# Optimal Instrumental Variables Estimation for ARMA 

Models

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In this paper a new class of Instrumental Variables estimators for linear processes and in particular ARMA models is developed. Previously, IV estimators based on lagged observations as instruments have been used to account for unmodelled MA(q) errors in the estimation of the AR parameters. Here it is shown that these IV methods can be used to improve efficiency of linear time series estimators in the presence of unmodelled conditional heteroskedasticity. Moreover an IV estimator for both the AR and MA parts is developed. Estimators based on a Gaussian likelihood are inefficient members of the class of IV estimators analyzed here when the innovations are conditionally heteroskedastic.

Keywords: ARMA, conditional heteroskedasticity, instrumental variables, efficiency lowerbound, frequency domain.

JEL Classification: C13, C22

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## 1. Introduction ${ }^{2}$

This paper considers instrumental variables (IV) estimators for linear time series models. Efficient estimation in this framework has been studied by Hayashi and Sims (1983), Stoica, Söderström and Friedlander (1985; 1987a,b) and Hansen and Singleton (1991, 1996). In these papers efficient estimation of autoregressive roots under the presence of moving average errors has been analyzed. The moving average part of the model is not estimated but rather treated as a nuisance parameter. The class of instruments is restricted to linear functions of past observations. It is also assumed in this literature that the innovations are conditionally homoskedastic. Instrumental variables methods using overidentifying restrictions have also been applied to autoregressive (AR) models in the context of missing observations by Chen and Zadrozny (1998).

Here it is shown that the same type of IV estimators based on linear functions of past observations can be used to improve efficiency of estimators for linear time series models in the presence of unmodelled conditional heteroskedasticity. A consequence of the results of this paper is that standard estimators of linear process models based on Gaussian pseudo maximum likelihood (PML) functions are inefficient generalized method

[^1]of moments (GMM) estimators if the innovations are conditionally heteroskedastic. This means in particular that ordinary least squares (OLS) estimators for autoregressive models of order $\mathrm{p}(\mathrm{AR}(\mathrm{p}))$ are inefficient GMM estimators if the innovations are heteroskedastic.

In Kuersteiner (1999) a feasible efficient IV estimator for the AR(p) model with martingale innovations satisfying some additional moment restrictions is developed. In this paper we extend the previous results in two directions. First, the form of the optimal instrument is analyzed under more general assumptions about the innovation process, relaxing some of the restrictions on the fourth moments of the innovations that were previously maintained. Second, the class of models is extended to include general ARMA (p,q) processes that require minimization of a nonlinear criterion function. While Kuersteiner (1999) is mainly concerned with the semiparametric implementation of the IV procedure we focus on identification issues which are the most important aspect in extending the previous results to nonlinear models.

In addition, the paper extends the current literature in two directions. First, IV estimation is extended to general autoregressive-moving average (ARMA) models when the innovations are conditionally heteroskedastic. Second, for the class of IV estimators with linear instruments the paper derives exact functional forms of optimal filters of the type developed in Hansen and Singleton (1991) for a simpler estimation problem. It is shown how the filters depend on fourth order cumulants of the innovation distribution and the impulse response function of the underlying process. This formulation allows to give exact conditions on the distribution of the error process under which optimal instrumental variables estimators are feasible. A detailed analysis of the properties of the optimal weight
matrix is provided.
The results in this paper are presented for the case of martingale difference innovations driving the linear process. Alternatively similar formulas with the same efficiency implications could be obtained under the weaker assumption of white noise innovations. In this case the space of permissible instruments is generated by all linear combinations of past observations and the efficiency bounds would be identical to the bounds of Hansen (1985) and Hansen, Heaton and Ogaki (1988). In the case of martingale difference innovations Hansen's bounds are based on a larger class of instruments and are therefore tighter than the bounds obtained here.

The main technical difficulty in extending previous procedures to the estimation of the moving average case lies in the consistency proof. We give a general characterization of instrument processes that lead to consistent estimators. We then establish that the optimal instrument satisfies these criteria.

In this paper we do not focus on implementation issues. For most parts of the analysis it is assumed that the optimal instrument is known a priori. It is clear that in practice a procedure for estimating the weight matrix is needed. In Kuersteiner (1997) such a feasible procedure is developed under stronger assumptions on the martingale difference innovations. If these assumptions are satisfied then the procedures developed in Kuersteiner (1997) can be directly applied to the present context. Explicit formulas are provided for this case. We also discuss feasible versions of the optimal procedure under the more general conditions analyzed in this paper without giving proofs of feasibility. In this case the feasible estimator depends on a bandwidth parameter. Monte Carlo simulations for this
case are reported to give some guidance in the choice of the bandwidth parameter.
The paper is organized as follows. Section 2 introduces the assumptions about the innovation sequence and specifies the inference problem. Section 3 develops an instrumental variables estimator for estimation of linear process models and proves consistency and asymptotic normality of estimators for the ARMA class. In Section 4 it is shown how to factorize the asymptotic covariance matrix of this class of instrumental variables estimators in a way to obtain a lower bound. Section 5 uses the lowerbound to derive an explicit formulation of the optimal IV estimator depending on the data periodogram and an optimal frequency domain filter. Numerical Examples for the ARMA $(1,1)$ model and some Monte Carlo simulations are reported in Section 6. Proofs of some important lemmas are contained in Appendix A while the proofs of the results in the paper are contained in Appendix B.

## 2. Model Specification

The econometrician observes a finite stretch of data $\left\{y_{t}\right\}_{t=1}^{n}$ from a univariate process which is generated by the following mechanism

$$
\begin{equation*}
y_{t}=\sum_{j=0}^{\infty} c(\beta, j) \varepsilon_{t-j} \tag{1}
\end{equation*}
$$

for a given $\beta=\beta_{0} \in \mathbb{R}^{d}$ and $c(\beta, j): \mathbb{R}^{d} \times \mathbb{N} \rightarrow \mathbb{R}$. The parameter $\beta_{0}$ is unknown but the functions $c(., j)$ are known. We define the lag polynomial $C(\beta, L)=\sum_{j=0}^{\infty} c(\beta, j) L^{j}$ where $L$ is the lag operator and impose the identifying restriction $c(\beta, 0)=1$.

The innovations $\varepsilon_{t}$ are assumed to be a univariate martingale difference sequence. The
martingale difference property imposes restrictions on the fourth order cumulants. These restrictions can be conveniently summarized by defining the following function

$$
\sigma(s, r)=\left\{\begin{array}{ll}
E\left(\varepsilon_{t}^{2} \varepsilon_{t-|s|} \varepsilon_{t-|r|}\right) & \text { if } r \neq s  \tag{2}\\
E\left(\varepsilon_{t}^{2} \varepsilon_{t-s}^{2}\right)-\sigma^{4} & \text { if } r=s
\end{array} \text { for } r, s \in\{0, \pm 1, \pm 2, \ldots\}\right.
$$

where $\sigma^{4}=\left(E \varepsilon_{t}^{2}\right)^{2}$. It should be emphasized that $\sigma(s, r)$ is equal to the fourth order cumulant for $s, r>0$. We assume that we have a probability space $(\Omega, \mathcal{F}, P)$ with a filtration $\mathcal{F}_{t}$ of increasing $\sigma$-fields such that $\mathcal{F}_{t} \subseteq \mathcal{F}_{t+1} \subseteq \mathcal{F} \forall t$. The doubly infinite sequence of random variables $\left\{\varepsilon_{t}\right\}_{t=-\infty}^{\infty}$ generates the filtration $\mathcal{F}_{t}$ such that $\mathcal{F}_{t}=\sigma\left(\varepsilon_{t}, \varepsilon_{t-1}, \ldots\right)$. The assumptions on $\left\{\varepsilon_{t}\right\}_{t=-\infty}^{\infty}$ are summarized as follows:

Assumption A1. (i) $\varepsilon_{t}$ is strictly stationary and ergodic, (ii) $E\left(\varepsilon_{t} \mid \mathcal{F}_{t-1}\right)=0$ almost surely, (iii) $E\left(\varepsilon_{t}^{2}\right)=\sigma^{2}>0$, (iv) $\sum_{r=1}^{\infty} \sum_{s=1}^{\infty}|\sigma(s, r)|=B<\infty$, (v) $E\left(\varepsilon_{t}^{2} \varepsilon_{t-s}^{2}\right)>\underline{\alpha}$ some $\underline{\alpha}>0$ for all $s$.

Remark 1. Assumption A1(ii) could be relaxed to $E \varepsilon_{t} \varepsilon_{s}=0$ for $t \neq s$ at the cost of slightly more complicated expressions for the optimal instruments. Assumption (iii) guarantees that $\varepsilon_{t}$ has a nondegenerate distribution. Assumption (iv) limits the dependence in higher moments by imposing a summability condition on the fourth cumulants. The assumption is needed to prove invertibility of the infinite dimensional weight matrix of the optimal GMM estimator. Assumption (v) together with (iii) rules out degenerate joint distributions of $\varepsilon_{t}$ and $\varepsilon_{t-s}$.

Remark 2. It can be checked that processes in the autoregressive conditionally het-
eroskedastic family such as ARCH, GARCH, EGARCH as well as stochastic volatility models satisfy the assumptions, provided that $E \varepsilon_{t}^{4}<\infty$. It is well known from Milhoj (1985) or Nelson (1990) that this condition is satisfied only if additional restrictions limiting the temporal dependence of conditional variances and/or the innovation distribution are imposed.

Assumption (A1) implies that $\varepsilon_{t}^{2}$ is strictly stationary and ergodic and therefore covariance stationary. It should be emphasized that no assumptions about third moments are made. In particular this allows for skewness in the error process.

For the special case of an $\operatorname{ARMA}(\mathrm{p}, \mathrm{q})$ process, the lag polynomial has the familiar rational form

$$
\begin{equation*}
C(\beta, z)=\frac{\theta(z)}{\phi(z)} \tag{3}
\end{equation*}
$$

with $\theta(z)=1-\theta_{1} z-\ldots-\theta_{q} z^{q}$ and $\phi(z)=1-\phi_{1} z-\ldots-\phi_{p} z^{p}$ and $\beta^{\prime}=\left(\phi_{1}, \ldots, \phi_{p}, \theta_{1}, \ldots, \theta_{q}\right)$. Let $g_{y y}(\beta, \lambda)=\left|C\left(\beta, e^{i \lambda}\right)\right|^{2}$ where $|z|=\left(z z^{*}\right)^{1 / 2}$ for $z \in \mathbb{C}$ and $z^{*}$ is the complex conjugate of $z$. Under Assumption (A1), the spectrum of $y_{t}$ is given by $f_{y y}(\beta, \lambda)=\frac{\sigma^{2}}{2 \pi} g_{y y}(\beta, \lambda)$.

Further restrictions on $C\left(\beta, e^{i \lambda}\right)$ are needed to insure identification of the model and for consistency and asymptotic normality of the estimators. The necessary assumptions are discussed in Hannan (1973), Dunsmuir and Hannan (1976), and Deistler, Dunsmuir and Hannan (1978). As shown in these articles, a careful distinction between convergence of the parameters in $c(\beta, j)$ and the structural form parameters is needed. Consistency proofs typically establish convergence in the pointwise topology. An identification condition is then needed to obtain convergence in the quotient topology.

Some of the results of this paper are presented for the general formulation $C(\beta, z)$.

At some points however a specialization to the ARMA case is made in order to obtain sharper results. This is especially the case for the consistency proof. In that case abstract high level assumptions can be made precise for the specific functional form of the ARMA model.

In the general case the functions $c(\beta, j) \in C\left(\left[\mathbb{R}^{d} \times \mathbb{N}\right], \mathbb{R}\right)$ are restricted to satisfy the following additional constraints.

Assumption B1. Let $C(\beta, z)=\sum_{j=0}^{\infty} c(\beta, j) z^{j}$. The parameter space $\Theta$ is a closed subset of $\Theta^{\prime}$ defined by $\Theta^{\prime}=\left\{\left.\beta \in \mathbb{R}^{d}| | C(\beta, z)\right|^{-2} \neq 0\right.$ for $|z| \leq 1,|C(\beta, z)|^{2} \neq 0$ for $\left.|z| \leq 1\right\}$ where we assume that $\Theta^{\prime}$ is open in $\mathbb{R}^{d}$ and has a compact closure denoted by $\bar{\Theta}$. Assume $\beta_{0} \in \Theta$. The coefficients $c(\beta, j)$ are twice continuously differentiable in $\beta \in \Theta$ for all $j$ and $c(\beta, 0)=1$. We require for $\beta \in \Theta$ that $\sum_{j=0}^{\infty}|j||c(\beta, j)|<\infty$ and $\sum_{j=0}^{\infty}|j|\left|\frac{\partial}{\partial \beta} c(\beta, j)\right|<\infty$.

Assumption B2. For all $\beta, \beta_{0} \in \Theta, g_{y y}\left(\beta_{0}, \lambda\right) \neq g_{y y}(\beta, \lambda)$ whenever $\beta \neq \beta_{0}$ for some subsets $L \subset[-\pi, \pi]$ with nonzero Lebesgue measure.

Assumption B3. For a neighborhood $U$ of $\beta_{0}, U \subset \Theta_{0}, \partial^{2} g_{y y}(\beta, \lambda) / \partial \beta \partial \beta^{\prime}$ is continuous in $\lambda \in[-\pi, \pi]$ and $\beta \in U$.

Remark 3. Assumption (B1) implies that the functions $g_{y y}(\beta, \lambda)$ and $\partial g_{y y}(\beta, \lambda) / \partial \beta$ are Lipschitz continuous. The Lipschitz condition also implies that $g_{y y}^{-1}(\beta, \lambda)$ is Lipschitz continuous on $\Theta$ and therefore that $\frac{\partial}{\partial \beta} \ln g_{y y}(\beta, \lambda)$ is Lipschitz continuous on $\Theta$.

Remark 4. Assumption (B1) is stronger than C2.2 in Dunsmuir (1979) where on the other hand conditional homoskedasticity is assumed. The stronger summability restrictions are needed to justify approximations of the instruments based on the innovation sequence.

Remark 5. Assuming $\Theta$ to be compact is of little practical importance and is commonly done in the time series literature. See for example Hosoya and Taniguchi (1982), Kabaila (1980), Taniguchi (1983).

The assumptions specified here are sufficient to identify the parameters $\beta$ in $C\left(\beta, e^{i \lambda}\right)$. For specific functional forms of $C\left(\beta, e^{i \lambda}\right)$ the assumptions can be made more explicit. A leading example is the ARMA model where the identifiable subset of $\mathbb{R}^{d}$ can be described more accurately. The following Assumption is equivalent to the previous assumptions for the case of an ARMA model.

Assumption B4. Let $C(\beta, z)=\theta(z) / \phi(z)$. The parameter space $\Theta$ is a closed subset of $\Theta^{\prime}$ defined by $\Theta^{\prime}=\left\{\beta \in \mathbb{R}^{d} \mid \phi(z) \neq 0\right.$ for $|z| \leq 1, \theta(z) \neq 0$ for $|z| \leq 1, \theta(z), \phi(z)$ have no common zeros, $\theta_{q} \neq 0$ or $\left.\phi_{p} \neq 0\right\}$. Assume $\beta_{0} \in \Theta$.

Remark 6. Deistler, Dunsmuir and Hannan (1978) show that $\Theta$ defined in Assumption (B4) satisfies the conditions of Assumption (B1). It is easy to show that all ARMA models in $\Theta$ satisfy the summability and differentiability requirements of (B1).

In the following analysis of the IV estimator results will first be obtained for the general linear process case. It will then be shown that high level assumptions needed for these results are satisfied for the case when Assumptions (B1-B3) are specialized to (B4).

## 3. Instrumental Variables Estimators

In this section a class of instrumental variables estimators is introduced. The instruments are constructed from linear filters of lagged innovations $\varepsilon_{t}$. An alternative, equivalent for-
mulation would be to allow for linear filters of the observable process $y_{t}$. Estimators of this form have been proposed by Hayashi and Sims (1983), Stoica, Söderström and Friedlander (1985, 1987a,b) and Hansen and Singleton (1991). Efficiency of these procedures is achieved by exploiting all the moment conditions of the form $E \varepsilon_{t} \varepsilon_{s}=0$ for $t \neq s$. The innovations $\varepsilon_{t}$ are functions of the observable data $g\left(y_{t}, y_{t-1}, \ldots, \beta_{0}\right)=g_{t}\left(\beta_{0}\right)=\varepsilon_{t}$. If we stack a finite number $m$ of innovations in $\varepsilon_{t}^{m}=\left[\varepsilon_{t-1}, \ldots, \varepsilon_{t-m}\right]^{\prime}$ then a standard GMM estimator minimizes the population equivalent of $E\left(g(\beta) \varepsilon_{t}^{m}\right)^{\prime} \Omega_{m}^{-1} E\left(g(\beta) \varepsilon_{t}^{m}\right)$ where $\Omega_{m}$ is a suitable weight matrix. The first order conditions of this problem are given by $E\left(\frac{\partial g(\beta)}{\partial \beta^{\prime}} \varepsilon_{t}^{m}\right)^{\prime} \Omega_{m}^{-1} E\left(g(\beta) \varepsilon_{t}^{m}\right)=0$. We will later use the notation $P_{m}=E\left(\frac{\partial g(\beta)}{\partial \beta^{\prime}} \varepsilon_{t}^{m}\right)$ where $P_{m}$ is a $m \times d$ dimensional matrix. In the context of linear instrumental variables estimators $P_{m}$ corresponds to the matrix of covariances between regressors and instruments. Using this notation we can now set up an equivalent problem which is to solve the set of equations $E\left(g(\beta) P_{m}^{\prime} \Omega_{m}^{-1} \varepsilon_{t}^{m}\right)=0$. Letting $z_{t}=P_{m}^{\prime} \Omega_{m}^{-1} \varepsilon_{t}^{m}$ be the $d$ dimensional vector of instruments then leads to the formulation of an equivalent, exactly identified IV estimator that solves the population equivalent of $E\left(g(\beta) z_{t}\right)=0$. Clearly, the two formulations of the problem are equivalent as far as their first order asymptotic efficiency properties are concerned.

