for all $i, j, k, l$, are used in the last equality.

Straightforward computations (for the third one, use property (49)) now show that

$$
\begin{aligned}
& \quad-\frac{1}{2} T^{1 / 2}\left(\tilde{\pi}-\tilde{\pi}^{*}\right)^{\prime} \mathcal{J}_{T, \tilde{\pi} \tilde{\pi}} T^{1 / 2}\left(\tilde{\pi}-\tilde{\pi}^{*}\right)=\frac{1}{2}\left(\tilde{\pi}-\tilde{\pi}^{*}\right)^{\prime} \nabla_{\tilde{\pi} \tilde{\pi}^{\prime}}^{2} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)\left(\tilde{\pi}-\tilde{\pi}^{*}\right), \\
& -T^{1 / 2}\left(\tilde{\pi}-\tilde{\pi}^{*}\right)^{\prime} \mathcal{J}_{T, \tilde{\pi} \varpi \omega^{\prime}}^{\prime} T^{1 / 2} \alpha_{1} \alpha_{2} v(\varpi) \\
& =\left(\tilde{\pi}-\tilde{\pi}^{*}\right)^{\prime}\left[c_{11} \nabla_{\tilde{\pi} \varpi_{1} \varpi_{1}}^{3} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right) \varpi_{1}^{2}+\ldots+c_{q_{2}-1, q_{2}} \nabla_{\tilde{\pi} \varpi_{q_{2}-1} \varpi_{q_{2}}} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right) \varpi_{q_{2}-1} \varpi_{q_{2}}\right] \\
& =\frac{3}{3!} \sum_{i=1}^{q} \sum_{j=1}^{q_{2}} \sum_{k=1}^{q_{2}} \nabla_{\tilde{\pi}_{i} \varpi_{j} \varpi_{k}}^{3} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)\left(\tilde{\pi}_{i}-\tilde{\pi}_{i}^{*}\right) \varpi_{j} \varpi_{k}, \\
& - \\
& -\frac{1}{2} T \alpha_{1}^{2} \alpha_{2}^{2} v(\varpi)^{\prime} \mathcal{J}_{T, \varpi \varpi \varpi \varpi} v(\varpi)=\frac{1}{4!} \sum_{i=1}^{q_{2}} \sum_{j=1}^{q_{2}} \sum_{k=1}^{q_{2}} \sum_{l=1}^{q_{2}} \nabla_{\varpi_{i} \varpi_{j} \varpi_{k} \varpi_{l}}^{4} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right) \varpi_{i} \varpi_{j} \varpi_{k} \varpi_{l} .
\end{aligned}
$$

Therefore the fourth-order Taylor expansion of $L_{T}^{\pi}(\alpha, \tilde{\pi}, \varpi)$ in (38) can be written as a quadratic expansion given by

$$
\begin{equation*}
L_{T}^{\pi}(\alpha, \tilde{\pi}, \varpi)-L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)=S_{T}^{\prime} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)-\frac{1}{2}\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\right]^{\prime} \mathcal{J}_{T}\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\right]+R_{T}^{(1)}(\alpha, \tilde{\pi}, \varpi) . \tag{50}
\end{equation*}
$$

Next, define

$$
\mathcal{I}=\left[\begin{array}{cc}
\mathcal{I}_{\tilde{\pi} \tilde{\pi}} & \mathcal{I}_{\tilde{\pi} w \infty}^{\prime}  \tag{51}\\
\mathcal{I}_{\tilde{\pi} \omega m} & \mathcal{I}_{\varpi \omega w n}
\end{array}\right],
$$

where the matrices $\mathcal{I}_{\tilde{\pi} \tilde{\pi}}\left(\left(q_{1}+q_{2}\right) \times\left(q_{1}+q_{2}\right)\right)$, $\mathcal{I}_{\tilde{\pi} \varpi \omega}^{\prime}\left(\left(q_{1}+q_{2}\right) \times q_{\vartheta}\right)$, and $\mathcal{I}_{\varpi \omega \varpi \varpi}\left(q_{\vartheta} \times q_{\vartheta}\right)$ are defined as follows

$$
\begin{aligned}
\mathcal{I}_{\tilde{\pi} \tilde{\pi}} & =E\left[\frac{\nabla f_{t}\left(\tilde{\pi}^{*}\right)}{f_{t}\left(\tilde{\pi}^{*}\right)} \frac{\nabla^{\prime} f_{t}\left(\tilde{\pi}^{*}\right)}{f_{t}\left(\tilde{\pi}^{*}\right)}\right], \\
\mathcal{I}_{\tilde{\pi} \varpi \varpi}^{\prime} & =\left[c_{i j} E\left[\frac{\nabla f_{t}\left(\tilde{\pi}^{*}\right)}{f_{t}\left(\tilde{\pi}^{*}\right)} X_{t, i, j}^{*}\right]\right]_{(i, j) \in \mathfrak{I}}, \\
\mathcal{I}_{\varpi \varpi \varpi \varpi} & =\left[c_{i j} c_{k l} E\left[X_{t, i, j}^{*} X_{t, k, l}^{*}\right]\right]_{(i, j, k, l) \in \mathfrak{I} \times \mathfrak{I}},
\end{aligned}
$$

and where we have used the short-hand notation $X_{t, i, j}^{*}, i, j \in\left\{1, \ldots, q_{2}\right\}$ (see (44)). Finiteness of $\mathcal{I}$ follows from Lemma F.1. Now, defining

$$
R_{T}^{(2)}(\alpha, \tilde{\pi}, \varpi)=-\frac{1}{2}\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\right]^{\prime}\left(\mathcal{J}_{T}-\mathcal{I}\right)\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\right]
$$

and adding and subtracting terms, expansion (50) can be written as (39) with

$$
R_{T}(\alpha, \tilde{\pi}, \varpi)=R_{T}^{(1)}(\alpha, \tilde{\pi}, \varpi)+R_{T}^{(2)}(\alpha, \tilde{\pi}, \varpi) .
$$

## F. 3 Some more explicit derivatives and the verification of Assumption 5(iii)

Some more explicit derivative expressions. We will require more explicit expressions for the components of $s_{t}$ in (43) (see also (40)-(42) and (44)). Straightforward computation shows that (as before, $\nabla f_{t}(\cdot)$ denotes differentiation of $f_{t}(\cdot)$ in (3) with respect to $\left.\tilde{\phi}=\left(\tilde{\phi}_{0}, \tilde{\phi}_{1}, \ldots, \tilde{\phi}_{p}, \tilde{\sigma}_{1}^{2}\right)\right)$

$$
\frac{\nabla f_{t}(\tilde{\phi})}{f_{t}(\tilde{\phi})}=\nabla \log \left(f_{t}(\tilde{\phi})\right), \quad \frac{\nabla^{2} f_{t}(\tilde{\phi})}{f_{t}(\tilde{\phi})}=\nabla^{2} \log \left(f_{t}(\tilde{\phi})\right)+\nabla \log \left(f_{t}(\tilde{\phi})\right) \nabla^{\prime} \log \left(f_{t}(\tilde{\phi})\right)
$$

where $\left(\operatorname{as} \log \left(f_{t}(\tilde{\phi})\right)=-\frac{1}{2} \log (2 \bar{\pi})-\frac{1}{2} \log \left(\tilde{\sigma}_{1}^{2}\right)-\frac{1}{2} g_{t}^{2}(\tilde{\phi})\right.$ with $\left.g_{t}(\tilde{\phi})=\left[y_{t}-\left(\tilde{\phi}_{0}+\tilde{\phi}_{1} y_{t-1}+\cdots+\tilde{\phi}_{p} y_{t-p}\right)\right] / \tilde{\sigma}_{1}\right)$ $\nabla \log \left(f_{t}(\tilde{\phi})\right)=\left[\begin{array}{c}\frac{1}{\tilde{\sigma}_{1}} g_{t}(\tilde{\phi}) \\ \frac{1}{\tilde{\sigma}_{1}} \boldsymbol{y}_{t-1} g_{t}(\tilde{\phi}) \\ \frac{1}{2 \tilde{\sigma}_{1}^{2}}\left(g_{t}^{2}(\tilde{\phi})-1\right)\end{array}\right], \quad \nabla^{2} \log \left(f_{t}(\tilde{\phi})\right)=\left[\begin{array}{ccc}-\frac{1}{\tilde{\sigma}_{1}^{2}} & -\frac{1}{\tilde{\sigma}_{1}^{2}} \boldsymbol{y}_{t-1}^{\prime} & -\frac{1}{\tilde{\sigma}_{1}^{3}} g_{t}(\tilde{\phi}) \\ -\frac{1}{\tilde{\sigma}_{1}^{2}} \boldsymbol{y}_{t-1} & -\frac{1}{\tilde{\sigma}_{1}^{2}} \boldsymbol{y}_{t-1} \boldsymbol{y}_{t-1}^{\prime} & -\frac{1}{\tilde{\sigma}_{1}^{3}} g_{t}(\tilde{\phi}) \boldsymbol{y}_{t-1} \\ -\frac{1}{\tilde{\sigma}_{1}^{3}} g_{t}(\tilde{\phi}) & -\frac{1}{\tilde{\sigma}_{1}^{3}} g_{t}(\tilde{\phi}) \boldsymbol{y}_{t-1}^{\prime} & -\frac{1}{2 \tilde{\sigma}_{1}^{4}}\left(2 g_{t}^{2}(\tilde{\phi})-1\right)\end{array}\right]$, so that

$$
\begin{align*}
& \frac{\nabla f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)}=\left[\begin{array}{c}
\frac{1}{\tilde{\sigma}_{1}^{*}} \varepsilon_{t} \\
\frac{1}{\tilde{\sigma}_{1}^{*}} \boldsymbol{y}_{t-1} \varepsilon_{t} \\
\frac{\tilde{\sigma}_{1}^{* 2}}{2 \varepsilon_{1}}\left(\varepsilon_{t}^{2}-1\right)
\end{array}\right], \\
& \frac{\nabla^{2} f_{t}\left(\pi^{*}\right)}{f_{t}\left(\pi^{*}\right)}=\left[\begin{array}{ccc}
-\frac{1}{\tilde{\sigma}_{1}^{* 2}} & -\frac{1}{\tilde{\sigma}_{1}^{* 2}} \boldsymbol{y}_{t-1}^{\prime} & -\frac{1}{\tilde{\sigma}_{1}^{* 3}} \varepsilon_{t} \\
-\frac{1}{\tilde{\sigma}_{1}^{* 2}} \boldsymbol{y}_{t-1} & -\frac{1}{\tilde{\sigma}_{1}^{* 2}} \boldsymbol{y}_{t-1} \boldsymbol{y}_{t-1}^{\prime} & -\frac{1}{\tilde{\sigma}_{1}^{* 3}} \boldsymbol{y}_{t-1} \varepsilon_{t} \\
-\frac{1}{\tilde{\sigma}_{1}^{* 3}} \varepsilon_{t} & -\frac{1}{\tilde{\sigma}_{1}^{* 3}} \boldsymbol{y}_{t-1}^{\prime} \varepsilon_{t} & -\frac{1}{2 \tilde{\sigma}_{1}^{* 4}}\left(2 \varepsilon_{t}^{2}-1\right)
\end{array}\right]+\left[\begin{array}{c}
\frac{1}{\tilde{\sigma}_{1}^{*}} \varepsilon_{t} \\
\frac{1}{\tilde{\sigma}_{1}^{*}} \boldsymbol{y}_{t-1} \varepsilon_{t} \\
\frac{1}{2 \tilde{\sigma}_{1}^{* 2}}\left(\varepsilon_{t}^{2}-1\right)
\end{array}\right]\left[\begin{array}{c}
\frac{1}{\tilde{\sigma}_{1}^{*}} \varepsilon_{t} \\
\frac{1}{\tilde{\sigma}_{1}^{*}} \boldsymbol{y}_{t-1} \varepsilon_{t} \\
\frac{1}{2 \tilde{\sigma}_{1}^{* 2}}\left(\varepsilon_{t}^{2}-1\right)
\end{array}\right]^{\prime} \\
&=\left[\begin{array}{ccc}
\frac{1}{\tilde{\sigma}_{1}^{* 2}}\left(\varepsilon_{t}^{2}-1\right) & \frac{1}{\tilde{\sigma}_{1}^{* 2}} \boldsymbol{y}_{t-1}^{\prime}\left(\varepsilon_{t}^{2}-1\right) & \frac{1}{2 \tilde{\sigma}_{1}^{* 3}}\left(\varepsilon_{t}^{3}-3 \varepsilon_{t}\right) \\
\frac{1}{\tilde{\sigma}_{1}^{* 2}} \boldsymbol{y}_{t-1}\left(\varepsilon_{t}^{2}-1\right) & \frac{1}{\tilde{\sigma}_{1}^{* 2}} \boldsymbol{y}_{t-1} \boldsymbol{y}_{t-1}^{\prime}\left(\varepsilon_{t}^{2}-1\right) & \frac{1}{2 \tilde{\sigma}_{1}^{* 3}} \boldsymbol{y}_{t-1}^{2}\left(\varepsilon_{t}^{3}-3 \varepsilon_{t}\right) \\
\frac{1}{2 \tilde{\sigma}_{1}^{* 3}}\left(\varepsilon_{t}^{3}-3 \varepsilon_{t}\right) & \frac{1}{2 \tilde{\sigma}_{1}^{* 3}} \boldsymbol{y}_{t-1}^{\prime}\left(\varepsilon_{t}^{3}-3 \varepsilon_{t}\right) & \frac{1}{4 \tilde{\sigma}_{1}^{* 4}}\left(\varepsilon_{t}^{4}-6 \varepsilon_{t}^{2}+3\right)
\end{array}\right] . \tag{52}
\end{align*}
$$

