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# MAXIMUM LIKELIHOOD INFERENCE IN WEAKLY IDENTIFIED DSGE MODELS 

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# Maximum Likelihood Inference in <br> Weakly Identified DSGE Models. <br> By Isaiah Andrews ${ }^{1}$ and Anna Mikusheva ${ }^{2} 3$ 


#### Abstract

This paper examines the problem of weak identification in maximum likelihood, motivated by problems with estimation and inference a multi-dimensional, non-linear DSGE model. We suggest a test for a simple hypothesis concerning the full parameter vector which is robust to weak identification. We also suggest a test for a composite hypothesis regarding a sub-vector of parameters. The suggested test is shown to be asymptotically exact when the nuisance parameter is strongly identified, and in some cases when the nuisance parameter is weakly identified. We pay particular attention to the question of how to estimate Fisher's information, and make extensive use of martingale theory.

Key words: weak identification, maximum likelihood, score test, $C(\alpha)$ - test


## 1 Introduction

Recent years have witnessed the rapid growth of the empirical literature on the highly parameterized micro-founded macro models known as Dynamic Stochastic General Equilibrium (DSGE) models. A number of papers in this literature have considered estimating these models by maximum likelihood (see for example Ingram, Kocherlakota and Savin (1994), Ireland (2004), Lindé (2005), and McGrattan, Rogerson and Wright(1997)). More recently, Bayesian estimation has become increasing popular, in large part due to the difficulty of maximum likelihood estimation in many DSGE models. As Fernéndez-Villaverde (2010) points out in his survey of DSGE estimation, "likelihoods of DSGE models are full of local maxima and minima and of nearly flat surfaces... the standard errors of the estimates are notoriously difficult to compute and their asymptotic distribution a poor approximation to the small sample one." The bad behavior of maximum likelihood estimation has fueled growing concerns about poor identification in many DSGE models

[^0](see Canova and Sala (2009), Guerron-Quintana, Inoue and Kilian (2009), and Iskrev (2009)).

In this paper, we consider the problem of weak identification in dynamic models estimated by maximum likelihood. Weak identification arises when the amount of information in the data about some parameter or group of parameters is small, and is generally modeled in such a way that information about parameters accumulates slowly along some dimensions. This leads to the breakdown of the usual asymptotics for maximum likelihood, with the asymptotic distributions for the maximum likelihood estimator and the standard LR, LM, and Wald statistics providing a poor approximation to their finite sample behavior. This is distinct from loss of point identification, and we assume throughout that the models we consider are point identified, and thus that changing the value of any parameter changes the distribution of the data, though the effect will be small for some parameters.

We focus on the problem of testing and confidence set construction in this context. In our view there are two main approaches to inference in models where identification may be weak. One is to create a two-step procedure, where one first differentiates (via a pre-test) between weakly and strongly identified models and then chooses a procedure based on the test result. We take the other approach. Rather than looking for a test for weak identification as such, we instead attempt to construct a test for parameters which is robust to weak identification. The ideal procedure should satisfy two conditions. First, it should control size well if identification is weak, and second, it should be asymptotically equivalent to the classical MLE tests if identification is strong. If such a procedure exists, it renders pretests unnecessary and, in general, inferior given the size problems endemic to multiple testing procedures.

We view this approach as analogous to the modern treatment of testing in the presence of potential heteroskedasticity. While in the past it was common to use pretests for hetroskedatsicity, current empirical practice is to simply use standard errors (such as those of White (1980)) which are correct asymptotically regardless of whether or not the data is heteroskedastic. Likewise, in weak instrumental variable regression (weak IV) there are tests available which have correct asymptotic size under the weak identification and (at least for the case of one endogenous variable) at least as much power as classical procedures under strong identification. Unlike the case of hetereoskedasticity, where weighted least squares could potentially improve precision, in weak IV the outcome of a
pretest cannot be used to increase power, so there is even less reason to use a pretest-based procedure.

We construct a robust test in two steps. First, we suggest a test for a simple hypothesis on the full parameter vector. This test is robust to weak identification and is asymptotically equivalent to the classical Lagrange Multiplier (LM) test when identification is strong. The assumptions needed for this result are extremely weak, and cover a large number of cases, including weak IV and weakly identified exponential family models, for example VARs with weakly identified structural parameters. The proofs for this test make extensive use of martingale theory, particularly the fact that the score (i.e. the gradient of the log likelihood) is a martingale when evaluated at the true parameter value.

Next, we turn to the problem of testing a subset of parameters without restricting the remaining parameters. Creation of such tests is critical for the construction of confidence sets, given that the common practice in applied work is to report a separate confidence interval for each element of the parameter vector. Constructing a test satisfying our first requirement, that is one which controls size well under weak identification, is straightforward using our test for the full parameter vector and the projection method. However, simultaneously satisfying the second condition, asymptotic equivalence to classical tests under strong identification, is a much more challenging problem which (to the best of our knowledge) has not been fully solved even for many simpler models.

The test which we suggest for a subset of parameters is asymptotically equivalent to Neyman's $C(\alpha)$ test when identification is strong. We show that the suggested test has the $\chi^{2}$ asymptotic distribution so long as the nuisance parameter (i.e. the part of the parameter vector which we are not testing) is strongly identified, without any assumption about the identification of the tested parameter. We also show that the suggested test has correct asymptotic size in some cases where the nuisance parameter is weakly identified. In particular we consider the case of an exponential family model where part of the nuisance parameter is weakly identified and enters linearly while no assumption is made on the identification of the tested parameter. As a special case we examine weak IV with one endogenous variable when the nuisance parameter is weakly identified.

In addition to these theoretical results, we report simulation results showing that our proposed test maintains size well in a simple nonlinear model and is conservative in weak IV with more than one endogenous variable. We also show the applicability of our results
to a basic DSGE model.

Relation to the Literature on Weak Identification The literature on weak identification is quite large. The most-studied and best-understood case is that of weak instrumental variables estimation. For a comprehensive survey of the literature on this topic, see Stock, Wright, and Yogo (2002). The weak identification framework was generalized to GMM by Stock and Wright (2000), who represented weak identification using an asymptotic embedding in which the objective function becomes flat along some dimensions as the sample grows. While we make use of a similar embedding to demonstrate the applicability of our assumptions, it is not necessary for our results, and we remain quite agnostic about the process generating the data. An alternative embedding for weak identification is introduced in Andrews and Chen (2009).

Making use of their embedding, Stock and Wright (2000) introduce tests for GMM which are robust to weak identification. They consider two types of test: a test for the full parameter vector (i.e. for a simple hypothesis) and a test for a sub-parameter for the case where the nuisance parameter is well-identified. Kleibergen and Mavroeidis (2009) suggest adaptations of the Stock and Wright (2000) S and Kleibergen (2005) KLM tests for a sub-parameter for the case when the nuisance parameter is weakly identified, which yield conservative tests asymptotically. While the statistics we consider are in many ways similar to those considered by Stock and Wright (2000), Kleibergen (2005), and Kleibergen and Mavroeidis (2009), their results do not in general apply to the context we consider as the variance of the moment condition (the score of the log likelihood) becomes degenerate asymptotically, violating one of their assumptions.

The issue of weak identification in DSGE models was first introduced by Canova and Sala (2009), who pointed out that the objective functions implied by many DSGE models are nearly flat in some directions. A weak identification-robust inference procedure for DSGE models based on likelihood analysis was introduced by Guerron-Quintana Inoue and Killian (2009). Their method makes extensive use of projection for constructing confidence sets which, given the high dimension of the parameter space in many DSGE models, has the potential to introduce a substantial amount of conservativeness in many applications. Another paper on weak identification in DSGE models is Iskrev (2008), which attempts to asses the quality of identification in DSGE models by considering the degeneracy of the Hessian of the log likelihood. There are also a few papers discussing
point-identification in DSGE models, which are unrelated to our paper as we assume point-identification. We refer the interested reader to Komunjer and Ng (2009) for an example of this literature.

Relation to the Classical MLE Literature The other major literature to which our paper is connected is the classical Statistics literature on maximum likelihood. This classical literature began in the i.i.d. context and was generalized considerably by Le Cam (see Le Cam and Yang (2000)), allowing the use of MLE in a wide array of problems, including with dependent data. The application of ML to dependent data was further explored by a number of other authors, including Silvey (1961), Crowder (1976), Heijmans and Magnus (1986) and Jeganathan (1995). Our approach is particularly informed by the strand of this literature which focuses on the martingale properties of the log likelihood and their implications for the asymptotics of the MLE, and especially by Bhat (1974) and Hall and Heyde (1980).

The weakly identified dynamic models which we consider differ from those in this classical literature in that the normalized second derivative of the log likelihood may not converge to a constant (or, if normalized to converge to a constant, may be singular asymptotically). As a result, these models fall outside of the classes considered by the previous literature (to take a non-dynamic example, it can be shown that the standard weak IV model is not Locally Asymptotically Quadratic, and thus is not subject to the results of Le Cam). Some additional complications in the DSGE context include the fact that the parameter space is in general quite large and that analytic expressions for the log likelihood are in general unavailable, though the likelihood can be evaluated numerically.

Structure of the paper Section 2 introduces our notation as well as some results from martingale theory; it also discusses the difference between two alternative measures of information. Section 3 suggests a test for the full parameter vector. Section 4 discusses the problem of testing a composite hypothesis about a sub-parameter, and introduces a statistic for such a test. Section 5 proves that our sub-vector test is valid when the nuisance parameter is strongly identified without any assumption on the identification of the tested parameter. Section 6 shows that this result can be extended to some cases when the nuisance parameter is weakly identified. Simulations supporting our theoretical results are provided in Section 7.

Throughout the rest of the paper, $I d_{k}$ is the $k \times k$ identity matrix, $\mathbb{I}\{\cdot\}$ is the indicatorfunction, $[\cdot]_{T}$ stands for quadratic variation of a martingale and $[\cdot, \cdot]_{T}$ - for joint quadratic variation of two martingales, $\Rightarrow$ denotes weak convergence (convergence in distribution), while $\rightarrow^{p}$ stands for convergence in probability.

## 2 Martingale Methods in Maximum Likelihood Theory

### 2.1 Setup

Let $X_{T}$ be the data available at time $T$. In general, we assume that $X_{T}=\left(x_{1}, \ldots, x_{T}\right)$. Let $\mathcal{F}_{t}$ be a sigma-algebra generated by $X_{t}=\left(x_{1}, \ldots, x_{t}\right)$. We assume that the log likelihood of the model,

$$
\ell\left(X_{T} ; \theta\right)=\log f\left(X_{T} ; \theta\right)=\sum_{t=1}^{T} \log f\left(x_{t} \mid \mathcal{F}_{t-1} ; \theta\right)
$$

is known up to the parameter $\theta$, which has true value $\theta_{0}$. We further assume that $\ell\left(X_{T} ; \theta\right)$ is twice continuously differentiable with respect to $\theta$, and that the class of likelihood gradients $\left\{\frac{\partial}{\partial \theta} \ell\left(X_{T} ; \theta\right): \theta \in \Theta\right\}$ and the class of second derivatives $\left\{\frac{\partial^{2}}{\partial \partial \partial \theta^{\prime}} \ell\left(X_{T} ; \theta\right)\right\}$ are both locally dominated integrable.

Our main object of study will be the score function,

$$
S_{T}(\theta)=\frac{\partial}{\partial \theta} \ell\left(X_{T}, \theta\right)=\sum_{t=1}^{T} \frac{\partial}{\partial \theta} \log f\left(x_{t} \mid \mathcal{F}_{t-1} ; \theta\right)
$$

where $s_{t}(\theta)=S_{t}(\theta)-S_{t-1}(\theta)=\frac{\partial}{\partial \theta} \log f\left(x_{t} \mid \mathcal{F}_{t-1} ; \theta\right)$ is the increment of the score. Under the assumption that we have correctly specified the model, the expectation of $s_{t}\left(\theta_{0}\right)$ conditional on all information up to $t-1$ is equal to zero,

$$
\begin{equation*}
E\left(s_{t}\left(\theta_{0}\right) \mid \mathcal{F}_{t-1}\right)=0 \quad \text { a.s. } \tag{1}
\end{equation*}
$$

This in turn implies that the score taken at the true parameter value, $S_{t}\left(\theta_{0}\right)$, is a martingale with respect to filtration $\mathcal{F}_{t}$. One way to view (1) is as a generalization of the first informational equality, which in i.i.d. models states that $E\left[s_{t}\left(\theta_{0}\right)\right]=0$, to the dynamic
context. To derive this equality, note that $s_{t}\left(\theta_{0}\right)=\frac{1}{f\left(x_{t} \mid \mathcal{F}_{t-1} ; \theta_{0}\right)} \frac{\partial}{\partial \theta} f\left(x_{t} \mid \mathcal{F}_{t-1} ; \theta_{0}\right)$,

$$
E\left(s_{t}\left(\theta_{0}\right) \mid \mathcal{F}_{t-1}\right)=\int s_{t}\left(\theta_{0}\right) f\left(x_{t} \mid \mathcal{F}_{t-1} ; \theta_{0}\right) d x_{t}=\int \frac{\partial}{\partial \theta} f\left(x_{t} \mid \mathcal{F}_{t-1} ; \theta_{0}\right) d x_{t}=0
$$

This observation is due to Silvey (1961).
Similarly, the second informational equality also generalizes to the dependent case. In the i.i.d. case, this equality states that we can calculate Fisher's information using either the Hessian of the log likelihood or the outer product of the score, i.e.

$$
\begin{equation*}
\mathfrak{I}\left(\theta_{0}\right)=-E\left(\frac{\partial^{2}}{\partial \theta \partial \theta^{\prime}} \log f\left(x_{t} ; \theta_{0}\right)\right)=E\left(\frac{\partial}{\partial \theta} \log f\left(x_{t} ; \theta_{0}\right) \frac{\partial}{\partial \theta^{\prime}} \log f\left(x_{t} ; \theta_{0}\right)\right) . \tag{2}
\end{equation*}
$$

Fisher's information plays a key role in the classical asymptotics for maximum likelihood, as it is directly related to the asymptotic variance of the MLE, and (2) suggests two different ways of estimating it which are asymptotically equivalent in the classical context. To generalize (2) to the dynamic context, following Barndorff-Nielsen and Sorensen (1991), we introduce two measures of information based on observed quantities:

- Observed information: the negative of Hessian of log-likelihood,

$$
I_{T}(\theta)=-\frac{\partial^{2}}{\partial \theta \partial \theta^{\prime}} \ell\left(X_{T} ; \theta\right)=\sum_{t=1}^{T} i_{t}(\theta)
$$

where $i_{t}(\theta)=-\frac{\partial^{2}}{\partial \theta \partial \theta^{\prime}} \log f\left(x_{t} \mid X_{t-1} ; \theta\right)$;

- Incremental observed information: the quadratic variation of the score of the log likelihood,

$$
J_{T}(\theta)=[S(\theta)]_{T}=\sum_{t=1}^{T} s_{t}(\theta) s_{t}^{\prime}(\theta),
$$

where as before $s_{t}(\theta)$ is the increment of $S_{T}(\theta)$.
Using these definitions, let $A_{T}(\theta)=J_{T}(\theta)-I_{T}(\theta)$ be the difference between the two measures of observed information. The second informational equality implies that $A_{T}\left(\theta_{0}\right)$ is a martingale with respect to $\mathcal{F}_{t}$. Specifically, the increment of $A_{T}\left(\theta_{0}\right)$ is $a_{t}\left(\theta_{0}\right)=$ $A_{t}\left(\theta_{0}\right)-A_{t-1}\left(\theta_{0}\right)$,

$$
a_{t}\left(\theta_{0}\right)=\frac{\partial^{2}}{\partial \theta \partial \theta^{\prime}} \log f\left(x_{t} \mid X_{t-1} ; \theta_{0}\right)+\frac{\partial}{\partial \theta} \log f\left(x_{t} \mid X_{t-1} ; \theta_{0}\right) \frac{\partial}{\partial \theta^{\prime}} \log f\left(x_{t} \mid X_{t-1} ; \theta_{0}\right),
$$

and an argument similar to that for the first informational equality gives us that $E\left(a_{t} \mid \mathcal{F}_{t-1}\right)=$ 0 a.s.