The advantage of this transformation lies in the fact that in our context of ARMA models the matrix $P_{m}$ can be estimated $\sqrt{n}$ consistently irrespective of the size of $m$. It is shown in Kuersteiner $(1997,1999)$ that under additional restrictions on the weight matrix $\Omega_{m}$ it is possible to set $m=n$ or in other words to let the number of moment conditions grow at the same rate as the sample size. This is not the case for the usual implementation
of GMM where $E\left(g(\beta) \varepsilon_{t}^{m}\right)$ is replaced by a sample average.
We now discuss the estimation problem for a general class of linear instruments. Introduce the space of absolutely summable sequences $l^{1}$ such that $x \in l^{1}$ if $\sum\left|x_{j}\right|<\infty$ for $x=\left\{x_{j}\right\}_{j=1}^{\infty}$. Define the set $\mathcal{A}$ of sequences of vectors $a_{j} \in \mathbb{R}^{d}$ such that

$$
\mathcal{A}=\left\{a=\left\{a_{j}\right\}_{j=1}^{\infty}: a_{j} \in \mathbb{R}^{d},\left\{\left[a_{j}\right]_{k}\right\}_{j=1}^{\infty} \in l^{1} \text { for all } 1 \leq k \leq d\right\}
$$

where $[.]_{k}$ denotes the $k$-th element of a vector. We define $z_{t}(\omega): \Omega \rightarrow \mathbb{R}^{d}$ for all $t$ as

$$
z_{t}=\sum_{k=1}^{\infty} a_{k} \varepsilon_{t-k} \text { a.s. }
$$

for $a \in \mathcal{A}$, $a$ fixed. The instruments satisfy the orthogonality condition

$$
\begin{equation*}
E\left[\left(C^{-1}\left(\beta_{0}, L\right) y_{t}\right) z_{t}\right]=0 \tag{4}
\end{equation*}
$$

since $C^{-1}\left(\beta_{0}, L\right) y_{t}=\varepsilon_{t}$ from (1). The estimator based on this condition is constructed in the time domain. If $C^{-1}\left(\beta_{0}, L\right)$ is of infinite order, as is the case for $M A(q)$ models, a sample analog to (4) needs to be based on an approximation. Such an approximation can be conveniently analyzed in the frequency domain. It should be stressed however that the estimator is set up in time domain. Let the expansion of the polynomial $C^{-1}(\beta, L)$ be $C^{-1}(\beta, L)=\sum_{j=0}^{\infty} \tilde{c}_{j}^{\beta} L^{j}$. The sample analog of the moment restriction is then given by

$$
\begin{equation*}
G_{n}(\beta, a)=\frac{1}{n} \sum_{t=1}^{n} z_{t} \sum_{j=0}^{t-1} \tilde{c}_{j}^{\beta} y_{t-j}=0 \tag{5}
\end{equation*}
$$

for $a \in \mathcal{A}$. The criterion function is indexed by $a$ to emphasize the fact that each choice of an instrument results in a different estimator. From (4) we see that $z_{t}$ has to be approximated as well. Discussion of this issue will be delayed until Section 5 where an optimal instrument is considered. For the time being it is therefore assumed that $z_{t}$ is known.

In the frequency domain the analog of (4) is

$$
\int_{-\pi}^{\pi} C^{-1}\left(\beta_{0}, e^{-i \lambda}\right) f_{y z}(\lambda) d \lambda=0
$$

where $f_{y z}(\lambda)=\sum_{j=-\infty}^{\infty} \gamma_{y z}(j) e^{-i \lambda j}$ and $\gamma_{y z}(j)=E y_{t} z_{t-j}$. We set

$$
G(\beta, a)=(2 \pi)^{-1} \int_{-\pi}^{\pi} C^{-1}\left(\beta, e^{-i \lambda}\right) f_{y z}(\lambda) d \lambda
$$

Note that $f_{y z}(\lambda)$ typically is a complex vector valued function $f_{y z}(\lambda):[-\pi, \pi] \rightarrow \mathbb{C}^{d}$. Also note that $\int_{-\pi}^{\pi} C^{-1}\left(\beta, e^{-i \lambda}\right) f_{y z}(\lambda) d \lambda$ is real valued.

We introduce discrete Fourier transforms of the data defined as $\omega_{n, y}(\lambda)=\frac{1}{\sqrt{n}} \sum_{t=1}^{n} y_{t} e^{-i t \lambda}$ and for the instrument as $\omega_{n, z}(\lambda)=\frac{1}{\sqrt{n}} \sum_{t=1}^{n} z_{t} e^{-i t \lambda}$. The cross periodogram is $I_{n, y z}(\lambda)=$ $\omega_{n, y}(\lambda) \omega_{n, z}(-\lambda)$. It is easy to check that $G_{n}(\beta, a)$ defined in (5) is identical to

$$
G_{n}(\beta, a)=(2 \pi)^{-1} \int_{-\pi}^{\pi} C^{-1}\left(\beta, e^{-i \lambda}\right) I_{n, y z}(\lambda) d \lambda
$$

We follow Hansen (1982) in defining the estimator $\beta_{n}$. Unless otherwise stated all conditions are for $a \in \mathcal{A}$, $a$ fixed

Assumption C1. The sequence of estimators $\beta_{n} \in \mathbb{R}^{d}$ is defined by

$$
\begin{equation*}
\beta_{n}=\underset{\beta \in \Theta}{\arg \min }\left\|G_{n}(\beta, a)\right\|^{2} . \tag{6}
\end{equation*}
$$

Assumption C2. Let the sets $B_{k}\left(\beta_{0}\right)$ for $k=1,2, \ldots$ form a countable local base ${ }^{3}$ around $\beta_{0}$. The sets $B_{k}\left(\beta_{0}\right)$ can be taken as the set of balls with rational radius centered at $\beta_{0}$. Let $z_{t}=\sum_{k=1}^{\infty} a_{k} \varepsilon_{t-k}$ a.s. where $\varepsilon_{t-k}$ satisfies Assumption (A1). Let $\mathcal{A}^{*} \subseteq \mathcal{A}$ be the set of all sequences $\left\{a_{k}\right\}_{k=1}^{\infty}$ such that

$$
\mathcal{A}^{*}=\left\{\left.a \in \mathcal{A}\right|_{\beta \in B_{k}\left(\beta_{0}\right)^{C} \cap \Theta}\left\|G\left(\beta_{n}, a\right)\right\|>0 \text { for } k=1,2, \ldots\right\}
$$

where $B_{k}\left(\beta_{0}\right)^{C}$ are the complements of $B_{k}\left(\beta_{0}\right)$. Assume that $\mathcal{A}^{*} \neq \emptyset$.

Remark 7. Assumption (C1) is the definition of the estimator. We show in the consistency proof that $\left\|G_{n}\left(\beta_{n}, a\right)\right\|^{2}=0$ almost surely is implied by the assumptions on $G_{n}$. (C2) is a familiar identification condition for global identification. Assumption (C2) specializes (B2) to the case of IV estimators. The difference is that identification of an entire class of estimators indexed by a needs to be guaranteed. Identification depends on the choice of the instrument or $a \in \mathcal{A}$ and does not hold for all $a \in \mathcal{A}$ (see Remark 10 for an example). We therefore define the subset $\mathcal{A}^{*}$ of instruments that satisfy the identification condition.

[^2]This imposes restrictions on $z_{t}$ or $a$. A complete description of the set $\mathcal{A}^{*}$ is possible for a given parametric class $C(\beta, z)$. A characterization will be given for the ARMA case.

Lemma 3.1. Assume (A1), (B1-B3), (C1-C2). Let $z_{t}=\lim _{m \rightarrow \infty} A_{m}^{\prime} \varepsilon_{t}^{m}$ a.s. with $A_{m}^{\prime}=$ $\left[a_{1}, \ldots, a_{m}\right],\left\{a_{k}\right\}_{k=1}^{\infty} \in \mathcal{A}^{*}$ and $\varepsilon_{t}^{m}=\left[\varepsilon_{t-1}, \ldots, \varepsilon_{t-m}\right]^{\prime}$. Then $\beta_{n} \rightarrow \beta_{0}$ almost surely.

Consistency of the IV estimator depends both on restrictions on the parameter space and the instruments $z_{t}$. Assumption (C2) restricts the class of allowable instruments. The conditions given are necessarily high level without further knowledge regarding the function $C(\beta, L)$. For practical purposes it is however important to characterize the set of instruments $\mathcal{A}^{*}$ leading to consistent estimators. In the case of an $\operatorname{ARMA}(\mathrm{p}, \mathrm{q})$ model it is possible to give conditions on the sequences $a \in \mathcal{A}^{*}$. This is done in the next proposition.

Proposition 3.2. Assume $C(\beta, L)=\theta_{0}(L) / \phi_{0}(L)$ is an $\operatorname{ARMA}(p, q)$ lag operator and the parameter space $\Theta$ satisfies Assumption (B4). Let

$$
S=\left\{x=\left[x_{1}, \ldots\right] \in l^{2}: \phi_{0}(L) x=0 \text { for } x_{j}, j>d,\left[x_{1}, \ldots, x_{d}\right]^{\prime}=\kappa, \kappa \in \mathbb{R}^{d}\right\}
$$

be the set of solutions to the difference equation $\phi_{0}(L) x=0$ with $d$ initial conditions к. Define $\operatorname{ker} A^{\prime}=\left\{x \in l^{2}: A^{\prime} x=0\right\}$ for $A=\left[a_{1}, \ldots\right]^{\prime}$ and $a \in \mathcal{A}$. Let $a \in \mathcal{A}$ with $A_{d}=\left[a_{1}, \ldots, a_{d}\right]^{\prime}$ where $d=p+q$. If $p=0$ and $A_{d}$ is nonsingular then $a \in \mathcal{A}^{*}$. If $0<p$, $A=\left[a_{1}, \ldots .\right]^{\prime}$ is of full column rank and $\operatorname{ker} A^{\prime} \cap S=0$ then $a \in \mathcal{A}^{*}$.

Remark 8. Proposition (3.2) shows that ARMA models can be consistently estimated by instrumental variables techniques provided that the instruments satisfy the specified restrictions.

Remark 9. The usual conditions for consistency in $I V$ estimation are $E \varepsilon_{t} z_{t}=0$ and $E \frac{\partial g_{t}(\beta)}{\partial \beta} z_{t}$ is of full rank. For linear models these two conditions are equivalent to Assumption (C2). In our context C2 may hold even if $E \frac{\partial g_{t}(\beta)}{\partial \beta} z_{t}$ is of reduced rank. An example is an $\operatorname{ARMA}(1,1)$ model with instruments $z_{t}$ defined by $a_{1}=[1,0]^{\prime}, a_{2}=\left[-\theta_{0}^{-1}, 1\right]^{\prime}$ and $a_{3}=\left[0,-\theta_{0}^{-1}\right]^{\prime}$ with $a_{j}=0$ for $j>3$. Then $A$ is of full row rank, ker $A^{\prime} \cap S=0$ but $E \frac{\partial g_{t}(\beta)}{\partial \beta} z_{t}$ is of reduced rank. On the other hand $E \frac{\partial g_{t}(\beta)}{\partial \beta} z_{t}$ being full rank implies $A^{\prime}$ to be of full row rank and $A^{\perp} \cap S=0$. In other words the identification conditions given here are weaker than standard conditions would imply.

Remark 10. If the instruments are replaced by $a_{1}=[1,0]^{\prime}, a_{2}=\left[-\phi_{0}^{-1}, 1\right]$ and $a_{3}=$ $\left[0,-\phi_{0}^{-1}\right]$ in the $\operatorname{ARMA}(1,1)$ example then $\operatorname{ker} A^{\prime} \cap S \neq 0$ which shows that $A^{\prime}$ being full row rank is not sufficient for identification and clearly $\mathcal{A}^{*} \neq \mathcal{A}$.

We now state additional assumptions that are sufficient to establish a result for the limiting distribution of $\sqrt{n}\left(\beta_{n}-\beta_{0}\right)$. Introduce the notation $\dot{\eta}(\beta, \lambda)=\partial \ln C\left(\beta, e^{-i \lambda}\right) / \partial \beta$ and $b_{k}=(2 \pi)^{-1} \int \dot{\eta}\left(\beta_{0}, \lambda\right) e^{i k \lambda} d \lambda$. It follows immediately that $b_{-k}=0$ and $b_{0}=0$. Let $l_{a}(\lambda)=\sum_{k=1}^{\infty} a_{k} e^{-i \lambda k}$ and define the matrices $P_{m}^{\prime}=\left[b_{1}, \ldots, b_{m}\right], A_{m}^{\prime}=\left[a_{1}, \ldots, a_{m}\right]$ and

$$
\Omega_{m}=\left[\begin{array}{lll}
\alpha_{1,1} & \cdots & \alpha_{1, m}  \tag{7}\\
\vdots & \ddots & \vdots \\
& & \\
\alpha_{m, 1} & \cdots & \alpha_{m, m}
\end{array}\right]
$$

where

$$
\alpha_{s, r}=\left\{\begin{array}{ll}
\sigma(s, r) & \text { if } s \neq r  \tag{8}\\
\sigma(r, r)+\sigma^{4} & \text { if } s=r
\end{array} .\right.
$$

It is easy to check that $\lim _{m} P_{m}^{\prime} A_{m}=(2 \pi)^{-1} \int \dot{\eta}\left(\beta_{0}, \lambda\right) l_{a}(-\lambda)^{\prime} d \lambda$. Also note that using our earlier notation $\lim _{m} \sigma^{2} P_{m}^{\prime} A_{m}=E \frac{\partial g_{t}(\beta)}{\partial \beta} z_{t}$. The following additional conditions are needed to prove the existence of a limiting distribution of $\beta_{n}$.

Assumption D1. Define $\mathcal{A}^{* *} \subseteq \mathcal{A}$ as $\mathcal{A}^{* *}=\left\{a \in \mathcal{A} \mid \operatorname{det} \int \dot{\eta}\left(\beta_{0}, \lambda\right) l_{a}(-\lambda)^{\prime} d \lambda \neq 0\right\}$. Assume that $\mathcal{A}^{*} \cap \mathcal{A}^{* *} \neq \emptyset$.

Remark 11. Assumption (D1) guarantees that there is an instrument $a \subseteq \mathcal{A}$ that satisfies both the identification condition $C 2$ and $\operatorname{det} \int \dot{\eta}\left(\beta_{0}, \lambda\right) l_{a}(-\lambda)^{\prime} d \lambda \neq 0$. Assumption (D1) corresponds to Assumption 3.4 in Hansen (1982). Note that for linear models there is no difference between the identification assumptions and Assumption (D1) while this is not true for nonlinear models. In general C2 does not imply D1.

The next lemma shows that for the ARMA model D1 does indeed imply C2 while we have seen in Remark 9 that the reverse is not true.

Lemma 3.3. Assume $C(\beta, L)=\theta_{0}(L) / \phi_{0}(L)$ is an $\operatorname{ARMA}(p, q)$ lag operator and the parameter space $\Theta$ satisfies Assumption (B4) and that $a \in \mathcal{A}$. Then $a \in \mathcal{A}^{* *} \Rightarrow a \in \mathcal{A}^{*}$.

The limiting distribution of the instrumental variables estimator is stated in the next theorem. For notational efficiency define $\lim _{m \rightarrow \infty} \sigma^{-4}\left(P_{m}^{\prime} A_{m}\right)^{-1} A_{m}^{\prime} \Omega_{m} A_{m}\left(A_{m}^{\prime} P_{m}\right)^{-1}=$ $\sigma^{-4}\left(P^{\prime} A\right)^{-1} A^{\prime} \Omega A\left(A^{\prime} P\right)^{-1}$. This notation will be justified in the next section in terms of operators on infinite dimensional spaces.