Similar formulas hold for $\mathrm{n}_{p}\left(\boldsymbol{y}_{t-1} ; \tilde{\phi}\right)$. As $\log \left(\mathrm{n}_{p}\left(\boldsymbol{y}_{t-1} ; \tilde{\phi}\right)\right)=-\frac{p}{2} \log (2 \bar{\pi})-\frac{1}{2} \log \left(\operatorname{det}\left(\boldsymbol{\Gamma}_{1, p}\right)\right)-\frac{1}{2}\left(\boldsymbol{y}_{t-1}-\right.$ $\left.\mu_{1} \mathbf{1}_{p}\right)^{\prime} \boldsymbol{\Gamma}_{1, p}^{-1}\left(\boldsymbol{y}_{t-1}-\mu_{1} \mathbf{1}_{p}\right)$ (see Section 2.2), we obtain, for each $i=1, \ldots, p+2$, (cf. Magnus and Neudecker (1999, p. 325)) ${ }^{11}$
$\nabla_{i} \log \left(\mathrm{n}_{p}\left(\boldsymbol{y}_{t-1} ; \tilde{\phi}\right)\right)=\frac{1}{2} \operatorname{tr}\left(\frac{\partial \boldsymbol{\Gamma}_{1, p}^{-1}}{\partial \tilde{\phi}_{i}} \boldsymbol{\Gamma}_{1, p}\right)+\left(\boldsymbol{y}_{t-1}-\mu_{1} \mathbf{1}_{p}\right)^{\prime} \boldsymbol{\Gamma}_{1, p}^{-1}\left(\frac{\partial \mu_{1}}{\partial \tilde{\phi}_{i}} \mathbf{1}_{p}\right)-\frac{1}{2}\left(\boldsymbol{y}_{t-1}-\mu_{1} \mathbf{1}_{p}\right)^{\prime} \frac{\partial \boldsymbol{\Gamma}_{1, p}^{-1}}{\partial \tilde{\phi}_{i}}\left(\boldsymbol{y}_{t-1}-\mu_{1} \mathbf{1}_{p}\right)$, where (see Section 2.2 for the notation)

$$
\frac{\partial \mu_{1}}{\partial \tilde{\phi}}=\left[\begin{array}{c}
(\tilde{\phi}(1))^{-1} \\
\tilde{\phi}_{0}(\tilde{\phi}(1))^{-2} \mathbf{1}_{p} \\
0
\end{array}\right]
$$

For the expression of $\frac{\partial \boldsymbol{\Gamma}_{1, p}^{-1}}{\partial \hat{\phi}_{i}}$, first note that $\boldsymbol{\Gamma}_{1, p}^{-1}$ can be expressed as (see, e.g., Galbraith and Galbraith (1974)) $\boldsymbol{\Gamma}_{1, p}^{-1}=\frac{1}{\tilde{\sigma}_{1}^{2}}\left(U^{\prime} U-V^{\prime} V\right)$ with $U$ and $V$ being $p \times p$ Toeplitz matrices given by

$$
U=\left[\begin{array}{cccc}
1 & & & \\
-\tilde{\phi}_{1} & \ddots & & \\
\vdots & \ddots & \ddots & \\
-\tilde{\phi}_{p-1} & \cdots & -\tilde{\phi}_{1} & 1
\end{array}\right], \quad V=\left[\begin{array}{cccc}
\tilde{\phi}_{p} & & & \\
\tilde{\phi}_{p-1} & \ddots & & \\
\vdots & \ddots & \ddots & \\
\tilde{\phi}_{1} & \cdots & \tilde{\phi}_{p-1} & \tilde{\phi}_{p}
\end{array}\right]
$$

Thus $\frac{\partial \boldsymbol{\Gamma}_{1, p}^{-1}}{\partial \tilde{\phi}_{i}}$ equals a zero matrix when differentiating with respect to $\tilde{\phi}_{0},-\frac{1}{\tilde{\sigma}_{1}^{2}} \boldsymbol{\Gamma}_{1, p}^{-1}$ when differentiating with respect to $\tilde{\sigma}_{1}^{2}$, and

$$
\frac{\partial \boldsymbol{\Gamma}_{1, p}^{-1}}{\partial \tilde{\phi}_{i}}=\frac{1}{\tilde{\sigma}_{1}^{2}}\left(\frac{\partial U^{\prime}}{\partial \tilde{\phi}_{i}} U+U^{\prime} \frac{\partial U}{\partial \tilde{\phi}_{i}}-\frac{\partial V^{\prime}}{\partial \tilde{\phi}_{i}} V-V^{\prime} \frac{\partial V}{\partial \tilde{\phi}_{i}}\right)
$$

[^8]when differentiating with respect to the autoregressive parameters. To summarize,
\[

\frac{\nabla \mathrm{n}_{p}(\tilde{\phi})}{\mathrm{n}_{p}(\tilde{\phi})}=\nabla \log \left(\mathrm{n}_{p}\left(\boldsymbol{y}_{t-1} ; \tilde{\phi}\right)\right)=\left[$$
\begin{array}{l}
d_{1}\left(\boldsymbol{y}_{t-1} ; \tilde{\phi}\right)  \tag{53}\\
d_{2}\left(\boldsymbol{y}_{t-1} ; \tilde{\phi}\right) \\
d_{3}\left(\boldsymbol{y}_{t-1} ; \tilde{\phi}\right)
\end{array}
$$\right],
\]

where (note that $\operatorname{tr}\left(\frac{\partial \boldsymbol{\Gamma}_{1, p}^{-1}}{\partial \dot{\phi}_{i}} \boldsymbol{\Gamma}_{1, p}\right)=\frac{\partial v e c\left(\boldsymbol{\Gamma}_{1, p}^{-1}\right)^{\prime}}{\partial \operatorname{d}} \operatorname{vec}\left(\boldsymbol{\Gamma}_{1, p}\right)$ )

$$
\begin{aligned}
d_{1}\left(\boldsymbol{y}_{t-1} ; \tilde{\phi}\right) & =(\tilde{\phi}(1))^{-1}\left(\boldsymbol{y}_{t-1}-\mu_{1} \mathbf{1}_{p}\right)^{\prime} \boldsymbol{\Gamma}_{1, p}^{-1} \mathbf{1}_{p}, \\
d_{2}\left(\boldsymbol{y}_{t-1} ; \tilde{\phi}\right) & =\frac{1}{2} \frac{\partial v e c\left(\boldsymbol{\Gamma}_{1, p}^{-1}\right)^{\prime}}{\partial\left(\tilde{\phi}_{1}, \ldots, \tilde{\phi}_{p}\right)} \operatorname{vec}\left(\boldsymbol{\Gamma}_{1, p}\right)+\tilde{\phi}_{0}(\tilde{\phi}(1))^{-2} \mathbf{1}_{p}\left(\mathbf{1}_{p}^{\prime} \boldsymbol{\Gamma}_{1, p}^{-1}\left(\boldsymbol{y}_{t-1}-\mu_{1} \mathbf{1}_{p}\right)\right) \\
& -\frac{1}{2} \frac{\partial v e c\left(\boldsymbol{\Gamma}_{1, p}^{-1}\right)^{\prime}}{\partial\left(\tilde{\phi}_{1}, \ldots, \tilde{\phi}_{p}\right)}\left(\left(\boldsymbol{y}_{t-1}-\mu_{1} \mathbf{1}_{p}\right) \otimes\left(\boldsymbol{y}_{t-1}-\mu_{1} \mathbf{1}_{p}\right)\right), \\
d_{3}\left(\boldsymbol{y}_{t-1} ; \tilde{\phi}\right) & =-\frac{p}{2 \tilde{\sigma}_{1}^{2}}+\frac{1}{2 \tilde{\sigma}_{1}^{2}}\left(\boldsymbol{y}_{t-1}-\mu_{1} \mathbf{1}_{p}\right)^{\prime} \boldsymbol{\Gamma}_{1, p}^{-1}\left(\boldsymbol{y}_{t-1}-\mu_{1} \mathbf{1}_{p}\right)
\end{aligned}
$$

(first and last scalars, middle one $p \times 1$ ). Therefore

$$
\frac{\nabla \mathrm{n}_{p}\left(\pi^{*}\right)}{\mathrm{n}_{p}\left(\pi^{*}\right)}=\nabla \log \left(\mathrm{n}_{p}\left(\boldsymbol{y}_{t-1} ; \pi^{*}\right)\right)=\left[\begin{array}{l}
d_{1}\left(\boldsymbol{y}_{t-1} ; \pi^{*}\right) \\
d_{2}\left(\boldsymbol{y}_{t-1} ; \pi^{*}\right) \\
d_{3}\left(\boldsymbol{y}_{t-1} ; \pi^{*}\right)
\end{array}\right] .
$$

Based on the preceding derivations, the derivatives appearing in (40)-(42) can now be expressed as

$$
\begin{aligned}
& \nabla_{\beta} l_{t}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right)=\frac{1}{\tilde{\sigma}_{1}^{*}} \varepsilon_{t} \\
& \nabla_{\pi} l_{t}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right)=\left[\begin{array}{c}
\frac{1}{\tilde{\sigma}_{1}^{*}} \\
\frac{1}{2 \tilde{\sigma}_{1}^{* 2}}\left(\varepsilon_{t-1}^{2}-1\right)
\end{array}\right] \quad((p+1) \times 1) \\
& \nabla_{w \omega^{\prime}}^{2} l_{t}^{\pi}\left(\alpha, \beta^{*}, \pi^{*}, 0\right)=\alpha(1-\alpha)\left\{\left[\begin{array}{ll}
d_{2}\left(\boldsymbol{y}_{t-1} ; \pi^{*}\right) \frac{1}{\tilde{\sigma}_{1}^{*}} \boldsymbol{y}_{t-1}^{\prime} \varepsilon_{t} & d_{2}\left(\boldsymbol{y}_{t-1} ; \pi^{*}\right) \frac{1}{2 \tilde{\sigma}^{* 2}}\left(\varepsilon_{t}^{2}-1\right) \\
d_{3}\left(\boldsymbol{y}_{t-1} ; \pi^{*}\right) \frac{1}{\tilde{\sigma}_{1}^{*}} \boldsymbol{y}_{t-1}^{\prime} \varepsilon_{t} & d_{3}\left(\boldsymbol{y}_{t-1} ; \pi^{*}\right) \frac{1}{2 \tilde{\sigma}_{1}^{* 2}}\left(\varepsilon_{t}^{2}-1\right)
\end{array}\right]\right. \\
& +\left[\begin{array}{cc}
\frac{1}{\tilde{\sigma}_{1}^{*}} \boldsymbol{y}_{t-1} \varepsilon_{t} d_{2}^{\prime}\left(\boldsymbol{y}_{t-1} ; \pi^{*}\right) & \frac{1}{\tilde{\sigma}_{1}^{*}} \boldsymbol{y}_{t-1} \varepsilon_{t} d_{3}\left(\boldsymbol{y}_{t-1} ; \pi^{*}\right) \\
\frac{1}{2 \tilde{\sigma}_{1}^{* 2}}\left(\varepsilon_{t}^{2}-1\right) d_{2}^{\prime}\left(\boldsymbol{y}_{t-1} ; \pi^{*}\right) & \frac{1}{2 \tilde{\sigma}_{1}^{* 2}}\left(\varepsilon_{t}^{2}-1\right) d_{3}\left(\boldsymbol{y}_{t-1} ; \pi^{*}\right)
\end{array}\right] \\
& \left.+\left[\begin{array}{cc}
\frac{1}{\tilde{\sigma}^{* 2}} \boldsymbol{y}_{t-1} \boldsymbol{y}_{t-1}^{\prime}\left(\varepsilon_{t}^{2}-1\right) & \frac{1}{2 \tilde{\sigma}_{1}^{* 3}} \boldsymbol{y}_{t-1}\left(\varepsilon_{t}^{3}-3 \varepsilon_{t}\right) \\
\frac{1}{2 \tilde{\sigma}_{1}^{3}} \boldsymbol{y}_{t-1}^{\prime}\left(\varepsilon_{t}^{3}-3 \varepsilon_{t}\right) & \frac{1}{4 \tilde{\sigma}_{1}^{* 4}}\left(\varepsilon_{t}^{4}-6 \varepsilon_{t}^{2}+3\right)
\end{array}\right]\right\} .
\end{aligned}
$$