In the classical context, $I_{T}$ and $J_{T}$ are asymptotically equivalent, which plays a key role in the asymptotics of maximum likelihood. In the i.i.d. case, for example, the law of large numbers implies that $\frac{1}{T} I_{T}\left(\theta_{0}\right) \rightarrow^{p}-E\left(\frac{\partial^{2}}{\partial \theta \partial \theta^{\prime}} \log f\left(x_{t}, \theta_{0}\right)\right)=\Im\left(\theta_{0}\right)$ and $\frac{1}{T} J_{T}\left(\theta_{0}\right) \rightarrow^{p}$ $E\left(\frac{\partial}{\partial \theta} \log f\left(x_{t}, \theta_{0}\right) \frac{\partial}{\partial \theta^{\prime}} \log f\left(x_{t}, \theta_{0}\right)\right)=\Im\left(\theta_{0}\right)$. As a result of this asymptotic equivalence, the classical literature in the i.i.d. context uses these two measures of information more or less interchangeably.

The classical literature in the dependent context makes use of a similar set of conditions to derive the asymptotic properties of the MLE, focusing in particular on the asymptotic negligibility of $A_{T}\left(\theta_{0}\right)$ relative to $J_{T}\left(\theta_{0}\right)$. For example, Hall and Heyde (1980), show that for $\theta$ scalar, if $J_{T}\left(\theta_{0}\right) \rightarrow \infty$ a.s. and in addition,

$$
\lim \sup _{T \rightarrow \infty} J_{T}\left(\theta_{0}\right)^{-1}\left|A_{T}\left(\theta_{0}\right)\right|<1 \text { a.s. }
$$

then the MLE for $\theta$ is strongly consistent. If moreover, $J_{T}\left(\theta_{0}\right)^{-1} I_{T}\left(\theta_{0}\right) \rightarrow 1$ a.s., then the ML estimator is asymptotically normal and $J_{T}\left(\theta_{0}\right)^{\frac{1}{2}}\left(\hat{\theta}-\theta_{0}\right) \Rightarrow N(0,1)$.

We depart from this classical approach in that we consider weak identification. Weak identification arises when information is small along some dimension, which we model by using an embedding such that Fisher's information is degenerate asymptotically. Similar embeddings have been used to study weak identification in other contexts, including the Weak Instrument asymptotics introduced by Staiger and Stock (1997), and the Weak GMM asymptotics of Stock and Wright (2000). In such an embedding the difference between our two measures of information is important, and $A_{T}\left(\theta_{0}\right)$ is no longer negligible asymptotically compared to observed incremental information $J_{T}$ as demonstrated in the weak IV example below.

### 2.2 Weak IV Example

We assume a reduced form model with normal errors:

$$
\left\{\begin{array}{c}
y_{t}=\beta \pi^{\prime} z_{t}+u_{t}  \tag{3}\\
x_{t}=\pi^{\prime} z_{t}+v_{t}
\end{array},\binom{u_{t}}{v_{t}} \sim \text { i.i.d. } N\left(0, I d_{2}\right),\right.
$$

We take $z_{t}$ to be a $k$-dimensional set of instruments, while $\beta$ is the parameter of interest and $\pi$ is a $k \times 1$ vector of nuisance parameters. Our assumption that the errors have known covariance matrix equal to $I d_{2}$ is not restrictive, since $u_{t}$ and $v_{t}$ are reduced form (rather than structural) errors, and thus are well-estimable. The analysis is done conditional on the instruments $z_{t}$, and for simplicity we assume that the data generating process for $z_{t}$ is such that it satisfies a law of large numbers. Following the approach laid out by Staiger and Stock (1997), we represent weak identification by modeling $\pi$ as local to zero, that is $\pi=\frac{1}{\sqrt{T}} C$, so $\pi$ is drifting to zero as the sample grows.

Let $Y=\left(y_{1}, \ldots, y_{T}\right)^{\prime}, X=\left(x_{1}, \ldots, x_{T}\right)^{\prime}$ be $T \times 1$ and $Z=\left(z_{1}, \ldots, z_{T}\right)^{\prime}$ be $T \times k$. In this model, we have the following log-likelihood:

$$
\ell_{T}(\beta, \pi)=\text { const }-\frac{1}{2}(Y-\beta Z \pi)^{\prime}(Y-\beta Z \pi)-\frac{1}{2}(X-Z \pi)^{\prime}(X-Z \pi) .
$$

The score is

$$
S_{\beta}(\theta)=\pi^{\prime} Z^{\prime}(Y-\beta-Z \pi) ; S_{\pi}(\theta)=\beta Z^{\prime}(Y-\beta-Z \pi)+Z^{\prime}(X-Z \pi)
$$

Finally, the two measures of information are:

$$
\begin{gathered}
I_{T}\left(\theta_{0}\right)=-\frac{\partial^{2}}{\partial \theta \partial \theta^{\prime}} l_{T}=\left(\begin{array}{cc}
\pi^{\prime} Z^{\prime} Z \pi & \beta \pi^{\prime} Z^{\prime} Z-U^{\prime} Z \\
\beta Z^{\prime} Z \pi-Z^{\prime} U & \left(1+\beta^{2}\right) Z^{\prime} Z
\end{array}\right) ; \\
J_{T}\left(\theta_{0}\right)=[S]_{T}=\left(\begin{array}{cc}
\pi^{\prime} \sum_{t} u_{t}^{2} Z_{t} Z_{t}^{\prime} \pi & \pi^{\prime} \sum_{t} u_{t}\left(\beta u_{t}+v_{t}\right) Z_{t} Z_{t}^{\prime} \\
\sum_{t} u_{t}\left(\beta u_{t}+v_{t}\right) Z_{t} Z_{t}^{\prime} \pi & \sum_{t}\left(\beta u_{t}+v_{t}\right)^{2} Z_{t} Z_{t}^{\prime}
\end{array}\right) .
\end{gathered}
$$

Using the weak instrument embedding $\pi=\frac{1}{\sqrt{T}} C$, we can use normalizing matrix $K_{T}=$ $\operatorname{diag}\left(1, \frac{1}{\sqrt{T}}\right)$ to get a non-trivial limit for both information matrices:

$$
\begin{gathered}
K_{T} J_{T}\left(\theta_{0}\right) K_{T} \rightarrow^{p}\left(\begin{array}{cc}
C^{\prime} Q_{Z} C & \beta C^{\prime} Q_{Z} \\
\beta Q_{Z} C & \left(1+\beta^{2}\right) Q_{Z}
\end{array}\right) \\
K_{T} I_{T}\left(\theta_{0}\right) K_{T} \rightarrow^{p}\left(\begin{array}{cc}
C^{\prime} Q_{Z} C & \beta C^{\prime} Q_{Z}-\xi \\
\beta Q_{Z} C-\xi & \left(1+\beta^{2}\right) Q_{Z}
\end{array}\right) .
\end{gathered}
$$

To derive these expressions we have used a law of large numbers, $\frac{1}{T} Z Z^{\prime} \rightarrow^{p} Q_{Z}$, and
a central limit theorem, $\frac{1}{\sqrt{T}} Z^{\prime} U \Rightarrow \xi=N\left(0, Q_{Z}\right)$. Notice that, under weak instrument asymptotics, there is a difference between the two information matrices (i.e. the addition of the term $-\xi$ to the off-diagonal elements of $I_{T}$ ), whereas for the strong IV case ( $\pi \neq 0$ and fixed) we have that $J_{T}^{-1} I_{T} \rightarrow^{p} I d_{2}$.

The difference between the two measures of information can be used to construct a test to detect weak identification. A potential test should compare the two observed informations at the true parameter value. As argued in the introduction, however, tests of identification are less useful than weak identification robust procedures so we do not pursue such tests here.

White (1982) shows in the context of quasi-MLE that the two measures of information may be asymptotically different if the likelihood is misspecified. As we point out above, even if the model is correctly specified the two informations may differ if identification is weak. While we are aware of one strand of the classical statistical literature which explores the difference between these different information measures, the literature on so-called non-ergodic models, these models are usually part of the LAMN (locally asymptotically mixed-normal) class, whereas the types of models which we consider in this paper are not in general LAMN.

## 3 Test for Full Parameter Vector

In this section, we suggest a test for a simple hypothesis on the full parameter vector, $H_{0}: \theta=\theta_{0}$, which is robust to weak identification. To allow for the possibility of an embedding such as weak IV, we consider a so-called scheme of series. In a scheme of series we assume that we have a series of experiments indexed by the sample size: the data $X_{T}$ of sample size T is generated by distribution $f_{T}\left(X_{T} ; \theta_{0}\right)$, which may change as T grows. We assume that in the definition of all quantities in the previous section there is a silent index $T$. For example, the log-likelihood is $\ell_{T}(\theta)=\sum_{t=1}^{T} \log f_{T}\left(x_{T, t} \mid X_{T, t-1} ; \theta\right)$, where the data is $X_{T}=\left(x_{T, 1}, \ldots, x_{T, T}\right)$ and $X_{T, t}=\left(x_{T, 1}, \ldots, x_{T, t}\right)$. All scores and information matrices also have this implied index T ; for each fixed T the score $S_{T, t}$ is a process indexed by $t, S_{T, t}\left(\theta_{0}\right)=\frac{\partial}{\partial \theta} \log f_{T}\left(X_{T, t} ; \theta_{0}\right)=\sum_{j=1}^{t} s_{T, j}\left(\theta_{0}\right)$, and is a martingale with respect to the sigma-field $\mathcal{F}_{T, t}$ generated by $X_{T, t}$. All other statistics are defined correspondingly. In this context, we introduce our first assumption:

Assumption 1 Assume that there exists a sequence of constant matrices $K_{T}$ such that:
(a) $\sum_{t=1}^{T} E\left(K_{T}\left|s_{t, T}\left(\theta_{0}\right)\right| \mathbb{I}\left\{\left|K_{T} s_{t, T}\left(\theta_{0}\right)\right|>\delta\right\} \mid \mathcal{F}_{t-1}\right) \rightarrow 0$
(b) $\sum_{t=1}^{T} K_{T} s_{t, T}\left(\theta_{0}\right) s_{t, T}\left(\theta_{0}\right)^{\prime} K_{T}=K_{T} J\left(\theta_{0}\right) K_{T} \rightarrow^{p} \Sigma$, where $\Sigma$ is constant positivedefinite matrix

## Discussion of Assumption 1

Assumption 1 (a) is a classical infinitesimality (or limit negligibility) condition. We can, if we prefer, replace it with a version of Linderberg's condition:

$$
\sum_{t=1}^{T} E\left(\left\|K_{T} s_{t, T}\right\|^{2} \mathbb{I}\left\{\left\|K_{T} s_{t, T}\left(\theta_{0}\right)\right\|>\delta\right\} \mid \mathcal{F}_{t-1}\right) \rightarrow 0
$$

although this condition is stronger than 1 (a). Assumption 1 (b) imposes the ergodicity of the quadratic variation $J_{T}\left(\theta_{0}\right)$ of martingale $S_{T}\left(\theta_{0}\right)$, which rules out some potentially interesting models including persistent (unit root) processes and non-ergodic models.

Assumption 1 is trivially satisfied for the weak IV model we consider in section 2.2, and can also be checked for an exponential family with weak identification. In particular, consider an exponential family with joint density of the form

$$
\begin{equation*}
f_{T}\left(X_{t} \mid \theta\right)=h\left(X_{T}\right) \exp \left\{\eta_{T}(\theta)^{\prime} \sum_{t=1}^{T} H\left(x_{t}\right)-T A_{T}\left(\eta_{T}(\theta)\right)\right\} . \tag{4}
\end{equation*}
$$

Here, $\eta$ is a $p$-dimensional reduced form parameter, while $\sum_{t=1}^{T} H\left(x_{t}\right)$ is a $p$-dimensional sufficient statistic. Model (4) covers VAR models with $\eta$ being a set of reduced form VAR coefficients and $x_{t}=\left(Y_{t}^{\prime}, \ldots, Y_{t-p}^{\prime}\right)^{\prime}$, where $Y_{t}$ is a vector of data observed at time $t$, and the sufficient statistics are the sample autocovariances of the $Y_{t}$. Fernéndez-Villaverde et al. (2007) discuss the relationship between linearized DGSE models and VARs.

Suppose that we can partition structural coefficient $\theta$ into sub-vectors $\alpha$ and $\beta, \theta=$ $(\alpha, \beta)$. We consider an embedding similar to that of Stock and Wright (2000) for weak GMM, which we use to model $\beta$ as weakly identified. In particular, we assume that

$$
\eta_{T}(\theta)=m(\alpha)+\frac{1}{\sqrt{T}} \tilde{m}(\alpha, \beta)
$$

where $\frac{\partial}{\partial \beta} m\left(\alpha_{0}\right)$ and $\frac{\partial}{\partial \theta} \tilde{m}\left(\alpha_{0}, \beta_{0}\right)$ are matrices of full rank $\left(\operatorname{dim}(\theta)=k=k_{\alpha}+k_{\beta} \leq p\right)$. This means that while $\theta$ is identified for any fixed $T$, the likelihood is close to flat in direc-
tions corresponding to $\beta$. Assumption 1 is trivially satisfied for $K_{T}=\left(\begin{array}{cc}\frac{1}{\sqrt{T}} I d_{k_{\alpha}} & 0 \\ 0 & I d_{k_{\beta}}\end{array}\right)$ so long as the infinitesimality condition holds for the sequence $\left\{\frac{1}{\sqrt{T}} H\left(x_{t}\right)\right\}_{t=1}^{T}$ and a law of large numbers holds for $H\left(x_{t}\right) H\left(x_{t}\right)^{\prime}$ (i.e. $\left.\frac{1}{T} \sum_{t=1}^{T} H\left(x_{t}\right) H\left(x_{t}\right)^{\prime} \rightarrow^{p} E\left[H\left(x_{t}\right) H\left(x_{t}\right)^{\prime}\right]\right)$. For example, if $x_{t}$ is i.i.d. a sufficient condition for Assumption 1 in this embedding is that $H\left(x_{t}\right)$ posses at least two finite moments.

The following theorem is a direct corollary of the multivariate martingale Central Limit Theorem (see Theorem 8, ch. 5 in Liptser and Shiryayev (1989))

Theorem 1 If Assumption 1 holds, then

$$
K_{T} S_{T}\left(\theta_{0}\right) \Rightarrow N(0, \Sigma)
$$

and

$$
\begin{equation*}
L M\left(\theta_{0}\right)=S_{T}\left(\theta_{0}\right) J_{T}\left(\theta_{0}\right)^{-1} S_{T}\left(\theta_{0}\right) \Rightarrow \chi_{k}^{2}, \tag{5}
\end{equation*}
$$

where $k=\operatorname{dim}\left(\theta_{0}\right)$.

Remark. A weaker form of Assumption 1 is sufficient for statement (5). In particular, we may allow $\Sigma$ in Assumption 1 (b) to be an almost surely positive definite random matrix, rather than being constant. This is the so-called non-ergodic case, statistical examples of which can be found in Basawa and Koul (1979).

Statement (5) of Theorem 1 suggests a test for simple hypotheses about the whole parameter vector $\theta$. Unlike the classical ML Wald and LR tests, the derivation of the asymptotic distribution of this statistic uses no assumptions about the strength of identification. The statistic is a special form of the classical LM (score) test, which is formulated as:

$$
L M=\frac{1}{T} S_{T}\left(\theta_{0}\right)^{\prime} \hat{\mathcal{I}}^{-1} S_{T}\left(\theta_{0}\right)
$$

where $\hat{\mathcal{I}}$ is any consistent estimator of Fisher's information. Our suggested statistic plugs in $\frac{1}{T} J_{T}\left(\theta_{0}\right)=\frac{1}{T}\left[S\left(\theta_{0}\right)\right]_{T}$ for this estimator. It is important to note that while the true Fisher information is asymptotically degenerate under weak identification, the appropriately defined LM statistic (as in (5)) nevertheless achieves a $\chi^{2}$ distribution
asymptotically. It is likewise important to note that the LM statistic calculated with other estimators of Fisher's information (for example $\frac{1}{T} I_{T}\left(\theta_{0}\right)$ ) is not necessarily robust to weak identification, as can be seen in the example of weak IV. It is also a bad idea to estimate the information matrix using an estimator of $\theta$, i.e. to use $\frac{1}{T} J_{T}(\hat{\theta})$. All of these alternative formulations deliver asymptotically equivalent tests in strongly identified models, but this equivalence fails under weak identification.

## 4 Test for a Subset of Parameters

### 4.1 The Problem

In applied economics, it is very common to report confidence sets for estimates as separate confidence intervals for each one-dimensional sub-parameter in the (often quite multidimensional) parameter vector $\theta$. Current standards require that each such confidence interval be valid, that is, it should have at least $95 \%$ coverage asymptotically (assuming the typical $95 \%$ confidence level). These one-dimensional confidence sets need not be valid jointly: if $\operatorname{dim}(\theta)=k$, the k-dimensional rectangle formed by the Cartesian product of the 1-dimensional confidence intervals need not have $95 \%$ asymptotic coverage. Going the other direction, if one has a $95 \%$ confidence set for $\theta$ and projects it on the one-dimensional subspaces corresponding to the individual sub-parameters, the resulting confidence sets for the one-dimensional parameters will of course be valid. However, confidence sets obtained in such a manner (usually called the projection method) tend to be conservative.