Theorem 3.4. Assume (A1), (B1-B3), (C1, C2) and (D1). Let $z_{t}=\lim _{m \rightarrow \infty} A_{m}^{\prime} \varepsilon_{t}^{m}$ with $A_{m}^{\prime}=\left[a_{1}, \ldots, a_{m}\right],\left\{a_{k}\right\}_{k=1}^{\infty} \in \mathcal{A}^{*} \cap \mathcal{A}^{* *}$ and $\varepsilon_{t}^{m}=\left[\varepsilon_{t-1}, \ldots, \varepsilon_{t-m}\right]^{\prime}$. Then the estimator
defined by $\beta_{n}=\arg \min \left\|G_{n}\left(\beta_{n}, a\right)\right\|^{2}$ has a limiting distribution given by

$$
\sqrt{n}\left(\beta_{n}-\beta_{0}\right) \xrightarrow{d} N\left(0, \sigma^{-4}\left(P^{\prime} A\right)^{-1} A^{\prime} \Omega A\left(A^{\prime} P\right)^{-1}\right)
$$

Remark 12. If $\beta_{n}$ is obtained from minimizing a Gaussian PML criterion function then the asymptotic covariance matrix is $\sigma^{-4}\left(P^{\prime} P\right)^{-1} P^{\prime} \Omega P\left(P^{\prime} P\right)^{-1}$. Such an estimator therefore corresponds to an IV estimator where $A=P$. This shows that Gaussian estimators have the interpretation of inefficient IV or GMM estimators when the innovations are conditionally heteroskedastic.

The main result of the paper will now be developed in two steps. We first obtain a lower bound for the covariance matrix

$$
\begin{equation*}
\sigma^{-4}\left(P^{\prime} A\right)^{-1} A^{\prime} \Omega A\left(A^{\prime} P\right)^{-1} \tag{9}
\end{equation*}
$$

in the next section. This lower bound is then used to construct an optimal instrumental variables estimator.

## 4. Covariance Matrix Lowerbound

Finding a lower bound for (9) poses certain technical difficulties having to do with the infinite dimensional nature of the instrument space. We investigate the properties of the fourth order cumulant matrix $\Omega_{m}$, first by holding $m$ fixed and then by looking at a related infinite dimensional problem. In particular we establish that the infinite dimensional operator $\Omega$, associated with $\Omega_{m}$ in a way to be defined, has a well behaved inverse.

Invertibility of $\Omega_{m}$ for all $m$ is not enough to show that $\Omega$ is invertible. We briefly review the theory of invertible operators (see Gohberg and Goldberg, 1980), p.65. For two Banach spaces $B_{1}$ and $B_{2}$ denote the set of bounded linear operators mapping $B_{1}$ into $B_{2}$ by $L\left(B_{1}, B_{2}\right)$. Then $A \in L\left(B_{1}, B_{2}\right)$ is invertible if there exists an operator $A^{-1} \in$ $L\left(B_{2}, B_{1}\right)$ such that $A^{-1} A x=x$ for all $x \in B_{1}$ and $A A^{-1} y=y$ for all $y \in B_{2}$. Let ker $A=\left\{x \in B_{1}: A x=0\right\}$ and $\operatorname{Im} A=\left\{A x: x \in B_{1}\right\}$. Then $A$ is invertible if ker $A=\{0\}$ and $\operatorname{Im} A=B_{2}$.

Following Hanani, Netanyahu and Reichaw (1968) we now choose $B_{1}, B_{2}$ as linear spaces whose points are sequences of real numbers denoted by $x=\left\{x_{1}, x_{2}, \ldots\right\}$ and $y=$ $\left\{y_{1}, y_{2}, \ldots\right\}$. Define the norm $\|x\|_{2}=\left(\sum_{i}\left|x_{i}\right|^{2}\right)^{1 / 2}$. Then $B$ is the space of all sequences that are bounded under the $\|\cdot\|_{2}$ norm and is denoted by $l^{2}$. An operator $A: l^{2} \mapsto l^{2}$ is defined by the infinite dimensional matrix $A=\left(a_{i, j}\right), i, j=1,2, \ldots$ such that $y=A x \in l^{2}$ for all $x \in l^{2}$. The operator $A$ can be defined element by element as $y_{i}=\sum_{j}^{\infty} a_{i, j} x_{j}$ for all $i$. The operator $A$ is invertible if the only solution to $A x=0$ is $x=\{0,0, \ldots$.$\} and$ $\operatorname{Im} A=l^{2}$. Note that $l^{2}$ is a Hilbert space with inner product $\langle x, y\rangle=\sum_{j}^{\infty} x_{j} y_{j}$. From Theorem 11.4 in Gohberg and Goldberg (1980) it follows $\operatorname{Ker} A^{\perp}=\operatorname{Im} A$ for a self adjoint operator $A$. It is thus enough to show $\operatorname{Ker} A=0$ for $A: l^{2} \rightarrow l^{2}, A$ selfadjoint.

Consider now the following infinite dimensional operator associated with $\Omega_{m}$. Define the operator $\Omega$ component-wise by its image for all $x \in l^{2}$ by $b_{i}=\lim _{m \rightarrow \infty} \sum_{j}^{m} \alpha_{i, j} x_{j}$ where $\alpha_{i, j}$ is defined in (8). In other words $\Omega$ is the infinite dimensional matrix such that any left upper corner sub matrix of dimension $m \times m$ has the same elements as $\Omega_{m}$.

Lemma 4.1. Let $\Omega_{m}$ be defined as in (7). Then $\Omega_{m}^{-1}$ exists for all $m, \Omega \in L\left(l^{2}, l^{2}\right)$ and
$\Omega^{-1}$ exists.

Proof. See Appendix B

Remark 13. The fact that the image of $\Omega$ is square summable, i.e. $\Omega x \in l^{2}$, depends on the summability properties of $\sigma(k, l)$. The interpretation of the summability condition is that the instruments $\varepsilon_{t}$ become unrelated in their fourth cumulants as the time spread between them increases.

By the Closed Graph Theorem (Gohberg and Goldberg (1980), Theorem X.4.2) it also follows that $\Omega^{-1}$ is bounded, i.e., $\left\|\Omega^{-1}\right\|=\sup _{\|x\|_{2} \leq 1}\left\|\Omega^{-1} x\right\|_{2}<\infty$. Thus $\sup _{i, j}\left|\omega_{i, j}\right|<\infty$ where $\left[\Omega^{-1}\right]_{i, j}=\omega_{i, j}$.

Next, we need to establish properties of the matrix $\Omega_{m}^{-1}$ as $m$ tends to infinity. In particular we want to establish that the inverse $\Omega_{m}^{-1}$ approximates $\Omega^{-1}$ as $m \rightarrow \infty$.

Lemma 4.2. Let $\Omega_{m}$ be as defined in (7). Define $\Omega_{m}^{-1}$ such that $\Omega_{m}^{-1} \Omega_{m}=I_{m}$ and $\Omega_{m} \Omega_{m}^{-1}=I_{m} \forall m$. Let

$$
\Omega_{m}^{*}=\left[\begin{array}{ll}
\Omega_{m} & 0  \tag{10}\\
0 & \sigma^{4} I
\end{array}\right]
$$

where I stands for an infinite dimensional identity matrix. Then $\Omega_{m}^{*-1}$ exists and $\left\|\Omega_{m}^{*-1}-\Omega^{-1}\right\| \rightarrow$ 0 as $m \rightarrow \infty$.

Proof. See Appendix B

Remark 14. A consequence of the convergence of $\Omega_{m}^{*-1}$ to $\Omega^{-1}$ in the operator norm is that $\left\|P^{\prime}\left(\Omega_{m}^{*-1}-\Omega^{-1}\right)\right\|_{2} \rightarrow 0$ in the $l^{2}$ norm. In other words $\operatorname{Var}\left(z_{t}^{m}-z_{t}\right) \rightarrow 0$ as $m \rightarrow \infty$
where $z_{t}^{m}=P^{\prime} \Omega_{m}^{*-1} \varepsilon_{t}^{\infty}$ and $z_{t}=P^{\prime} \Omega^{-1} \varepsilon_{t}^{\infty}$. In this sense Lemma (4.2) provides an algorithm to approximate the infinite dimensional inverse $\Omega^{-1}$. For $m$ fixed, the finite dimensional inverse of $\Omega_{m}$ is computed and used to construct $\Omega_{m}^{*-1}$. The resulting instrument then converges in a mean squared sense.

We define the $d$ dimensional product of sequence spaces $l_{d}^{2}=l^{2} \times \ldots \times l^{2}$. Define the infinite dimensional matrix $P=\left[b_{1}, \ldots\right]^{\prime}$ by stacking elements of the sequence $\left\{b_{k}\right\}_{k=1}^{\infty} \in l_{d}^{2}$. Introduce notation for the reverse operation of extracting a sequence from the rows of a matrix by defining $b(P):=\left\{b_{k}\right\}_{k=1}^{\infty}$. Define the matrix $\Xi=\left(P^{\prime} \Omega^{-1} P\right)^{-1}$.

Using this notation we can state our next theorem which establishes a lower bound for the covariance matrix (9).

Theorem 4.3. For any $a \in \mathcal{A}$ let $A^{\prime}=\left[a_{1}, \ldots\right]$ and $P$ and $\Omega$ as previously defined. If $a(A) \in \mathcal{A}^{* *}$ then the matrix $\left(P^{\prime} A\right)^{-1} A^{\prime} \Omega A\left(A^{\prime} P\right)^{-1}$ satisfies

$$
\left(P^{\prime} A\right)^{-1} A^{\prime} \Omega A\left(A^{\prime} P\right)^{-1}-\left(P^{\prime} \Omega^{-1} P\right)^{-1} \geq 0
$$

where $\geq 0$ stands for positive semi-definite.

Proof. See Appendix B

Remark 15. If $a \in \mathcal{A}^{*} \cap \mathcal{A}^{* *}$ then $\left(P^{\prime} A\right)^{-1} A^{\prime} \Omega A\left(A^{\prime} P\right)^{-1}$ is the asymptotic covariance matrix of an estimator based on $a$. However, it is important to point out that the lowerbound is for IV estimators in the class of all instruments which are linear functions of the innovation process and have an innovation filter in $\mathcal{A}^{* *}$. The construction of the lower
bound does not involve consistency restrictions for the instruments. In order to construct an efficient estimator in practice it has to be established that the optimal instrument does in fact satisfy consistency restrictions.

## 5. Optimal Instrumental Variables Estimators

Theorem (4.3) immediately leads to the construction of an efficient IV estimator. The optimal instrument is determined by the linear filter $A^{\prime}=P^{\prime} \Omega^{-1}$. It is not a priori true that the optimal filter also results in a consistent estimator. However, for important parametric examples such as the ARMA class this is indeed the case.

Theorem 5.1. Assume $C(\beta, L)=\theta(L) / \phi(L)$ and the parameter space $\Theta$ satisfies Assumption (B4). If $A^{\prime}=P^{\prime} \Omega^{-1}$ then the sequence $a=a\left(P^{\prime} \Omega^{-1}\right)$ defined by the rows of $A$ satisfies $a \in \mathcal{A}^{*} \cap \mathcal{A}^{* *}$.

Theorem (5.1) together with Theorem (3.4) and Theorem (4.3) establish that the IV estimator for the ARMA model constructed with instruments based on $A^{\prime}=P^{\prime} \Omega^{-1}$ achieves a lowerbound of the same type as in Hansen and Singleton (1991) but under the weaker martingale difference sequence assumptions on $\varepsilon_{t}$ detailed in Assumption (A1). This result is summarized in the following Corollary.

Corollary 5.2. Assume (A1), $C(\beta, L)=\theta_{0}(L) / \phi_{0}(L)$ is an $\operatorname{ARMA}(p, q)$ lag operator, the parameter space $\Theta$ satisfies Assumption (B4) and $z_{t}=\lim _{m \rightarrow \infty} P_{m}^{\prime} \Omega_{m}^{-1} \varepsilon_{t}^{m}$. Then

$$
\sqrt{n}\left(\beta_{n}-\beta_{0}\right) \xrightarrow{d} N\left(0, \sigma^{-4}\left(P^{\prime} \Omega^{-1} P\right)^{-1}\right)
$$

Feasible versions of the optimal IV procedure have to be based on approximations of the optimal instrument $z_{t}$. Such approximations proceed in two steps. First, the infinite number of unobserved innovations is replaced by a finite number of proxies based on the observed sample and constructed from $\hat{\varepsilon}_{t}=y_{t}-\sum_{j=1}^{t-1} c\left(\beta_{0}, j\right) y_{t-j}$ for $t=1, \ldots, n$ where $\hat{\varepsilon}_{1}=y_{1}$. This leads to a pseudo feasible estimator $\tilde{\beta}$ based on the instruments $\hat{z}_{t}=\sum_{j=1}^{t-1} a_{j} \hat{\varepsilon}_{t-j}$. A fully feasible estimator denoted by $\tilde{\beta}(\hat{a})$ is obtained by substituting $\beta_{0}$ for a first stage consistent estimator $\hat{\beta}$ in the construction of $\hat{\varepsilon}_{t}$ and by replacing the weights $a_{j}$ by estimated quantities $\hat{a}_{j}$. Gaussian PMLE procedures which are consistent but inefficient in our context can be used to generate first stage estimators $\hat{\beta}$.

The empirical analog of the moment restriction for the pseudo feasible estimator becomes

$$
\begin{equation*}
\tilde{G}_{n}(\tilde{\beta}, a)=\frac{1}{n} \sum_{t=1}^{n} \hat{z}_{t} \sum_{j=0}^{t-1} \tilde{c}_{j}^{\beta} y_{t-j}=0 . \tag{11}
\end{equation*}
$$

It is shown in the proof of Corollary (5.4) that $\sup _{\beta \in \Theta}\left|\tilde{G}_{n}^{F}(\beta, a)-\tilde{G}_{n}(\tilde{\beta}, a)\right|=O_{p}\left(n^{-1}\right)$ where

$$
\begin{equation*}
\tilde{G}_{n}^{F}(\beta, a)=(2 \pi)^{-1} \int_{-\pi}^{\pi} C^{-1}\left(\beta, e^{-i \lambda}\right) h\left(\beta_{0}, \lambda\right) I_{n, y y}(\lambda) d \lambda \tag{12}
\end{equation*}
$$

Here $I_{n, y y}(\lambda)$ is the data periodogram and the filter $h(\lambda):[-\pi, \pi] \rightarrow \mathbb{C}^{d}$ is defined as $h\left(\beta_{0}, \lambda\right)=l_{\psi}(-\lambda) C^{-1}\left(\beta_{0}, e^{i \lambda}\right)$ with

$$
l_{\psi}(\lambda)=\sum_{j=1}^{\infty} a_{j} e^{-i \lambda j}
$$

The coefficients of the optimal instrument are given by

$$
a_{j}=\sum_{k=1}^{\infty} b_{k} \omega_{k j}
$$

where $b_{k}$ is the Fourier coefficient of the derivative of the log spectral density of $y_{t}$ and $\omega_{k j}$ is the $k j$-th entry of the inverse $\Omega^{-1}$. The $b_{k}$ coefficients have simple interpretations in special parametric models. In the case of an $A R(p)$ model for example they are equivalent to the impulse response function and can therefore be computed easily. It can also be noted that the Gaussian estimators are obtained by setting $a_{j}=b_{j}$.

It is shown in Kuersteiner (1997) that a sufficient condition for the validity of the approximation $\hat{z}_{t}$ for $z_{t}$ is that the coefficients of the instruments satisfy

$$
\begin{equation*}
\sum_{j=1}^{\infty} j\left|\left[a_{j}\right]_{k}\right|<\infty \text { for } k=1, \ldots, d \tag{13}
\end{equation*}
$$

The following lemma shows that under strengthened summability restrictions on the fourth order cumulants Condition (13) is satisfied for the optimal instrumental variables estimator of the $\operatorname{ARMA}(p, q)$ model. The Corollary summarizes the result that based on Theorem (5.3) the instruments can be approximated without affecting the first order asymptotic distribution.

Theorem 5.3. Assume $C(\beta, L)=\theta(L) / \phi(L)$ and the parameter space $\Theta$ satisfies $A s$ sumption (B4). Strengthen Assumption (A1v) to $\sum_{r=1}^{\infty} \sum_{s=1}^{\infty} s|\sigma(s, r)|=B<\infty$. By symmetry this implies $\sum_{r=1}^{\infty} \sum_{s=1}^{\infty} r|\sigma(s, r)|=B<\infty$. If $A^{\prime}=P^{\prime} \Omega^{-1}$ then $a=a\left(P^{\prime} \Omega^{-1}\right)$ satisfies (13).

Proof. See Appendix B

Corollary 5.4. Assume that the conditions of Theorem 5.3 hold together with Assumption A1. Let $\tilde{\beta} \in \Theta$ satisfy $\tilde{G}_{n}(\tilde{\beta}, a)=o_{p}\left(n^{-1 / 2}\right)$. Then $\sqrt{n}\left(\tilde{\beta}-\beta_{n}\right)=o_{p}(1)$.