Verification of Assumption 5(iii). Finiteness of $\mathcal{I}$ was already established in Supplementary Appendix F.2. For positive definiteness, it suffices to show that the components of the vector $s_{t}$ are linearly independent. Note that for linear independence, it does not matter if the order of the elements is changed or if some of the elements are multiplied by nonzero constants. Therefore, making use of the explicit expressions given above, it suffices to show that the components of the vector $\tilde{s}_{t}=\left(\tilde{s}_{t, 1}, \tilde{s}_{t, 2}, \tilde{s}_{t, 3}, \tilde{s}_{t, 4}, \tilde{s}_{t, 5}, \tilde{s}_{t, 6}\right)$ (where the dimensions of the six components are $1, p, 1, p(p+1) / 2, p, 1$,
respectively) are linearly independent, where

$$
\left[\begin{array}{c}
\tilde{s}_{t, 1} \\
\tilde{s}_{t, 2} \\
\tilde{s}_{t, 3} \\
\tilde{s}_{t, 4}^{*} \\
\tilde{s}_{t, 5} \\
\tilde{s}_{t, 6}
\end{array}\right]=\left[\begin{array}{c}
\frac{1}{\tilde{\sigma}_{1}^{*}} \boldsymbol{y}_{t-1} \varepsilon_{t} \\
\frac{\tilde{\sigma}_{1}^{* 2}}{2\left(\varepsilon_{t}^{2}-1\right)} \\
v e c h\left[d_{2}\left(\boldsymbol{y}_{t-1} ; \pi^{*}\right) \boldsymbol{y}_{t-1}^{\prime}+\boldsymbol{y}_{t-1} d_{2}^{\prime}\left(\boldsymbol{y}_{t-1} ; \pi^{*}\right) \frac{1}{\tilde{\sigma}_{1}^{*}} \varepsilon_{t}+v e c h\left[\boldsymbol{y}_{t-1} \boldsymbol{y}_{t-1}^{\prime}\right] \frac{1}{\tilde{\sigma}_{3}^{* 2}}\left(\varepsilon_{t}^{2}-1\right)\right. \\
d_{2}\left(\boldsymbol{y}_{t-1} ; \pi^{*}\right) \frac{1}{2 \tilde{\sigma}_{1}^{* 2}}\left(\varepsilon_{t}^{2}-1\right)+\boldsymbol{y}_{t-1} d_{3}\left(\boldsymbol{y}_{t-1} ; \pi^{*}\right) \frac{1}{\tilde{\sigma}_{1}^{*}} \varepsilon_{t}+\frac{1}{2 \tilde{\sigma}_{13}^{* 3}} \boldsymbol{y}_{t-1}\left(\varepsilon_{t}^{3}-3 \varepsilon_{t}\right) \\
d_{3}\left(\boldsymbol{y}_{t-1} ; \pi^{*}\right) \frac{1}{\tilde{\sigma}_{1}^{* 2}}\left(\varepsilon_{t}^{2}-1\right)+\frac{1}{4 \tilde{\sigma}_{1}^{* 4}}\left(\varepsilon_{t}^{4}-6 \varepsilon_{t}^{2}+3\right)
\end{array}\right] .
$$

To this end, suppose that $c^{\prime} \tilde{s}_{t}=\left(c_{1}, c_{2}, c_{3}, c_{4}, c_{5}, c_{6}\right)^{\prime}\left(\tilde{s}_{t, 1}, \tilde{s}_{t, 2}, \tilde{s}_{t, 3}, \tilde{s}_{t, 4}, \tilde{s}_{t, 5}, \tilde{s}_{t, 6}\right)=0$ (with the dimension of $c$ and its subvectors chosen conformably). Note that the only random quantities $\tilde{s}_{t}$ depends on are $\boldsymbol{y}_{t-1}$ and $\varepsilon_{t}$ which are independent. First, as the term $\varepsilon_{t}^{4}$ only appears in $\tilde{s}_{t, 6}$, the equality $E\left[c^{\prime} \tilde{s}_{t} \mid \varepsilon_{t}\right]=0$ can be expressed as $c_{6} \varepsilon_{t}^{4} /\left(4 \tilde{\sigma}_{1}^{* 4}\right)+P_{3}\left(\varepsilon_{t}\right)=0$, where $P_{3}\left(\varepsilon_{t}\right)$ is a third-order polynomial in $\varepsilon_{t}$. As the components of the vector $\left(\varepsilon_{t}, \varepsilon_{t}^{2}, \varepsilon_{t}^{3}, \varepsilon_{t}^{4}\right)$ are linearly independent (this clearly follows from normality), it follows that $c_{6}=0$. Next, basic properties of the standard normal distribution imply that $E\left[c^{\prime} \tilde{s}_{t}\left(\varepsilon_{t}^{3}-3 \varepsilon_{t}\right) \mid \boldsymbol{y}_{t-1}\right]=c_{5}^{\prime} \frac{1}{2 \tilde{\sigma}_{1}^{* 3}} \boldsymbol{y}_{t-1} E\left[\left(\varepsilon_{t}^{3}-3 \varepsilon_{t}\right)^{2}\right]=0$, so that necessarily $c_{5}=0$ (as the components of $\boldsymbol{y}_{t-1}$ are linearly independent and $E\left[\left(\varepsilon_{t}^{3}-3 \varepsilon_{t}\right)^{2}\right]>0$ ). Next note that (as $c_{5}=0$, $c_{6}=0$ )

$$
0=E\left[c^{\prime} \tilde{s}_{t}\left(\varepsilon_{t}^{2}-1\right) \mid \boldsymbol{y}_{t-1}\right]=c_{3} E\left[\left(\varepsilon_{t}^{2}-1\right)^{2}\right] /\left(2 \tilde{\sigma}_{1}^{* 2}\right)+c_{4}^{\prime} \operatorname{vech}\left[\boldsymbol{y}_{t-1} \boldsymbol{y}_{t-1}^{\prime}\right] E\left[\left(\varepsilon_{t}^{2}-1\right)^{2}\right] / \tilde{\sigma}_{1}^{* 2}
$$

so that $c_{4}^{\prime} \operatorname{vech}\left[\boldsymbol{y}_{t-1} \boldsymbol{y}_{t-1}^{\prime}\right]=-c_{3} / 2$. As the components of $\operatorname{vech}\left[\boldsymbol{y}_{t-1} \boldsymbol{y}_{t-1}^{\prime}\right]$ are linearly independent (as $\operatorname{vech}\left(\boldsymbol{y}_{t-1} \boldsymbol{y}_{t-1}^{\prime}\right)=D_{p}^{+} \operatorname{vec}\left(\boldsymbol{y}_{t-1} \boldsymbol{y}_{t-1}^{\prime}\right)=D_{p}^{+}\left(\boldsymbol{y}_{t-1} \otimes \boldsymbol{y}_{t-1}\right)$, with $D_{p}^{+}$denoting the Moore-Penrose inverse of the duplication matrix $D_{p}, \operatorname{Cov}\left[\operatorname{vech}\left(\boldsymbol{y}_{t-1} \boldsymbol{y}_{t-1}^{\prime}\right)\right]=D_{p}^{+} \operatorname{Cov}\left(\boldsymbol{y}_{t-1} \otimes \boldsymbol{y}_{t-1}\right) D_{p}^{+\prime}$; because $D_{p}^{+}$is of full row rank and $\operatorname{Cov}\left(\boldsymbol{y}_{t-1} \otimes \boldsymbol{y}_{t-1}\right)$ has rank $p(p+1) / 2$, see Thm 4.3(v) of Magnus and Neudecker (1979), $\operatorname{Cov}\left[\operatorname{vech}\left(\boldsymbol{y}_{t-1} \boldsymbol{y}_{t-1}^{\prime}\right)\right]$ is positive definite), it necessarily follows that $c_{4}=0$ and $c_{3}=0$. Finally, as only $c_{1}$ and $c_{2}$ may be nonzero, $E\left[c^{\prime} \tilde{s}_{t} \varepsilon_{t} \mid \boldsymbol{y}_{t-1}\right]=c_{1} \frac{1}{\tilde{\sigma}_{1}^{*}}+c_{2}^{\prime} \frac{1}{\tilde{\sigma}_{1}^{*}} \boldsymbol{y}_{t-1}=0$, from which $c_{2}=0$ and $c_{1}=0$ follow (as the components of $\boldsymbol{y}_{t-1}$ are linearly independent).

## F. 4 Verification of Assumption 5(iv)

First consider $R_{T}^{(1)}(\alpha, \tilde{\pi}, \varpi)$. Of the quantities on the right hand side of (48), the first two are equal to zero because $\nabla_{\varpi} l_{t}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)=0$ and $\nabla_{\tilde{\pi} \omega^{\prime}} l_{t}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)=0$; for the other eight quantities, Lemma F. 4 provides upper bounds that aid in bounding them. Now, to verify Assumption 5(iv) (for $R_{T}^{(1)}(\alpha, \tilde{\pi}, \varpi)$ ), let $\left\{\gamma_{T}, T \geq 1\right\}$ be an arbitrary sequence of (non-random) positive scalars such that $\gamma_{T} \rightarrow 0$ as $T \rightarrow \infty$. Condition $\left\|(\beta, \pi, \varpi)-\left(\beta^{*}, \pi^{*}, 0\right)\right\| \leq \gamma_{T}$ (appearing in Assumption 5(iv)), together with the properties $\|\boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)\|^{2}=\left\|\tilde{\pi}-\tilde{\pi}^{*}\right\|^{2}+\alpha_{1}^{2} \alpha_{2}^{2}\|v(\varpi)\|^{2}$ and $\|v(\varpi)\| \leq\|\varpi\|^{2}$, implies that

$$
\|\boldsymbol{\theta}(\alpha, \beta, \pi, \varpi)\|^{1 / 2} \leq\left\|\tilde{\pi}-\tilde{\pi}^{*}\right\|^{1 / 2}+\alpha_{1}^{1 / 2} \alpha_{2}^{1 / 2}\|v(\varpi)\|^{1 / 2} \leq \gamma_{T}^{1 / 2}+\alpha_{1}^{1 / 2} \alpha_{2}^{1 / 2} \gamma_{T} .
$$