Using our proposed test of the full parameter vector, which is robust to weak identification, we have the option to produce robust confidence sets for sub-parameters via the projection method. This approach has been used many times in the literature, for example by Dufour and Taamouti (2005) for weak IV and Guerron-Quintana, Inoue, and Killian (2009) for DSGE. The typical DSGE model has a large number of parameters to estimate (often between 20 and 60 ), which makes projection less attractive as the degree of conservativeness may be very high, which in turn makes the resulting confidence sets less informative.

For some intuition on the source of this conservativeness, imagine for a moment that we are concerned with a two-dimensional parameter $\theta=\left(\theta_{1}, \theta_{2}\right)$, and have a t-
statistic for each $\theta_{i}$. Suppose, moreover, that these two statistics are asymptotically normal and asymptotically independent of each other. We can construct a confidence set for each parameter in two ways: the first and most commonly used is to invert the t-test for the corresponding sub-parameter, which is equivalent to using the the squared t-statistic and $\chi_{1}^{2}$ critical values and yields $C_{1, \theta_{i}}=\left\{\theta_{i}: \frac{\left(\hat{\theta}_{i}-\theta_{i}\right)^{2}}{\sigma_{i}^{2}} \leq \chi_{1,95}^{2}\right\}$. As an alternative, one may construct a joint confidence set for $\theta$, which in this case will be an ellipse $C_{2, \theta}=\left\{\theta: \frac{\left(\hat{\theta}_{1}-\theta_{1}\right)^{2}}{\sigma_{1}^{2}}+\frac{\left(\hat{\theta}_{2}-\theta_{2}\right)^{2}}{\sigma_{2}^{2}} \leq \chi_{2,95}^{2}\right\}$, and then use the projection method to obtain $C_{2, \theta_{1}}=\left\{\theta_{1}: \exists \theta_{2}\right.$ s.t. $\left.\left(\theta_{1}, \theta_{2}\right) \in C_{2, \theta}\right\}$ (and likewise for $\theta_{2}$ ). One can notice that $C_{2, \theta_{i}}$ ultimately uses the same t-statistic as $C_{1, \theta_{i}}$, but compares this statistic to the critical value of a $\chi_{2}^{2}$ rather than a $\chi_{1}^{2}$. As a result, in this example the projection method produces unnecessarily wide (and conservative) confidence sets for each sub-parameter.

The projection method, when applied to strongly identified models, produces a less powerful test than classical MLE. Thus, when using projection it is natural to combine it with a pre-test procedure which first discriminates between weakly and strongly identified models and then, based on the results of the test, uses either classical MLE or the projection method. There are two obstacles to such an approach: first, we are unaware of procedures for effectively discriminating between weak and strong identification in maximum likelihood. Second, the size properties of two-step testing procedures are notoriously difficult to asses. Our approach is different, and instead constructs a test which maintains correct asymptotic size under weak identification, but which is equivalent to the classical MLE tests under strong identification.

We are aware of a number of papers dealing with this issue in the context of weak identification. In particular, Stock and Wright (2000) prove that for GMM, under some assumptions, if $\theta=(\alpha, \beta)$ and $\alpha$ is well-identified then it is possible to test the hypothesis $H_{0}: \beta=\beta_{0}$ by comparing the GMM objective function, minimized with respect to $\alpha$, to the critical values of a $\chi_{p-k_{\alpha}}^{2}$ distribution, where $p$ is the number of moment conditions used and $k_{\alpha}=\operatorname{dim}(\alpha)$. Their result shows that it is possible to reduce the degrees of freedom for projection-based confidence sets in weak GMM provided the nuisance parameter is well identified.

Kleibergen and Mavroeidis (2009) prove that it is possible to extend this result to some models where the nuisance parameter may not be well identified. They consider a test statistic, called $H\left(\theta_{0}\right)$ here, for testing the simple hypothesis $H_{0}: \theta=\theta_{0}$ (they use the Anderson-Rubin and IV-LM tests). Assume again that $\theta=(\alpha, \beta)$, and that the
hypothesis of interest is $H_{0}: \beta=\beta_{0}$. Kleibergen and Mavroeidis (2009) demonstrate that one can again use the quantiles of a $\chi_{p-k_{\alpha}}^{2}$ as critical values. This test is similar if identification is strong, and somewhat conservative if identification is weak. In this paper we consider a class of models which, as discussed above, differs from those in the weak GMM literature in that the variance of the moment conditions may be degenerate asymptotically and take a different approach.

### 4.2 Classical LM Tests for Composite Hypotheses

We assume that $\theta=(\alpha, \beta)$. We are interested in testing the composite hypothesis $H_{0}: \beta=\beta_{0}$, treating $\alpha$ as a nuisance parameter. The classical theory for maximum likelihood considers two LM tests for such a setting: Rao's score test and Neyman's $C(\alpha)$-test.

Let $S_{T}(\theta)=\left(S_{\alpha}(\theta)^{\prime}, S_{\beta}(\theta)^{\prime}\right)^{\prime}, \mathcal{I}(\theta)=\left(\begin{array}{cc}\mathcal{I}_{\alpha \alpha} & \mathcal{I}_{\alpha \beta} \\ \mathcal{I}_{\alpha \beta}^{\prime} & \mathcal{I}_{\beta \beta}\end{array}\right)$ be Fisher's information, and $\hat{\theta}_{0}$ be the restricted ML estimator of $\theta$, under the restriction $\hat{\beta}=\beta_{0}$. Assume, in addition, that all martingales introduced in Section 2 are divided into sub-matrices corresponding to $\alpha$ and $\beta$. Rao's score test is based on the statistic

$$
\text { Rao }=S_{T}\left(\hat{\theta}_{0}\right)^{\prime} \mathcal{I}\left(\hat{\theta}_{0}\right)^{-1} S_{T}\left(\hat{\theta}_{0}\right),
$$

where $\hat{\theta}_{0}$ is the restricted ML estimator of $\theta$ under $H_{0}$.
Neyman's $C(\alpha)$ test was developed as a locally asymptotically most powerful (LAMP) test for composite hypotheses in the classical ML model (see Akritas (1987)). The statistic is defined as

$$
C(\alpha)=\left.\left(S_{\beta}-\mathcal{I}_{\alpha \beta}^{\prime} \mathcal{I}_{\alpha \alpha}^{-1} S_{\alpha}\right)^{\prime} \mathcal{I}_{\beta \beta, \alpha}^{-1}\left(S_{\beta}-\mathcal{I}_{\alpha \beta}^{\prime} \mathcal{I}_{\alpha \alpha}^{-1} S_{\alpha}\right)\right|_{\theta=\left(\hat{\alpha}, \beta_{0}\right)},
$$

where $\hat{\alpha}$ is any $\sqrt{T}$ consistent estimator of $\alpha$, and $\mathcal{I}_{\beta \beta, \alpha}=\mathcal{I}_{\beta \beta}-\mathcal{I}_{\beta \alpha} \mathcal{I}_{\alpha \alpha}^{-1} \mathcal{I}_{\alpha \beta}$.
Kocherlakota and Kocherlakota (1991) show that the two statistics are the same if one takes $\hat{\alpha}$ in Neyman's $C(\alpha)$ test to be the restricted MLE. If the classical ML assumptions are satisfied then both statistics are distributed $\chi_{k_{\beta}}^{2}$ asymptotically. In this paper, we suggest a statistic which is asymptotically equivalent to both Rao's score and Neyman's $C(\alpha)$ if the classical ML assumptions are satisfied. In particular, we consider the same

LM statistic defined in (5) but evaluated at $\theta=\left(\hat{\alpha}, \beta_{0}\right)$, where $\hat{\alpha}$ is the restricted MLE, that is, the solution to equation

$$
\begin{equation*}
S_{\alpha}\left(\hat{\alpha}, \beta_{0}\right)=0 \tag{6}
\end{equation*}
$$

One can easily see that

$$
\begin{equation*}
\widetilde{L M}\left(\beta_{0}\right)=L M\left(\hat{\alpha}, \beta_{0}\right)=\left.S_{\beta}^{\prime}\left(J_{\beta \beta}-J_{\beta \alpha} J_{\alpha \alpha}^{-1} J_{\beta \alpha}^{\prime}\right)^{-1} S_{\beta}\right|_{\theta=\left(\hat{\alpha}, \beta_{0}\right)} \tag{7}
\end{equation*}
$$

## 5 Test for a Subset of Parameters- Strong Identification

In this section, we establish that if $\alpha$ is strongly identified then the statistic defined in (7) has a $\chi_{k_{\beta}}^{2}$ distribution asymptotically, regardless of the strength of identification of $\beta$.

### 5.1 How We Define Strong Identification of $\alpha$

When we test $H_{0}: \beta=\beta_{0}$, under the null $\alpha$ is the only unknown parameter. We call $\alpha$ strongly identified if it satisfies the assumptions below, which guarantee that the restricted ML estimate of $\alpha$ is consistent and asymptotically normal. We adapt Baht's (1974) result on the consistency and asymptotic normality of the MLE for time series to show that this is the case.

Let $A_{\alpha \alpha, T}=J_{\alpha \alpha, T}-I_{\alpha \alpha, T}$, where the last two quantities are the sub-matrices of $J_{T}\left(\theta_{0}\right)$ and $I_{T}\left(\theta_{0}\right)$ corresponding to $\alpha$.

Assumption 2 Assume that matrix $K_{T}$ from Assumption 1 is diagonal and $K_{\alpha, T}$ and $K_{\beta, T}$ are the sub-matrices of $K_{T}$ corresponding to $\alpha$ and $\beta$, respectively.
(a) $K_{\alpha, T} A_{\alpha \alpha, T} K_{\alpha, T} \rightarrow^{p} 0$.
(b) for any $\delta>0$ we have $\sup _{\left|\alpha_{1}-\alpha_{0}\right|<\delta K_{\alpha, T}}\left|\left(I_{\alpha \alpha}\left(\alpha_{1}, \beta_{0}\right)-I_{\alpha \alpha}^{0}\right)\left(I_{\alpha \alpha}^{0}\right)^{-1}\right| \rightarrow^{p} 0$
(c) $K_{\alpha, T} \rightarrow 0$ as $T \rightarrow \infty$.

Lemma 1 If Assumptions 1 and 2 are satisfied, then the restricted MLE $\hat{\alpha}\left(\beta_{0}\right)$ is consistent for $\alpha$ and

$$
\begin{equation*}
K_{\alpha, T}\left(\hat{\alpha}-\alpha_{0}\right)=K_{\alpha, T} J_{\alpha, \alpha, T}^{-1} S_{\alpha, T}+o_{p}(1) \Rightarrow N\left(0, \Sigma_{\alpha \alpha}^{-1}\right) \tag{8}
\end{equation*}
$$

## Discussion of Assumption 2

Assumption 2(a) may be reformulated as

$$
J_{\alpha \alpha, T}^{-1} I_{\alpha \alpha, T} \rightarrow^{p} I d_{k_{\alpha}}
$$

which requires that the two information matrices be the same asymptotically. We mentioned a condition of this nature in our discussion of weak identification in section 2. One approach to checking 2(a) in many contexts is to establish a law of large numbers for $A_{\alpha \alpha, T}$. Indeed, $A_{\alpha \alpha, T}$ is a martingale of the form

$$
A_{\alpha \alpha, T}=\sum_{t=1}^{T} \frac{1}{f\left(x_{t} \mid X_{t-1}, \theta_{0}\right)} \frac{\partial^{2}}{\partial \alpha \partial \alpha^{\prime}} f\left(x_{t} \mid X_{t-1}, \theta_{0}\right) .
$$

If the terms $\frac{1}{f\left(x_{t} \mid X_{t-1}, \theta_{0}\right)} \frac{\partial^{2}}{\partial \alpha \partial \alpha^{\prime}} f\left(x_{t} \mid X_{t-1}, \theta_{0}\right)$ are uniformly integrable and $K_{\alpha, T}$ converges to zero no slower than $\frac{1}{\sqrt{T}}$, then the martingale law of large numbers gives us Assumption 2(a).

Assumption 2(b) is an assumption on the smoothness of the log-likelihood. We can reformulate it using a third type of martingale, corresponding to the third derivatives:

$$
\begin{equation*}
\Lambda_{\alpha \alpha \alpha_{i}, T}=\sum_{t=1}^{T} \frac{1}{f\left(x_{t} \mid X_{t-1}, \theta_{0}\right)} \frac{\partial^{3}}{\partial \alpha_{i} \partial \alpha \partial \alpha^{\prime}} f\left(x_{t} \mid X_{t-1}, \theta_{0}\right) \tag{9}
\end{equation*}
$$

For all $i, \Lambda_{\alpha \alpha \alpha_{i}, T}$ is a martingale so long as we can interchange differentiation and integration of the log-likelihood function three times. An alternative to Assumption 2(b) is

Assumption 2( $b^{\prime}$ ) for any $i: K_{\alpha_{i}, T} K_{\alpha, T} \Lambda_{\alpha_{i} \alpha \alpha, T} K_{\alpha, T} \rightarrow^{p} 0$
Lemma 2 Assumptions 1, 2(a) and 2(b') imply assumption 2(b).
Finally, Assumption 2(c) implies that information about $\alpha$ accumulates as the sample size increases. This assumption is critical for consistency, but turns out to be less
important for asymptotic normality, so long as $2(\mathrm{a})$ and 2 (b) hold. In weak IV, for example, if we assume that $\pi$ is known and local to zero (i.e. $\pi=C / \sqrt{T}$ ), the restricted MLE for $\beta$ is asymptotically normal, even though it is not consistent. The corresponding normalization is $K_{\beta, T}=1$.

### 5.2 Result

As we show in section 3, to test a simple hypothesis about the whole parameter vector it is enough to have a CLT for the score function. Kleibergen and Mavroeidis (2009) impose a stronger assumption for the their test of a subset of parameters, namely that the CLT also hold for the derivative of the moment condition (in fact, they impose a functional CLT). For our test of a subset of parameters, we likewise need an additional assumption, specifically a CLT on the derivative of the score, which is directly related to the martingale $A_{T}$ (the difference of the two information matrices).

Assumption 3 Consider the sequence of martingales $M_{T}=\left(S_{T}\left(\theta_{0}\right)^{\prime} \text {, vec }\left(A_{\alpha, \beta, T}\left(\theta_{0}\right)\right)\right)^{\prime}=$ $\sum_{t=1}^{T} m_{t, T}$. Assume that there exists a sequence of non-stochastic diagonal matrices $K_{M, T}$ such that:
(a) $\sum_{t=1}^{T} E\left(K_{M, T}\left|m_{t, T}\right| \mathbb{I}\left\{\left|K_{M, T} m_{t, T}\right|>\delta\right\} \mid \mathcal{F}_{t-1}\right) \rightarrow 0$
(b) $\sum_{t=1}^{T} K_{M, T} m_{t, T} m_{t, T}^{\prime} K_{M, T} \rightarrow^{p} \Sigma_{M}$, where $\Sigma_{M}$ is a constant matrix whose sub-matrix $\Sigma$ corresponding to the martingale $S_{T}$ is positive definite.

Let us define the martingales associated with the third derivative of likelihood function

$$
\begin{equation*}
\Lambda_{\alpha_{i} \alpha_{j} \beta_{n}}=\sum_{t=1}^{T} \frac{1}{f\left(y_{t} \mid Y_{t-1}, \theta_{0}\right)} \cdot \frac{\partial^{3} f\left(y_{t} \mid Y_{t-1}, \theta_{0}\right)}{\partial \alpha_{i} \partial \alpha_{j} \partial \beta_{n}} . \tag{10}
\end{equation*}
$$

If we can interchange integration and differentiation three times then each entry of $\Lambda_{\alpha \alpha \beta, T}$ is a martingale. We also use the fourth-order martingales

$$
\Gamma_{\alpha_{i} \alpha_{j} \alpha_{m} \beta_{n}}=\sum_{t=1}^{T} \frac{1}{f\left(y_{t} \mid Y_{t-1}, \theta_{0}\right)} \cdot \frac{\partial^{4} f\left(y_{t} \mid Y_{t-1}, \theta_{0}\right)}{\partial \alpha_{i} \partial \alpha_{j} \partial \alpha_{m} \partial \beta_{n}}
$$

For the proof of the theorem below we will also need the following assumptions:
Assumption 4 (a) $\lim _{T \rightarrow \infty} K_{\alpha, T} K_{\alpha \beta, T}^{-1} K_{\beta, T}=C$ where $C$ is some finite matrix (which may be zero).
(b) $K_{\alpha_{i}, T} K_{\alpha_{j}, T} K_{\beta_{n}, T} \sqrt{\left[\Lambda_{\alpha_{i} \alpha_{j} \beta_{n}}\right]_{T}} \rightarrow^{p} 0$ for any $i, j, n$.
(c) $K_{\alpha_{i}, T} K_{\alpha_{j}, T} K_{\alpha_{m}, T} K_{\beta_{n}, T} \sqrt{\left[\Gamma_{\alpha_{i} \alpha_{j} \alpha_{m} \beta_{n}}\right]_{T}} \rightarrow^{p} 0$ for any $i, j, m, n$.