## Proof. See Appendix B

Feasible versions of the optimal estimator are then obtained by replacing $\tilde{G}_{n}(\beta, a)$ by $\tilde{G}_{n}(\beta, \hat{a})$ where in $\tilde{G}_{n}(\beta, \hat{a})$ we replace $h\left(\beta_{0}, \lambda\right)$ by $\hat{l}_{\psi}(\lambda) C^{-1}\left(\hat{\beta}, e^{i \lambda}\right)$ and $\hat{\beta}$ is a consistent first stage estimate. The challenging part is to estimate $\hat{l}_{\psi}(\lambda)$ consistently. Here we only discuss a special case where $\Omega_{m}$ is diagonal. This case with the additional restrictions $\alpha_{j, k}=0$ for $j \neq k$ on the moments of $\varepsilon_{t}$ has been analyzed in Kuersteiner (1997) in the context of estimating an $\operatorname{AR}(\mathrm{p})$ model. The restriction $\alpha_{j, k}=0$ is satisfied for GARCH processes with symmetric innovation distributions and stochastic volatility models. Under these circumstances it is possible to estimate $\hat{l}_{\psi}(\lambda)$ consistently without the need to introduce bandwidth parameters controlling the number of instruments. The simplification comes from the fact that in that particular case $\Omega^{-1}$ is diagonal such that $a_{j}=b_{j} / \alpha_{j, j}$. An estimate of the optimal instrument is obtained from

$$
\begin{equation*}
\hat{z}_{t}=\sum_{j=1}^{t-1} \hat{b}_{j} / \hat{\alpha}_{j, j} \hat{\varepsilon}_{t-j} \tag{14}
\end{equation*}
$$

where $\hat{b}_{j}=(2 \pi)^{-1} \int_{-\pi}^{\pi} \dot{\eta}(\hat{\beta}, \lambda) e^{i \lambda j} d \lambda$. To define $\hat{\alpha}_{j, j}$ introduce a truncation sequence $d_{n}=$ $c n^{-1 / 2+\nu}$ for some $0<\nu<1 / 2$ and some constant $c>0$. Then define $\hat{\alpha}_{j, j}=\hat{\alpha}_{j, j}^{*}$ if $\hat{\alpha}_{j, j}^{*}>d_{n}$ and $\hat{\alpha}_{j, j}=d_{n}$ otherwise where $\hat{\alpha}_{j, j}^{*}=\frac{1}{n} \sum_{t=1+j}^{n} \hat{\varepsilon}_{t}^{2} \hat{\varepsilon}_{t-j}^{2}$. In simulation experiments this truncation parameter does not affect the outcome and in practice $d_{n}=1$ works well. In
order to prove the result formally we need to require that additional higher order moments exists.

Assumption E1. Let $c_{\varepsilon . . . \varepsilon}\left(t_{1}, \ldots, t_{k-1}\right)$ be the $k$-th order cumulant of the error process $\varepsilon_{t}$ and assume that

$$
\sum_{t_{1}} \cdots \sum_{t_{k-1}}\left(1+\left|t_{j}\right|\right)\left|c_{\varepsilon \ldots \varepsilon}\left(t_{1}, \ldots, t_{k-1}\right)\right|<\infty, \text { for all } j=1, \ldots, k-1 \text { and } k=2,3, . ., 8
$$

Theorem 5.5. Assume (A1), $\alpha_{j, k}=0$ for $j \neq k$ and $(E 1), C(\beta, L)=\theta_{0}(L) / \phi_{0}(L)$ is an ARMA $(p, q)$ lag operator, the parameter space $\Theta$ satisfies Assumption (B4). Let $\tilde{\beta}(\hat{a}) \in \Theta$ satisfy $\tilde{G}_{n}(\tilde{\beta}, \hat{a})=o_{p}\left(n^{-1 / 2}\right)$ with $\hat{z}_{t}$ defined in (14). Then $\sqrt{n}\left(\tilde{\beta}(\hat{a})-\beta_{n}\right)=o_{p}(1)$.

Proof. See Appendix B
In the more general case where $\alpha_{j, k} \neq 0$ for $j \neq k$ the elements $\omega_{k j}$ can be estimated from a sample analog of the approximation matrix $\Omega_{m}^{*}$ defined in (10). Denote by $\omega_{k j}^{*}$ the $k, j$-th element of the inverse $\Omega_{m}^{*-1}$ and let $z_{t}^{*}=\sum_{j=1}^{t-1} \sum_{k=1}^{t-1} b_{k} \omega_{k j}^{*} \hat{\varepsilon}_{t-j}$ be the instrument based on the approximation and denote by $\tilde{G}_{n}\left(\beta, a^{*}(m)\right)$ the corresponding criterion function where $a^{*}(m)$ indicates that instruments depend on $m$. Let $\tilde{\beta}^{*}=\arg \min \left\|\tilde{G}_{n}\left(\beta, a^{*}(m)\right)\right\|^{2}$. From $\hat{\Omega}_{m}=n^{-1} \sum_{t=m+1}^{n} \hat{\varepsilon}_{t}^{2} \hat{\varepsilon}_{t}^{m} \hat{\varepsilon}_{t}^{m^{\prime}}$ we form the $n \times n$ matrix

$$
\hat{\Omega}_{m}^{*-1}=\left[\begin{array}{cc}
\hat{\Omega}_{m}^{-1} & 0 \\
0 & I_{n-m}\left(\hat{\sigma}^{4}\right)^{-1}
\end{array}\right]
$$

and obtain estimates $\hat{\omega}_{k j}^{*}$ of $\omega_{k j}$ where $\hat{\omega}_{k j}^{*}$ is the $k, j$-th element of $\hat{\Omega}_{m}^{*-1}$. Here $I_{n-m}$ is an $n-m$ dimensional identity matrix and $\hat{\sigma}^{4}=\left(\frac{1}{n} \sum_{t=1}^{n} \hat{\varepsilon}_{t}^{2}\right)^{2}$. We then form the truncated
estimate of the $a_{j}$ coefficients by setting $\hat{a}_{j}(m)=\sum_{k=1}^{n} \hat{b}_{k} \hat{\omega}_{k j}^{*}$. The feasible estimator resulting from minimizing $\left\|\tilde{G}_{n}(\beta, \hat{a}(m))\right\|^{2}$ is denoted by $\tilde{\beta}^{*}(\hat{a}(m))$. In the next section we use Monte Carlo experiments to evaluate some selection rules for $m$.

## 6. Numerical Examples

In this section we take the case of an $\operatorname{ARMA}(1,1)$ model and analyze its asymptotic efficiency properties for the case when $\varepsilon_{t}$ is an $\operatorname{ARCH}(1)$ process. We also show plots of the likelihood contour. We finally conduct a small Monte Carlo experiment to explore the performance of the IV procedures in finite samples.

The model we investigate is a univariate process $y_{t}$ defined as the stationary solution to

$$
\begin{equation*}
y_{t}=\phi y_{t-1}+\varepsilon_{t}-\theta \varepsilon_{t-1} . \tag{15}
\end{equation*}
$$

The innovations $\varepsilon_{t}$ are an $\operatorname{ARCH}(1)$ process generated by $\varepsilon_{t}=u_{t} h_{t}^{-1 / 2}$ and $h_{t+1}=\gamma_{0}+\gamma_{1} \varepsilon_{t}^{2}$ where $u_{t} \sim N(0,1)$ is an iid sequence of random variables. We assume that $\gamma_{0}>0$ and set it to . 1 for all numerical calculations. In this case $\beta=[\phi, \theta]^{\prime}$ and $\Theta^{\prime}=\{\beta| | \phi|<1,|\theta|<1\}$. The parameter $\alpha_{1}$ satisfies $0 \leq \gamma_{1}<\sqrt{1 / 3}$ which guarantees that $E \varepsilon_{t}^{4}<\infty$. It is shown in Milhoj (1985) that $\sigma^{4}=\left(\gamma_{0} / 1-\gamma_{1}\right)^{2}, \operatorname{Cov}\left(\varepsilon_{t}^{2}, \varepsilon_{t-i}^{2}\right)=2 \gamma_{0}^{2} \gamma_{1}^{i} /\left[\left(1-\gamma_{1}\right)^{2}\left(1-3 \gamma_{1}^{2}\right)\right]$ for $i>0$ and $E\left(\varepsilon_{t}^{2} \varepsilon_{t-k} \varepsilon_{t-j}\right)=0$ for $j \neq k$. It follows that $\Omega_{m}=\operatorname{diag}\left(\alpha_{1,1}, \ldots, \alpha_{m, m}\right)$ where $\alpha_{k, k}=\operatorname{Cov}\left(\varepsilon_{t}^{2}, \varepsilon_{t-k}^{2}\right)+\sigma^{4}$.

The polynomial $C(\beta, L)=(1-\theta L) /(1-\phi L)$ has $\log$ derivative $\dot{\eta}(\beta, \lambda)=\left[e^{-i \lambda} /(1-\right.$ $\left.\left.\phi e^{-i \lambda}\right),-e^{-i \lambda} /\left(1-\theta e^{-i \lambda}\right)\right]^{\prime}$. It follows that $b_{j}=\left[\phi^{j-1},-\theta^{j-1}\right]^{\prime}$.

For $\gamma_{1}=0$ it follows that $\operatorname{Cov}\left(\varepsilon_{t}^{2}, \varepsilon_{t-k}^{2}\right)=0$ and $\alpha_{k, k}=\sigma^{4}$ for $k>0$ since $\varepsilon_{t}=\gamma_{0} u_{t}$ is
iid in this case. The optimal instrument then is $z_{t}=\lim _{m \rightarrow \infty} P_{m}^{\prime} \varepsilon_{t}^{m}$ a.s. and the limiting distribution of $\beta_{n}$ is $\sqrt{n}\left(\beta_{n}-\beta_{0}\right) \rightarrow N\left(0,\left(P^{\prime} P\right)^{-1}\right)$ which is the same as the limiting distribution of the maximum likelihood estimator. Figure 1 shows a contour plot of the function - $\arctan \|G(\beta, a)\|^{2}$ which has the shape of a long flat valley with a unique global minimum at $\beta_{0}$ where $-\arctan \|G(\beta, a)\|^{2}=0$.

Figure 1

When $\gamma_{1}>0$ the optimal instrument is given by

$$
z_{t}=\left[\sum_{k=1}^{\infty} \phi_{0}^{k-1} / \alpha_{k, k} \varepsilon_{t-k},-\sum_{k=1}^{\infty} \theta_{0}^{k-1} / \alpha_{k, k} \varepsilon_{t-k}\right]^{\prime}
$$

In Figure 2 we plot the asymptotic relative efficiency $\sigma^{2}\left(\phi_{P M L}\right) / \sigma^{2}\left(\phi_{I V}\right)$ of the PML estimator which reaches .7 around $\phi=-.9$ and improves as $\phi$ gets closer to zero. The overall shape and magnitude of the efficiency improvement of IV over PML is very similar to the pure $\operatorname{AR}(1)$ model analyzed in Kuersteiner (1999). In Figure 3 we plot the same efficiency gains for the parameter $\theta$ when $\phi$ is varied over $[-1,1]$. Note that here the shape of the efficiency curve is much less symmetric.

Figure 2,3

We now report a Monte Carlo experiment to investigate how well the asymptotic efficiency properties of the IV procedure hold in finite samples. We generate samples of size
$n=2^{k}$ for $k=7,8, \ldots, 10$ from Model (15) with $\operatorname{ARCH}(1)$ innovations.
Starting values are $y_{0}=0, h_{0}=0$ and $\varepsilon_{0}=0$. In each sample the first 500 observations are discarded to eliminate dependence on initial conditions. Small sample properties are evaluated for different values of $\beta, \gamma_{1} \in[0,1)$. For GARCH processes it was shown in Milhoj (1985) and Bollerslev (1986) that asymptotic normality established in previous chapters only obtains for a subset of values for $\gamma_{1}$. Nevertheless, simulation results are reported for parametrizations outside this range in order to analyze the robustness of the proposed IV procedure to departures from the assumptions.

The parameter $\beta$ is estimated in several different ways. The PML estimator based on a Gaussian likelihood is denoted by $\hat{\beta}_{O L S}$ and was computed using the NAG Fortran routine G13AFF, Mark 18. This routine forces the parameters to be in $\Theta$. The optimal instrumental variables estimator is obtained from the consistent first stage estimator $\hat{\beta}_{O L S}$ as $\tilde{\beta}(\hat{a})=\arg \min \left\|\tilde{G}_{n}(\beta, \hat{a})\right\|^{2}$. We also compute the more general IV estimator $\tilde{\beta}(\hat{a}(m))=\arg \min \left\|\tilde{G}_{n}(\beta, \hat{a}(m))\right\|^{2}$ where $m$ is set to $5,10, \sqrt{n}, n^{2 / 5}$. Both problems are solved numerically with the constrained nonlinear optimization routine nag_nlp_sol of the NAG fl90 (release 3) library of Fortran 90 subroutines using $\hat{\beta}_{O L S}$ as a starting value.

We can compare the asymptotic gains reported in Figure 1 with the empirical efficiency of the estimators $\tilde{\beta}(\hat{a})$ and $\tilde{\beta}(\hat{a}(m))$ based on 1, 000 replications for sample sizes ranging from 128 to 1024 . The results are summarized in Tables 1-16. The results for $\gamma_{1}=.5$ in Tables 1-8 can be directly compared to the theoretical efficiency gain calculations. For the estimator $\tilde{\beta}(\hat{a})$ based on the diagonal weight matrix, reported in Tables 1-4, the sample sizes needed to achieve efficiency gains are relatively large with 1024 observations. As
expected, the strongest gains are achieved for $|\phi|>.5$. The moving average parameter $\theta$ is generally less well estimated by this method than the autoregressive parameter $\phi$. For parameter values near the point of non-identification where $\theta=\phi$ the IV procedure performs relatively poorly compared to Gaussian pseudo likelihood estimators. Overall the performance of the IV estimator is superior for large sample sizes and at well chosen points in the parameter space. For smaller sample sizes and at less favorable points in the parameter space the Gaussian estimators tend to dominate but IV performs reasonable even under these circumstances.

In Tables 5-8 we report the performance of the unrestricted IV procedures $\tilde{\beta}(\hat{a}(m))$. These procedures perform better than $\tilde{\beta}(\hat{a})$ and achieve gains over the Gaussian estimator at favorable points in the parameter space even in small samples of 128 observations. In larger samples they almost uniformly dominate Gaussian estimators, with the exception of points very close to $\theta=\phi$. As has to be expected from the theoretical calculations the estimators for $\theta$ again fare slightly worse than for $\phi$. The overall preferred choice for $m$ is $m=n^{2 / 5}$.

In Tables 9-16 we report results for $\gamma_{1}=.9$. Strictly speaking these experiments are not covered by the theoretical results because $\gamma_{1}>\sqrt{1 / 3}$ implies that $E \varepsilon_{t}^{4}$ does not exist. The IV estimators nevertheless perform well and in fact achieve even larger efficiency gains. This is especially true for $\tilde{\beta}(\hat{a}(m))$ which achieves variance reductions of up to $60 \%$ in large samples. There is however an outlier problem for this estimator. This is particularly evident from Table 14 where the mean absolute error for $\tilde{\beta}(\hat{a}(m))$ is improved over $\beta_{O L S}$ while the same is not always true for the variance measure. Looking at the inter quantile
range however shows that the distribution of $\tilde{\beta}(\hat{a}(m))$ is more concentrated than for $\beta_{O L S}$ at least as long as $|\phi|>.5$.

The simulations show that even if the true weight matrix $\Omega$ is diagonal it is better not to impose this restriction and to estimate the full non-diagonal weight matrix. The recommendation that can be given based on the simulation results is to use $\tilde{\beta}(\hat{a}(m))$ with $m=n^{2 / 5}$. The choice of $m$ reported here is clearly limited by the small scale of the Monte Carlo experiment and no claim is made that this choice is optimal from a theoretical point of view. The simulations however also seem to indicate that the performance of $\tilde{\beta}(\hat{a}(m))$ is not very sensitive to the choice of $m$ at least in the cases considered.

## 7. Conclusions

In this paper we have analyzed the instrumental variables estimator for stationary linear process models and ARMA models in particular. It was shown that a GMM estimator based on the infinite number of moment conditions $E \varepsilon_{t} \varepsilon_{s}$ can be constructed. The maintained hypothesis in this paper is that the innovations $\varepsilon_{t}$ are martingale difference sequences. The overidentified GMM estimator can be conveniently represented in the form of an exactly identified IV estimator. It is shown that under the additional restrictions used in Kuersteiner (1999) a feasible version of the GMM estimator is available even when the number of moment conditions used in estimation is the same as the number of observations.

The procedures proposed in this paper are optimal in the class of GMM estimators based on instruments that are linear in the observed data. While inclusion of nonlinear instruments is possible in principle and would improve asymptotic efficiency it creates
difficult issues of implementation. Linear instruments have the advantage of being approximately ordered by their time lag in terms of their importance. This ordering is lost when nonlinear instruments are included and a more sophisticated selection procedure has to be implemented to make any such procedure feasible.

The proposed procedures are developed for univariate time series models. Extensions to the class of vector ARMA models could be proved along similar lines. If $h_{0}(L) y_{t}=g_{0}(L) \varepsilon_{t}$ in the notation of Dunsmuir and Hannan (1976) then IV procedures would be constructed based on the moment conditions $E \varepsilon_{t} \otimes \varepsilon_{s}$. Consistent estimates then have to be based on instruments that are not orthogonal to any process $g^{-1}(L) h(L) y_{t}$ in the VARMA class.