This, Lemma F.4, and the fact that $\alpha$ is bounded away from zero and one on $A$, imply that for some sequence $\left\{\tilde{\gamma}_{T}, T \geq 1\right\}$ of (non-random) positive scalars such that $\tilde{\gamma}_{T} \rightarrow 0$ as $T \rightarrow \infty$ and for some
finite $C$,

$$
\begin{align*}
& \sup _{\alpha \in A} \sup _{\substack{(\tilde{\pi}, \varpi) \in B \times \Pi_{\alpha},\left\|(\tilde{\pi}, \varpi)-\left(\tilde{\pi}^{*}, 0\right)\right\| \leq \gamma_{T}}} \frac{\left|R_{T}^{(1)}(\alpha, \tilde{\pi}, \varpi)\right|}{\left(1+\left\|T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\right\|\right)^{2}} \\
& \leq \tilde{\gamma}_{T} \sum_{i, j, k} \sup _{\alpha \in A}\left[\left|T^{-1} \nabla_{\tilde{\pi}_{i} \tilde{\pi}_{j} \tilde{\pi}_{k}}^{3} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)\right|+\left|T^{-1} \nabla_{\tilde{\pi}_{i} \tilde{\pi}_{j} \varpi_{k}}^{3} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)\right|+\left|T^{-1 / 2} \nabla_{\varpi_{i} \varpi_{j} \varpi_{k}}^{3} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)\right|\right] \\
& +\tilde{\gamma}_{T} \sum_{i, j, k, l} \sup _{\alpha \in A} \sup _{\substack{(\tilde{\pi}, \varpi) \in B \times \Pi_{\alpha},\left\|(\tilde{\pi}, \varpi)-\left(\tilde{\pi}^{*}, 0\right)\right\| \leq \gamma_{T}}}\left[\left|T^{-1} \nabla_{\tilde{\pi}_{i} \tilde{\pi}_{j} \tilde{\pi}_{k} \tilde{\pi}_{l}}^{4} L_{T}^{\pi}(\alpha, \tilde{\pi}, \varpi)\right|+\left|T^{-1} \nabla_{\tilde{\pi}_{i} \tilde{\pi}_{j} \tilde{\pi}_{k} \varpi_{l}}^{4} L_{T}^{\pi}(\alpha, \tilde{\pi}, \varpi)\right|\right. \\
& \left.+\left|T^{-1} \nabla_{\tilde{\tilde{\pi}_{i}}}^{4} \tilde{\pi}_{j} \varpi_{k} \varpi_{l} L_{T}^{\pi}(\alpha, \tilde{\pi}, \varpi)\right|+\left|T^{-1} \nabla_{\tilde{\pi}_{i} \varpi_{j} \varpi_{k} \varpi_{l}}^{4} L_{T}^{\pi}(\alpha, \tilde{\pi}, \varpi)\right|\right] \\
& +C \sum_{i, j, k, l} \sup _{\alpha \in A} \sup _{\substack{(\tilde{\pi}, \varpi) \in B \times \Pi_{\alpha},\left\|(\tilde{\pi}, \varpi)-\left(\tilde{\pi}^{*}, 0\right)\right\| \leq \gamma_{T}}}\left|T^{-1}\left(\nabla_{\varpi_{i} \varpi_{j} \varpi_{k} \varpi_{l}}^{4} L_{T}^{\pi}(\alpha, \tilde{\pi}, \varpi)-\nabla_{\varpi_{i} \varpi_{j} \varpi_{k} \varpi_{l}}^{4} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)\right)\right|, \tag{54}
\end{align*}
$$

where the summations above are understood to contain counterparts of each term in (48). As the data is assumed to be generated by a linear autoregression (Assumption 1), the $y_{t}$ 's form a stationary and ergodic process. Moreover, as the reparameterized log-likelihood of the GMAR model is four times continuously differentiable (see Assumption 4 and its verification), also the $\nabla_{\tilde{\pi}_{i} \tilde{\pi}_{j} \tilde{\pi}_{k}}^{3} l_{t}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)$ 's form a stationary and ergodic process (for any $i, j, k$ ). An analogous result holds for all the third and fourth partial derivatives of $l_{t}^{\pi}(\alpha, \tilde{\pi}, \varpi)$ appearing on the majorant side of (54).

Now, Lemma F.2(iii) together with the ergodic theorem implies that $\sup _{\alpha \in A}\left[\left|T^{-1} \nabla_{\tilde{\pi}_{i} \tilde{\pi}_{j} \tilde{\pi}_{k}}^{3} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)\right|\right]=$ $O_{p}(1)$ (for any $\left.i, j, k\right)$. Similarly, $\sup _{\alpha \in A}\left[\left|T^{-1} \nabla_{\tilde{\pi}_{i} \tilde{\pi}_{j} \varpi_{k}}^{3} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)\right|\right]=O_{p}(1)$ (for any $\left.i, j, k\right)$. Expression of $\nabla_{\varpi_{i} \varpi_{j} \varpi_{k}}^{3} l_{t}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)$ in Supplementary Appendix F.1, Lemmas F. 1 and F.3, and the compactness of A, imply that $\sup _{\alpha \in A}\left|T^{-1 / 2} \nabla_{\varpi_{i} \varpi_{j} \varpi_{k}}^{3} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)\right| \leq C\left|T^{-1 / 2} \sum_{t=1}^{T} M D S_{t, i, j, k}\left(\pi^{*}\right)\right|$ for some finite $C$ and for some square integrable martingale difference sequence $M D S_{t, i, j, k}\left(\pi^{*}\right)$. Moreover, for any $i, j, k$, the last upper bound is $O_{p}(1)$ by an appropriate central limit theorem (Billingsley (1961)).

As for the fourth-order partial derivatives appearing on the majorant side of (54), Lemma F.2(iv) and a uniform law of large numbers for stationary and ergodic processes (Ranga Rao (1962)) imply that $\sup _{\alpha \in A} \sup _{(\tilde{\pi}, \varpi) \in B \times \Pi_{\alpha},\left\|(\tilde{\pi}, \varpi)-\left(\tilde{\pi}^{*}, 0\right)\right\| \leq \gamma_{T}}\left|T^{-1} \nabla_{\tilde{\pi}_{i} \tilde{\pi}_{j} \tilde{\pi}_{k} \tilde{\pi}_{l}} L_{T}^{\pi}(\alpha, \tilde{\pi}, \varpi)\right|=O_{p}(1)$ (for any $\left.i, j, k, l\right)$. The next three terms in (54) can be handled similarly. As for the last term on the majorant side of (54),

$$
\begin{aligned}
& \sup _{\alpha \in A} \sup _{\substack{(\tilde{\pi}, \varpi) \in B \times \Pi_{\alpha},\left\|(\tilde{\pi}, \varpi)-\left(\tilde{\pi}^{*}, 0\right)\right\| \leq \gamma_{T}}}\left|T^{-1}\left(\nabla_{\varpi_{i} \varpi_{j} \varpi_{k} \varpi_{l}}^{4} L_{T}^{\pi}(\alpha, \tilde{\pi}, \varpi)-\nabla_{\varpi_{i} \varpi_{j} \varpi_{k} \varpi_{l}}^{4} L_{T}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)\right)\right| \\
& \quad \leq 2 \sup _{\alpha \in A} \sup _{\substack{(\tilde{\pi}, \varpi) \in B \times \Pi_{\alpha},}} \mid T^{-1} \nabla_{\varpi_{i} \varpi_{j} \varpi_{k} \varpi_{l}}^{4} L_{T}^{\pi}(\alpha, \tilde{\pi}, \varpi)-\left(\tilde{\pi}^{*}, 0\right) \| \leq \gamma_{T}
\end{aligned}
$$

where the dominant side is $o_{p}(1)$ (again relying on Lemma F.2(iv) and a uniform LLN). To conclude, the upper bound in (54) is $\tilde{\gamma}_{T} O_{p}(1)+C o_{p}(1)=o_{p}(1)$. This completes the verification of Assumption $5(\mathrm{iv})$ for the term $R_{T}^{(1)}(\alpha, \tilde{\pi}, \varpi)$.

Now consider $R_{T}^{(2)}(\alpha, \tilde{\pi}, \varpi)=-\frac{1}{2}\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\right]^{\prime}\left(\mathcal{J}_{T}-\mathcal{I}\right)\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\right]$. We will below show that (a) $\mathcal{J}_{T} \xrightarrow{p} \mathcal{J}$ as $T \rightarrow \infty$, where the matrix $\mathcal{J}$ will be specified below (and $\mathcal{J}_{T}, \mathcal{J}$ do not depend on $\alpha$ ). Write $(-2$ times $) R_{T}^{(2)}(\alpha, \tilde{\pi}, \varpi)$ as

$$
\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\right]^{\prime}\left(\mathcal{J}_{T}-\mathcal{J}\right)\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\right]+\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\right]^{\prime}(\mathcal{J}-\mathcal{I})\left[T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\right]
$$

We will below also show that (b) the latter term above equals zero. The validity of Assumption 5(iv) for the term $R_{T}^{(2)}(\alpha, \tilde{\pi}, \varpi)$ follows from results (a) and (b) (together with usual properties of the Euclidean
norm).
To prove claim (a), we first define the matrix $\mathcal{J}$ as

$$
\mathcal{J}=\left[\begin{array}{cc}
\mathcal{J}_{\pi \pi} & \mathcal{J}_{\pi \varpi \varpi}^{\prime} \\
\mathcal{J}_{\pi \omega \varpi} & \mathcal{J}_{\omega w \omega \varpi}
\end{array}\right]
$$

where the matrices $\mathcal{J}_{\tilde{\pi} \tilde{\pi}}\left(\left(q_{1}+q_{2}\right) \times\left(q_{1}+q_{2}\right)\right), \mathcal{J}_{\tilde{\pi} \varpi \varpi}^{\prime}\left(\left(q_{1}+q_{2}\right) \times q_{\vartheta}\right)$, and $\mathcal{J}_{\omega \varpi \omega \varpi}\left(q_{\vartheta} \times q_{\vartheta}\right)$ are defined as

$$
\begin{aligned}
\mathcal{J}_{\tilde{\pi} \tilde{\pi}} & =E\left[\frac{\nabla f_{t}^{*}}{f_{t}^{*}} \frac{\nabla^{\prime} f_{t}^{*}}{f_{t}^{*}}\right] \\
\mathcal{J}_{\tilde{\pi} \varpi \varpi}^{\prime} & =E\left[\left[c_{i j} \frac{\nabla f_{t}^{*}}{f_{t}^{*}} X_{t, i, j}^{*}\right]_{(i, j) \in \mathcal{I}}\right] \\
\mathcal{J}_{\varpi \varpi \varpi \varpi} & =\frac{1}{3}\left[c_{i j} c_{k l}\left(E\left[X_{t, i, j}^{*} X_{t, k, l}^{*}\right]+E\left[X_{t, i, k}^{*} X_{t, j, l}^{*}\right]+E\left[X_{t, i, l}^{*} X_{t, j, k}^{*}\right)\right]_{(i, j, k, l) \in \mathfrak{J} \times \mathcal{I}}\right.
\end{aligned}
$$

where the $X_{t, i, j}^{*}\left(i, j \in\left\{1, \ldots, q_{2}\right\}\right)$ are as in (44). Finiteness of $\mathcal{J}$ follows from Lemma F.1.
Now consider the convergence result $\mathcal{J}_{T} \xrightarrow{p} \mathcal{J}$ for each block at a time. For the top-left block, from Supplementary Appendix F. 1 we have $\nabla_{\tilde{\pi} \tilde{\pi}} l_{t}^{\pi *}=\frac{\nabla^{2} f_{t}^{*}}{f_{t}^{*}}-\frac{\nabla f_{t}^{*}}{f_{t}^{*}} \frac{{ }^{\prime} f_{t}^{*}}{f_{t}^{*}}$ so that ergodic theorem and Lemmas F. 1 and F. 3 (latter ensuring the first term on the right-hand side of the previous equation has zero expectation) imply that $\mathcal{J}_{T, \tilde{\pi} \tilde{\pi}}=-T^{-1} \nabla_{\tilde{\tilde{\pi}} \tilde{\pi}^{2}}^{2} L_{T}^{\pi *} \xrightarrow{p} \mathcal{J}_{\tilde{\pi} \tilde{\pi}}$.

For the off-diagonal block, consider the expression of $\nabla_{\tilde{\pi}_{i} \varpi_{j} \omega_{k}}^{3} l_{t}^{\pi *}$ in Supplementary Appendix F.1. Lemma F. 3 ensures that of the ten summands in this expression, only the second, fourth, and sixth ones have non-zero expectation. Therefore the ergodic theorem and Lemma F. 1 imply that

$$
\mathcal{J}_{T, \tilde{\pi} \varpi \varpi}^{\prime}=-T^{-1} \frac{1}{\alpha_{1} \alpha_{2}}\left[c_{11} \nabla_{\tilde{\pi} \omega_{1} \varpi_{1}}^{3} L_{T}^{\pi *}: \cdots: c_{q_{2}-1, q_{2}} \nabla_{\tilde{\pi} \omega_{q_{2}-1} \varpi_{q_{2}}}^{3} L_{T}^{\pi *}\right] \xrightarrow{p} \mathcal{J}_{\tilde{\pi} \omega \varpi}^{\prime} .
$$

Lastly, for the bottom-right block, consider the expression of $\nabla_{w_{i} \varpi_{j} \omega_{k} \omega_{l}}^{4} l t{ }_{t}^{\pi *}$ in Supplementary Appendix F.1. Lemma F. 3 reveals that the terms in this expression that have non-zero expectation can be expressed as

$$
-\alpha_{1}^{2} \alpha_{2}^{2}\left[X_{t, i, j}^{*} X_{t, k, l}^{*}+X_{t, i, k}^{*} X_{t, j, l}^{*}+X_{t, i, l}^{*} X_{t, j, k}^{*}\right] .
$$

Therefore the ergodic theorem and Lemma F. 1 imply that

$$
\mathcal{J}_{T, \varpi \varpi \varpi \varpi}=-T^{-1} \frac{8}{4!} \frac{1}{\alpha_{1}^{2} \alpha_{2}^{2}}\left[c_{i j} c_{k l} \nabla_{\varpi_{i} \varpi_{j} \varpi_{k} \varpi_{l}}^{4} L_{T}^{\pi_{*}^{*}}\right]_{(i, j, k, l) \in \mathfrak{J} \times \mathfrak{J}} \xrightarrow{p} \mathcal{J}_{\varpi \varpi \varpi \varpi} .
$$