## Discussion of Assumption 4

Assumptions 4(b) and (c) state that the higher order derivatives with respect to $\alpha$ are not important for the analysis, and are used primarily in the proofs. For example, Assumption 4(c) implies that

$$
K_{\alpha, T}^{2} K_{\beta_{i}, T} \sup _{\left|\alpha^{*}-\alpha_{0}\right|<\delta K_{\alpha, T}}\left|I_{\alpha \alpha \beta_{i}}^{*}-I_{\alpha \alpha \beta_{i}}^{0}\right| \rightarrow^{p} 0 \text { for all } i .
$$

If $\alpha$ is strongly identified, then Assumptions 4(b) and (c) generally hold, and can be checked using some law of large numbers, since the normalization $K_{\alpha, T}^{2}$ or $K_{\alpha, T}^{3}$ converges to zero very quickly. Finally, Assumption 4 holds trivially for weak IV, as well as for the exponential family case discussed in section 3.

Theorem 2 If Assumptions 2, 3 and 4 are satisfied then under the null $H_{0}: \beta=\beta_{0}$ we have

$$
\widetilde{L M}\left(\beta_{0}\right) \Rightarrow \chi_{k_{\beta}}^{2}
$$

### 5.3 How Our Result Differs from the Previous Literature

As discussed above, Stock and Wright (2000) develop a framework for weakly identified GMM and construct a test for the hypothesis $H_{0}: \beta=\beta_{0}$ when the nuisance parameter $\alpha$ is strongly identified (Theorem 3 in Stock and Wright (2000)). They consider GMM with moment condition $\operatorname{Em}\left(x_{t}, \alpha, \beta\right)=0$ and construct a statistic based on

$$
S(\theta)=\left(\frac{1}{\sqrt{T}} \sum_{t=1}^{T} m\left(x_{t} ; \theta\right)\right)^{\prime} W_{T}(\theta)\left(\frac{1}{\sqrt{T}} \sum_{t=1}^{T} m\left(x_{t} ; \theta\right)\right)
$$

where $W_{T}(\theta)$ is a consistent estimator of the variance of the moment condition. They show that, for $\hat{\alpha}=\arg \min _{\alpha} S W\left(\alpha, \beta_{0}\right)$, their statistic $S\left(\hat{\alpha}, \beta_{0}\right)$ has an asymptotic $\chi^{2}$ distribution with degrees of freedom equal to $p-k_{\alpha}$, where $p=\operatorname{dim}\left(m\left(x_{t}, \theta\right)\right)$ and $k_{\alpha}=\operatorname{dim}(\alpha)$.

Kleibergen (2005) considers an alternative statistic based on the LM test for GMM and proves that this statistic, minimized over $\alpha$, is also the basis of a valid test of
$H_{0}: \beta=\beta_{0}$ when $\alpha$ is strongly identified. In our context, however, if we use the score of the log-likelihood as the GMM moment condition the system is just-identified and Kleibergen's KLM statistic is equal to Stock and Wright's S statistic.

Our result, though of a similar flavor, is quite different and is not covered by these previous results. First, the weak ML model does not satisfy the assumptions in the above mentioned papers. Specifically, if we consider ML estimation as GMM using the moment condition $E S_{T}\left(\theta_{0}\right)=0$, the variance matrix of our moment condition (information matrix) is directly linked to identification. In particular, the matrix $W_{T}(\theta)$ (to use Stock and Wright's notation) becomes degenerate asymptotically, which is ruled out by the assumptions of Stock and Wright (2000), Kleibergen (2005), and Kleibergen and Mavroeidis (2009). Second, we apply a different principle to go from a test of the full parameter vector to a test for a subset of parameters. In the above mentioned papers the authors minimize the statistic over the nuisance parameter, while we plug in the restricted MLE. In fact, in our context minimizing the statistic over the nuisance parameter does not necessarily lead to a $\chi^{2}$ distribution, as illustrated in the following weak IV example.

## Weak IV Example (cont.)

Consider the weak IV model (3) and consider the LM statistic for $L M(\pi, \beta)$ for testing the whole parameter vector $\theta=(\pi, \beta)$, defined as in section 2.3. Suppose we wish to test the composite hypothesis $H_{0}: \beta=\beta_{0}$ by considering the concentrated statistic:

$$
\begin{equation*}
L M^{c}\left(\beta_{0}\right)=\min _{\pi} L M\left(\pi, \beta_{0}\right)=L M\left(\tilde{\pi}, \beta_{0}\right) . \tag{11}
\end{equation*}
$$

We can show (see proof in the Appendix) that

$$
L M^{c}\left(\beta_{0}\right)=\frac{\left(Q_{S}+Q_{T}\right)-\sqrt{\left(Q_{S}+Q_{T}\right)^{2}-4 Q_{S T}^{2}}}{2}
$$

where $Q_{S}, Q_{T}$, and $Q_{S T}$ are defined as in Andrews, Moreira, and Stock (2006) (again, see Appendix). If the instruments are weak, that is if $\pi=C / \sqrt{T}$, then the asymptotic distribution of $L M^{c}\left(\beta_{0}\right)$ is stochastically dominated by a $\chi_{1}^{2}$, and the resulting test is conservative.

## 6 Test for a Subset of Parameters- Weak Identification

In the previous section we show that our subset-test statistic $\widetilde{L M}\left(\beta_{0}\right)$ for the composite hypothesis $H_{0}: \beta=\beta_{0}$ is asymptotically $\chi^{2}$ when the nuisance parameter $\alpha$ is strongly identified, without any assumptions about the identification of $\beta$.

This result can be extended somewhat to cases when $\alpha$ is weakly identified. Below are two such examples.

### 6.1 Weak IV Case

Here we consider a weak IV model with one endogenous variable, when the hypothesis tested is one about $\pi$, that is, $H_{0}: \pi=\pi_{0}$, while the weakly identified parameter $\beta$ is treated as a nuisance parameter.

As in section 2.2, we consider the model:

$$
\left\{\begin{array}{c}
Y=\beta_{0} Z \pi_{0}+U_{0} \\
X=Z \pi_{0}+V_{0}
\end{array} ; \quad\binom{u_{t}}{v_{t}} \sim N\left(0, I d_{2}\right)\right.
$$

For simplicity we consider a slightly different version of the quadratic variation of S, namely expected quadratic variation.

$$
\widetilde{J}=\langle S\rangle=\sum_{t=1}^{T} E\left(s_{t} s_{t}^{\prime} \mid \mathcal{F}_{t-1}\right)=\left(\begin{array}{cc}
\pi^{\prime} Z^{\prime} Z \pi & \beta \pi^{\prime} Z^{\prime} Z \\
\beta Z^{\prime} Z \pi & \left(1+\beta^{2}\right) Z^{\prime} Z
\end{array}\right)
$$

The difference between $J$ and $\widetilde{J}$ doesn't matter asymptotically as $J^{-1} \widetilde{J} \rightarrow^{p} I d_{k+1}$ uniformly over the strength of instruments.

According equation (7) our statistic of interest is

$$
\widetilde{L M}\left(\pi_{0}\right)=L M\left(\hat{\beta}, \pi_{0}\right),
$$

where $\hat{\beta}$ is the restricted ML estimator of $\beta$, and $L M\left(\beta, \pi_{0}\right)$ is defined as in (5) with the slight modification that $\widetilde{J}$ is used in place of $J$. A simple formula for $\operatorname{LM}\left(\beta, \pi_{0}\right)$ is given in (28). Note that $S_{\beta}\left(\hat{\beta}, \pi_{0}\right)=0$, and we can explicitly solve for $\hat{\beta}$ as

$$
\begin{equation*}
\pi_{0}^{\prime} Z^{\prime}\left(Y-\hat{\beta} Z \pi_{0}\right)=0 \Leftrightarrow \hat{\beta}=\frac{\pi_{0}^{\prime} Z^{\prime} Y}{\pi_{0}^{\prime} Z^{\prime} Z \pi_{0}} \tag{12}
\end{equation*}
$$

Simple calculations show that

$$
\begin{equation*}
L M\left(\hat{\beta}, \pi_{0}\right)=\left(\hat{\beta} \hat{U}+V_{0}\right)^{\prime} Z\left(\left(1+\hat{\beta}^{2}\right) Z^{\prime} Z-\frac{\hat{\beta}^{2} Z^{\prime} Z \pi_{0} \pi_{0}^{\prime} Z^{\prime} Z}{\pi_{0}^{\prime} Z^{\prime} Z \pi_{0}}\right)^{-1} Z^{\prime}\left(\hat{\beta} \hat{U}+V_{0}\right) \tag{13}
\end{equation*}
$$

where $\hat{U}=Y-\hat{\beta} Z \pi_{0}$.

Lemma 3 If $\pi_{0}=c / \sqrt{T}$ we have $\operatorname{LM}\left(\hat{\beta}, \pi_{0}\right) \Rightarrow \chi_{k}^{2}$
The idea of the proof is the following. Under the weak instruments embedding, $\hat{\beta}$ is not consistent but is asymptotically normal. We can show that $\left(Z^{\prime} Z\right)^{-1 / 2} Z^{\prime} \hat{U}, \hat{\beta}$ and $\left(Z^{\prime} Z\right)^{-1 / 2} Z^{\prime} V_{0}$ are asymptotically normal and asymptotically uncorrelated with each other. If we consider statistic $\operatorname{LM}\left(\hat{\beta}, \pi_{0}\right)$, conditional on $\hat{\beta}$ it becomes a correctly normalized quadratic form of an asymptotically normal $k$-dimensional random variable and thus conditionally asymptotically $\chi_{k}^{2}$. As a result, unconditional convergence holds as well.

### 6.2 Case Where Score is Linear in $\alpha$

The case considered in the previous subsection is interesting in that the nuisance parameter is weakly identified, but is somewhat trivial since the parameter tested is strongly identified. We can to a limited extent generalize this result to more interesting contexts. Below, we consider the problem of testing a hypothesis about a weakly identified parameter in an exponential family model. The nuisance parameter will be divided into two subsets, one of which is strongly identified while the other is weakly identified. We will make the very restrictive assumption that the weakly identified nuisance parameters enter linearly.

Assume that the experiment at time $T$ is generated by the exponential family (4). As already discussed, model (4) covers VAR models, and many linearized DGSE models can be represented as VARs (see Fernéndez-Villaverde et al. (2007)).

Assume that we are interested in structural parameters $\theta=\left(\alpha_{1}, \alpha_{2}, \beta\right)$, where the relation between the structural and reduced form parameters is given by

$$
\begin{equation*}
\eta_{T}(\theta)=m\left(\alpha_{1}\right)+\frac{1}{\sqrt{T}} n\left(\alpha_{1}, \beta\right) \alpha_{2}+\frac{1}{\sqrt{T}} r\left(\alpha_{1}, \beta\right) . \tag{14}
\end{equation*}
$$

We assume that the matrix $\left(\frac{\partial}{\partial \alpha_{1}} m\left(\alpha_{1}\right), n\left(\alpha_{1}, \beta\right), \frac{\partial}{\partial \beta} n\left(\alpha_{1}, \beta\right) \alpha_{2}+\frac{\partial}{\partial \beta} r\left(\alpha_{1}, \beta\right)\right)$ has full rank $k=\operatorname{dim}(\theta) \leq p$ and call this the rank assumption. That is, we assume that the structural parameters are identified, though only $\alpha_{1}$ is strongly identified (parameters $\alpha_{2}$ and $\beta$ are weakly identified).

We are interested in testing a composite hypothesis $H_{0}: \beta=\beta_{0}$, treating $\alpha=\left(\alpha_{1}, \alpha_{2}\right)$ as a nuisance parameter. We use $\widetilde{L M}\left(\beta_{0}\right)$ statistic as defined in (7).

Theorem 3 Assume that in model (4) and (14) which satisfies the rank assumption the following convergence holds at the true value of $\theta_{0}$ :
(a) $A_{T}(\eta) \rightarrow A(\eta)$, as $T \rightarrow \infty$ in a neighborhood of $\eta_{\infty}$ and the first four derivatives of $A_{T}$ at $\eta_{\infty}$ converge to those of $A(\cdot)$;
(b) $\frac{1}{T} \sum_{t=1}^{T} H\left(x_{t}\right) \rightarrow^{p} \dot{A}$;
(c) $\frac{1}{T} \sum_{t=1}^{T}\left(H\left(x_{t}\right)-\dot{A}\right)\left(H\left(x_{t}\right)-\dot{A}\right)^{\prime} \rightarrow^{p}-\frac{\partial^{2}}{\partial \eta \partial \eta^{\prime}} A\left(\eta_{\infty}\right)=-\ddot{A}$, where $\ddot{A}$ is a positivedefinite matrix;
(d) $\frac{1}{T} \sum_{t}\left(H\left(x_{t}\right)\right)^{3}=O_{p}(1)$.

Then under the null we have

$$
\widetilde{L M}\left(\beta_{0}\right) \Rightarrow \chi_{k_{\beta}}^{2} .
$$

## $7 \quad$ Simulation Results

We have a number of simulation results which both support our theoretical results and suggest directions for further research. We focus on simulation results from three models: a simple DSGE model based on Clarida, Gali, and Gertler (1999), a nonlinear extension of the standard weak IV model discussed earlier in this paper, and a weak IV model with two endogenous variables. In all cases, we simulate the behavior of our proposed statistics and compare the finite sample distributions of the statistics in question to their limiting distributions. In the DSGE example, we argue that estimation in the model behaves in a manner consistent with weak identification, and that our proposed statistics offer a substantial improvement over the usual Wald-based statistics for testing in this model. For the other two models, we use a standard specification for weak identification and show that our proposed tests have good properties in simulation.

### 7.1 DSGE Model

We consider a simple DSGE model based on Clarida, Gali and Gertler (1999). ${ }^{4}$ For this model, we first explore the properties of ML (and the usual ML-based test statistics), then discuss the properties of the information matrix, and finally explore the behavior of our proposed test statistics, both for the full parameter vector and for subsets of parameters.

The (log-linearized) equilibrium conditions for the model are

$$
\begin{gathered}
\beta E_{t} \pi_{t+1}+\kappa x_{t}-\pi_{t}+\varepsilon_{t}=0 \\
-\left[r_{t}-E_{t} \pi_{t+1}-r r_{t}^{*}\right]+E_{t} x_{t+1}-x_{t}=0 \\
r_{t}=\alpha r_{t-1}+(1-\alpha) \phi_{\pi} \pi_{t}+(1-\alpha) \phi_{x} x_{t}+u_{t} \\
r r_{t}^{*}=\rho \Delta a_{t}
\end{gathered}
$$

while the exogenous variables ( $\Delta a_{t}$ and $u_{t}$ ) evolve according to

$$
\begin{gathered}
\Delta a_{t}=\rho \Delta a_{t-1}+\varepsilon_{a, t} \\
u_{t}=\delta u_{t-1}+\varepsilon_{u, t} \\
{\left[\begin{array}{c}
\varepsilon_{t} \\
\varepsilon_{a, t} \\
\varepsilon_{u, t}
\end{array}\right] \sim i i d N\left(\left[\begin{array}{l}
0 \\
0 \\
0
\end{array}\right],\left[\begin{array}{ccc}
\sigma^{2} & 0 & 0 \\
0 & \sigma_{a}^{2} & 0 \\
0 & 0 & \sigma_{u}^{2}
\end{array}\right]\right)}
\end{gathered}
$$

The model has ten parameters: the discount rate $\beta$, the structural parameters $\kappa, \phi_{x}, \phi_{\pi}$, and $\alpha$, and the parameters describing the evolution of the exogenous variables. We calibrate the structural parameters at generally accepted values: $\beta=.99, \kappa=\frac{(1-\theta)(1+\phi)(1-\beta \theta)}{\theta} \approx$ .1717, $\phi_{x}=0, \phi_{\pi}=1.5$ and $\alpha=0$. For the parameters describing the exogenous variables, we choose $\rho=.2$ and $\delta=.2$, to introduce a degree of persistence while maintaining stationarity, and set $\sigma_{a}=1, \sigma_{u}=1, \sigma=1$. Using this model, we generate samples of size 300 and then discard the first 100 observations. We use only the last 200 observations from each simulation draw for the remainder of the analysis. Given well-documented problems with estimating $\beta$ in many models, for this point forward we also calibrate this

[^1]parameter at its true value, and conduct the analysis using the remaining 9 parameters. ${ }^{5}$

### 7.1.1 MLE Monte-Carlo Results

We begin by examining the behavior of the maximum likelihood estimator for the nine non-calibrated parameters in the model. We report histograms for the resulting estimates in Figure 1 (based on 500 Monte-Carlo draws), with the true value of each parameter reported in parentheses at the top of each subplot. As can be seen from the figure, the distribution of many of the estimates is quite far from the normal limiting distribution of the maximum likelihood estimator under the usual assumptions. Moreover, it appears that this non-normality is not purely the result of bad behavior on the part of one parameter: after experimenting with calibrating (to their true values) a number of different parameters, it appears that we need to exclude from the estimation (calibrate) at least three parameters before the distributions of the remaining parameters begin to appear well-approximated by normal distributions.