Simulation evidence for a univariate $\operatorname{ARMA}(1,1)$ model shows that the procedures do achieve, sometimes significant, efficiency gains in finite samples. This is especially true when the model is well identified and the autoregressive parameter is larger than .5 in absolute value.

## A. Appendix - Lemmas

Lemma A.1. Under Assumption (A1) for each $m \in\{1,2 \ldots\}$, $m$ fixed, the vector

$$
\frac{1}{\sqrt{n}} \sum_{t=1}^{n}\left[\varepsilon_{t} \varepsilon_{t-1}, \ldots, \varepsilon_{t} \varepsilon_{t-m}\right] \Rightarrow N\left(0, \Omega_{m}\right)
$$

with $\Omega_{m}$ defined in (7).

Proof. We note that individually all the terms $\varepsilon_{t} \varepsilon_{t-k}$ with $k \geq 1$ are martingale difference sequences (mds). Now define $Y_{t}^{\prime}=\left[\varepsilon_{t} \varepsilon_{t-1}, \ldots, \varepsilon_{t} \varepsilon_{t-m}\right]$. Then it is enough to show that for all $\ell \in \mathbb{R}^{m}$ such that $\ell^{\prime} \ell=1$ we have $\frac{1}{\sqrt{n}} \sum \ell^{\prime} \tilde{Y}_{t} \Rightarrow N(0,1)$ where now $\tilde{Y}_{t}=\Omega_{m}^{-1 / 2} Y_{t}$ and $\Omega_{m}=E Y_{t} Y_{t}^{\prime}$. Note that $\ell^{\prime} \tilde{Y}_{t}$ is a mds and a martingale CLT (see Hall and Heyde, 1980, Theorem 3.2, p.52) can be applied to the sum $\sum_{t} Y_{n t}=\frac{1}{\sqrt{n}} \sum_{t} \tilde{Y}_{t}$.

Lemma A.2. Let $\varepsilon_{t}$ satisfy Assumption (A1). Then $\nexists \alpha \in l^{2}, \alpha \neq 0$ such that $\sum_{i=0}^{\infty} \alpha_{i} \varepsilon_{t-i}=$ 0 a.s.

Proof. If $\exists \alpha \in l^{2}, \alpha \neq 0$ such that $\sum \alpha_{i} \varepsilon_{t-i}=0$ a.s assume without loss of generality $\alpha_{1} \neq 0$. If $\alpha_{i}=0$ for all $i=2,3, \ldots$ then $\sum \alpha_{i} \varepsilon_{t-i}=0$ a.s. is trivially contradicted. Now assume $\alpha_{i} \neq 0$ for at least one $i=2,3, \ldots$ such that $\varepsilon_{t-1}=-\alpha_{1}^{-1} \sum_{i=2}^{\infty} \alpha_{i} \varepsilon_{t-i}$ a.s. But then $E\left(\varepsilon_{t-1} \mid \mathcal{F}_{t-2}\right)=-\alpha_{1}^{-1} \sum_{i=2}^{\infty} \alpha_{i} \varepsilon_{t-i}$ a.s. so that $E\left(\varepsilon_{t-1} \mid \mathcal{F}_{t-2}\right) \neq 0$ with positive probability. This contradicts the martingale difference assumption.

Lemma A.3. Let $I_{n, y z}(\lambda)=\omega_{n, y}(\lambda) \omega_{n, z}(-\lambda) . I_{n, \varepsilon \varepsilon}(\lambda)$ is the periodogram of $\left\{\varepsilon_{1}, \ldots, \varepsilon_{n}\right\}$. Assume $\varepsilon_{t}$ satisfy Assumption (A1) and that $y_{t}=\sum_{j=0}^{\infty} \psi_{j} \varepsilon_{t-j}$ with $\psi(\lambda)=\sum_{j=0}^{\infty} \psi_{j} e^{-i \lambda j}$
such that $\sum_{j=0}^{\infty}|j|\left|\psi_{j}\right|<\infty$. Also let $z_{t}=\sum_{j=1}^{\infty} a_{j} \varepsilon_{t-j}$ with $a \in \mathcal{A}$. Let $\varsigma$ (.) be a function on $[-\pi, \pi] \rightarrow \mathbb{C}$ with absolutely summable Fourier coefficients $\left\{c_{k},-\infty<k<\infty\right\}$ such that $\varsigma(\lambda)=\sum_{j=-\infty}^{\infty} c_{j} e^{-i \lambda j}$. Then for any $\eta, \epsilon>0$

$$
P\left(\sqrt{n}(2 \pi)^{-1}\left|\int_{-\pi}^{\pi} I_{n, y z}(\lambda) \varsigma(\lambda) d \lambda-\int_{-\pi}^{\pi} I_{n, \varepsilon z}(\lambda) C\left(\beta_{0}, \lambda\right) \varsigma(\lambda) d \lambda\right|>\eta\right)<\epsilon
$$

as $n \rightarrow \infty$.

Proof. First an expression for $R_{n}(\lambda)=I_{n, y z}(\lambda)-I_{n, \varepsilon z}(\lambda) \psi(\lambda)$ is obtained. Using

$$
\begin{equation*}
\omega_{n, y}(\lambda)=\psi(\lambda) \omega_{n, \varepsilon}(\lambda)+n^{-1 / 2} \sum_{j=0}^{\infty} \psi_{j} e^{-i \lambda j} U_{n j}(\lambda) \tag{16}
\end{equation*}
$$

where $U_{n j}(\lambda)=\sum_{t=1-j}^{n-j} \varepsilon_{t} e^{-i \lambda t}-\sum_{t=1}^{n} \varepsilon_{t} e^{-i \lambda t}$ leads to

$$
R_{n}(\lambda):=I_{n, y z}(\lambda)-\psi(\lambda) I_{n, \varepsilon z}(\lambda)=\omega_{z}(-\lambda) n^{-1 / 2} \sum_{j=0}^{\infty} \psi_{j} e^{-i \lambda j} U_{n j}(\lambda)
$$

Note that $(2 \pi)^{-1} \int R_{n}(\lambda) \varsigma(\lambda) d \lambda=n^{-1} \sum_{k=1}^{\infty} \sum_{l=0}^{\infty} \sum_{r=1}^{l} \sum_{m=-\infty}^{\infty} a_{k} \psi_{l} c_{m} \varepsilon_{r+j+m-k}\left(\varepsilon_{r-l}-\varepsilon_{n-l+r}\right)$. Then using the Markov inequality it is enough to consider

$$
E \sqrt{n}\left|(2 \pi)^{-1} \int_{-\pi}^{\pi} R_{n}(\lambda) \varsigma(\lambda) d \lambda\right| \leq 2 \sup _{k} \alpha_{k}^{1 / 2} n^{-1 / 2} \sum_{k=1}^{\infty} \sum_{l=0}^{\infty} \sum_{m=-\infty}^{\infty}\left|a_{k} \psi_{l} c_{m}\right||l| \rightarrow 0
$$

since the last term is bounded from $\sum_{k=1}^{\infty}\left|a_{k}\right|<\infty$ and $\sum_{l=0}^{\infty}|l|\left|\psi_{l}\right|<\infty$.

Lemma A.4. Let $I_{n, \varepsilon z}(\lambda)=\omega_{n, \varepsilon}(\lambda) \omega_{n, z}(-\lambda)$. Assume $\varepsilon_{t}$ satisfy Assumption (A1) and
let $z_{t}=\sum_{j=0}^{\infty} a_{j} \varepsilon_{t-j}$ with $a \in \mathcal{A}$. Then for any $\ell \in \mathbb{R}^{d}$ such that $\ell^{\prime} \ell=1$,

$$
n^{1 / 2}(2 \pi)^{-1} \int_{-\pi}^{\pi} \ell^{\prime} I_{n, \varepsilon z}(\lambda) d \lambda \xrightarrow{d} N\left(0, \sum_{l=1}^{\infty} \sum_{k=1}^{\infty} \alpha_{k, l} \ell^{\prime} a_{k} a_{l}^{\prime} \ell\right) .
$$

Proof. First note that $(2 \pi)^{-1} \int_{-\pi}^{\pi} I_{n, \varepsilon z}(\lambda) d \lambda=n^{-1} \sum_{t=1}^{n} \varepsilon_{t} z_{t}$ such that $E n^{1 / 2}(2 \pi)^{-1} \int_{-\pi}^{\pi} I_{n, \varepsilon z}(\lambda) d \lambda=$ 0 . It also follows that $\varepsilon_{t} z_{t}$ is a martingale difference sequence. However $z_{t}=\sum_{k=1}^{\infty} a_{k} \varepsilon_{t-k}$ such that a direct application of Lemma (A.1) is not possible.

For a fixed $m$ we introduce $z_{t}^{m}=\sum_{k=1}^{m} a_{k} \varepsilon_{t-k}$ such that $\omega_{n, z}^{m}(\lambda)=n^{-1 / 2} \sum_{t=1}^{k} z_{t}^{m} e^{-i \lambda k}$ and $I_{n, \varepsilon z}^{m}(\lambda)=\omega_{n, \varepsilon}(\lambda) \omega_{n, z}^{m}(-\lambda)$. From Billingsley (1968, Theorem 4.2) it is enough to show that for all $\epsilon>0$,

$$
\lim _{m \rightarrow \infty} \limsup _{n \rightarrow \infty} P\left\{\left|n^{1 / 2} \int_{-\pi}^{\pi} \ell^{\prime}\left(I_{n, \varepsilon z}^{m}(\lambda)-I_{n, \varepsilon z}(\lambda)\right) d \lambda\right| \geq \epsilon\right\}=0
$$

where

$$
n^{1 / 2}(2 \pi)^{-1} \int_{-\pi}^{\pi} \ell^{\prime}\left(I_{n, \varepsilon z}^{m}(\lambda)-I_{n, \varepsilon z}(\lambda)\right) d \lambda=n^{-1 / 2} \sum_{t=1}^{n} \sum_{k>m}^{\infty} \ell^{\prime} a_{k} \varepsilon_{t} \varepsilon_{t-k}
$$

Since $E a_{k} \varepsilon_{t} \varepsilon_{t-k}=0$ it is enough to consider

$$
n^{-1} E\left(\sum_{t=1}^{n} \sum_{k>m}^{\infty} \ell^{\prime} a_{k} \varepsilon_{t} \varepsilon_{t-k}\right)^{2}=n^{-1} \sum_{t=1}^{n} \sum_{j>m}^{\infty} \sum_{k>m}^{\infty} \ell^{\prime} a_{k} a_{l}^{\prime} \ell \alpha_{k, l} \rightarrow 0 \text { as } m \rightarrow \infty .
$$

Applying Lemma (A.1) then gives the result.

Lemma A.5. Assumption (B1) implies that $\tilde{c}(\beta, j)=(2 \pi)^{-1} \int C^{-1}(\beta, \lambda) e^{i \lambda j} d \lambda$ satisfies $\sum j|\tilde{c}(\beta, j)|<\infty$ for all $\beta \in \Theta$.

Proof. Since $C^{-1}(\beta, \pi)=C^{-1}(\beta,-\pi)$ it follows from integration by parts that

$$
\begin{equation*}
|\tilde{c}(\beta, j)|=j^{-1}\left|(2 \pi)^{-1} \int \partial C^{-1}(\beta, \lambda) / \partial \lambda e^{i \lambda j} d \lambda\right| . \tag{17}
\end{equation*}
$$

From $\partial C^{-1}(\beta, \lambda) / \partial \lambda=C^{-2}(\beta, \lambda) \partial C(\beta, \lambda) / \partial \lambda$ and the fact that $C(\beta, \lambda)$ satisfies $\sum j|c(\beta, j)|<$ $\infty$ it follows that $\partial C^{-1}(\beta, \lambda) / \partial \lambda$ has absolutely summable Fourier coefficients. Rearranging (17) and summing over $j$ then gives the result.

Lemma A.6. Assume (A1), (B1-B3), (C1-C2). Let $z_{t}=\lim _{m \rightarrow \infty} A_{m}^{\prime} \varepsilon_{t}^{m}$ with $A_{m}^{\prime}=$ $\left[a_{1}, \ldots, a_{m}\right],\left\{a_{k}\right\}_{k=1}^{\infty} \in \mathcal{A}^{*}$ and $\varepsilon_{t}^{m}=\left[\varepsilon_{t-1}, \ldots, \varepsilon_{t-m}\right]^{\prime}$. Then for any convergent sequence $\beta_{n} \in \Theta$ with $\beta_{n} \rightarrow \beta \in \Theta$ there exists an event $E$ with probability one such that for all outcomes in $E, G_{n}\left(\beta_{n}, a\right) \rightarrow G(\beta, a)$.

Proof. Without loss of generality assume that $z_{t}$ takes values in $\mathbb{R}$. Let $E y_{t} z_{s}=$ $\gamma_{y z}(t-s)$, and $\operatorname{cum}\left(y_{t}, z_{s}, y_{q}, z_{r}\right)=c_{y y z z}(t-s, t-q, t-r)$. Then, from Assumption (A1) and the proof of Theorem 2.8.1 in Brillinger (1981) it follows that $\sum_{j}\left|\gamma_{y z}(j)\right|<\infty$ and $\sum_{s, q, r}\left|c_{y y z z}(s, q, r)\right|<\infty$. For each $\epsilon>0$ there exists an $n_{0}<\infty$ and $\delta>0$ such that $\left\|\beta_{n}-\beta\right\|<\delta$ and

$$
\sup _{\left\|\beta^{\prime}-\beta\right\|<\delta} \sup _{\lambda}\left|C^{-1}\left(\beta^{\prime}, \lambda\right)-C^{-1}(\beta, \lambda)\right|<\epsilon
$$

for $n>n_{0}$ by continuity of $C^{-1}(\beta, \lambda)$ at $\beta \in \Theta$. For $\beta^{\prime}$ such that $\left\|\beta^{\prime}-\beta\right\|<\delta$ the lag polynomial $C^{-1}\left(\beta^{\prime}, z\right)$ has an expansion with coefficients $\tilde{c}\left(\beta^{\prime}, j\right)$ such that $\sum_{j=1}^{\infty} j\left|\tilde{c}\left(\beta^{\prime}, j\right)\right|<$ $\infty$. We will use the short hand notation $\tilde{c}^{\prime}=\tilde{c}\left(\beta^{\prime}, j\right)$. Let $X_{n}(\beta)=G_{n}(\beta, a)-E G_{n}(\beta, a)$ and define $X_{n}=\sup _{\left\|\beta^{\prime}-\beta\right\|<\delta}\left|X_{n}\left(\beta^{\prime}\right)\right|$. Since $E G_{n}\left(\beta^{\prime}, a\right) \rightarrow G\left(\beta^{\prime}, a\right)$ and $\left|G\left(\beta^{\prime}, a\right)-G(\beta, a)\right| \leq$ $\epsilon \int\left|f_{y z}(\lambda)\right| d \lambda$ uniformly for all $\beta^{\prime}$ such that $\left\|\beta^{\prime}-\beta\right\|<\delta$ it is enough to show that $X_{n} \rightarrow 0$
almost surely. Thus letting $X_{n}(j)=\sum_{t=1}^{n-j} y_{t} z_{t+j}-\gamma_{y z}(-j)$

$$
X_{n} \leq \sup _{\left\|\beta^{\prime}-\beta\right\|<\delta} n^{-1} \sum_{j=0}^{n}\left|\tilde{c}_{j}^{\prime}\right|\left|X_{n}(j)\right| \leq K_{0} n^{-1}\left(\sum_{j=0}^{n} j^{-2}\left|X_{n}(j)\right|^{2}\right)^{1 / 2}
$$

where $K_{0}=\sup _{\left\|\beta^{\prime}-\beta\right\|<\delta}\left(\sum_{j=0}^{\infty}\left|\tilde{c}_{j}^{\prime}\right| j\right)$. We consider

$$
E X_{n}^{2} \leq K_{0}^{2} n^{-2} \sum_{j=0}^{n} j^{-2}\left(E X_{n}(j)^{2}\right)
$$

Since

$$
E X_{n}(j)^{2} \leq n \sum_{k=-\infty}^{\infty}\left|\gamma_{y y}(k) \gamma_{z z}(k)+\gamma_{y z}(k) \gamma_{y z}(k)\right|+n \sum_{j, k, l=-\infty}^{\infty}\left|c_{y y z z}(j, k, l)\right|
$$

for all $j$ there is a $K_{1}$ such that $E X_{n}^{2} \leq K_{2} n^{-1}$ where $K_{2}=\frac{\pi^{2}}{6} K_{1} K_{0}^{2}$. For $n / 2 \leq n_{1}<n$ consider $X_{n, n_{1}}=\sup _{\left\|\beta^{\prime}-\beta\right\|<\delta}\left|X_{n}\left(\beta^{\prime}\right)-X_{n_{1}}\left(\beta^{\prime}\right)\right|$ such that

$$
\begin{aligned}
X_{n, n_{1}} \leq & K_{0}\left(n-n_{1}\right)\left(n n_{1}\right)^{-1}\left(\sum_{j=0}^{n_{1}} j^{-2}\left|X_{n_{1}}(j)\right|^{2}\right)^{1 / 2} \\
& +K_{0} n^{-1}\left(\sum_{j=0}^{n} j^{-2}\left(\sum_{t=\max \left(n_{1}-j, 1\right)}^{n-j} y_{t} z_{t+j}-\gamma_{y z}(-j)\right)^{2}\right)^{1 / 2}
\end{aligned}
$$

Now

$$
K_{0}^{2}\left(n-n_{1}\right)^{2}\left(n n_{1}\right)^{-2} E \sum_{j=0}^{n_{1}} j^{-2}\left|X_{n_{1}}(j)\right|^{2} \leq K_{2}\left(n-n_{1}\right) n^{-2}
$$

and

$$
K_{0}^{2} n^{-2} \sum_{j=0}^{n} j^{-2} E\left(\sum_{t=\max \left(n_{1}-j, 1\right)}^{n-j} y_{t} z_{t+j}-\gamma_{y z}(-j)\right)^{2} \leq K_{2} n^{-2}\left(n-n_{1}\right)
$$

together with $E(Y+Z)^{2} \leq E Y^{2}+2\left(E Y^{2} E Z^{2}\right)^{1 / 2}+E Z^{2}$ implies that $E X_{n, n_{1}}^{2} \leq K_{2} n^{-2}(n-$ $n_{1}$ ). It now follows from Lemma 3 in Gaposhkin (1980) that $X_{n} \rightarrow 0$ almost surely. Let this event be $E$. From $\left|G_{n}\left(\beta_{n}, a\right)-G_{n}(\beta, a)\right| \leq X_{n}$ for all $n>n_{0}$ the result follows.