This completes the proof of claim (a).
To prove claim (b), first note that from the definitions of $\mathcal{J}$ and $\mathcal{I}$ (see (51)) it can be seen that only the bottom-right blocks of $\mathcal{J}$ and $\mathcal{I}$ differ. Therefore, as $T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)=\left(T^{1 / 2}(\tilde{\pi}-\right.$ $\left.\tilde{\pi}^{*}\right), T^{1 / 2}\left(\alpha_{1} \alpha_{2} v(\varpi)\right)$ ), claim (b) holds if $T\left(\alpha_{1} \alpha_{2}\right)^{2} v(\varpi)^{\prime}\left(\mathcal{J}_{\varpi \varpi \varpi \varpi}-\mathcal{I}_{\varpi \varpi \varpi \varpi}\right) v(\varpi)=0$ where

$$
\begin{aligned}
\mathcal{J}_{\varpi \varpi \varpi \varpi} & =\frac{1}{3}\left[c_{i j} c_{k l}\left(E\left[X_{t, i, j}^{*} X_{t, k, l}^{*}\right]+E\left[X_{t, i, k}^{*} X_{t, j, l}^{*}\right]+E\left[X_{t, i, l}^{*} X_{t, j, k}^{*}\right]\right)\right]_{(i, j, k, l) \in \mathfrak{I} \times \mathfrak{I}}, \\
\mathcal{I}_{\varpi \varpi \varpi \varpi} & =\left[c_{i j} c_{k l} E\left[X_{t, i, j}^{*} X_{t, k, l}^{*},\right]_{(i, j, k, l) \in \mathfrak{I} \times \mathfrak{I}} .\right.
\end{aligned}
$$

Note that the scalars $A_{i j k l}=E\left[X_{t, i, j}^{*} X_{t, k, l}^{*}\right]$ satisfy $A_{i j k l}=A_{j i k l}$ and $A_{i j k l}=A_{i j l k}$ for all $i, j, k, l$ so
that using property (49) we obtain

$$
\begin{aligned}
v(\varpi)^{\prime} \mathcal{J}_{\varpi \varpi \varpi \varpi} v(\varpi) & =\frac{1}{3} v(\varpi)^{\prime}\left[c_{i j} c_{k l}\left(E\left[X_{t, i, j}^{*} X_{t, k, l}^{*}\right]+E\left[X_{t, i, k}^{*} X_{t, j, l}^{*}\right]+E\left[X_{t, i, l}^{*} X_{t, j, k}^{*}\right]\right)\right]_{(i, j, k, l) \in \mathfrak{J} \times \mathfrak{I}} v(\varpi) \\
& =\frac{1}{3} \frac{1}{4} \sum_{i=1}^{q_{2}} \sum_{j=1}^{q_{2}} \sum_{k=1}^{q_{2}} \sum_{l=1}^{q_{2}}\left(E\left[X_{t, i, j}^{*} X_{t, k, l}^{*}\right]+E\left[X_{t, i, k}^{*} X_{t, j, l}^{*}\right]+E\left[X_{t, i, l}^{*} X_{t, j, k}^{*}\right]\right) \varpi_{i} \varpi_{j} \varpi_{k} \varpi_{l} \\
& =\frac{1}{4} \sum_{i=1}^{q_{2}} \sum_{j=1}^{q_{2}} \sum_{k=1}^{q_{2}} \sum_{l=1}^{q_{2}} E\left[X_{t, i, j}^{*} X_{t, k, l}^{*}\right] \varpi_{i} \varpi_{j} \varpi_{k} \varpi_{l} \\
& =v(\varpi)^{\prime} \mathcal{I}_{\varpi \varpi \varpi \varpi} v(\varpi) .
\end{aligned}
$$

This completes the proof of claim (b).
Therefore, the verification of Assumption 5(iv) is done.

## F. 5 Additional Lemmas

The following four lemmas contain results needed in the proofs. Note that the first and the third lemma are not specific to the examples in this paper, whereas the second and fourth lemmas concern only the GMAR example. In the first lemma, $\mathrm{n}_{p+1}(\tilde{\phi})=\mathrm{n}_{p+1}\left(y_{t}, \boldsymbol{y}_{t-1} ; \tilde{\phi}\right)$ denotes the $(p+1)$-dimensional density function of an $\operatorname{AR}(p)$ process based on parameter value $\tilde{\phi}$ evaluated at ( $y_{t}, \boldsymbol{y}_{t-1}$ ); cf. equations (6)-(8) for the $p$-dimensional counterpart $\mathrm{n}_{p}(\tilde{\phi})=\mathrm{n}_{p}\left(\boldsymbol{y}_{t-1} ; \tilde{\phi}\right)$.

Lemma F.1. For any $i, j, k, l \in\{1, \ldots, p+2\}$ and any positive $r$, the following moments are all finite:
(i) $E\left[\sup _{\tilde{\phi} \in \tilde{\Phi} \mid}\left|\nabla_{i} f_{t}(\tilde{\phi}) / f_{t}(\tilde{\phi})\right|^{r}\right], E\left[\sup _{\tilde{\phi} \in \tilde{\Phi} \mid}\left|\nabla_{i j}^{2} f_{t}(\tilde{\phi}) / f_{t}(\tilde{\phi})\right|^{r}\right], \ldots, E\left[\sup _{\tilde{\phi} \in \tilde{\Phi} \mid}\left|\nabla_{i j k l}^{4} f_{t}(\tilde{\phi}) / f_{t}(\tilde{\phi})\right|^{r}\right]$,
(ii) $E\left[\sup _{\tilde{\phi} \in \tilde{\Phi} \mid}\left|\nabla_{i} \mathrm{n}_{p}(\tilde{\phi}) / \mathrm{n}_{p}(\tilde{\phi})\right|^{r}\right], E\left[\sup _{\tilde{\phi} \in \tilde{\Phi}}\left|\nabla_{i j}^{2} \mathrm{n}_{p}(\tilde{\phi}) / \mathrm{n}_{p}(\tilde{\phi})\right|^{r}\right], \ldots, E\left[\sup _{\tilde{\phi} \in \tilde{\Phi}}\left|\nabla_{i j k l}^{4} \mathrm{n}_{p}(\tilde{\phi}) / \mathrm{n}_{p}(\tilde{\phi})\right|^{r}\right]$,
(iii) $E\left[\sup _{\tilde{\phi} \in \tilde{\Phi}}\left|\nabla_{i} \mathrm{n}_{p+1}(\tilde{\phi})\right|^{r}\right], E\left[\sup _{\tilde{\phi} \in \tilde{\Phi}}\left|\nabla_{i j}^{2} \mathrm{n}_{p+1}(\tilde{\phi})\right|^{r}\right], \ldots, E\left[\sup _{\tilde{\phi} \in \tilde{\Phi}}\left|\nabla_{i j k l}^{4} \mathrm{n}_{p+1}(\tilde{\phi})\right|^{r}\right]$.

Lemma F.2. In the GMAR example the following hold, where each of (the scalars) $z_{1}, z_{2}, z_{3}, z_{4}$ is a 'placeholder' for any of $\tilde{\pi}_{i}, \tilde{\pi}_{j}, \tilde{\pi}_{k}, \tilde{\pi}_{l}\left(i, j, k, l \in\left\{1, \ldots, q_{1}+q_{2}\right\}\right)$ or $\varpi_{i}, \varpi_{j}, \varpi_{k}, \varpi_{l}(i, j, k, l \in$ $\left\{1, \ldots, q_{2}\right\}$ ):

> (i) $E\left[\sup _{\alpha \in A} \sup _{(\tilde{\pi}, \varpi) \in B \times \Pi_{\alpha}}\left|\nabla_{z_{1}} l_{t}^{\pi}(\alpha, \tilde{\pi}, \varpi)\right|\right]<\infty$,
> (ii) $E\left[\sup _{\alpha \in A} \sup _{(\tilde{\pi}, \varpi) \in B \times \Pi_{\alpha}}\left|\nabla_{z_{1} z_{2}}^{2} l_{t}^{\pi}(\alpha, \tilde{\pi}, \varpi)\right|\right]<\infty$,
> (iii) $E\left[\sup _{\alpha \in A} \sup _{(\tilde{\pi}, \varpi) \in B \times \Pi_{\alpha}}\left|\nabla_{z_{1} z_{2} z_{3}}^{3} l_{t}^{\pi}(\alpha, \tilde{\pi}, \varpi)\right|\right]<\infty$,
> (iv) $E\left[\sup _{\alpha \in A} \sup _{(\tilde{\pi}, \varpi) \in B \times \Pi_{\alpha}}\left|\nabla_{z_{1} z_{2} z_{3} z_{4}}^{4} l_{t}^{\pi}(\alpha, \tilde{\pi}, \varpi)\right|\right]<\infty$.

Lemma F.3. For any $i, j, k, l \in\{1, \ldots, p+2\}$,

$$
E\left[\left.\frac{\nabla_{i} f_{t}^{*}}{f_{t}^{*}} \right\rvert\, \boldsymbol{y}_{t-1}\right]=E\left[\left.\frac{\nabla_{i j}^{2} f_{t}^{*}}{f_{t}^{*}} \right\rvert\, \boldsymbol{y}_{t-1}\right]=E\left[\left.\frac{\nabla_{i, j, k}^{3} f_{t}^{*}}{f_{t}^{*}} \right\rvert\, \boldsymbol{y}_{t-1}\right]=E\left[\left.\frac{\nabla_{i, j k}^{4} f_{t}^{*}}{f_{t}^{*}} \right\rvert\, \boldsymbol{y}_{t-1}\right]=0 .
$$

Lemma F.4. In the GMAR example the following hold for all $\alpha \in A,(\tilde{\pi}, \varpi) \in B \times \Pi_{\alpha}$, $T$, and $i, j, k, l \in\left\{1, \ldots, q_{1}+q_{2}\right\}$ (subindex in $\tilde{\pi}$ ) or $i, j, k, l \in\left\{1, \ldots, q_{2}\right\}$ (subindex in $\varpi$ ):
(i) $T\left|\tilde{\pi}_{i}-\tilde{\pi}_{i}^{*}\right|\left|\tilde{\pi}_{j}-\tilde{\pi}_{j}^{*}\right|\left|\tilde{\pi}_{k}-\tilde{\pi}_{k}^{*}\right| \leq\left(1+\left\|T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\right\|\right)^{2}\left\|\tilde{\pi}-\tilde{\pi}^{*}\right\|$,
(ii) $T\left|\tilde{\pi}_{i}-\tilde{\pi}_{i}^{*}\right|\left|\tilde{\pi}_{j}-\tilde{\pi}_{j}^{*}\right|\left|\varpi_{k}\right| \leq\left(1+\left\|T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\right\|\right)^{2}\|\varpi\|$,

$$
\text { (iii) } \quad T^{1 / 2}\left|\varpi_{i}\right|\left|\varpi_{j}\right|\left|\varpi_{k}\right| \leq\left(1+\left\|T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\right\|\right)^{2}\left(\alpha_{1} \alpha_{2}\right)^{-3 / 2}\|\boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\|^{1 / 2},
$$

(iv) $T\left(\tilde{\pi}_{i}-\tilde{\pi}_{i}^{*}\right)\left(\tilde{\pi}_{j}-\tilde{\pi}_{j}^{*}\right)\left(\tilde{\pi}_{k}-\tilde{\pi}_{k}^{*}\right)\left(\tilde{\pi}_{l}-\tilde{\pi}_{l}^{*}\right) \leq\left(1+\left\|T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\right\|\right)^{2}\left\|\tilde{\pi}-\tilde{\pi}^{*}\right\|^{2}$,
(v) $T\left(\tilde{\pi}_{i}-\tilde{\pi}_{i}^{*}\right)\left(\tilde{\pi}_{j}-\tilde{\pi}_{j}^{*}\right)\left(\tilde{\pi}_{k}-\tilde{\pi}_{k}^{*}\right) \varpi_{l} \leq\left(1+\left\|T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\right\|\right)^{2}\left\|\tilde{\pi}-\tilde{\pi}^{*}\right\|\|\varpi\|$,
(vi) $T\left(\tilde{\pi}_{i}-\tilde{\pi}_{i}^{*}\right)\left(\tilde{\pi}_{j}-\tilde{\pi}_{j}^{*}\right) \varpi_{k} \varpi_{l} \leq\left(1+\left\|T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\right\|\right)^{2}\|\varpi\|^{2}$,
(vii) $T\left(\tilde{\pi}_{i}-\tilde{\pi}_{i}^{*}\right) \varpi_{j} \varpi_{k} \varpi_{l} \leq\left(1+\left\|T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\right\|\right)^{2}\left(\alpha_{1} \alpha_{2}\right)^{-3 / 2}\|\boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\|^{1 / 2}$,
(viii) $T \varpi_{i} \varpi_{j} \varpi_{k} \varpi_{l} \leq\left(1+\left\|T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\right\|\right)^{2}\left(\alpha_{1} \alpha_{2}\right)^{-2}$.