While the results in Figure 1 show that the usual asymptotics for the ML estimator provide a poor approximation to its finite-sample distribution in this model, our theoretical results focus on questions of inference rather than estimation, so we also look at the behavior of the usual maximum likelihood tests for this model. We consider each of the trinity of classical tests (LR, Wald, and LM) in turn, focusing on tests of the full parameter vector. Specifically, we test the hypothesis $H_{0}: \theta=\theta_{0}$, where $\theta$ is the vector consisting of all parameters other than $\beta$ and $\theta_{0}$ is the true value. Under the usual assumptions for ML, all of these statistics should have a $\chi_{9}^{2}$ distribution asymptotically. In simulations, however, the distribution of these statistics appears quite far from a $\chi^{2}$. To illustrate this fact, in Table 1 we list the size of a number of classical test statistics which, under classical assumptions, should have asymptotic size $5 \%$ or $10 \%$ (for the left and right columns, respectively, based on 2000 simulations). These sizes were generated by calculating the appropriate test statistic in simulation and comparing it to the 95th (or 90th) percentile of a $\chi_{9}^{2}$ distribution. The LM statistic listed in Table 1 is calculated as $L M\left(\theta_{0}\right)=S\left(\theta_{0}\right)^{\prime} I^{-1}\left(\theta_{0}\right) S\left(\theta_{0}\right)$ where $I\left(\theta_{0}\right)=-\ddot{\ell}\left(\theta_{0}\right)$ is the observed information (rather than with J as our LM statistic, $\widetilde{L M}\left(\theta_{0}\right)=S\left(\theta_{0}\right)^{\prime} J^{-1}\left(\theta_{0}\right) S\left(\theta_{0}\right)$ ). Table 1 also lists four variations on the Wald statistic, corresponding to different estimators of the

[^2]asymptotic variance used in $\left(\hat{\theta}-\theta_{0}\right) \hat{V}^{-1}\left(\hat{\theta}-\theta_{0}\right)$. In particular, Wald $(I(\hat{\theta}))$ is the usual Wald statistic which uses the inverse of the observed information, evaluated at $\hat{\theta}$, to estimate the asymptotic variance. Wald $\left(I\left(\theta_{0}\right)\right)$, on the other hand, evaluates the observed information at the true parameter value. Likewise, Wald $(J(\hat{\theta}))$ and Wald $\left(J\left(\theta_{0}\right)\right)$ use $J^{-1}$ as the estimator of the asymptotic variance, calculated at $\hat{\theta}$ and $\theta_{0}$ respectively.

As can be seen in Table 1, the LR statistic is fairly conservative, with type I error less than half the desired asymptotic size. All versions of the Wald statistic which we consider severely overreject. Finally, the usual LM statistic (calculated using the negative hessian) somewhat overrejects at the $5 \%$ level and underrejects at the $10 \%$ level, and its CDF is very poorly approximated by a that of a $\chi_{9}^{2}$. Taken together, these results strongly suggest that the usual approaches to ML estimation and inference are poorly behaved when applied to this model.

### 7.1.2 Behavior of the Information Matrix

Having examined the behavior of the usual ML estimator and tests in this model, we can also look directly at the properties of the information matrix. The embedding which we use to describe weak identification in the exponential family case implies a singular information matrix in the limit. Thus, intuitively, if we think that there are problems of weak identification in this model, we would expect the information matrix to be in some sense close to singular. While we have not formalized this idea, examination of the eigenvalues of Fisher's information matrix $E\left[-\ddot{\ell}\left(\theta_{0}\right)\right]$ (which we calculate by simulation) confirms, consistent with the intuition of weak identification, that the information is "small" in this model, with the smallest eigenvalue (.08) far smaller than the largest (15644.04).

In Section 2 we associated weak identification with the difference between two information measures $A_{T}\left(\theta_{0}\right)$ being large compared to $J_{T}\left(\theta_{0}\right)$. We point out that observed incremental information $J_{T}$ is almost surely positive-definite by construction, while $A_{T}$ is a mean zero random matrix. If $A_{T}$ is negligible compared to $J_{T}$, then the observed information $I_{T}\left(\theta_{0}\right)=J_{T}\left(\theta_{0}\right)-A_{T}\left(\theta_{0}\right)$ is positive definite for majority of realizations. We can check positive-definiteness of $I_{T}\left(\theta_{0}\right)$ directly in simulations. Considering the observed information evaluated at the true value $\left(I_{T}\left(\theta_{0}\right)=-\ddot{\ell}\left(\theta_{0}\right)\right)$, we see that it has at least one negative eigenvalue in over $95 \%$ of simulation draws, and at least two negative
eigenvalues in over $40 \%$ of simulation draws (based on 2000 simulations). While this falls far short of a formal test for weak identification, it is consistent with the idea that weak identification is the source of the bad behavior of ML estimation in this model.

### 7.1.3 LM Test for Full Parameter Vector

We now turn to the weak identification-robust statistics discussed earlier in this paper. We begin by considering the behavior of the the test for the full parameter vector described in section 3. As the reader will recall, this statistic is $\widetilde{L M}(\theta)=S(\theta)^{\prime} J(\theta)^{-1} S(\theta)$, and under appropriate assumptions we have that $\widetilde{L M}\left(\theta_{0}\right) \rightarrow_{d} \chi_{k}^{2}$ under $H_{0}: \theta=\theta_{0}$. In Figure 2, we plot the CDF of the simulated distribution of $\widetilde{L M}\left(\theta_{0}\right)$, together with a $\chi_{9}^{2}$. If we use $\chi_{9}^{2}$ critical values to construct a test based on this statistic, a $5 \%$ test rejects $9.84 \%$ of the time, while a $10 \%$ test rejects $16.68 \%$ of the time: though this shows that the test based on $\widetilde{L M}$ and $\chi_{9}^{2}$ critical values is not exact, the $\chi^{2}$ approximation is far better for $\widetilde{L M}$ than for the usual Wald or LM statistics.

### 7.1.4 Subset Tests

Finally, we simulate tests for subsets of parameters. Specifically, as before we consider a partition of the parameter vector, $\theta=(a, b)$, and consider the problem of testing $b=b_{0}$ without any restrictions on $a$. In this context, we simulate two tests. One is based on the $\widetilde{L M}$ statistic evaluated at $\left(\hat{a}, b_{0}\right)$ for $\hat{a}$ the ML estimator, which we have discussed extensively in this paper. The other is based on $\min _{a} \widetilde{L M}\left(a, b_{0}\right)$, suggested Stock and Wright (2000) for GMM when $a$ is strongly identified. As discussed above, Kleibergen and Mavroeidis (2009) argue that (under under similar assumptions to Stock and Wright) when a is weakly identified the asymptotic distribution of this statistic is dominated by that of a $\chi_{p}^{2}$, where p is the dimension of $b$. For both approaches, and for several subsets of parameters, we simulate the distribution of the statistic and then construct tests using quantiles from the $\chi_{p}^{2}$ distribution as critical values.

We first consider testing the six parameters other than $\alpha, \rho$, and $\delta$, (so we have $a=(\alpha, \rho, \delta)$ and $\left.b=\left(\phi_{x}, \phi_{\pi}, \kappa, \sigma_{a}, \sigma_{u}, \sigma\right)\right)$. The size of $5 \%$ and $10 \%$ tests based on these statistics using asymptotic $\left(\chi_{6}^{2}\right)$ critical values are given in Table 2. As can be seen, while the $\chi_{6}^{2}$ distribution does not provide a perfect approximation to the distribution of either statistic, it is fairly close. Both statistics tend to over-reject, so since the test based on
$\min _{a} \widetilde{L M}\left(a, b_{0}\right)$ is more conservative by construction it performs somewhat better.
We next consider testing the six parameters other than $\phi_{x}, \phi_{\pi}$ and $\kappa$ (so $a=\left(\phi_{x}, \phi_{\pi}, \kappa\right)$, while $\left.b=\left(\alpha, \rho, \delta, \sigma_{a}, \sigma_{u}, \sigma\right)\right)$. The results from this simulation are given in Table 3. Again, the tests over-reject compared to their asymptotic size, and the problem is actually somewhat worse than for the previous subset of parameters considered, although still not too severe.

Finally, we may be interested in testing only one parameter at a time (for example to generate confidence sets). We report results for Taylor Rule parameters $\phi_{x}$ and $\phi_{\pi}$. Based on 1000 simulations, the results for $\widetilde{L M}\left(\hat{\alpha}, \beta_{0}\right)$ (reported in Table 4) are similar to those in the other parameter subsets: the test over-rejects, although not severely. Interestingly, when we consider the minimized statistic, the tests we receive for $\phi_{x}$ and $\phi_{\pi}$ separately appear conservative, rather than over-rejecting as they did for larger subsets.

### 7.2 Nonlinear Weak IV

In section 5 we prove that, provided $\alpha$ is well identified, under appropriate assumptions $\widetilde{L M}\left(\hat{\alpha}, \beta_{0}\right)$ converges to a $\chi_{k_{\beta}}^{2}$ distribution asymptotically, where $k_{\beta}$ is the dimension of $\beta$. As shown in section 6 , for the exponential family model where $\alpha$ is weakly identified but enters linearly we again have that $\widetilde{L M}\left(\hat{\alpha}, \beta_{0}\right)$ converges to a $\chi_{k_{\beta}}^{2}$. To understand the extent to which the result relies on the fact that $\alpha$, the nuisance parameter, enters the expression linearly, we here consider a variation on the usual weak IV model in which $\beta$ enters the equation for $Y$ nonlinearly. In particular, the model is:

$$
\begin{gathered}
Y=\pi\left(\beta^{2} Z^{2}+\beta Z\right)+U \\
X=\pi Z+V
\end{gathered}
$$

with $\beta, \pi$ scalar and $\binom{u_{t}}{v_{t}} \sim \operatorname{iid} N(0, I)$. As usual with weak IV, we take the firststage parameter to zero as the sample size grows, $\pi=\frac{c}{\sqrt{T}}$. The log-likelihood for this model is $\ell(\theta)=$ const $-\frac{1}{2} \sum\left(y_{t}-\pi\left(\beta^{2} z_{t}^{2}+\beta z_{t}\right)\right)^{2}-\frac{1}{2} \sum\left(x_{t}-\pi z_{t}\right)$. As before, let $s$ be the increment of the score vector $S=\frac{\partial}{\partial \theta} \ell$ and let $J$ be its quadratic variation, $J=[S]=\sum s_{t} s_{t}^{\prime}$. We consider testing $H_{0}: \pi=\pi_{0}$ using $\widetilde{L M}\left(\pi_{0}, \hat{\beta}_{M L}\right)$, and are interested in whether this statistic has a $\chi_{1}^{2}$ distribution asymptotically. While we do not have any theoretical results for this case, we have run a number of simulations, which
suggest that a $\chi_{1}^{2}$ is a reasonable approximation to the distribution of this statistic. In particular, we set $\beta=1$ and, $c=.01$, and consider $T=100$ and $T=10,000$. For each value of T, we simulate 10,000 Monte-Carlo draws, and calculate the size of asymptotic $5 \%$ and $10 \%$ tests (using critical values based on a $\chi_{1}^{2}$ ) for sample sizes 100 and 10,000 , which we report in Table 5 . We also plot the CDF of $\widetilde{L M}\left(\pi_{0}, \hat{\beta}_{M L}\right)$, together with that of a $\chi_{1}^{2}$, in Figure 3. These simulation results show that the distribution of $\widetilde{L M}\left(\pi_{0}, \hat{\beta}_{M L}\right)$ is close to a $\chi_{1}^{2}$ in this model, suggesting that it may be possible to extend our theoretical results to this context.

### 7.3 Weak IV for Two Endogenous Variables

Another possible application of our proposed statistics is to the case of weak IV with more than one endogenous variable. Here, we consider the case of IV with two endogenous variables, using the model

$$
Y=Z \pi \beta+U
$$

$$
\begin{gathered}
X=Z \pi+V \\
\text { where } \pi=\left[\begin{array}{cc}
\pi_{11} & \pi_{12} \\
\pi_{21} & \pi_{22}
\end{array}\right], \beta=\left[\begin{array}{l}
\beta_{1} \\
\beta_{2}
\end{array}\right], V=\left[\begin{array}{cc}
v_{11} & v_{12} \\
\vdots & \vdots \\
v_{T 1} & v_{T 2}
\end{array}\right] \text { and }\left[\begin{array}{c}
u_{t} \\
v_{1 t} \\
v_{2 t}
\end{array}\right] \sim N\left(\left[\begin{array}{l}
0 \\
0 \\
0
\end{array}\right], \Sigma\right) .
\end{gathered}
$$

We consider the problem of testing $\beta_{1}=\beta_{1}^{0}$ without imposing any restrictions on $\beta_{2}$ or $\pi$. To do this, we calculate $\widetilde{L M}\left(\beta_{1}^{0}, \hat{\beta}_{2}, \hat{\pi}\right)$, where $\left(\hat{\beta}_{2}, \hat{\pi}\right)$ is the restricted ML estimator of $\left(\beta_{2}, \pi\right)$. As for nonlinear weak IV we have no theoretical results on this model, but have simulated this statistic for $\pi=\frac{1}{\sqrt{T}}\left[\begin{array}{cc}1 & \frac{1}{3} \\ \frac{1}{3} & 1\end{array}\right], \Sigma=\left[\begin{array}{ccc}1 & .4 & .7 \\ .4 & 1 & .5 \\ .7 & .5 & 1\end{array}\right]$, and $T$ (total number of observations) equal to 100 and 10,000 (in both cases using 10,000 MonteCarlo draws). ${ }^{6}$ In Table 6 we report the size of what would be asymptotic $5 \%$ and $10 \%$ tests if the $\widetilde{L M}$ statistic converged to a $\chi_{1}^{2}$ asymptotically, while in Figure 10 we plot the distribution of the statistic for $T=100$. At both sample sizes (and for a number of other parameter values not shown), the statistic is dominated by a $\chi_{1}^{2}$, and as a result

[^3]the statistic under-rejects. These results suggest that it may be possible to extend our results to show that the $\widetilde{L M}$ statistic is dominated by a $\chi_{k}^{2}$ in the case of weak IV with more than one endogenous variable.