## B. Appendix - Proofs

Proof of Lemma 3.1 From the definition of $\beta_{n}$ it follows that

$$
\begin{equation*}
0 \leq \liminf _{n}\left\|G_{n}\left(\beta_{n}, a\right)\right\|^{2} \leq \limsup _{n}\left\|G_{n}\left(\beta_{n}, a\right)\right\|^{2} \leq \limsup _{n}\left\|G_{n}\left(\beta_{0}, a\right)\right\|^{2} . \tag{18}
\end{equation*}
$$

From Lemma (A.6) it follows that $G_{n}\left(\beta_{0}, a\right) \rightarrow G\left(\beta_{0}, a\right)=0$ almost surely. Thus

$$
\begin{equation*}
\underset{n}{\limsup }\left\|G_{n}\left(\beta_{n}, a\right)\right\|^{2}=\lim _{n}\left\|G_{n}\left(\beta_{n}, a\right)\right\|^{2}=0 \text { almost surely. } \tag{19}
\end{equation*}
$$

Let $E$ be the probability one event in Lemma (A.6). Now consider the sequence $\beta_{n} \in \Theta$. If $\beta_{n}$ does not converge to $\beta_{0}$ then by compactness of $\Theta$ there exists a subsequence $\beta_{n_{k}}$ such that $\beta_{n_{k}} \rightarrow \beta \in \Theta$. By Lemma (A.6) and Assumption (C2) $\liminf _{k}\left\|G_{n_{k}}\left(\beta_{n_{k}}, a\right)\right\|^{2}>0$ a.s. contradicting (19). Therefore $\beta_{n} \rightarrow \beta_{0}$.

Proof of Proposition 3.2 We only prove that Assumption (C2) holds. We first note
that $f_{y z}(\lambda)=C\left(\beta_{0}, e^{-i \lambda}\right) l_{a}(-\lambda)$ where $l_{a}(\lambda)=\sum_{k=1}^{\infty} a_{k} e^{-i \lambda k}$ such that

$$
G(\beta, a)=(2 \pi)^{-1} \int_{-\pi}^{\pi} \psi\left(\beta, e^{-i \lambda}\right) l_{a}(-\lambda) d \lambda
$$

with $\psi\left(\beta, e^{-i \lambda}\right)=C^{-1}\left(\beta, e^{-i \lambda}\right) C\left(\beta_{0}, e^{-i \lambda}\right)$. It is clear that $\psi\left(\beta_{0}, e^{-i \lambda}\right)=1$ so that $G\left(\beta_{0}, a\right)=$ 0.

We need to show that for $C\left(\beta, e^{-i \lambda}\right)=\theta\left(e^{-i \lambda}\right) / \phi\left(e^{-i \lambda}\right)$ there is no other $\beta \in \Theta$ such that $G(\beta, a)=0$. The orthogonality conditions can be written as

$$
\begin{equation*}
(2 \pi)^{-1} \int_{-\pi}^{\pi}\left(\psi\left(\beta, e^{-i \lambda}\right)-1\right) l_{a}(-\lambda) d \lambda=0 . \tag{20}
\end{equation*}
$$

We want to show that the only function $\psi\left(\beta, e^{-i \lambda}\right)-1:[-\pi, \pi] \rightarrow \mathbb{C}$ satisfying this condition is

$$
\begin{equation*}
\psi\left(\beta, e^{-i \lambda}\right)-1 \equiv 0 . \tag{21}
\end{equation*}
$$

If the assumptions of Proposition (3.2) hold then the only value $\beta$ for which $\psi\left(\beta, e^{-i \lambda}\right)-1 \equiv$ 0 is $\beta_{0}$.

Now showing that $\psi\left(\beta, e^{-i \lambda}\right)-1 \equiv 0$ is equivalent to showing $\phi\left(e^{-i \lambda}\right) \theta_{0}\left(e^{-i \lambda}\right) / \phi_{0}\left(e^{-i \lambda}\right)-$ $\theta\left(e^{-i \lambda}\right) \equiv 0$ since the polynomial $\theta\left(e^{-i \lambda}\right)$ is not zero for any $\lambda \in[-\pi, \pi]$ for $\beta \in \Theta$. It is more convenient to rewrite this equation as

$$
\left(\phi\left(e^{-i \lambda}\right)-\phi_{0}\left(e^{-i \lambda}\right)\right) C\left(\beta_{0}, e^{-i \lambda}\right)-\left(\theta\left(e^{-i \lambda}\right)-\theta_{0}\left(e^{-i \lambda}\right)\right) \equiv 0 .
$$

Here $C\left(\beta_{0}, e^{-i \lambda}\right)=\theta_{0}\left(e^{-i \lambda}\right) / \phi_{0}\left(e^{-i \lambda}\right)$ is the lag polynomial of an $\operatorname{ARMA}(p, q)$ with a one
sided Fourier expansion $\sum_{j=0}^{\infty} c_{j} e^{-i \lambda j}$.
For $j \geq \max (p, q+1)-p$ the coefficients $c_{j}$ satisfy the well known restriction

$$
\begin{equation*}
c_{j}-\phi_{0,1} c_{j-1}-\ldots-\phi_{0, p} c_{j-p}=0 . \tag{22}
\end{equation*}
$$

We define the infinite dimensional matrix $C$ with $p$ rows as $C=\left[c,\left[0, c^{\prime}\right]^{\prime},\left[0_{2}, c^{\prime}\right]^{\prime}, \ldots,\left[0_{p-1}, c^{\prime}\right]^{\prime}\right]$ with $c^{\prime}=\left[c_{0}, c_{1}, \ldots\right]$ and $0_{k}$ is the $k$-dimensional column vector of zeros then Condition (20) has a matrix representation

$$
\begin{equation*}
A^{\prime} C\left(\phi-\phi_{0}\right)-A_{q}^{\prime}\left(\theta-\theta_{0}\right)=0 \tag{23}
\end{equation*}
$$

which can be written as $R \delta=0$ where $R$ is the $d \times d$ matrix $R=A^{\prime} D$ where

$$
D=\left[\begin{array}{cc} 
& -I \\
C, & \\
& 0
\end{array}\right]
$$

with 0 an $\infty \times q$ dimensional matrix of zeros and $\delta=\beta-\beta_{0}$. We need to show that ker $R=0$ which follows if $R$ is of full rank. We can distinguish two cases. If $p=0$ then $R=A_{q}^{\prime}$ and $\delta=\left(\theta-\theta_{0}\right)$ such that $\delta=0$ if $A_{q}$ is of full rank. If $p>0$ then $C$ contains $p$ linearly independent vectors in $l^{2}$ which are also linearly independent of $[-I, 0]^{\prime}$. So $D$ has full column rank. It is a finite rank operator mapping $\mathbb{R}^{d}$ into the $d$-dimensional subspace $\operatorname{Im} D$ of $l^{2}$. Since $l^{2}$ is a Hilbert space this subspace is closed and has an orthogonal complement $(\operatorname{Im} D)^{\perp}$ (see Gohberg and Goldberg, p.205). The finite rank operator $A^{\prime}$ maps $l^{2}$ into $\mathbb{R}^{d}$. If $\operatorname{Im} D \cap \operatorname{ker} A^{\prime}=0$ then $\operatorname{Im} D=\left(\operatorname{ker} A^{\prime}\right)^{\perp}$ since $l^{2}=\operatorname{ker} A^{\prime} \oplus\left(\operatorname{ker} A^{\prime}\right)^{\perp}$
where $\oplus$ is the direct sum. Then, by theorem II.11.4 in Gohberg and Goldberg (1980), $\operatorname{Im} A^{\prime}=(\operatorname{ker} A)^{\perp}=\mathbb{R}^{d}$ where the last equality is due to $\operatorname{ker} A=0$ since $A$ is of full column rank. It follows that $\left\{A^{\prime} x \mid x \in \operatorname{Im} D\right\}=\mathbb{R}^{d}$ and $A^{\prime} x=0$ for $x \in \operatorname{Im} D$ if and only if $x=0$. But this means that $\operatorname{Im} R=\mathbb{R}^{d}$ showing that $R$ is of full rank.

Finally, we show that $\operatorname{Im} D$ is the space of all the solutions $x=\left[x_{1}, \ldots\right]$ to $\phi_{0}(L) x=0$ for $x_{j}, j>d$ with $d=p+q$ initial conditions determining $x_{1}, \ldots, x_{d}$. To see this note that $c=\left\{c_{j}\right\}_{j=0}^{\infty}$ is the solution to $\phi_{0}(L) c=0$ for $c_{j}, j>\max (p, q+1)-p$ which has general form $c_{j}=\sum_{i=1}^{k} \sum_{n=0}^{r_{i}} \kappa_{i n} j^{n} \xi_{i}^{-j}$ where $\xi_{i}, i=1, \ldots, k$ are the distinct zeros of $\phi_{0}(L)$ with multiplicity $r_{i}$. The first $\max (p, q+1)-p$ coefficients $c_{j}$ are determined by the $\max (p, q+1)$ boundary conditions implied by $C\left(\beta_{0}, L\right)$. The $l^{2}$ sequence $D \delta$ then has $j$-th element $\tilde{c}_{j}=\sum_{i=1}^{p} \delta_{i} c_{j-i}$ with $c_{j-m}=\sum_{i=1}^{k} \sum_{n=0}^{r_{i}} \sum_{l=0}^{n}\binom{n}{l} \tilde{\kappa}_{i n} j^{l} m^{n-l} \xi_{i}^{-j-m}$ such that the $p$ coefficients $\tilde{\kappa}_{i n}$ of $\tilde{c}_{j}$ can be set arbitrarily to satisfy $p$ initial conditions. The remaining $q$ initial conditions can be set by appropriately choosing $\delta_{p+1}, \ldots, \delta_{d}$.

Remark 16. For finite dimensional matrices it is known from Corollary 6.2 in Marsaglia and Styan (1974) that $\operatorname{rank}\left(A^{\prime} D\right)=\operatorname{rank}(D)-\operatorname{dim}\left(\operatorname{ker} A^{\prime} \cap \operatorname{Im} D\right)$. Our proof extends this result to finite rank operators on Hilbert spaces when $A$ and $D$ are of identical and full column rank.

Proof of Lemma 3.3: Remember that $\int_{-\pi}^{\pi} \dot{\eta}\left(\beta_{0}, \lambda\right) l_{a}(-\lambda) d \lambda=A^{\prime} P$ with $P=\left[b_{1}, b_{2}, \ldots\right]^{\prime}$ where $b_{k}=(2 \pi)^{-1} \int_{-\pi}^{\pi} \partial \ln C\left(\beta_{0}, e^{-i \lambda}\right) / \partial \beta e^{i \lambda k} d \lambda$. For $C\left(\beta_{0}, e^{-i \lambda}\right)=\theta_{0}(\lambda) / \phi_{0}(\lambda)$ we have

$$
\partial \ln C\left(\beta_{0}, e^{-i \lambda}\right) / \partial \beta=\left[\frac{e^{-i \lambda}}{\phi_{0}(\lambda)}, \cdots, \frac{e^{-i \lambda p}}{\phi_{0}(\lambda)}, \frac{e^{-i \lambda}}{\theta_{0}(\lambda)}, \cdots, \frac{e^{-i \lambda q}}{\theta_{0}(\lambda)}\right]^{\prime} .
$$

Define the expansions of $\phi_{0}^{-1}(z)=\sum_{j=0}^{\infty} \psi_{\phi, j} z^{j}$ and $\theta_{0}^{-1}(z)=\sum_{j=0}^{\infty} \psi_{\theta, j} z^{j}$. The coefficients in the expansion satisfy the difference equation $\psi_{\phi, j}-\phi_{0,1} \psi_{\phi, j-1}-\ldots-\phi_{0, p} \psi_{\phi, j-p}=0$ which has $p$ linearly independent solutions. A similar expression holds for $\psi_{\theta, j}$. Set $\psi_{\phi, j}=\psi_{\theta, j}=$ 0 for $j<0$. Then $b_{k}=\left[\psi_{\phi, k-1}, \cdots, \psi_{\phi, k-p}, \psi_{\theta, k-1}, \cdots, \psi_{\theta, k-q}\right]^{\prime}$. Any set of $d=p+q$ vectors $b_{k_{1}}, b_{k_{2}}, \ldots, b_{k_{d}}$ is linearly independent because of the linear independence of the solutions to $\phi_{0}(L) x=0$ and $\theta_{0}(L) x=0$ together with the requirement that $\phi(L)$ and $\theta(L)$ have no common zeros and that $\phi_{p} \neq 0$ or $\theta_{q} \neq 0$. Thus $P$ has full column rank.

When $p>0$ we distinguish two cases. For the case where $q=0, \operatorname{Im} P=\operatorname{Im} D$ where $D$ was defined in the proof of Proposition (3.2). This implies that ker $P^{\prime} \cap \operatorname{Im} D=0$ since $\operatorname{Im} P=\left(\operatorname{ker} P^{\prime}\right)^{\perp}$.

When $q>0$ then $\operatorname{Im} P=S_{1}$ where

$$
S_{1}=\left\{x=\left[x_{1}, \ldots\right] \in l^{2}: \phi_{0}(L) \theta_{0}(L) x=0 \text { for } x_{j}, j>d,\left[x_{1}, \ldots, x_{d}\right]^{\prime}=\kappa, \kappa \in \mathbb{R}^{d}\right\}
$$

while $\operatorname{Im} D=S$. Since $\phi_{0}(L) x=0 \Rightarrow \phi_{0}(L) \theta_{0}(L) x=0$ it follows that $\operatorname{Im} P \supset \operatorname{Im} D$. But $\operatorname{Im} P=\left(\operatorname{ker} P^{\prime}\right)^{\perp}$ and $\operatorname{ker} P^{\prime} \cap\left(\operatorname{ker} P^{\prime}\right)^{\perp}=0$ which implies that $\operatorname{Im} D \cap \operatorname{ker} P^{\prime}=0$. To see that $\operatorname{Im} P=S_{1}$ note that the $j$-th element $c_{j}$ in $\operatorname{Im} P$ is $c_{j}=\sum_{i=1}^{p} \delta_{i} \psi_{\phi, j-i}+$ $\sum_{i=1}^{q} \delta_{p+i} \psi_{\theta, j-1}$ for $\delta \in \mathbb{R}^{d}$. Since $\psi_{\phi, j-m}=\sum_{i=1}^{k} \sum_{n=0}^{r_{i}} \sum_{l=0}^{n}\binom{n}{l} \kappa_{i n} j^{l} m^{n-l} \xi_{\phi, i}^{-j-m}$ from the general solution with a similar expression for $\psi_{\theta, j-m}$ it follows that $\delta_{m} \psi_{\phi, j-m}+\delta_{p+w} \psi_{\theta, j-w}$ is the general solution of a difference equation with roots $\xi_{\phi, i}$ and $\xi_{\theta, i}$ which is the same as $\phi_{0}(L) \theta_{0}(L) x=0$. Since there are $d$ free parameters $\delta_{i}, c_{j}$ can be made to satisfy $d$ initial conditions.

We now show that $\mathcal{A}^{* *} \subset \mathcal{A}^{*}$. First let $p=0$. Then $a \in \mathcal{A}^{* *} \operatorname{implies} \operatorname{rank}\left(A^{\prime} P\right)=q$.

Assume that $\operatorname{rank} A<q$. Then $\operatorname{dim}\left(\operatorname{Im} A^{\prime}\right)<q$ because $\operatorname{ker} A \neq 0$. This contradicts $A^{\prime} P$ to be of full rank. Thus $A_{q}$ is of full rank and $a \in \mathcal{A}^{*}$.