## F. 6 Proofs of Lemmas F.1-F. 4

Proof of Lemma F.1. Writing $g_{t}(\tilde{\phi})=\left[y_{t}-\left(\tilde{\phi}_{0}+\tilde{\phi}_{1} y_{t-1}+\cdots+\tilde{\phi}_{p} y_{t-p}\right)\right] / \tilde{\sigma}_{1}$ and recalling the definition of $f_{t}(\tilde{\phi})$ we can write $f_{t}(\tilde{\phi})=\tilde{\sigma}_{1}^{-1} \mathfrak{N}\left(g_{t}(\tilde{\phi})\right)$ where $\mathfrak{N}(\cdot)$ denotes the density function of a standard normal random variable. Recall also that derivatives of $\mathfrak{N}(\cdot)$ can be expressed using (one version of) Hermite polynomials $H_{n}(x)$ as $\frac{d^{n}}{d x^{n}} \mathfrak{N}(x)=(-1)^{n} H_{n}(x) \mathfrak{N}(x)$. Using the chain rule for differentiation repeatedly, it can therefore be seen that each of the functions $\nabla_{i} f_{t}(\tilde{\phi}) / f_{t}(\tilde{\phi}), \nabla_{i j}^{2} f_{t}(\tilde{\phi}) / f_{t}(\tilde{\phi})$, $\nabla_{i j k}^{3} f_{t}(\tilde{\phi}) / f_{t}(\tilde{\phi})$, and $\nabla_{i j k l}^{4} f_{t}(\tilde{\phi}) / f_{t}(\tilde{\phi})$ can be expressed as a sum of terms each of which is a product involving Hermite polynomials $H_{n}\left(g_{t}(\tilde{\phi})\right)$ and powers of derivatives of $g_{t}(\tilde{\phi})$ (and functions of $\tilde{\phi}$ ). Thus, each of these functions is a polynomial in terms of $y_{t}, y_{t-1}, \ldots, y_{t-p}$. As the $y_{t}$ 's are generated by a stationary linear Gaussian $\operatorname{AR}(p)$ model, they possess moments of all orders, implying (together with the definition of $\tilde{\Phi}$, implying in particular that $\tilde{\sigma}_{1}$ is bounded away from zero on $\tilde{\Phi}$ ) the finiteness of the moments listed in part (i) of the lemma.

As for part (ii), note that $\mathbf{n}_{p}(\tilde{\phi})$ can be expressed as $\mathbf{n}_{p}(\tilde{\phi})=g_{1}(\tilde{\phi}) \mathfrak{N}\left(g_{2, t}(\tilde{\phi})\right)$ for some function $g_{1}(\tilde{\phi})$ not depending on the $y_{t}$ 's and $g_{2, t}(\tilde{\phi})$ the square root of a second-order polynomial in $y_{t-1}, \ldots, y_{t-p}$. Therefore the finiteness of the moments listed in the part (ii) follows using similar arguments as above (noting that the definition of $\tilde{\Phi}$ implies that the determinant of the covariance matrix appearing in $\mathrm{n}_{p}(\tilde{\phi})$ is bounded away from zero on $\left.\tilde{\Phi}\right)$.

Finally, for part (iii), similar arguments, together with the observation that $\mathrm{n}_{p+1}(\tilde{\phi})$ is bounded on $\tilde{\Phi}$, yield the desired result.

Proof of Lemma F.2. To prove (i), first consider the derivatives with respect to $\varpi$. From the formulas in Supplementary Appendix F. 7 we obtain

$$
\begin{aligned}
\nabla_{\varpi} l_{t}^{\pi}(\alpha, \tilde{\pi}, \varpi) & =\alpha_{1, t} \alpha_{2, t}\left(\frac{D_{\tilde{\phi}, \varpi}^{(1) \prime} \nabla \mathrm{n}_{p}(\tilde{\phi})}{\mathrm{n}_{p}(\tilde{\phi})}-\frac{D_{\tilde{\varphi}, \varpi}^{(1) \prime} \nabla \mathrm{n}_{p}(\tilde{\varphi})}{\mathrm{n}_{p}(\tilde{\varphi})}\right) \frac{f_{t}(\tilde{\phi})}{f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)}+\frac{D_{\tilde{\phi}, \varpi}^{(1) \prime} \nabla f_{t}(\tilde{\phi})}{f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)} \alpha_{1, t} \\
& -\alpha_{1, t} \alpha_{2, t}\left(\frac{D_{\tilde{\phi}, \varpi}^{(1) \prime} \nabla \mathrm{n}_{p}(\tilde{\phi})}{\mathrm{n}_{p}(\tilde{\phi})}-\frac{D_{\tilde{\varphi}, \varpi}^{(1) \prime} \nabla \mathrm{n}_{p}(\tilde{\varphi})}{\mathrm{n}_{p}(\tilde{\varphi})}\right) \frac{f_{t}(\tilde{\varphi})}{f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)}+\frac{D_{\tilde{\varphi}, \varpi}^{(1) \prime} \nabla f_{t}(\tilde{\varphi})}{f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)}\left(1-\alpha_{1, t}\right)
\end{aligned}
$$

where $\tilde{\phi}$ and $\tilde{\varphi}$ are understood as functions of $(\alpha, \tilde{\pi}, \varpi)$ (i.e., $\tilde{\phi}=\left(\beta, \pi+\alpha_{2} \varpi\right)$ and $\tilde{\varphi}=\left(\beta, \pi-\alpha_{1} \varpi\right)$ ). Note that whenever $\alpha \in A$ and $(\tilde{\pi}, \varpi) \in B \times \Pi_{\alpha}, \tilde{\phi} \in \tilde{\Phi}$ and $\tilde{\varphi} \in \tilde{\Phi}$. Also note that over $\alpha \in A$ and
$(\tilde{\pi}, \varpi) \in B \times \Pi_{\alpha}$, the quantities

$$
\left|\alpha_{1, t}\right|,\left|\alpha_{2, t}\right|, \quad\left\|D_{\tilde{\phi}, \varpi}^{(1)}\right\|,\left\|D_{\tilde{\varphi}, \varpi}^{(1)}\right\|,\left|f_{t}(\tilde{\phi}) / f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)\right|,\left|f_{t}(\tilde{\varphi}) / f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)\right|
$$

are all bounded by finite constants. Therefore $E\left[\sup _{\alpha \in A} \sup _{(\tilde{\pi}, \varpi) \in N_{\tilde{\sim}}\left(\tilde{\pi}^{*}, 0\right)}\left\|\nabla_{\varpi} l_{t}^{\pi}(\alpha, \tilde{\pi}, \varpi)\right\|\right]<\infty$ as long as $E\left[\sup _{\tilde{\phi} \in \tilde{\Phi}}\left\|\nabla f_{t}(\tilde{\phi}) / f_{t}(\tilde{\phi})\right\|\right]<\infty$ and $E\left[\sup _{\tilde{\phi} \in \tilde{\Phi}}\left\|\nabla \mathrm{n}_{p}(\tilde{\phi}) / \mathrm{n}_{p}(\tilde{\phi})\right\|\right]<\infty$, which is ensured by Lemma F.1. The argument for $\nabla_{\tilde{\pi}} l_{t}^{\pi}(\alpha, \tilde{\pi}, \varpi)$ is entirely similar and is omitted.

To prove (ii)-(iv), entirely similar arguments can be used. Tedious calculations (details omitted) show that the finiteness of the required moments is ensured by the finiteness of the moments in Lemma F.1(i) and (ii).

Proof of Lemma F.3. For the first two derivatives, the stated result follows directly from the expressions of $\nabla f_{t}^{*} / f_{t}^{*}$ and $\nabla^{2} f_{t}^{*} / f_{t}^{*}$ in (52). The results for the third and fourth derivatives can be obtained with straightforward calculation.

Proof of Lemma F.4. First recall that $\left\|T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\right\|^{2}=T\left\|\tilde{\pi}-\tilde{\pi}^{*}\right\|^{2}+T \alpha_{1}^{2} \alpha_{2}^{2}\|v(\varpi)\|^{2}$. (i) By an elementary inequality, $T\left|\tilde{\pi}_{i}-\tilde{\pi}_{i}^{*}\right|\left|\tilde{\pi}_{j}-\tilde{\pi}_{j}^{*}\right|\left|\tilde{\pi}_{k}-\tilde{\pi}_{k}^{*}\right| \leq T\left\|\tilde{\pi}-\tilde{\pi}^{*}\right\|^{3}$ and therefore the result follows by adding nonnegative terms on the majorant side of this inequality. Parts (ii) and (iv)-(vi) are shown similarly. (vii) As $\boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)=\left(\tilde{\pi}-\tilde{\pi}^{*}, \alpha_{1} \alpha_{2} v(\varpi)\right)$, each of the terms $\left|\tilde{\pi}_{i}-\tilde{\pi}_{i}^{*}\right|, \alpha_{1} \alpha_{2} \varpi_{j}^{2}, \alpha_{1} \alpha_{2} \varpi_{k}^{2}$, and $\alpha_{1} \alpha_{2} \varpi_{l}^{2}$ are dominated by $\|\boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\|$. Therefore
$T\left|\tilde{\pi}_{i}-\tilde{\pi}_{i}^{*}\right|\left|\varpi_{j}\right|\left|\varpi_{k}\right|\left|\varpi_{l}\right| \leq T\left(\alpha_{1} \alpha_{2}\right)^{-3 / 2}\|\boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\|^{5 / 2} \leq\left(\alpha_{1} \alpha_{2}\right)^{-3 / 2}\left(1+\left\|T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\right\|\right)^{2}\|\boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\|^{1 / 2}$
where the second inequality holds because nonnegative terms were added to the majorant side. (viii) Similarly as in the previous part,

$$
T\left|\varpi_{i}\right|\left|\varpi_{j}\left\|\varpi_{k}| | \varpi_{l} \mid \leq\left(\alpha_{1} \alpha_{2}\right)^{-2}\right\| T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi) \|^{2} \leq\left(\alpha_{1} \alpha_{2}\right)^{-2}\left(1+\left\|T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\right\|\right)^{2}\right.
$$

Finally, for (iii) we, similarly as above but scaling with $T^{1 / 2}$ instead of $T$, obtain
$T^{1 / 2}\left|\varpi_{i}\right|\left|\varpi_{j}\left\|\varpi_{k} \mid \leq T^{1 / 2}\left(\alpha_{1} \alpha_{2}\right)^{-3 / 2}\right\| \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi)\left\|^{3 / 2} \leq\left(\alpha_{1} \alpha_{2}\right)^{-3 / 2}\left(1+\left\|T^{1 / 2} \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi \alpha)\right\|\right)^{2}\right\| \boldsymbol{\theta}(\alpha, \tilde{\pi}, \varpi) \|^{1 / 2}\right.$,
which completes the proof.

## F. 7 Partial derivatives of the reparameterized log-likelihood function (continued)

Note that $l_{t}^{\pi}(\alpha, \beta, \pi, \varpi)=\log \left[f_{2, t}^{\pi}(\alpha, \beta, \pi, \varpi)\right]$ with

$$
\begin{aligned}
f_{2, t}^{\pi}(\alpha, \beta, \pi, \varpi) & =\alpha_{1, t}^{\pi}(\alpha, \beta, \pi, \varpi) f_{t}\left(\beta, \pi+\alpha_{2} \varpi\right)+\left(1-\alpha_{1, t}^{\pi}(\alpha, \beta, \pi, \varpi)\right) f_{t}\left(\beta, \pi-\alpha_{1} \varpi\right), \\
\alpha_{1, t}^{\pi}(\alpha, \beta, \pi, \varpi) & =\alpha_{1, t}^{G}\left(\alpha,\left(\beta, \pi+\alpha_{2} \varpi\right),\left(\beta, \pi-\alpha_{1} \varpi\right)\right) .
\end{aligned}
$$

For the sake of brevity, but with slight abuse of notation, we will write these as

$$
\begin{aligned}
f_{2, t}^{\pi}(\alpha, \beta, \pi, \varpi) & =\alpha_{1, t}^{\pi}(\alpha, \beta, \pi, \varpi) f_{t}(\tilde{\phi})+\left(1-\alpha_{1, t}^{\pi}(\alpha, \beta, \pi, \varpi)\right) f_{t}(\tilde{\varphi}) \\
\alpha_{1, t}^{\pi}(\alpha, \beta, \pi, \varpi) & =\frac{\alpha \mathrm{n}_{p}(\tilde{\phi})}{\alpha \mathrm{n}_{p}(\tilde{\phi})+(1-\alpha) \mathrm{n}_{p}(\tilde{\varphi})}
\end{aligned}
$$

where $\tilde{\phi}$ and $\tilde{\varphi}$ are understood as functions of $(\alpha, \beta, \pi, \varpi)$, that is, $\tilde{\phi}=\left(\beta, \pi+\alpha_{2} \varpi\right)$ and $\tilde{\varphi}=(\beta, \pi-$ $\left.\alpha_{1} \varpi\right)$.