## 8 References

D.W.K. Andrews, and Xu Chen, 2009, Estimation and Inference with Weak, semi-strong and strong identification, unpublished manuscript
D.W.K. Andrews , M. Moreira and J. Stock, 2006, "Optimal Two-Sided Invariant Similar Tests for Instrumental Variables Regression," Econometrica, 74(3), 715-752.
M. Akritas (1988),"An Asymptotic Deviration of Neyman's $C(\alpha)$ Test", Statistics and Probablity Letters, 6, pp.363-367
O.E. Barndorff-Nielsen and M. Sorensen (1991) "Information quantities in non-classical settings", Computational Statiatics and Data Analysis, 12, pp.143-158
I.V. Basawa and H.L. Koul "Asymptotic tests of composite hypotheses for non-ergodic type stochastic processes" Stochastic Processes and their Applications 9(1979) 291-305

Bhat B.R. "On the Method of Maximum-Likelihood for Dependent Observations" Journal of the Royal Statistical Society. Series B (Methodological), Vol. 36, No. 1(1974), pp. 48-53

Canova, Fabio, and Luca Sala (2009), "Back to Square One: Identification Issues in DSGE Models," Journal of Monetary Economics, 56, 431-449.
D. Chant "On Asymptotic Tests of Composite Hypotheses in Nonstandard Conditions" : Biometrika, Vol. 61, No. 2 (Aug., 1974), pp. 291-298
S. Chaudhuri and E. Zivot (2009)"A New Method of Projection-based Inference in GMM with Weakly Identified Nuisance Parameters", Unpublished Manuscript
R. Clarida, J. Gali and M. Gertler (1999), "The Science of Monetary Policy: A New Keynesian Perspective" Journal of Economic Literature, 37 (4), pp 1661-1707
M. J. Crowder "Maximum Likelihood Estimation for Dependent Observations" Journal of the Royal Statistical Society. Series B (Methodological), Vol. 38, No. 1(1976), pp. 45-53
J.M. Dufour and M. Taamouti (2005) "Projection-Based Statistical Inference in Linear Structural Models with Possibly Weak Instruments", Econometrica, 73(4), pp. 1351-1365
J. Fernéndez-Villaverde(2010) " The Econometrics of DSGE Models." SERIES: Journal of the Spanish Economic Association 1, 3-49

Fernéndez-Villaverde J., J. F. Rubio-Ramírez, T. J. Sargent, and M. W. Watson, (2007) "ABCs (and Ds) of Understanding VARs," American Economic Review, vol. 97(3), pp. 1021-1026.
P. Guerron-Quintana, A. Inoue and L. Kilian "Frequentist Inference in Weakly Identified DSGE Models", unpublished manuscript, 2009
P. Hall and C.C. Heyde (1980)"Martingale Limit Theory and its Application", Academic Press
R.D.H. Heijmans and J.R. Magnus (1986)"Consistent Maximum-Likelihood Estimation with Dependent Observations," Journal of Econometrics, 32, pp.253-285
B. F. Ingram , N. R. Kocherlakota and N. E. Savin (1994) "Explaining business cycles: A multiple-shock approach", Journal of Monetary Economics, 34, Issue 3, pp. 415-428
P. N. Ireland "Technology Shocks in the New Keynesian Model" The Review of Economics and Statistics, Vol. 86, No. 4 (2004), pp. 923-936

Iskrev, Nikolay (2008), "How Much Do We Learn from the Estimation of DSGE Models? A Case Study of Identification Issues in a New Keynesian Business Cycle Model," unpublished manuscript, University of Michigan.
P. Jeganathan (1995)"Some Aspect of Asymptotic Theory with Applications to Time Series Models", Econometric Theory, 11(5), 818-887
F. Kleibergen (2005) "Testing Parameters in GMM without Assuming that They Are Identified" Econometrica,73, 1103-1123.
F. Kleibergen and S. Mavroeidis (2009)"Weak Instrument Robust Tests in GMM and the New Keynesian Phillips Curve", Journal of Business \& Economic Statistics,27, 293311.
F. Kleibergen and S. Mavroeidis (2009)"Inference on subsets of parameters in GMM without assuming identification", unpublished manuscript
S. Kocherlakota and K. Kocherlakota "Neyman's $C(\alpha)$ test and Rao's efficient score test for composite hypotheses" Statistics E Brobability Letters 11 (1991) 491-493

Komunjer, Ivana, and Serena Ng (2009), "Dynamic Identification of DSGE Models," unpublished manuscript, University of California, San Diego, and Columbia University.

Le Cam L. and G.L. Yang (2000)"Asymptotics in Statistics: Some Basic Concepts", Springer Series in Statistics
J. Lindé(2005)"Estimating New-Keynesian Phillips curves: A full information maximum likelihood approach" Journal of Monetary Economics 52(6), Pages 1135-1149
R. Liptser and A. Shiryayev (1989)"Theory of Martingales," Springer
E. R. McGrattan, R. Rogerson and R. Wright(1997) "An Equilibrium Model of the business cycle with household production and fiscal policy", International Economic Review, 38(2) pp.267-290.

Moon, H.R., and Frank Schorfheide (2009), "Bayesian and Frequentist Inference in Partially Identified Models," unpublished manuscript, University of Southern California and University of Pennsylvania.
S.D. Silvey (1961)"A Note on Maximum-Likelihood in the Case of Dependent Random Variables" Journal of the Royal Statistical Society. Series B (Methodological), Vol. 23, No. 2(1961), pp. 444-452

Staiger, D., and Stock, J. H. (1997), "Instrumental Variables Regression With Weak Instruments," Econometrica, 65, 557-586.

Stock, J. H., and Wright, J. H. (2000), "GMM With Weak Identification," Econometrica, 68, 1055-1096.

Stock, J.H., Wright, J. H. and M. Yogo (2002)"A Survey of Weak Instruments and Weak Identification in Generalized Method of Moments" Journal of Business \& Economic statistics, 20(4), pp. 518-529
H. White (1980) "A Heteroskedasticity-Consistent Covariance Matrix Estimator and a Direct Test for Heteroskedasticity" Econometrica, 48, pp. 817-838.
H. White (1982) "Maximum Likelihood Estimation in Misspecified Models" Econometrica, 50, pp. 1-25.

## 9 Appendix with Proofs

As before, we assume that the normalizing matrix $K_{M, T}$ is diagonal. In order to simplify notation when dealing with objects in three or more dimensions, we will somewhat abuse our notation by using $K_{M, T}$ as if it were a sequence of constants. For example, $K_{\alpha, T}^{2} K_{\beta, T} \Lambda_{\alpha \alpha \beta} \rightarrow^{p} 0$ means that for all indexes $i, j, n$ we have $K_{\alpha_{i}, T} K_{\alpha_{j}, T} K_{\beta_{n}, T} \Lambda_{\alpha_{i} \alpha_{j} \beta_{n}} \rightarrow^{p}$ 0. In the Taylor expansions used in the proof for Theorem 2, the expansion is assumed to be for each entry of the expanded matrix. In addition, there is often a silent summation over the indices in $\alpha$, as will be clear in context.

## Proof of Lemma 1

The proof follows closely the argument of Bhat (1974), starting with the Taylor expansion:

$$
0=S_{\alpha}\left(\hat{\alpha}, \beta_{0}\right)=S_{\alpha}^{0}-I_{\alpha \alpha}^{0}\left(\hat{\alpha}-\alpha_{0}\right)-\left(I_{\alpha \alpha}\left(\alpha^{*}, \beta_{0}\right)-I_{\alpha \alpha}^{0}\right)\left(\hat{\alpha}-\alpha_{0}\right)
$$

where $\alpha^{*}$ is a convex combination of $\hat{\alpha}$ and $\alpha_{0}$. As usual, we may consider different $\alpha^{*}$ for different rows of $I_{\alpha \alpha}$. Assumption 2(b) helps to control the last term of this expansion, while Assumption 2(a) allows us to substitute $J_{\alpha \alpha, T}$ for $I_{\alpha \alpha, T}$ in the second term. Assumption 1 gives the CLT for $K_{\alpha, T} S_{\alpha, T}$.

Lemma 4 Let $M_{T}=\sum_{t=1}^{T} m_{t}$ be a multi-dimensional martingale with respect to sigmafield $\mathcal{F}_{t}$, and let $[X]_{t}$ be its quadratic variation. Assume that there is a sequence of diagonal matrices $K_{T}$ such that
(a) $K_{T}[M]_{T} K_{T} \rightarrow \Sigma$, where $\Sigma$ is some finite matrix;
(b) for any $\varepsilon>0$ we have that $\sum_{t} E\left(K_{T}\left|m_{t}\right| \mathbb{I}\left\{K_{T}\left|m_{t}\right|>\varepsilon\right\}\right) \rightarrow 0$ as $T \rightarrow \infty$.

Let $m_{i, t}$ be $i$-th component of $m_{t}$, and $K_{i, T}$ the $i$-th diagonal element of $K_{T}$. For any $i, j, l$ :

$$
K_{i, T} K_{j, T} K_{l, T} \sum_{t=1}^{T} m_{i, t} m_{j, t} m_{l, t} \rightarrow^{p} 0
$$

## Proof of Lemma 4

Take any $\varepsilon>0$.

$$
\begin{array}{r}
\left|K_{i, T} K_{j, T} K_{l, T} \sum_{t=1}^{T} m_{i, t} m_{j, t} m_{l, t}\right| \leq \max _{t}\left|K_{i, T} m_{i, t}\right|\left(K_{j, T} K_{l, T} \sum_{t=1}^{T} m_{j, t} m_{l, t}\right) \leq \\
\leq \max _{t}\left|K_{i, T} m_{i, t}\right|\left(K_{j, T} K_{l, T}\left[M_{j}, M_{l}\right]_{T}\right)
\end{array}
$$

Condition (a) implies that $K_{j, T} K_{l, T}\left[M_{j}, M_{l}\right]_{T} \rightarrow^{p} \Sigma_{j, l}$ is bounded in probability.

$$
\begin{array}{r}
E\left(\max _{t}\left|K_{i, T} m_{i, t}\right|\right) \leq \varepsilon+E\left(K_{i, T} \max _{t}\left|m_{i, t}\right| \mathbb{I}\left\{\left|K_{i, T} m_{i, t}\right|>\varepsilon\right\}\right) \leq \\
\leq \varepsilon+\sum_{t} E\left(K_{i, T}\left|m_{i, t}\right| \mathbb{I}\left\{\left|K_{i, T} m_{i, t}\right|>\varepsilon\right\}\right)
\end{array}
$$

the last term converges to 0 according to condition (b).

## Proof of Lemma 2

Notice first that

$$
\begin{equation*}
-\frac{\partial}{\partial \alpha} I_{\alpha \alpha}\left(\theta_{0}\right)=-3\left[A_{\alpha \alpha}, S_{\alpha}\right]_{T}+2 \sum_{t=1}^{T} s_{\alpha, t}^{3}+\Lambda_{\alpha \alpha \alpha} \tag{15}
\end{equation*}
$$

where $\Lambda_{\alpha \alpha \alpha}$ is a martingale defined in (9), and $[\cdot, \cdot]_{T}$ is a joint quadratic variation of any two martingales.

Denote by $f_{t}=f\left(x_{t} \mid X_{t-1} ; \theta\right)$ the (valid) pdf, while $f_{\alpha, t}, f_{\alpha \alpha, t}$ etc. are its partial derivatives with respect to $\alpha$. Notice that the increments of martingales $S_{\alpha, T}, A_{\alpha \alpha, T}$ and $\Lambda_{\alpha \alpha \alpha}$ are $s_{\alpha, t}=\frac{f_{\alpha, t}}{f_{t}}, a_{\alpha \alpha, t}=\frac{f_{\alpha \alpha, t}}{f_{t}}$, and $\lambda_{\alpha \alpha \alpha, t}=\frac{f_{\alpha \alpha \alpha, t}}{f_{t}}$ respectively. By definition

$$
-\frac{\partial}{\partial \alpha} I_{\alpha \alpha}=\frac{\partial^{3}}{\partial \alpha^{3}} \sum_{t}^{T} \log f_{t}=\sum_{t}^{T} \frac{f_{\alpha \alpha \alpha, t}}{f_{t}}-3 \sum_{t}^{T} \frac{f_{\alpha \alpha, t}}{f_{t}} \frac{f_{\alpha, t}}{f_{t}}+2 \sum_{t}^{T}\left(\frac{f_{\alpha, t}}{f_{t}}\right)^{3}
$$

so (15) follows.
Given Assumptions 1 and (2a) we have that $K_{\alpha, T} I_{\alpha \alpha} K_{\alpha, T} \rightarrow \Sigma$. Now consider the quantity of interest from Assumption (2b)

$$
\left|K_{\alpha, T}^{2}\left(I_{\alpha \alpha}\left(\alpha_{1}, \beta_{0}\right)-I_{\alpha \alpha}^{0}\right)\right|=K_{\alpha, T}^{2}\left|\frac{\partial}{\partial \alpha} I_{\alpha \alpha}\right|\left|\alpha_{1}-\alpha_{0}\right| \leq \delta K_{\alpha, T}^{3}\left|\frac{\partial}{\partial \alpha} I_{\alpha \alpha}\right| .
$$

It suffices to show that $K_{\alpha, T}^{3}\left|\frac{\partial}{\partial \alpha} I_{\alpha \alpha}\right| \rightarrow^{p} 0$, using identity (15). Assumption (2b') implies that the last term converges to zero in probability. Lemma 4 implies that the second term is negligible. And finally, Assumption (2a) gives us that the first term also converges to zero in probability.

## Proof of Theorem 2

We denote by super-script 0 quantities evaluated at $\theta_{0}=\left(\alpha_{0}, \beta_{0}\right)$. According to martingale CLT, Assumption 3 implies that

$$
\begin{equation*}
\left(K_{\alpha, T} S_{\alpha}^{0}, K_{\beta, T} S_{\beta}^{0}, K_{\alpha \beta, T} A_{\alpha \beta}^{0}\right) \Rightarrow\left(\xi_{\alpha}, \xi_{\beta}, \xi_{\alpha \beta}\right), \tag{16}
\end{equation*}
$$

where $\xi$ 's are jointly normal with variance matrix $\Sigma_{M}$.

We Taylor expand $S_{\beta}\left(\hat{\alpha}, \beta_{0}\right)$, keeping in mind that $I_{\beta \alpha}^{0}=-\frac{\partial^{2}}{\partial \alpha \partial \beta} \ell\left(\alpha_{0}, \beta_{0}\right)$, and receive

$$
K_{\beta, T} S_{\beta}\left(\hat{\alpha}, \beta_{0}\right)=K_{\beta, T} S_{\beta}^{0}-K_{\beta, T} I_{\beta \alpha}^{0}\left(\hat{\alpha}-\alpha_{0}\right)+\frac{1}{2} K_{\beta, T}\left(\hat{\alpha}-\alpha_{0}\right)^{\prime}\left(I_{\alpha \alpha \beta}^{0}\right)\left(\hat{\alpha}-\alpha_{0}\right)+\tilde{R}
$$

with residual

$$
\tilde{R}=K_{\beta, T} \frac{1}{2}\left(\hat{\alpha}-\alpha_{0}\right)^{\prime}\left(I_{\alpha \alpha \beta}^{*}-I_{\alpha \alpha \beta}^{0}\right)\left(\hat{\alpha}-\alpha_{0}\right)
$$

where $I_{\alpha \alpha \beta}^{0}=\frac{\partial^{3}}{\partial \alpha^{2} \partial \beta} \ell\left(\alpha_{0}, \beta_{0}\right), I_{\alpha \alpha \beta}^{*}=\frac{\partial^{3}}{\partial \alpha^{2} \partial \beta} \ell\left(\alpha^{*}, \beta_{0}\right)$ and $\alpha^{*}$ is again a point between $\hat{\alpha}$ and $\alpha_{0}$. From Lemma 1 we have that $\hat{\alpha}$ is $K_{\alpha, T}$-consistent. As a result, Assumption 4 (c) makes the Taylor residual negligible:

$$
K_{\beta, T} S_{\beta}\left(\hat{\alpha}, \beta_{0}\right)=K_{\beta, T} S_{\beta}^{0}-K_{\beta, T} I_{\beta \alpha}^{0}\left(\hat{\alpha}-\alpha_{0}\right)+K_{\beta, T} \frac{1}{2}\left(\hat{\alpha}-\alpha_{0}\right)^{\prime}\left(I_{\alpha \alpha \beta}^{0}\right)\left(\hat{\alpha}-\alpha_{0}\right)+o_{p}(1) .
$$

We plug asymptotic statement (8) into this equation and get
$K_{\beta, T} S_{\beta}\left(\hat{\alpha}, \beta_{0}\right)=K_{\beta, T} S_{\beta}^{0}-K_{\beta, T} I_{\beta \alpha}^{0}\left(I_{\alpha \alpha}^{0}\right)^{-1} S_{\alpha}^{0}+\frac{1}{2} K_{\beta, T} S_{\alpha}^{0^{\prime}}\left(I_{\alpha \alpha}^{0}\right)^{-1}\left(I_{\alpha \alpha \beta}^{0}\right)\left(I_{\alpha \alpha}^{0}\right)^{-1} S_{\alpha}^{0}+o_{p}(1)$.