For $p>0$ it follows by the same argument that the row rank of $A^{\prime}$ has to be full. To show that ker $A^{\prime} \cap S=0$ assume that $\operatorname{ker} A^{\prime} \cap S \neq 0$. We have shown that $S=\operatorname{Im} D$ and $\operatorname{Im} D \subseteq \operatorname{Im} P$. This implies $\operatorname{ker} A^{\prime} \cap \operatorname{Im} P \neq 0$. But then $\exists x \in \mathbb{R}^{d}, x \neq 0$ such that $P x \in \operatorname{ker} A^{\prime}$ thus $A^{\prime} P x=0$. This contradicts $A^{\prime} P$ being full rank.

Proof of Theorem 3.4: Let $M_{n}(\beta, a)=G_{n}(\beta, a)-G(\beta, a)$. We use a mean value expansion for $M_{n}(\beta, a)-M_{n}\left(\beta_{0}, a\right)=\frac{\partial}{\partial \beta^{\prime}} M_{n}\left(\beta^{+}, a\right)\left(\beta-\beta_{0}\right)$ with $\left\|\beta^{+}-\beta_{0}\right\| \leq\left\|\beta-\beta_{0}\right\|$. Then

$$
\begin{aligned}
\sup _{\left\|\beta-\beta_{0}\right\|<\delta}\left\|M_{n}(\beta, a)-M\left(\beta_{0}, a\right)\right\| \leq & \delta n^{-1} \sum_{j=0}^{n-1}\left\|\frac{\partial}{\partial \beta} \tilde{c}_{j}^{\beta^{+}}-\frac{\partial}{\partial \beta} \tilde{c}_{j}\right\|\left\|X_{n}(j)\right\| \\
& +\delta \sum_{j=n}^{\infty}\left\|\frac{\partial}{\partial \beta} \tilde{c}_{j}^{\beta^{+}}-\frac{\partial}{\partial \beta} \tilde{c}_{j}\right\|\left\|\gamma_{y z}(-j)\right\|
\end{aligned}
$$

where the second term is $O\left(n^{-1}\right)$. Note that $\frac{\partial}{\partial \beta} C^{-1}(\beta, \lambda)=C^{-2}(\beta, \lambda) \frac{\partial}{\partial \beta} C(\beta, \lambda)$ such that $j\left\|\frac{\partial}{\partial \beta} \tilde{c}_{j}^{\beta}\right\|$ is summable by Assumption (B1). It then follows from arguments similar to the proof of Lemma (A.6) that $E \sup _{\left\|\beta-\beta_{0}\right\|<\delta}\left\|M_{n}(\beta, a)-M\left(\beta_{0}, a\right)\right\| \leq K \delta / \sqrt{n}$ for some constant $K$. For any sequence $\delta_{n} \rightarrow 0$ the Markov inequality then implies that $\sup _{\left\|\beta-\beta_{0}\right\|<\delta_{n}} \sqrt{n}\left\|M_{n}(\beta, a)-M\left(\beta_{0}, a\right)\right\|=o_{p}(1)$. From Theorem 3.2.5 in van der Vaart and Wellner (1996) it follows that $\left\|\beta-\beta_{0}\right\|=O_{p}\left(n^{-1 / 2}\right)$. Following the proof of Theorem 3.3 in Pakes and Pollard (1989) we now consider

$$
\left\|G_{n}\left(\beta_{n}\right)-G_{n}\left(\beta_{0}\right)-\left[\frac{\partial}{\partial \beta_{i}} G_{n}\left(\beta_{n}\right)^{\prime}\right]\left(\beta_{n}-\beta_{0}\right)\right\|
$$

$$
\leq\left\|M_{n}\left(\beta_{n}, a\right)-M\left(\beta_{0}, a\right)\right\|+\left\|G\left(\beta_{n}\right)-\left[\frac{\partial}{\partial \beta_{i}} G_{n}\left(\beta_{n}^{i}\right)^{\prime}\right]\left(\beta_{n}-\beta_{0}\right)\right\|
$$

Pick a sequence $\delta_{n}$ such that $P\left(\left\|\beta_{n}-\beta_{0}\right\| \geq \delta_{n}\right) \rightarrow 0$. Then for any $\epsilon, \eta>0$ there exists an $n$ such that

$$
\begin{aligned}
& P\left(\sqrt{n}\left\|M_{n}\left(\beta_{n}, a\right)-M\left(\beta_{0}, a\right)\right\|\right.>\epsilon) \leq P\left(\sup _{\left\|\beta-\beta_{0}\right\|<\delta_{n}} \sqrt{n}\left\|M_{n}\left(\beta_{n}, a\right)-M\left(\beta_{0}, a\right)\right\|>\epsilon\right) \\
&+P\left(\left\|\beta_{n}-\beta_{0}\right\| \geq \delta_{n}\right) \leq \eta
\end{aligned}
$$

The term $\left\|G\left(\beta_{n}\right)-\left[\frac{\partial}{\partial \beta_{i}} G_{n}\left(\beta_{n}^{i}\right)^{\prime}\right]\left(\beta_{n}-\beta_{0}\right)\right\|=o_{p}\left(n^{-1 / 2}\right)$ by a mean value expansion argument of $G\left(\beta_{n}\right)$ around $\beta_{0}$ and the fact that convergence of $\left[\frac{\partial}{\partial \beta_{i}} G_{n}\left(\beta_{n}^{i}\right)^{\prime}\right]$ can be shown by the same arguments as convergence of $G_{n}(\beta, a)$ noting that both $y_{t}$ and $z_{t}$ are strictly stationary and $\partial \ln C\left(\beta, e^{-i \lambda}\right) / \partial \beta$ is uniformly continuous on $[-\pi, \pi] \times U$ for $U \subset \Theta, U$ compact, $\beta_{0} \in \Theta$. A set $U$ with these properties exists by local compactness of the parameter space. The details are omitted.

We have thus established that

$$
\left[\frac{\partial}{\partial \beta} G_{n}\left(\beta_{n}\right)\right] \sqrt{n} G_{n}\left(\beta_{n}\right)=\left(M+o_{p}(1)\right)\left[\sqrt{n} G_{n}\left(\beta_{0}\right)+\left[\frac{\partial}{\partial \beta_{i}} G_{n}\left(\beta_{n}^{i}\right)^{\prime}\right] \sqrt{n}\left(\beta_{n}-\beta_{0}\right)\right]+o_{p}(1)
$$

where $M=\sigma^{2}(2 \pi)^{-1} \int_{-\pi}^{\pi} \partial \ln C\left(\beta_{0}, e^{-i \lambda}\right) / \partial \beta l_{a}(\lambda) d \lambda$. Next, turn to

$$
\sqrt{n} G_{n}\left(\beta_{0}\right)=\sqrt{n}(2 \pi)^{-1} \int_{-\pi}^{\pi} C^{-1}\left(\beta_{0}, e^{-i \lambda}\right) I_{n, y z}(\lambda) d \lambda
$$

From Lemma (A.3) it follows that $\sqrt{n} \int_{-\pi}^{\pi} C^{-1}\left(\beta_{0}, e^{-i \lambda}\right) I_{n, y z}(\lambda) d \lambda-\sqrt{n} \int_{-\pi}^{\pi} I_{n, \varepsilon z}(\lambda) d \lambda=$
$o_{p}(1)$. Using Lemma (A.4) then shows that $\sqrt{n} G_{n}\left(\beta_{0}\right) \xrightarrow{d} N\left(0, \sum_{l=1}^{\infty} \sum_{k=1}^{\infty} \alpha_{k, l} a_{k} a_{l}^{\prime}\right)$ where it should be noted that $\sum_{l=1}^{\infty} \sum_{k=1}^{\infty} \alpha_{k, l} a_{k} a_{l}^{\prime}=\lim A_{m}^{\prime} \Omega_{m} A_{m}$. The result now follows from $\partial \ln C\left(\beta_{0}, e^{-i \lambda}\right) / \partial \beta=\sum_{k=1}^{\infty} b_{k} e^{-i \lambda k}$ such that $(2 \pi)^{-1} \int_{-\pi}^{\pi} \partial \ln C\left(\beta_{0}, e^{-i \lambda}\right) / \partial \beta l_{a}(\lambda)^{\prime} d \lambda=$ $\sum_{k=1}^{\infty} b_{k} a_{k}^{\prime}$.

Proof of Lemma 4.1 From Assumption (A1) it is clear that $\Omega x \in l^{2}$ for all $x \in l^{2}$. It remains to show that $\operatorname{ker} \Omega=0$. Assume there is $x \in l^{2}$ such that $x \neq\{0,0, \ldots\}$ and $\Omega x=0$. Then also $x^{\prime} \Omega x=0$ which can be written as $E\left(\sum_{i=1}^{\infty} x_{i} \varepsilon_{t} \varepsilon_{t-i}\right)^{2}=0$. But this is only possible if $\sum x_{i} \varepsilon_{t} \varepsilon_{t-i}=0$ with probability one. Now $\sum x_{i} \varepsilon_{t} \varepsilon_{t-i}=0$ a.s. if $\varepsilon_{t} \varepsilon_{t-i}=0$ a.s. or the functions $\varepsilon_{t-i}$ are linearly dependent a.s. which is ruled out by Lemma (A.2).

On the other hand if $\varepsilon_{t} \varepsilon_{t-i}=0$ a.s. for all $i$ then $\varepsilon_{t}^{2} \varepsilon_{t-i}^{2}=0$ a.s. But then $E\left(\varepsilon_{t}^{2} \varepsilon_{t-i}^{2}\right)=0$ for all $i$ which contradicts Assumption (A1). Therefore $\Omega x=0$ can only hold if $x=0$. Thus ker $\Omega=0$. Symmetry of $\Omega$ now implies that $\operatorname{Im} \Omega=l^{2}$ therefore $\Omega^{-1}$ exists and is bounded on $l^{2}$.

Finally, it follows at once from before that $x_{m}^{\prime} \Omega_{m} x_{m}=E\left(\sum x_{i, m} \varepsilon_{t} \varepsilon_{t-i}\right)^{2}>0$ where the inequality is strict by Assumption (A1). So $\Omega_{m}$ is positive definite such that $\lambda_{j}^{m}>0$ $\forall j, m$. This shows that $\Omega_{m}$ has full rank.

Proof of Lemma 4.2 By Assumption (A1) we know that $\sum \sum|\sigma(k, l)|<B$ thus $\sum_{k}|\sigma(k, l)|<B$ for any $l$. Therefore for any fixed $l, \sigma(k, l) \rightarrow 0$ as $k \rightarrow \infty$. This holds also if the roles of $k$ and $l$ are reversed. Also $\sum_{k}|\sigma(k, k)|<B$ such that $\sigma(k, k) \rightarrow 0$ as $k \rightarrow \infty$. Define the infinite dimensional matrices $S_{12}^{m}, S_{21}^{m}$ and $S_{22}^{m}$ according to the
following partition

$$
\Omega=\left[\begin{array}{ll}
\Omega_{m} & S_{12}^{m} \\
S_{21}^{m} & S_{22}^{m}
\end{array}\right]
$$

Then $\operatorname{tr}\left(S_{12}^{m} S_{12}^{m^{\prime}}\right)=\sum_{l=m+1}^{\infty} \sum_{k=1}^{m}|\sigma(k, l)|^{2} \rightarrow 0, \operatorname{tr}\left(S_{21}^{m} S_{21}^{m^{\prime}}\right) \rightarrow 0$ and $\operatorname{tr}\left(S_{22}^{m}-\sigma^{4} I\right)\left(S_{22}^{m}-\right.$ $\left.\sigma^{4} I\right)^{\prime} \rightarrow 0$ as $m \rightarrow \infty$. Define the infinite dimensional approximation matrix

$$
\Omega_{m}^{*}=\left[\begin{array}{ll}
\Omega_{m} & 0 \\
0 & \sigma^{4} I
\end{array}\right]
$$

Clearly $\Omega_{m}^{*-1}$ exists $\forall m$ by Lemma (4.1) and the partitioned inverse formula. We now have

$$
\left(\Omega^{-1}-\Omega_{m}^{*-1}\right)=\Omega_{m}^{*-1}\left(\Omega-\Omega_{m}^{*}\right) \Omega^{-1}
$$

such that

$$
\left\|\Omega^{-1}-\Omega_{m}^{*-1}\right\| \leq\left\|\Omega_{m}^{*-1}\right\|\left\|\Omega-\Omega_{m}^{*}\right\|\left\|\Omega^{-1}\right\| .
$$

where $\|$.$\| is the matrix norm defined by \|A\|=\sup _{\|x\|_{2} \leq 1}\|A x\|_{2}$. First show that $\left\|\Omega_{m}^{*-1}\right\|$ is bounded. By the partitioned inverse formula

$$
\Omega_{m}^{*-1}=\left[\begin{array}{ll}
\Omega_{m}^{-1} & 0 \\
0 & \sigma^{-4} I
\end{array}\right]
$$

such that $\left\|\Omega_{m}^{*-1}\right\| \leq\left\|\Omega_{m}^{-1}\right\|+\sigma^{-4}$. We have shown in Lemma (4.1) that $0<\min _{x^{\prime} x=1} x^{\prime} \Omega x$ which together with $\min _{x^{\prime} x=1} x^{\prime} \Omega x \leq \min _{x_{m}^{\prime} x_{m}=1} x_{m}^{\prime} \Omega_{m} x_{m} \forall m$ implies that the smallest eigenvalue $\lambda_{1}^{m}$ of $\Omega_{m}$ is bounded away from zero uniformly in $m$. Then by a familiar
inequality for all $x \in \mathbb{R}^{m} x^{\prime} \Omega_{m}^{-1} x / x^{\prime} x \leq 1 / \lambda_{1}^{m}<\infty \forall m$ such that $\left\|\Omega_{m}^{-1}\right\|<\infty$ since for finite $m$ all norms are equivalent. Also $\left\|\Omega^{-1}\right\|<\infty$ by Lemma (4.1) and

$$
\begin{aligned}
\left\|\Omega-\Omega_{m}^{*}\right\| & =\sup _{\|x\| \leq 1}\left(\sum_{k=1}^{m}\left|\sum_{l=m+1}^{\infty} \sigma(k, l) x_{l}\right|^{2}+\sum_{k=m+1}^{\infty}\left|\sum_{l=1}^{\infty} \sigma(k, l) x_{l}\right|^{2}\right)^{1 / 2} \\
& \leq \sup _{\|x\| \leq 1} \sum_{k=1}^{m} \sum_{l=m+1}^{\infty}|\sigma(k, l)|\left|x_{l}\right|+\sup _{\|x\| \leq 1} \sum_{k=m+1}^{\infty} \sum_{l=1}^{\infty}|\sigma(k, l)|\left|x_{l}\right| \\
& \leq 2 \sum_{l=m+1}^{\infty} \sum_{k=1}^{\infty}|\sigma(k, l)| \rightarrow 0 \text { as } m \rightarrow \infty
\end{aligned}
$$

Thus $\left\|\Omega^{-1}-\Omega_{m}^{*-1}\right\| \rightarrow 0$ as $m \rightarrow \infty$

Proof of Theorem 4.3 For all $m$ fixed it follows from standard results that

$$
\left(P_{m}^{\prime} A_{m}\right)^{-1}\left(A_{m}^{\prime} \Omega_{m} A_{m}\right)\left(A_{m}^{\prime} P_{m}\right)^{-1}-\left(P_{m}^{\prime} \Omega_{m}^{-1} P_{m}\right)^{-1} \geq 0
$$

But since for any sequence $\left\{x_{m}\right\}$ such that $x_{m} \geq 0$ for all $m$ it follows that $\liminf _{m} x_{m} \geq 0$ the above inequality also holds in the limit. Since both $b(P) \in l^{1}$ and $a(A) \in l^{1}$ it follows from a bounded convergence argument that $\lim _{m} P_{m}^{\prime} A_{m}$ exists and is finite. If $a \in \mathcal{A}^{* *}$ then the inverse exists as well. The same arguments can be used to show that $\lim _{m} A_{m}^{\prime} \Omega_{m} A_{m}$ exists and is finite.

Finally note that

$$
\lim _{m} P_{m}^{\prime} \Omega_{m}^{-1} P_{m}=P^{\prime} \Omega^{-1} P
$$

by Lemma (4.2) since $P_{m} \in l^{2}$.