The following notation will be helpful:

$$
\begin{aligned}
& D_{\tilde{\phi}, \tilde{\pi}}^{(1)}=\frac{\partial\left(\beta, \pi+\alpha_{2} \varpi\right)}{\partial \tilde{\pi}^{\prime}}=I_{1+q_{2}} \\
& D_{\tilde{\phi}, \varpi}^{(1)}=\frac{\partial\left(\beta, \pi+\alpha_{2} \varpi\right)}{\partial \varpi^{\prime}}=\left[\begin{array}{c}
0 \\
\alpha_{2} I_{q_{2}}
\end{array}\right] \quad\left(\left(1+q_{2}\right) \times q_{2}\right) \\
& D_{\tilde{\varphi}, \tilde{\pi}}^{(1)}=\frac{\partial\left(\beta, \pi-\alpha_{1} \varpi\right)}{\partial \tilde{\pi}^{\prime}}=I_{1+q_{2}} \\
& D_{\tilde{\varphi}, \varpi}^{(1)}=\frac{\partial\left(\beta, \pi-\alpha_{1} \varpi\right)}{\partial \varpi^{\prime}}=\left[\begin{array}{c}
0 \\
-\alpha_{1} I_{q_{2}}
\end{array}\right] \quad\left(\left(1+q_{2}\right) \times q_{2}\right)
\end{aligned}
$$

First-order partial derivatives. With straightforward differentation we obtain

$$
\begin{aligned}
& \nabla_{\tilde{\pi}} l_{t}^{\pi}(\alpha, \tilde{\pi}, \varpi)=\frac{\nabla_{\tilde{\pi}} f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)}{f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)} \\
& \nabla_{\varpi} l_{t}^{\pi}(\alpha, \tilde{\pi}, \varpi)=\frac{\nabla_{\varpi} f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)}{f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)}
\end{aligned}
$$

with

$$
\begin{aligned}
\nabla_{\tilde{\pi}} f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi) & =\nabla_{\tilde{\pi}} \alpha_{1, t} f_{t}(\tilde{\phi})+\nabla f_{t}(\tilde{\phi}) \alpha_{1, t}-\nabla_{\tilde{\pi}} \alpha_{1, t} f_{t}(\tilde{\varphi})+\nabla f_{t}(\tilde{\varphi})\left(1-\alpha_{1, t}\right) \\
\nabla_{\varpi} f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi) & =\nabla_{\varpi} \alpha_{1, t} f_{t}(\tilde{\phi})+D_{\tilde{\phi}, \varpi}^{(1) \prime} \nabla f_{t}(\tilde{\phi}) \alpha_{1, t}-\nabla_{\varpi} \alpha_{1, t} f_{t}(\tilde{\varphi})+D_{\tilde{\varphi}, \varpi}^{(1) \prime} \nabla f_{t}(\tilde{\varphi})\left(1-\alpha_{1, t}\right)
\end{aligned}
$$

and

$$
\begin{aligned}
\nabla_{\tilde{\pi}} \alpha_{1, t} & =\frac{\alpha_{1} \nabla \mathrm{n}_{p}(\tilde{\phi})}{\alpha_{1} \mathrm{n}_{p}(\tilde{\phi})+\alpha_{2} \mathrm{n}_{p}(\tilde{\varphi})}-\alpha_{1, t} \frac{\alpha_{1} \nabla \mathrm{n}_{p}(\tilde{\phi})+\alpha_{2} \nabla \mathrm{n}_{p}(\tilde{\varphi})}{\alpha_{1} \mathrm{n}_{p}(\tilde{\phi})+\alpha_{2} \mathrm{n}_{p}(\tilde{\varphi})} \\
\nabla_{\varpi} \alpha_{1, t} & =\frac{\alpha_{1} D_{\tilde{\phi}, \varpi}^{(1) \prime} \nabla \mathrm{n}_{p}(\tilde{\phi})}{\alpha_{1} \mathrm{n}_{p}(\tilde{\phi})+\alpha_{2} \mathrm{n}_{p}(\tilde{\varphi})}-\alpha_{1, t} \frac{\alpha_{1} D_{\tilde{\phi}, \varpi}^{(1) \prime} \nabla \mathrm{n}_{p}(\tilde{\phi})+\alpha_{2} D_{\tilde{\varphi}, \varpi}^{(1) \prime} \nabla \mathrm{n}_{p}(\tilde{\varphi})}{\alpha_{1} \mathrm{n}_{p}(\tilde{\phi})+\alpha_{2} \mathrm{n}_{p}(\tilde{\varphi})}
\end{aligned}
$$

where simplification leads to

$$
\begin{aligned}
\nabla_{\tilde{\pi}} \alpha_{1, t} & =\alpha_{1, t} \frac{\nabla \mathrm{n}_{p}(\tilde{\phi})}{\mathrm{n}_{p}(\tilde{\phi})}-\alpha_{1, t}\left(\alpha_{1, t} \frac{\nabla \mathrm{n}_{p}(\tilde{\phi})}{\mathrm{n}_{p}(\tilde{\phi})}+\alpha_{2, t} \frac{\nabla \mathrm{n}_{p}(\tilde{\varphi})}{\mathrm{n}_{p}(\tilde{\varphi})}\right) \\
& =\alpha_{1, t} \alpha_{2, t}\left(\frac{\nabla \mathrm{n}_{p}(\tilde{\phi})}{\mathrm{n}_{p}(\tilde{\phi})}-\frac{\nabla \mathrm{n}_{p}(\tilde{\varphi})}{\mathrm{n}_{p}(\tilde{\varphi})}\right) \\
\nabla_{\varpi} \alpha_{1, t} & =\alpha_{1, t} \frac{D_{\tilde{\phi}, \varpi}^{(1) \prime} \nabla \mathrm{n}_{p}(\tilde{\phi})}{\mathrm{n}_{p}(\tilde{\phi})}-\alpha_{1, t}\left(\alpha_{1, t} \frac{D_{\tilde{\phi}, \varpi}^{(1) \prime} \nabla \mathrm{n}_{p}(\tilde{\phi})}{\mathrm{n}_{p}(\tilde{\phi})}+\alpha_{2, t} \frac{D_{\tilde{\varphi}, \varpi}^{(1) \prime} \nabla \mathrm{n}_{p}(\tilde{\varphi})}{\mathrm{n}_{p}(\tilde{\varphi})}\right) \\
& =\alpha_{1, t} \alpha_{2, t}\left(\frac{D_{\tilde{\phi}, \varpi}^{(1) \prime} \nabla \mathrm{n}_{p}(\tilde{\phi})}{\mathrm{n}_{p}(\tilde{\phi})}-\frac{D_{\tilde{\varphi}, \varpi}^{(1) \prime} \nabla \mathrm{n}_{p}(\tilde{\varphi})}{\mathrm{n}_{p}(\tilde{\varphi})}\right)
\end{aligned}
$$

Evaluated at $(\alpha, \tilde{\pi}, \varpi)=\left(\alpha, \tilde{\pi}^{*}, 0\right)$ we get

$$
\nabla_{\tilde{\pi}} \alpha_{1, t}^{*}=0, \quad \nabla_{\varpi} \alpha_{1, t}^{*}=\alpha_{1} \alpha_{2} \frac{\nabla_{(2, \ldots, p+2)} \mathrm{n}_{p}\left(\tilde{\pi}^{*}\right)}{\mathrm{n}_{p}\left(\tilde{\pi}^{*}\right)}
$$

$$
f_{2, t}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)=f_{t}\left(\tilde{\pi}^{*}\right), \quad \nabla_{\tilde{\pi}} f_{2, t}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)=\nabla f_{t}\left(\tilde{\pi}^{*}\right), \quad \nabla_{\varpi} f_{2, t}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)=0
$$

so that

$$
\nabla_{\tilde{\pi}} l_{t}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)=\frac{\nabla f_{t}\left(\tilde{\pi}^{*}\right)}{f_{t}\left(\tilde{\pi}^{*}\right)}, \quad \quad \nabla_{\varpi} l_{t}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)=0
$$

Second-order partial derivatives With straightforward differentation we obtain

$$
\begin{aligned}
\nabla_{\tilde{\pi} \tilde{\pi}}^{2} l_{t}^{\pi}(\alpha, \tilde{\pi}, \varpi) & =\frac{\nabla_{\tilde{\pi} \tilde{\pi}^{\prime}}^{2} f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)}{f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)}-\frac{\nabla_{\tilde{\pi}} f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)}{f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)} \frac{\nabla_{\tilde{\pi}^{\prime}} f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)}{f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)} \\
\nabla_{\tilde{\pi} \varpi^{\prime}}^{2} l_{t}^{\pi}(\alpha, \tilde{\pi}, \varpi) & =\frac{\nabla_{\tilde{\pi} \varpi^{\prime}}^{2} f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)}{f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)}-\frac{\nabla_{\tilde{\pi}} f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)}{f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)} \frac{\nabla_{\varpi^{\prime}} f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)}{f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)} \\
\nabla_{\varpi \varpi^{\prime}}^{2} l_{t}^{\pi}(\alpha, \tilde{\pi}, \varpi) & =\frac{\nabla_{\varpi \varpi^{\prime}}^{2} f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)}{f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)}-\frac{\nabla_{\varpi} f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)}{f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)} \frac{\nabla_{\varpi^{\prime}} f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)}{f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi)}
\end{aligned}
$$

with

$$
\begin{aligned}
\nabla_{\tilde{\pi} \tilde{\pi}^{\prime}}^{2} f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi) & =\nabla_{\tilde{\pi} \tilde{\pi}^{\prime}}^{2} \alpha_{1, t} f_{t}(\tilde{\phi})+\nabla_{\tilde{\pi}} \alpha_{1, t} \nabla^{\prime} f_{t}(\tilde{\phi}) \\
& +\nabla f_{t}(\tilde{\phi}) \nabla_{\tilde{\pi}^{\prime}} \alpha_{1, t}+\alpha_{1, t} \nabla^{2} f_{t}(\tilde{\phi}) \\
& -\nabla_{\tilde{\pi} \tilde{\pi}^{\prime}}^{2} \alpha_{1, t} f_{t}(\tilde{\varphi})-\nabla_{\tilde{\pi}} \alpha_{1, t} \nabla^{\prime} f_{t}(\tilde{\varphi}) \\
& -\nabla f_{t}(\tilde{\varphi}) \nabla_{\tilde{\pi}^{\prime}} \alpha_{1, t}+\left(1-\alpha_{1, t}\right) \nabla^{2} f_{t}(\tilde{\varphi}) \\
\nabla_{\tilde{\pi} \varpi^{\prime}}^{2} f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi) & =\nabla_{\tilde{\pi} \varpi^{\prime}}^{2} \alpha_{1, t} f_{t}(\tilde{\phi})+\nabla_{\tilde{\pi}} \alpha_{1, t} \nabla^{\prime} f_{t}(\tilde{\phi}) D_{\tilde{\phi}, \varpi}^{(1)} \\
& +\nabla f_{t}(\tilde{\phi}) \nabla_{\varpi^{\prime}} \alpha_{1, t}+\alpha_{1, t} \nabla^{2} f_{t}(\tilde{\phi}) D_{\tilde{\phi}, \varpi}^{(1)} \\
& -\nabla_{\tilde{\pi} \varpi^{\prime}}^{2} \alpha_{1, t} f_{t}(\tilde{\varphi})-\nabla_{\tilde{\pi}} \alpha_{1, t} \nabla^{\prime} f_{t}(\tilde{\varphi}) D_{\tilde{\varphi}, \varpi}^{(1)} \\
& -\nabla f_{t}(\tilde{\varphi}) \nabla_{\varpi^{\prime}} \alpha_{1, t}+\left(1-\alpha_{1, t}\right) \nabla^{2} f_{t}(\tilde{\varphi}) D_{\tilde{\varphi}, \varpi}^{(1)} \\
\nabla_{\varpi \varpi^{\prime}}^{2} f_{2, t}^{\pi}(\alpha, \tilde{\pi}, \varpi) & =\nabla_{\varpi \varpi^{\prime}}^{2} \alpha_{1, t} f_{t}(\tilde{\phi})+\nabla_{\varpi} \alpha_{1, t} \nabla^{\prime} f_{t}(\tilde{\phi}) D_{\tilde{\phi}, \varpi}^{(1)} \\
& +\alpha_{1, t} D_{\tilde{\phi}, \varpi}^{(1) \prime} \nabla^{2} f_{t}(\tilde{\phi}) D_{\tilde{\phi}, \varpi}^{(1)}+D_{\tilde{\phi}, \varpi}^{(1) \prime} \nabla f_{t}(\tilde{\phi}) \nabla_{\varpi^{\prime}} \alpha_{1, t} \\
& -\nabla_{\varpi \varpi \varpi^{\prime}}^{2} \alpha_{1, t} f_{t}(\tilde{\varphi})-\nabla_{\varpi} \alpha_{1, t} \nabla^{\prime} f_{t}(\tilde{\varphi}) D_{\tilde{\varphi}, \varpi}^{(1)} \\
& +\left(1-\alpha_{1, t}\right) D_{\tilde{\varphi}, \varpi}^{(\tilde{\varphi}) \prime} \nabla^{2} f_{t}(\tilde{\varphi}) D_{\tilde{\varphi}, \varpi}^{(1)}-D_{\tilde{\varphi}, \varpi}^{(\tilde{\varphi}) \prime} \nabla f_{t}(\tilde{\varphi}) \nabla_{\varpi^{\prime}} \alpha_{1, t}
\end{aligned}
$$