Reacall that by definition $I_{\beta \alpha}^{0}=J_{\beta \alpha}^{0}-A_{\beta \alpha}^{0}$. We use this substitution in the equation above, and receive:

$$
\begin{align*}
K_{\beta, T} S_{\beta}\left(\hat{\alpha}, \beta_{0}\right)=K_{\beta, T} S_{\beta}^{0}- & K_{\beta, T} J_{\beta \alpha}^{0}\left(I_{\alpha \alpha}^{0}\right)^{-1} S_{\alpha}^{0}+K_{\beta, T} A_{\beta \alpha}^{0}\left(I_{\alpha \alpha}^{0}\right)^{-1} S_{\alpha}^{0}+ \\
& +\frac{1}{2} K_{\beta, T} S_{\alpha}^{0^{\prime}}\left(I_{\alpha \alpha}^{0}\right)^{-1}\left(I_{\alpha \alpha \beta}^{0}\right)\left(I_{\alpha \alpha}^{0}\right)^{-1} S_{\alpha}^{0}+o_{p}(1) . \tag{17}
\end{align*}
$$

One can notice that we have the following informational equality:

$$
I_{\alpha \alpha \beta}^{0}=-\left[A_{\alpha \alpha}^{0}, S_{\beta}^{0}\right]-2\left[A_{\alpha \beta}^{0}, S_{\alpha}^{0}\right]+2 \sum_{t=1}^{T} s_{\alpha, t}^{2} s_{\beta, t}+\Lambda_{\alpha \alpha \beta} .
$$

It can be obtained in the same manner as (15). Assumption $4(\mathrm{~b})$ implies that $K_{\beta, T} K_{\alpha, T}^{2} \Lambda_{\alpha \alpha \beta} \rightarrow{ }^{p}$ 0. Assumption 2(a) and Assumption 3 together imply that $\frac{K_{\alpha, T}^{2}}{K_{\alpha \alpha, T}} \rightarrow 0$. Using Assumption 2(a) and Lemma 4, we notice that

$$
\begin{equation*}
K_{\beta, T} K_{\alpha, T}^{2} I_{\alpha \alpha \beta}^{0}=-2 K_{\beta, T} K_{\alpha, T}^{2}\left[A_{\alpha \beta}^{0}, S_{\alpha}^{0}\right]+o_{p}(1) \tag{18}
\end{equation*}
$$

According to Assumption 4 (a), $2 K_{\beta, T} K_{\alpha, T}^{2}\left[A_{\alpha \beta}^{0}, S_{\alpha}^{0}\right]$ is asymptotically bounded so $K_{\beta, T} K_{\alpha, T}^{2} I_{\alpha \alpha \beta}^{0}=$ $O_{p}(1)$. According to Assumption 2(a) $K_{\alpha, T} I_{\alpha \alpha}^{0} K_{\alpha, T}=K_{\alpha, T} J_{\alpha \alpha} K_{\alpha, T}+o_{p}(1)$, so Assump-
tion 4 (a) implies that $K_{\alpha, T} A_{\alpha \beta} K_{\beta, T}$ is bounded. Taken together, these statements imply that we can substitute $J_{\alpha \alpha}^{0}$ for $I_{\alpha \alpha}^{0}$ everywhere in (17). Doing so gives us:

$$
\begin{gathered}
K_{\beta, T} S_{\beta}\left(\hat{\alpha}, \beta_{0}\right)=K_{\beta, T} S_{\beta}^{0}-K_{\beta, T} J_{\beta \alpha}^{0}\left(J_{\alpha \alpha}^{0}\right)^{-1} S_{\alpha}^{0}+K_{\beta, T} A_{\alpha \beta}^{0}\left(J_{\alpha \alpha}^{0}\right)^{-1} S_{\alpha}^{0}+ \\
+\frac{1}{2} K_{\beta, T} S_{\alpha}^{0^{\prime}}\left(J_{\alpha \alpha}^{0}\right)^{-1}\left(I_{\alpha \alpha \beta}^{0}\right)\left(J_{\alpha \alpha}^{0}\right)^{-1} S_{\alpha}^{0}+o_{p}(1)
\end{gathered}
$$

Or

$$
\begin{equation*}
K_{\beta, T} S_{\beta}\left(\hat{\alpha}, \beta_{0}\right)=K_{\beta, T} S_{\beta}^{0}-K_{\beta, T} J_{\beta \alpha}^{0}\left(J_{\alpha \alpha}^{0}\right)^{-1} S_{\alpha}^{0}+D\left(J_{\alpha \alpha}^{0} K_{\alpha, T}\right)^{-1} S_{\alpha}^{0}+o_{p}(1) \tag{19}
\end{equation*}
$$

where

$$
D=K_{\alpha, T} K_{\beta, T} A_{\alpha \beta}^{0}+\frac{1}{2} K_{\alpha, T} K_{\beta, T}\left(I_{\alpha \alpha \beta}^{0}\right)\left(J_{\alpha \alpha}^{0}\right)^{-1} S_{\alpha}^{0^{\prime}}
$$

Notice that $D$ is asymptotically normal (though it may have zero variance, i.e. it may converge to zero) and asymptotically independent of $K_{\alpha, T} S_{\alpha}^{0}$. Indeed, using (??) we have:

$$
\begin{array}{r}
D=\frac{K_{\alpha, T} K_{\beta, T}}{K_{\alpha \beta, T}} K_{\alpha \beta, T} A_{\alpha \beta}^{0}+\left(K_{\alpha, T}^{2} K_{\beta, T}\left[A_{\alpha \beta}^{0}, S_{\alpha}^{0}\right]+o_{p}(1)\right)\left(K_{\alpha, T} J_{\alpha \alpha}^{0} K_{\alpha, T}\right)^{-1} K_{\alpha, T} S_{\alpha}^{0^{\prime}}= \\
=\frac{K_{\alpha, T} K_{\beta, T}}{K_{\alpha \beta, T}}\left(K_{\alpha \beta, T} A_{\alpha \beta}^{0}-\left(K_{\alpha, T} K_{\alpha \beta, T}\left[A_{\alpha \beta}^{0}, S_{\alpha}^{0}\right]\right)\left(K_{\alpha, T} J_{\alpha \alpha}^{0} K_{\alpha, T}\right)^{-1} K_{\alpha, T} S_{\alpha}^{0^{\prime}}\right)+o_{p}(1) \Rightarrow \\
\Rightarrow C\left(\xi_{\alpha \beta}-\frac{\operatorname{cov}\left(\xi_{\alpha \beta}, \xi_{\alpha}\right)}{\operatorname{Var}\left(\xi_{\alpha}\right)} \xi_{\alpha}\right)
\end{array}
$$

where variables $\left(\xi_{\alpha}^{\prime}, \xi_{\alpha \beta}^{\prime}\right)=\lim \left(K_{\alpha, T} S_{\alpha}^{0^{\prime}}, K_{\alpha \beta, T} A_{\alpha \beta}^{0^{\prime}}\right)$ are as described at the beginning of the proof.

Plugging the last statement and (16) into equation (19) we have:

$$
K_{\beta, T} S_{\beta}\left(\hat{\alpha}, \beta_{0}\right) \Rightarrow \xi_{\beta}-\frac{\operatorname{cov}\left(\xi_{\beta}, \xi_{\alpha}\right)}{\operatorname{Var}\left(\xi_{\alpha}\right)} \xi_{\alpha}+C\left(\xi_{\alpha \beta}-\frac{\operatorname{cov}\left(\xi_{\alpha \beta}, \xi_{\alpha}\right)}{\operatorname{Var}\left(\xi_{\alpha}\right)} \xi_{\alpha}\right) \frac{\xi_{\alpha}}{\operatorname{Var}\left(\xi_{\alpha}\right)}
$$

Conditional on $\xi_{\alpha}, K_{\beta, T} S_{\beta}\left(\hat{\alpha}, \beta_{0}\right)$ is asymptotically normal with mean zero and condi-
tional variance

$$
\begin{align*}
& \quad \lim \operatorname{Var}\left(K_{\beta, T} S_{\beta}\left(\hat{\alpha}, \beta_{0}\right) \mid K_{\alpha, T} S_{\alpha}^{0}=\xi_{\alpha}\right)=\left(\operatorname{Var}\left(\xi_{\beta}\right)-\frac{\operatorname{cov}^{2}\left(\xi_{\beta}, \xi_{\alpha}\right)}{\operatorname{Var}\left(\xi_{\alpha}\right)}\right)+ \\
& +2 C \frac{\xi_{\alpha}}{\operatorname{Var}\left(\xi_{\alpha}\right)}\left(\operatorname{cov}\left(\xi_{\alpha \beta}, \xi_{\beta}\right)-\frac{\operatorname{cov}\left(\xi_{\alpha \beta}, \xi_{\alpha}\right) \operatorname{cov}\left(\xi_{\beta}, \xi_{\alpha}\right)}{\operatorname{Var}\left(\xi_{\alpha}\right)}\right)+ \\
& \quad+\left(C \frac{\xi_{\alpha}}{\operatorname{Var}\left(\xi_{\alpha}\right)}\right)^{2}\left(\operatorname{Var}\left(\xi_{\alpha \beta}\right)-\frac{\operatorname{cov}^{2}\left(\xi_{\alpha \beta}, \xi_{\alpha}\right)}{\operatorname{Var}\left(\xi_{\alpha}\right)}\right) \tag{20}
\end{align*}
$$

Now we turn to the inverse variance term in formula (7) for $\widetilde{L M}\left(\beta_{0}\right),\left.\left(J_{\beta \beta}-J_{\beta \alpha} J_{\alpha \alpha}^{-1} J_{\beta \alpha}^{\prime}\right)\right|_{\left(\hat{\alpha}, \beta_{0}\right)}$. Below we prove the following lemma:

Lemma 5 Under Assumptions of Theorem 1 we have:
(a) $K_{\beta, T}^{2} J_{\beta \beta}\left(\hat{\alpha}, \beta_{0}\right) \Rightarrow \operatorname{Var}\left(\xi_{\beta}\right)+2 C \operatorname{cov}\left(\xi_{\alpha \beta}, \xi_{\beta}\right) \frac{\xi_{\alpha}}{\operatorname{Var}\left(\xi_{\alpha}\right)}+C^{2} \operatorname{Var}\left(\xi_{\alpha \beta}\right)\left(\frac{\xi_{\alpha}}{\operatorname{Var}\left(\xi_{\alpha}\right)}\right)^{2}$
(b) $K_{\alpha, T} K_{\beta, T} J_{\alpha \beta}\left(\hat{\alpha}, \beta_{0}\right) \Rightarrow \operatorname{cov}\left(\xi_{\alpha}, \xi_{\beta}\right)+C \cdot \operatorname{cov}\left(\xi_{\alpha \beta}, \xi_{\alpha}\right) \frac{\xi_{\alpha}}{\operatorname{Var}\left(\xi_{\alpha}\right)}$
(c) $K_{\alpha, T}^{2} J_{\alpha \alpha}\left(\hat{\alpha}, \beta_{0}\right) \rightarrow^{p} \operatorname{Var}\left(\xi_{\alpha}\right)$

Lemma 5 implies that

$$
\begin{aligned}
& \left.K_{\beta, T}^{2}\left(J_{\beta \beta}-J_{\beta \alpha} J_{\alpha \alpha}^{-1} J_{\beta \alpha}^{\prime}\right)\right|_{\left(\hat{\alpha}, \beta_{0}\right)} \Rightarrow \\
& \Rightarrow \operatorname{Var}\left(\xi_{\beta}\right)+2 \operatorname{Cov}\left(\xi_{\alpha \beta}, \xi_{\beta}\right) \frac{\xi_{\alpha}}{\operatorname{Var}\left(\xi_{\alpha}\right)}+C^{2} \operatorname{Var}\left(\xi_{\alpha \beta}\right)\left(\frac{\xi_{\alpha}}{\operatorname{Var}\left(\xi_{\alpha}\right)}\right)^{2}- \\
& -\left(\operatorname{cov}\left(\xi_{\alpha}, \xi_{\beta}\right)+C \cdot \operatorname{cov}\left(\xi_{\alpha \beta}, \xi_{\alpha}\right) \frac{\xi_{\alpha}}{\operatorname{Var}\left(\xi_{\alpha}\right)}\right) \frac{1}{\operatorname{Var}\left(\xi_{\alpha}\right)}\left(\operatorname{cov}\left(\xi_{\alpha}, \xi_{\beta}\right)+C \cdot \operatorname{cov}\left(\xi_{\alpha \beta}, \xi_{\alpha}\right) \frac{\xi_{\alpha}}{\operatorname{Var}\left(\xi_{\alpha}\right)}\right) .
\end{aligned}
$$

Note that the last expression is the same as (20). That is, $\left.K_{\beta, T}^{2}\left(J_{\beta \beta}-J_{\beta \alpha} J_{\alpha \alpha}^{-1} J_{\beta \alpha}^{\prime}\right)\right|_{\left(\hat{\alpha}, \beta_{0}\right)}$ is asymptotically equal to the asymptotic variance of $K_{\beta, T} S_{\beta}\left(\hat{\alpha}, \beta_{0}\right)$ conditional on $\xi_{\alpha}$. As a result statistic $\widetilde{L M}\left(\beta_{0}\right)$, conditional on $\xi_{\alpha}$, is distributed $\chi_{k_{\beta}}^{2}$ asymptotically and thus is asymptotically $\chi_{k_{\beta}}^{2}$ unconditionally as well. This completes the proof of Theorem 2.

## Proof of Lemma 5

(a) We can Taylor expand $J_{\beta \beta}\left(\hat{\alpha}, \beta_{0}\right)$ as:

$$
\begin{equation*}
J_{\beta \beta}\left(\hat{\alpha}, \beta_{0}\right)=J_{\beta \beta}^{0}+\frac{\partial}{\partial \alpha} J_{\beta \beta}^{0}\left(\hat{\alpha}-\alpha_{0}\right)+\frac{1}{2} \frac{\partial^{2}}{\partial \alpha^{2}} J_{\beta \beta}^{0}\left(\hat{\alpha}-\alpha_{0}\right)^{2}+R, \tag{21}
\end{equation*}
$$

where $R=\frac{\partial^{3}}{\partial \alpha^{3}} J_{\beta \beta}^{0}\left(\alpha^{*}-\alpha_{0}\right)^{2}$. Consider the first term of the Taylor expansion above:

$$
\frac{\partial}{\partial \alpha} J_{\beta \beta}=\frac{\partial}{\partial \alpha} \sum_{t} s_{\beta, t} s_{\beta, t}=-2 \sum i_{\alpha \beta, t} s_{\beta, t}=2 \sum a_{\alpha \beta, t} s_{\beta, t}-2 \sum s_{\alpha, t} s_{\beta, t}^{2} .
$$

Using the statement of Lemma 4 and Assumption 4 (a) we have

$$
\begin{align*}
& K_{\alpha, T} K_{\beta, T}^{2} \frac{\partial}{\partial \alpha} J_{\beta \beta}=2 \frac{K_{\alpha, T} K_{\beta, T}}{K_{\alpha \beta, T}} K_{\alpha \beta, T} K_{\beta, T} \sum a_{\alpha \beta, t} s_{\beta, t}-2 K_{\alpha, T} K_{\beta, T}^{2} \sum s_{\alpha, t} s_{\beta, t}^{2} \rightarrow^{p} \\
& \rightarrow^{p} 2 \operatorname{Cov}\left(\xi_{\alpha \beta}, \xi_{\beta}\right) \tag{22}
\end{align*}
$$

Now let us consider the second derivative of $J_{\beta \beta}$ :

$$
\begin{array}{r}
\frac{\partial^{2}}{\partial \alpha^{2}} J_{\beta \beta}=2 \frac{\partial}{\partial \alpha} \sum a_{\alpha \beta, t} s_{\beta, t}-2 \frac{\partial}{\partial \alpha} \sum s_{\alpha, t} s_{\beta, t}^{2}= \\
=2 \sum \lambda_{\alpha \alpha \beta, t} s_{\beta, t}+2 \sum a_{\alpha \beta, t}^{2}-8 \sum a_{\alpha \beta, t} s_{\beta, t} s_{\alpha, t}-2 \sum a_{\alpha \alpha, t} s_{\beta, t}^{2}+6 \sum_{t} s_{\alpha, t}^{2} s_{\beta, t}^{2} .
\end{array}
$$