Proof of Theorem 5.1 We first show that $a(A) \in \mathcal{A}$ for $A^{\prime}=P^{\prime} \Omega^{-1}$. From Assumption
(A1) it follows that $\Omega$ maps $l^{1}$ into $l^{1}$. To see this write $\Omega=\Sigma+\sigma^{4} I$ where the matrix $\Sigma$ consists of elements $\sigma(k, l)$. For $x \in l^{1}$ we have $\Omega x=\Sigma x+\sigma^{4} x$ with $\Sigma x \in l^{1}$ because of the summability restrictions on $\sigma(k, l)$. From Lemma 4.1 we know that for $x \in l^{1} \subset l^{2}$ we have $\Omega^{-1} x \in l^{2}$. Assume $\Omega^{-1} x \notin l^{1}$. Then $x=\Omega \Omega^{-1} x=\Sigma \Omega^{-1} x+\sigma^{4} \Omega^{-1} x$. But $\Sigma \Omega^{-1} x \in l^{1}$. Thus $\left\|\sigma^{4} \Omega^{-1} x\right\|_{1}=\left\|x-\Sigma \Omega^{-1} x\right\|_{1} \leq\|x\|_{1}+\left\|\Sigma \Omega^{-1} x\right\|_{1}$ and $\|x\|_{1}$ becomes unbounded because of $\left\|\sigma^{4} \Omega^{-1} x\right\|_{1}$. But this contradicts the assumption that $x \in l^{1}$. It follows that the image of $l^{1}$ under $\Omega^{-1}$ is also in $l^{1}$ which in turn implies that $\sum_{k=1}^{\infty}\left|\omega_{l k}\right|<\infty$ for all $l$. This can be seen by considering the image under $\Omega^{-1}$ of the $l$-th unit vector. Since $P \in \mathcal{A}$ it now follows that $P^{\prime} \Omega^{-1} \in \mathcal{A}$.

In light of Lemma (3.3) we only need to show that $a(A) \in \mathcal{A}^{* *}$ which implies $a(A) \in \mathcal{A}^{*}$. The optimal instrument is defined by $A^{\prime}=P^{\prime} \Omega^{-1}$ or $A^{\prime} \Omega=P^{\prime}$. The row rank of $A^{\prime}$ is therefore the same as the column rank of $P$ which has full column rank, thus establishing that $A_{d}=\left[a_{1}, \ldots, a_{d}\right]$ is nonsingular.

Next, $\int \dot{\eta}(\beta, \lambda) l_{a}(-\lambda)^{\prime} d \lambda=P^{\prime} \Omega^{-1} P$ and $P^{\prime} \Omega^{-1} P=E\left(\varepsilon_{t}^{2} z_{t} z_{t}^{\prime}\right)$. Now, $\operatorname{det} P^{\prime} \Omega^{-1} P=$ $0 \Rightarrow \exists \ell \in \mathbb{R}^{d}, \ell \neq 0$ such that $\ell^{\prime} E\left(\varepsilon_{t}^{2} z_{t} z_{t}^{\prime}\right)=0 \Rightarrow \ell^{\prime} E\left(\varepsilon_{t}^{2} z_{t} z_{t}^{\prime}\right) \ell=0$. Then for $x_{t}:=\ell^{\prime} z_{t}$, $0=E\left(\varepsilon_{t}^{2} x_{t}^{2}\right)=E x_{t}^{2} E\left[\varepsilon_{t}^{2} \mid \mathcal{F}_{t-1}\right] \Rightarrow E\left[\varepsilon_{t}^{2} \mid \mathcal{F}_{t-1}\right]=0$ a.s. or $x_{t}^{2}=0$ a.s. Now, clearly $E\left[\varepsilon_{t}^{2} \mid \mathcal{F}_{t-1}\right]=0$ a.s. is ruled out by Assumption $(A 1)$. Then $x_{t}^{2}=0$ a.s. implies $x_{t}=$ $\ell^{\prime} z_{t}=0$ a.s. From Lemma (A.2) it follows that $z_{t}=0$ a.s. is impossible and we have shown before that the column rank of $A$ is full so that $\ell^{\prime} z_{t}=0$ a.s. is also impossible.

Proof of Theorem 5.3: We need to show that $\sum_{k=1}^{\infty}|k|\left|\left[a_{k}\right]_{j}\right|$ for $j=1, \ldots, d$ is bounded. Since $P \in \mathcal{A}$ we can write $P^{\prime}=P^{\prime} \Omega^{-1} \Omega=P^{\prime} \Omega^{-1}\left(\sigma^{4} I+\Sigma\right)$. Define the vector $\ell_{k}=k e_{k}$
where $e_{k}$ is the $k$-th unit vector. Then

$$
\begin{equation*}
P^{\prime} \ell_{k}=P^{\prime} \Omega^{-1}\left(\sigma^{4} I+\Sigma\right) \ell_{k} . \tag{24}
\end{equation*}
$$

Now, the sequence $\left\{P^{\prime} \ell_{k}\right\}_{k=1}^{\infty} \in \mathcal{A}$ and $\Sigma \ell_{k} \in l^{1}$ for all $k$. Therefore, by the fact that $a\left(P^{\prime} \Omega^{-1}\right) \in \mathcal{A}$ and by the summability assumption of Lemma (5.3), $\left\{P^{\prime} \Omega^{-1} \Sigma \ell_{k}\right\}_{k=1}^{\infty} \in \mathcal{A}$. From (24) we have

$$
\begin{aligned}
\left|P^{\prime} \Omega^{-1} \ell_{k} \sigma^{4}\right| & =\left|P^{\prime} \ell_{k}-P^{\prime} \Omega^{-1} \Sigma \ell_{k}\right| \\
& \leq\left|P^{\prime} \ell_{k}\right|+\left|P^{\prime} \Omega^{-1} \Sigma \ell_{k}\right| .
\end{aligned}
$$

where $|$.$| is a vector norm on \mathbb{R}^{d}$. Without loss of generality we use $|x|=\sup _{i}\left|x_{i}\right|$ for $x \in \mathbb{R}^{d}$. Summing over $k$ gives $\sigma^{4} \sum_{k=1}^{\infty}\left|P^{\prime} \Omega^{-1} \ell_{k}\right| \leq \sum_{k=1}^{\infty}\left(\left|P^{\prime} \ell_{k}\right|+\left|P^{\prime} \Omega^{-1} \Sigma \ell_{k}\right|\right)<\infty$. Note that $\left|P^{\prime} \Omega^{-1} \ell_{k}\right|=k\left|\sum_{l=1}^{\infty} b_{l} \omega_{l k}\right|$. This establishes the result

Proof of Corollary 5.4: We first need to establish consistency. For this we show that uniformly in $\Theta,\left\|\tilde{G}_{n}(\beta, a)-G_{n}(\beta, a)\right\| \xrightarrow{p} 0$. Note that $z_{t}-\hat{z}_{t}=\sum_{j=1}^{t-1} a_{j}\left(\varepsilon_{t-j}-\hat{\varepsilon}_{t-j}\right)+$ $\sum_{j=t}^{\infty} a_{j} \varepsilon_{t-j}$ and $\varepsilon_{t-j}-\hat{\varepsilon}_{t-j}=\sum_{l=t-j}^{\infty} c\left(\beta_{0}, l\right) y_{t-j-l}$. Without loss of generality assume $a_{j} \in \mathbb{R}$. Then

$$
\begin{aligned}
\left|\tilde{G}_{n}(\beta, a)-G_{n}(\beta, a)\right| \leq & n^{-1}\left|\sum_{t=1}^{n} \sum_{j=1}^{t-1} a_{j} \sum_{l=t-j}^{\infty} c\left(\beta_{0}, l\right) y_{t-j-l} \sum_{r=0}^{t-1} \tilde{c}_{r}^{\beta} y_{t-r}\right| \\
& +n^{-1}\left|\sum_{t=1}^{n} \sum_{j=t}^{\infty} a_{j} \varepsilon_{t-j} \sum_{r=0}^{t-1} \tilde{c}_{r}^{\beta} y_{t-r}\right|
\end{aligned}
$$

where, using $t \leq 2|j||l|$ on the relevant range of summation, the first term can be bounded
by

$$
\begin{aligned}
& 2 n^{-1} \sum_{t=1}^{n} \sum_{j=1}^{t-1} j\left|a_{j}\right| \sum_{l=t-j}^{\infty} t^{-1} l\left|c\left(\beta_{0}, l\right)\right| \sum_{r=0}^{t-1}\left|\tilde{c}_{r}^{\beta}\right|\left|y_{t-j-l} y_{t-r}-\gamma_{y y}(j+l-r)\right| \\
& +2 n^{-1} \sum_{t=1}^{n} \sum_{j=1}^{t-1} j\left|a_{j}\right| \sum_{l=t-j}^{\infty} t^{-1} l\left|c\left(\beta_{0}, l\right)\right| \sum_{r=0}^{t-1}\left|\tilde{c}_{r}^{\beta}\right|\left|\gamma_{y y}(j+l-r)\right| .
\end{aligned}
$$

Note that $E\left|y_{t-j-l} y_{t-r}-\gamma_{y y}(j+l-r)\right|$ is uniformly bounded in $t, j, l, r$ and the remaining terms in the expression are summable in $\Theta$. We can therefore bound the sum by

$$
K n^{-1} \sum_{t=1}^{n} t^{-1}<K n^{-1+\nu} \sum_{t=1}^{n} t^{-(1+\nu)}=O\left(n^{-1+\nu}\right)
$$

for $\nu \in(0,1 / 2)$ and some constant $K$. Thus by the Markov inequality the first term is $O_{p}\left(n^{-1}\right)$. The second term can be bounded by
$n^{-1} \sum_{t=1}^{n} t^{-1} \sum_{j=t}^{\infty} j\left|a_{j}\right| \sum_{r=0}^{t-1}\left|\tilde{c}_{r}^{\beta}\right|\left|y_{t-r} \varepsilon_{t-j}-\gamma_{y \varepsilon}(r-j)\right|+n^{-1} \sum_{t=1}^{n} t^{-1} \sum_{j=t}^{\infty} j\left|a_{j}\right| \sum_{r=0}^{t-1}\left|\tilde{c}_{r}^{\beta}\right|\left|\gamma_{y \varepsilon}(r-j)\right|$
where again $E\left|y_{t-r} \varepsilon_{t-j}-\gamma_{y \varepsilon}(r-j)\right|$ is uniformly bounded and $n^{-1} \sum_{t=1}^{n} t^{-1}=O\left(n^{-1+\nu}\right)$ for $\nu \in(0,1 / 2)$. This establishes $\sup _{\beta \in \Theta}\left|\tilde{G}_{n}(\beta, a)-G_{n}(\beta, a)\right| \xrightarrow{p} 0$ and thus $\tilde{\beta} \xrightarrow{p} \beta_{0}$. Next we show that $\sup _{\beta \in \Theta}\left|\tilde{G}_{n}(\beta, a)-\tilde{G}_{n}^{F}(\beta, a)\right|=O_{p}\left(n^{-1}\right)$. Note that

$$
\begin{aligned}
\tilde{G}_{n}^{F}(\beta, a) & =\frac{1}{n} \sum_{j=0}^{\infty} \sum_{k=1}^{\infty} \sum_{l=0}^{\infty} c_{l} \tilde{c}_{j}^{\beta} a_{k} \sum_{t=\max (k+l, j)}^{\min (n+k+l, n+j)} y_{t-j} y_{t-k-l} \\
& =G_{n}(\beta, a)+\frac{1}{n} \sum_{j=0}^{\infty} \sum_{k=1}^{\infty} \sum_{l=0}^{\infty} c_{l} \tilde{c}_{j}^{\beta} a_{k}\{\max (j, k+l)<n\}\left(\sum_{t=n}^{\min (n+k+l, n+j)} y_{t-j} y_{t-k-l}\right)
\end{aligned}
$$

$$
+\frac{1}{n} \sum_{j=0}^{\infty} \sum_{k=1}^{\infty} \sum_{l=0}^{\infty} c_{l} \tilde{c}_{j}^{\beta} a_{k}\{\max (j, k+l) \geq n\}\left(\sum_{t=\max (k+l, j)+1}^{\min (n+k+l, n+j)} y_{t-j} y_{t-k-l}\right)
$$

where the second term is uniformly bounded in expectation by

$$
\frac{1}{n} \sup _{t, s} E\left|y_{t} y_{s}\right| \sup _{\beta} \sum_{j=0}^{n} \sum_{k=1}^{n} \sum_{l=0}^{n} j l k\left|c_{l}\right|\left|\tilde{c}_{j}^{\beta}\right|\left|a_{k}\right|=O\left(n^{-1}\right)
$$

where we can exchange sup and $E$ by a dominated convergence argument since $j\left|\tilde{c}_{j}^{\beta}\right|$ is uniformly summable on $\Theta$. The third term can be uniformly bounded in expectation by

$$
\begin{aligned}
& \frac{1}{n} \sup _{t, s} E\left|y_{t} y_{s}\right| \sup _{\beta} \sum_{k=1}^{\infty} \sum_{l=0}^{\infty} \sum_{j=0}^{k+l}(k+l)\left|c_{l}\right|\left|\tilde{c}_{j}^{\beta}\right|\left|a_{k}\right|\{k+l \geq n\} \\
& +\frac{1}{n} \sup _{t, s} E\left|y_{t} y_{s}\right| \sup _{\beta} \sum_{k=1}^{\infty} \sum_{l=0}^{\infty} \sum_{j=n}^{\infty} j\left|c_{l}\right|\left|\tilde{c}_{j}^{\beta}\right|\left|a_{k}\right|
\end{aligned}
$$

which is also $O\left(n^{-1}\right)$. It then follows that $E \sup _{\beta}\left|\tilde{G}_{n}(\beta, a)-\tilde{G}_{n}^{F}(\beta, a)\right|=O\left(n^{-1}\right)$. The result follows by the Markov inequality. For the limiting distribution expand (12) as

$$
\sqrt{n}\left(\tilde{\beta}-\beta_{0}\right)=\left[\partial \tilde{G}_{n}^{F}\left(\beta^{+}, a\right) / \partial \beta\right]^{-1} \sqrt{n} \tilde{G}_{n}^{F}\left(\beta_{0}, a\right)+o_{p}(1)
$$

where $\left\|\beta^{+}-\beta_{0}\right\| \leq\left\|\tilde{\beta}-\beta_{0}\right\|$ where $\beta^{+}$varies across different rows in $\partial \tilde{G}_{n}^{F}\left(\beta^{+}, a\right)$ by the mean value theorem. From Lemma A. 2 and A. 3 in Kuersteiner (1999) it follows that $\sqrt{n} \tilde{G}_{n}^{F}\left(\beta_{0}, a\right) \xrightarrow{d} N\left(0, P^{\prime} \Omega^{-1} P\right)$ and by standard arguments

$$
\partial \tilde{G}_{n}^{F}\left(\beta^{+}, a\right) / \partial \beta \xrightarrow{p} \int \dot{\eta}\left(\beta_{0}, \lambda\right) l_{\psi}(-\lambda)^{\prime} d \lambda=P^{\prime} \Omega^{-1} P .
$$

Proof of Theorem 5.5: For consistency we establish that $\sup _{\beta \in \Theta}\left|\tilde{G}_{n}^{F}(\beta, a)-\tilde{G}_{n}^{F}(\beta, \hat{a})\right| \xrightarrow{p}$ 0. Here
$\sup _{\beta \in \Theta}\left|\tilde{G}_{n}^{F}(\beta, a)-\tilde{G}_{n}^{F}(\beta, \hat{a})\right| \leq \frac{1}{2 \pi} \sup _{\substack{\beta \in \Theta \\ \lambda \in[-\pi, \pi]}}\left|C^{-1}\left(\beta, e^{-i \lambda}\right)\right| \sup _{\lambda}\left|h\left(\beta_{0}, \lambda\right)-\hat{h}(\hat{\beta}, \lambda)\right| \int_{-\pi}^{\pi} I_{n, y y}(\lambda) d \lambda$
where $\sup _{\lambda}\left|h\left(\beta_{0}, \lambda\right)-h(\hat{\beta}, \lambda)\right|=o_{p}(1)$ by Theorem 5.1 in Kuersteiner (1999) such that consistency follows by standard arguments. The same arguments also lead to

$$
\sup _{\beta \in \Theta}\left|\frac{\partial}{\partial \beta} \tilde{G}_{n}^{F}(\beta, a)-\frac{\partial}{\partial \beta} \tilde{G}_{n}^{F}(\beta, \hat{a})\right| \xrightarrow{p} 0
$$

The limit theory is established by showing that $\sqrt{n}\left(\tilde{G}_{n}^{F}\left(\beta_{0}, a\right)-\tilde{G}_{n}^{F}\left(\beta_{0}, \hat{a}\right)\right)=o_{p}(1)$. The proof is essentially identical to the proof of Theorem 5.2 in Kuersteiner (1999) and is omitted.

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[^1]:    ${ }^{1}$ This paper is partly based on results in my Ph.D. dissertation at Yale University, 1997. I wish to thank Peter C.B. Phillips for continued encouragement and support. I have also benefited from comments by Donald Andrews, Jinyong Hahn, Jerry Hausman, Whitney Newey, Oliver Linton, Chris Sims and participants of the econometrics workshop at Yale and the University of Pennsylvania. I have also received very helpful comments from an associate editor and two referees that lead to a substantial improvement of the paper. All remaining errors are my own. Financial support from an Alfred P. Sloan Doctoral Dissertation Fellowship is gratefully acknowledged.

[^2]:    ${ }^{3} A$ collection of open subsets $\mathcal{B}$ of a space $X$ is called a base if for each open set $O \subset X$ and each $x \in O$ there is a set $B \in \mathcal{B}$ such that $x \in B \subset O$. A collection $\mathcal{B}_{x}$ of open sets containing a point $x$ is called a local base at $x$ if for each open set $O$ containing x there is a $B \in \mathcal{B}_{x}$ such that $x \in B \subset O$. Every metric space has a countable base at each point (see Royden (1988), p. 175).