For brevity, we omit the expressions of $\nabla_{\tilde{\pi} \tilde{\pi}^{\prime}}^{2} \alpha_{1, t}, \nabla_{\tilde{\pi} \varpi^{\prime}}^{2} \alpha_{1, t}$, and $\nabla_{\varpi \varpi^{\prime}}^{2} \alpha_{1, t}$. Evaluated at $(\alpha, \tilde{\pi}, \varpi)=$ $\left(\alpha, \tilde{\pi}^{*}, 0\right)$ we get

$$
\begin{aligned}
\nabla_{\tilde{\pi} \tilde{\pi}^{\prime}}^{2} f_{2, t}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right) & =\nabla^{2} f_{t}\left(\tilde{\pi}^{*}\right) \\
\nabla_{\tilde{\pi} \varpi^{\prime}}^{2} f_{2, t}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right) & =0 \\
\nabla_{\varpi \varpi^{\prime}}^{2} f_{2, t}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right) & =\nabla_{\varpi} \alpha_{1, t}^{*} \nabla_{(2, \ldots, p+2)}^{\prime} f_{t}\left(\tilde{\pi}^{*}\right)+\nabla_{(2, \ldots, p+2)} f_{t}\left(\tilde{\pi}^{*}\right) \nabla_{\varpi^{\prime}} \alpha_{1, t}^{*}+\alpha_{1} \alpha_{2} \nabla_{(2, \ldots, p+2)(2, \ldots, p+2)}^{2} f_{t}\left(\tilde{\pi}^{*}\right) \\
& =\nabla_{\varpi} \alpha_{1, t}^{*} \nabla_{(2, \ldots, p+2)}^{\prime} f_{t}\left(\tilde{\pi}^{*}\right)+\nabla_{(2, \ldots, p+2)} f_{t}\left(\tilde{\pi}^{*}\right) \nabla_{\varpi^{\prime}} \alpha_{1, t}^{*}+\alpha_{1} \alpha_{2} \nabla_{(2, \ldots, p+2)(2, \ldots, p+2)}^{2} f_{t}\left(\tilde{\pi}^{*}\right)
\end{aligned}
$$

so that

$$
\begin{aligned}
\nabla_{\tilde{\pi} \tilde{\pi}^{\prime}}^{2} l_{t}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)= & \frac{\nabla^{2} f_{t}\left(\tilde{\pi}^{*}\right)}{f_{t}\left(\tilde{\pi}^{*}\right)}-\frac{\nabla f_{t}\left(\tilde{\pi}^{*}\right)}{f_{t}\left(\tilde{\pi}^{*}\right)} \frac{\nabla^{\prime} f_{t}\left(\tilde{\pi}^{*}\right)}{f_{t}\left(\tilde{\pi}^{*}\right)} \\
\nabla_{\tilde{\pi} w^{\prime}}^{2} l_{t}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)= & 0 \\
\nabla_{\varpi \varpi \omega^{\prime}}^{2} l_{t}^{\pi}\left(\alpha, \tilde{\pi}^{*}, 0\right)= & \frac{\nabla_{\varpi} \alpha_{1, t}^{*} \nabla_{(2, \ldots, p+2)}^{\prime} f_{t}\left(\tilde{\pi}^{*}\right)+\nabla_{(2, \ldots, p+2)} f_{t}\left(\tilde{\pi}^{*}\right) \nabla_{\varpi^{\prime}} \alpha_{1, t}^{*}+\alpha_{1} \alpha_{2} \nabla_{(2, \ldots, p+2)(2, \ldots, p+2)}^{2} f_{t}\left(\tilde{\pi}^{*}\right)}{f_{t}\left(\tilde{\pi}^{*}\right)} \\
= & \alpha_{1} \alpha_{2}\left[\frac{\nabla_{(2, \ldots, p+2)} \mathrm{n}_{p}\left(\tilde{\pi}^{*}\right)}{\mathrm{n}_{p}\left(\tilde{\pi}^{*}\right)} \frac{\nabla_{(2, \ldots, p+2)}^{\prime} f_{t}\left(\tilde{\pi}^{*}\right)}{f_{t}\left(\tilde{\pi}^{*}\right)}+\frac{\nabla_{(2, \ldots, p+2)} f_{t}\left(\tilde{\pi}^{*}\right)}{f_{t}\left(\tilde{\pi}^{*}\right)} \frac{\nabla_{(2, \ldots, p+2)}^{\prime}}{\mathrm{n}_{p}\left(\tilde{\pi}_{p}\right)}\right. \\
& \left.\quad+\frac{\nabla_{(2, \ldots, p+2)(2, \ldots, p+2)}^{2} f_{t} f_{t}\left(\tilde{\pi}^{*}\right)}{f_{t}\left(\tilde{\pi}^{*}\right)}\right] .
\end{aligned}
$$

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[^1]:    ${ }^{1}$ Specifically, $\mu_{1}=\tilde{\phi}_{0} / \tilde{\phi}(1)$ and $\mu_{2}=\tilde{\varphi}_{0} / \tilde{\varphi}(1)$, whereas $\boldsymbol{\Gamma}_{m, p}(m=1,2)$ is a $p \times p$ symmetric Toeplitz matrix with $\gamma_{m, 0}=\operatorname{Cov}\left[\nu_{m, t}, \nu_{m, t}\right]$ along the main diagonal and $\gamma_{m, i}=\operatorname{Cov}\left[\nu_{m, t}, \nu_{m, t-i}\right], i=1, \ldots, p-1$, on the diagonals above and below the main diagonal. Here $\mu_{1}$ and $\mu_{2}$, as well as the elements of $\boldsymbol{\Gamma}_{1, p}$ and $\boldsymbol{\Gamma}_{2, p}$, are treated as functions of the parameters $\dot{\phi}$ and $\tilde{\varphi}$, respectively; for details of this dependence, see Lütkepohl (2005, eqn (2.1.39)).
    ${ }^{2}$ A major motivation for defining the mixing weights as in (6) is that an explicit form for the stationary distribution of $\left(y_{t}, \ldots, y_{t-p+1}\right)$, and hence of $y_{t}$, can be obtained and shown to be a mixture of normal distributions with known expectations and covariance matrices and with mixing weights $\alpha$ and $1-\alpha$ (for details, see Kalliovirta et al. (2015)).

[^2]:    ${ }^{3}$ Note that assuming the autoregressive parameters $\phi$ and $\varphi$ to have a common parameter space is made for simplicity and could be relaxed; for an example where such a relaxation would be needed, see the ACR model of Bec et al. (2008). ${ }^{4}$ See the remarks following Proposition 5 in KS (2012) and property (a) on p. 1633 of KS (2015).

[^3]:    ${ }^{5}$ Note that the parameter space for $\tilde{\phi}$ is the compact set $\tilde{\Phi}$ and not the entire stationarity region of a (causal) linear $\operatorname{AR}(p)$ model. Asymptotically, also the OLS estimator can be used.
    ${ }^{6}$ Set, e.g., $p=1, \tilde{\phi}_{0}=\tilde{\varphi}_{0}=0$, and $\tilde{\sigma}_{1}^{2} /\left(1-\tilde{\phi}_{1}^{2}\right)=\tilde{\sigma}_{2}^{2} /\left(1-\tilde{\varphi}_{1}^{2}\right)$; the last equality can hold even if $\left(\tilde{\phi}, \tilde{\sigma}_{1}^{2}\right) \neq\left(\tilde{\varphi}, \tilde{\sigma}_{2}^{2}\right)$.

[^4]:    ${ }^{7}$ In what follows, $\nabla f_{t}(\cdot)$ denotes differentiation of $f_{t}(\cdot)$ in (3) with respect to $\tilde{\phi}=\left(\tilde{\phi}_{0}, \tilde{\phi}_{1}, \ldots, \tilde{\phi}_{p}, \tilde{\sigma}_{1}^{2}\right)$.

[^5]:    ${ }^{8}$ Here and in what follows, a subset of the listed assumptions would sometimes suffice for the stated results.

[^6]:    ${ }^{9}$ This result could also be obtained from Andrews and Ploberger (1995, Sec 2.2, 2.4) as their assumptions 1-5 appear to be satisfied in the LMAR case.

[^7]:    ${ }^{10}$ To justify (49), partition the index set $\mathfrak{I}$ as $\mathfrak{I}=\left(\mathfrak{I}_{1}, \mathfrak{J}_{2}\right)$ with $\mathfrak{I}_{1}=\left((1,1),(2,2), \ldots\left(q_{2}, q_{2}\right)\right)$ and $\mathfrak{I}_{2}=$ $\left((1,2),(1,3), \ldots,\left(1, q_{2}\right),(2,3), \ldots,\left(q_{2}-1, q_{2}\right)\right)$. With straightforward algebra,
    $\left[B_{i j}\right]_{(i, j) \in \mathcal{J}}\left[c_{i j} c_{k l} A_{i j k l}\right]_{(i, j, k, l) \in \mathcal{J} \times \mathcal{J}}\left[B_{k l}\right]_{(k, l) \in \mathcal{J}}^{\prime}$
    $=\sum_{(i, j) \in \mathcal{J}} \sum_{(k, l) \in \mathcal{J}} c_{i j} c_{k l} A_{i j k l} B_{i j} B_{k l}$
    $=\frac{1}{4} \sum_{(i, j) \in \mathcal{I}_{1}} \sum_{(k, l) \in \mathcal{I}_{1}} A_{i j k l} B_{i j} B_{k l}+\frac{1}{2} \sum_{(i, j) \in \mathcal{I}_{1}} \sum_{(k, l) \in \mathcal{I}_{2}} A_{i j k l} B_{i j} B_{k l}+\frac{1}{2} \sum_{(i, j) \in \mathcal{J}_{2}} \sum_{(k, l) \in \mathcal{I}_{1}} A_{i j k l} B_{i j} B_{k l}+\sum_{(i, j) \in \mathcal{J}_{2}} \sum_{(k, l) \in \mathcal{I}_{2}} A_{i j k l} B_{i j} B_{k l}$
    $=\frac{1}{4}\left[\sum_{i=j} \sum_{k=l} A_{i j k l} B_{i j} B_{k l}+2 \sum_{i=j} \sum_{k<l} A_{i j k l} B_{i j} B_{k l}+2 \sum_{i<j} \sum_{k=l} A_{i j k l} B_{i j} B_{k l}+4 \sum_{i<j} \sum_{k<l} A_{i j k l} B_{i j} B_{k l}\right]$
    $=\frac{1}{4} \sum_{i=1}^{q_{2}} \sum_{j=1}^{q_{2}} \sum_{k=1}^{q_{2}} \sum_{l=1}^{q_{2}} A_{i j k l} B_{i j} B_{k l}$,

[^8]:    ${ }^{11}$ As before, $\nabla$ denotes differentiation with respect to $\tilde{\phi}=\left(\tilde{\phi}_{0}, \tilde{\phi}_{1}, \ldots, \tilde{\phi}_{p}, \tilde{\sigma}_{1}^{2}\right)$ and $\nabla_{i}, i=1, \ldots, p+2$, with respect to the $i$ th component of $\tilde{\phi}$.