We is interested in the limit of $K_{\alpha, T}^{2} K_{\beta, T}^{2} \frac{\partial^{2}}{\partial \alpha^{2}} J_{\beta \beta}\left(\hat{\alpha}, \beta_{0}\right)$. According to Lemma 4, $K_{\alpha, T}^{2} K_{\beta, T}^{2} \sum a_{\alpha \beta, t} s_{\beta, t} s_{\alpha, t} \rightarrow^{p} 0$ and $K_{\alpha, T}^{2} K_{\beta, T}^{2} \sum a_{\alpha \alpha, t} s_{\beta, t}^{2} \rightarrow^{p} 0$. Analogously we can show that $K_{\alpha, T}^{2} K_{\beta, T}^{2} \sum_{t} s_{\alpha, t}^{2} s_{\beta, t}^{2} \rightarrow^{p} 0$. Assumption 4 (b) implies that $K_{\alpha, T}^{2} K_{\beta, T}^{2} \sum \lambda_{\alpha \alpha \beta, t} s_{\beta, t} \rightarrow^{p}$ 0. Finally using Assumption 3(b) we get

$$
\begin{equation*}
K_{\alpha, T}^{2} K_{\beta, T}^{2} \frac{\partial^{2}}{\partial \alpha^{2}} J_{\beta \beta} \rightarrow^{p} 2 C^{2} \operatorname{Var}\left(\xi_{\alpha \beta}\right) . \tag{23}
\end{equation*}
$$

In the same way as above we can show that

$$
\begin{equation*}
K_{\alpha, T}^{3} K_{\beta, T}^{2} \frac{\partial^{3}}{\partial \alpha^{3}} J_{\beta \beta} \rightarrow^{p} 0 \tag{24}
\end{equation*}
$$

Putting the expressions for derivatives (22), (23) and (24) into equation (21), and also noticing that due to Lemma $1 K_{\alpha, T}^{-1}\left(\hat{\alpha}-\alpha_{0}\right) \Rightarrow \frac{\xi_{\alpha}}{\operatorname{Var}\left(\xi_{\alpha}\right)}$, we get statement (a) of Lemma 5.
(b) Again we use Taylor expansion:

$$
\begin{equation*}
J_{\alpha \beta}\left(\hat{\alpha}, \beta_{0}\right)=J_{\alpha \beta}^{0}+\frac{\partial}{\partial \alpha} J_{\alpha \beta}^{0}\left(\hat{\alpha}-\alpha_{0}\right)+\frac{1}{2} \frac{\partial^{2}}{\partial \alpha^{2}} J_{\alpha \beta}^{0}\left(\alpha^{*}-\alpha_{0}\right)^{2} . \tag{25}
\end{equation*}
$$

From assumption 3(b)

$$
\begin{equation*}
K_{\alpha, T} K_{\beta, T} J_{\alpha \beta}^{0} \rightarrow^{p} \operatorname{cov}\left(\xi_{\alpha}, \xi_{\beta}\right) \tag{26}
\end{equation*}
$$

Taking the derivative we see

$$
\frac{\partial}{\partial \alpha} J_{\alpha \beta}=\frac{\partial}{\partial \alpha} \sum_{t} s_{\alpha, t} s_{\beta, t}=\left[A_{\alpha \alpha}, S_{\beta}\right]+\left[A_{\alpha \beta}, S_{\alpha}\right]-2 \sum s_{\alpha, t}^{2} s_{\beta, t}
$$

According to Lemma $4 K_{\alpha, T}^{2} K_{\beta, T} \sum s_{\alpha, t}^{2} s_{\beta, t} \rightarrow 0$. Assumptions 2(a) and 3 imply that $\frac{K_{\alpha, T}^{2}}{K_{\alpha \alpha, T}} \rightarrow 0$, so $K_{\alpha, T}^{2} K_{\beta, T}\left[A_{\alpha \alpha}, S_{\beta}\right] \rightarrow^{p} 0$. We have

$$
\begin{equation*}
K_{\alpha, T}^{2} K_{\beta, T} \frac{\partial}{\partial \alpha} J_{\alpha \beta}=\frac{K_{\alpha, T}^{2} K_{\alpha \beta, T}}{K_{\alpha \beta, T}} K_{\alpha \beta, T} K_{\alpha, T}\left[A_{\alpha \beta}, S_{\alpha}\right]+o_{p}(1) \rightarrow^{p} C \cdot \operatorname{cov}\left(\xi_{\alpha \beta}, \xi_{\alpha}\right) \tag{27}
\end{equation*}
$$

Similarly, we can show that

$$
K_{\alpha, T}^{3} K_{\beta, T} \frac{\partial^{2}}{\partial \alpha^{2}} J_{\alpha \beta}^{0} \rightarrow 0
$$

Putting the last equation, together with (26) and (27), into (25) and using Lemma 1 we get statement (b) of Lemma 5.
(c) As before we use Taylor expansion

$$
\begin{aligned}
K_{\alpha, T}^{2} J_{\alpha \alpha}\left(\hat{\alpha}, \beta_{0}\right) & =K_{\alpha, T}^{2} J_{\alpha \alpha}^{0}+K_{\alpha, T}^{3} \frac{\partial}{\partial \alpha} J_{\alpha \alpha}^{0} \frac{\left(\alpha^{*}-\alpha_{0}\right)}{K_{\alpha, T}} \\
\frac{\partial}{\partial \alpha} J_{\alpha \alpha} & =2\left[A_{\alpha \alpha} S_{\alpha}\right]+2 \sum s_{\alpha, t}^{3} .
\end{aligned}
$$

By the same argument as before $K_{\alpha, T}^{3}\left[A_{\alpha \alpha} S_{\alpha}\right] \rightarrow^{p} 0$, and according to Lemma $4 K_{\alpha, T}^{3} \sum s_{\alpha, t}^{3} \rightarrow^{p}$ 0 . Given the result of Lemma 1 we arrive at statement (c).

## Proof of statement from Section 5.3

We consider model

$$
Y=\beta_{0} Z \pi_{0}+U_{0} ; \quad X=Z \pi_{0}+V_{0}
$$

Here $Y, X, U_{0}$ and $V_{0}$ are $n \times 1, Z$ is $n \times k, \pi_{0}$ is $k \times 1$, and $\beta_{0}$ is $1 \times 1$. Here we assume
that $u_{t}^{0}, v_{t}^{0}$ are i.i.d. with mean zero and known variance matrix $\left(\begin{array}{cc}1 & 0 \\ 0 & 1\end{array}\right)$.
The test we wish to perform is $H_{0}: \beta=\beta_{0}$, so $\pi$ is a nuisance parameter. In what follows all parameters and errors with zero indexes stand for their true values, while those without zero are calculated with values of $\pi$ which are not necessarily the truth. For example, $U=Y-\beta_{0} Z \pi$ is equal to $U_{0}$ only if $\pi=\pi_{0}$. We also use $\widetilde{V}=\beta_{0} U+V$ and $E=U-\beta_{0} V$. Notice that $E=E_{0}$ for all $\pi$.

The score (up to a constant multiplier) is $(k+1) \times 1$ vector $S_{T}=\left(S_{1}, S_{2}^{\prime}\right)^{\prime}$, where

$$
S_{1}=\pi^{\prime} Z^{\prime} U ; \quad S_{2}=\beta_{0} Z^{\prime} U+Z^{\prime} V=Z^{\prime} \widetilde{V}
$$

We consider a slightly different version of the $J$ matrix, using the expected quadratic variation $J=\sum_{t} E\left[s_{t} s_{t}^{\prime}\right]=\left(\begin{array}{cc}\pi^{\prime} Z^{\prime} Z \pi & \beta_{0} \pi^{\prime} Z^{\prime} Z \\ \beta_{0} Z^{\prime} Z \pi & \left(1+\beta_{0}^{2}\right) Z^{\prime} Z\end{array}\right)$. The LM statistic for the full parameter vector $\left(\beta_{0}, \pi\right)$ is

$$
\begin{equation*}
L M\left(\beta_{0}, \pi\right)=S^{\prime} J^{-1} S=\frac{1}{1+\beta_{0}^{2}}\left(\frac{\left(\pi^{\prime} Z^{\prime} E\right)^{2}}{\pi^{\prime} Z^{\prime} Z \pi}+\widetilde{V}^{\prime} Z\left(Z^{\prime} Z\right)^{-1} Z^{\prime} \widetilde{V}\right) \tag{28}
\end{equation*}
$$

Consider an alternative to our statistic introduced in (7). In particular, consider statistic $L M^{c}\left(\beta_{0}\right)$ introduced in equation (11). It is of the same type as the statistic considered in Kleibergen and Mavroeidis (2009). We need to minimize $L M\left(\beta_{0}, \pi\right)$ with respect to $\pi$. For this purpose we take first derivative with respect to $\pi$ :

$$
\frac{\partial L M}{\partial \pi}=\frac{1}{1+\beta_{0}^{2}}\left(2 \frac{\pi^{\prime} Z^{\prime} E}{\pi^{\prime} Z^{\prime} Z \pi} Z^{\prime} E-2 \frac{\left(\pi^{\prime} Z^{\prime} E\right)^{2}}{\left(\pi^{\prime} Z^{\prime} Z \pi\right)^{2}} Z^{\prime} Z \pi-2\left(1+\beta_{0}^{2}\right) Z^{\prime} \tilde{V}\right)
$$

So our first order condition is

$$
\begin{equation*}
\frac{\pi^{\prime} Z^{\prime} E}{\pi^{\prime} Z^{\prime} Z \pi} Z^{\prime} E-\frac{\left(\pi^{\prime} Z^{\prime} E\right)^{2}}{\left(\pi^{\prime} Z^{\prime} Z \pi\right)^{2}} Z^{\prime} Z \pi-\left(1+\beta_{0}^{2}\right) Z^{\prime} \tilde{V}=0 \tag{29}
\end{equation*}
$$

Note that if we multiply equation (29) by $\pi^{\prime}$ from the left the first two terms cancel out and

$$
\begin{equation*}
\pi^{\prime} Z^{\prime} \widetilde{V}=0 \tag{30}
\end{equation*}
$$

We next derive a formula for $\widetilde{V}^{\prime} Z\left(Z^{\prime} Z\right)^{-1} Z^{\prime} \tilde{V}$ from equation (29). As a first step we
multiply it by $\widetilde{V}^{\prime} Z\left(Z^{\prime} Z\right)^{-1}$

$$
\begin{aligned}
\left(1+\beta_{0}^{2}\right) \widetilde{V}^{\prime} Z\left(Z^{\prime} Z\right)^{-1} Z^{\prime} \widetilde{V} & =\frac{\pi^{\prime} Z^{\prime} E}{\pi^{\prime} Z^{\prime} Z \pi} \widetilde{V}^{\prime} Z\left(Z^{\prime} Z\right)^{-1} Z^{\prime} E-\frac{\left(\pi^{\prime} Z^{\prime} E\right)^{2}}{\left(\pi^{\prime} Z^{\prime} Z \pi\right)^{2}} \widetilde{V}^{\prime} Z\left(Z^{\prime} Z\right)^{-1} Z^{\prime} Z \pi= \\
& =\frac{\pi^{\prime} Z^{\prime} E}{\pi^{\prime} Z^{\prime} Z \pi} \widetilde{V}^{\prime} Z\left(Z^{\prime} Z\right)^{-1} Z^{\prime} E
\end{aligned}
$$

The last equality follows from (30). We plug the last equation into (28) and receive:

$$
L M^{c}\left(\beta_{0}\right)=\frac{1}{1+\beta_{0}^{2}}\left(\frac{\left(\tilde{\pi}^{\prime} Z^{\prime} E\right)^{2}}{\tilde{\pi}^{\prime} Z^{\prime} Z \tilde{\pi}}+\frac{1}{1+\beta_{0}^{2}} \frac{\tilde{\pi}^{\prime} Z^{\prime} E}{\tilde{\pi}^{\prime} Z^{\prime} Z \tilde{\pi}} \widetilde{V}^{\prime} Z\left(Z^{\prime} Z\right)^{-1} Z^{\prime} E\right)
$$

where $\tilde{\pi}$ is the minimizer of $\operatorname{LM}\left(\beta_{0}, \pi\right)$ and the point at which we consider FOC (29). Now let us introduce $\tilde{Y}=\beta_{0} Y+X=\left(1+\beta_{0}^{2}\right) Z \tilde{\pi}+\tilde{V}$, or $\tilde{V}=\tilde{Y}-\left(1+\beta_{0}^{2}\right) Z \tilde{\pi}$ and plug it in the last expression:

$$
\begin{equation*}
L M\left(\beta_{0}, \tilde{\pi}\right)=\frac{1}{\left(1+\beta_{0}^{2}\right)^{2}} \frac{\tilde{\pi}^{\prime} Z^{\prime} E}{\tilde{\pi}^{\prime} Z^{\prime} Z \tilde{\pi}} \widetilde{Y}^{\prime} Z\left(Z^{\prime} Z\right)^{-1} Z^{\prime} E . \tag{31}
\end{equation*}
$$

We adopt notation from Andrews, Moreira, and Stock (2006),

$$
Q_{S}=\frac{E^{\prime} Z\left(Z^{\prime} Z\right)^{-1} Z^{\prime} E}{1+\beta_{0}^{2}} ; Q_{S T}=\frac{E^{\prime} Z\left(Z^{\prime} Z\right)^{-1} Z^{\prime} \widetilde{Y}}{1+\beta_{0}^{2}} ; Q_{T}=\frac{\widetilde{Y}^{\prime} Z\left(Z^{\prime} Z\right)^{-1} Z^{\prime} \widetilde{Y}}{1+\beta_{0}^{2}}
$$

Let us also denote $\gamma=1+\beta_{0}^{2} ; a=\tilde{\pi}^{\prime} Z^{\prime} E ; b=\tilde{\pi}^{\prime} Z^{\prime} Z \tilde{\pi}$. We can re-write (29) as

$$
\begin{equation*}
\frac{a}{b} Z^{\prime} E-\frac{a^{2}}{b^{2}} Z^{\prime} Z \pi-\gamma Z^{\prime} \widetilde{Y}+\gamma^{2} Z^{\prime} Z \pi=0 \tag{32}
\end{equation*}
$$

Multiply equation (32) by $E^{\prime} Z\left(Z^{\prime} Z\right)^{-1}$ from the left and multiply (32) by $\widetilde{Y}^{\prime} Z\left(Z^{\prime} Z\right)^{-1}$ from the left and take a linear combination of the resulting expressions. We arrive at the following equation

$$
\left(\frac{a}{\gamma b}\right)^{2} Q_{S T}-\frac{a}{\gamma b}\left(Q_{T}+Q_{S}\right)+Q_{S T}=0
$$

As a result

$$
\frac{a}{\gamma b}=\frac{\left(Q_{S}+Q_{T}\right)-\sqrt{\left(Q_{S}+Q_{T}\right)^{2}-4 Q_{S T}^{2}}}{2 Q_{S T}} .
$$

Putting this into formula (31) we get

$$
L M^{c}\left(\beta_{0}\right)=\frac{\left(Q_{S}+Q_{T}\right)-\sqrt{\left(Q_{S}+Q_{T}\right)^{2}-4 Q_{S T}^{2}}}{2}
$$

Simulations show that under weak identification, the statistic is strictly dominated by $\chi_{1}^{2}$.

## Proof of Lemma 3.

First we prove that $\left(Z^{\prime} Z\right)^{-1 / 2} Z^{\prime} \hat{U}, \hat{\beta}$ and $\left(Z^{\prime} Z\right)^{-1 / 2} Z^{\prime} V_{0}$ are asymptotically normal and asymptotically uncorrelated with each other. Two $k$-vectors $\left(Z^{\prime} Z\right)^{-1 / 2} Z^{\prime} V_{0}$ and $\left(Z^{\prime} Z\right)^{-1 / 2} Z^{\prime} U_{0}$ have exact normal distribution and independent from each other (conditionally on $Z$, however, the whole analysis here is done conditionally on $Z$ ).

From the definition of $\hat{U}$ we can see:

$$
\left(Z^{\prime} Z\right)^{-1 / 2} Z^{\prime} \hat{U}=\left(Z^{\prime} Z\right)^{-1 / 2} Z^{\prime} U_{0}-\frac{\pi_{0}^{\prime} Z^{\prime} U_{0}}{\pi_{0}^{\prime} Z^{\prime} Z \pi_{0}}\left(Z^{\prime} Z\right)^{-1 / 2} Z^{\prime} \pi_{0}
$$

and $\hat{\beta}-\beta_{0}=\frac{\pi^{\prime} Z^{\prime} U_{0}}{\pi^{\prime} Z^{\prime} Z \pi}$. As we can see $\left(Z^{\prime} Z\right)^{-1 / 2} Z^{\prime} U_{0}$ is represented as a sum of two orthogonal components $\left(Z^{\prime} Z\right)^{-1 / 2} Z^{\prime} \hat{U}$ and $\frac{\pi_{0}^{\prime} Z^{\prime} U_{0}}{\pi_{0}^{\prime} Z^{\prime} Z \pi_{0}}\left(Z^{\prime} Z\right)^{-1 / 2} Z^{\prime} \pi_{0}$. The last one is normalized $\hat{\beta}$. So, asymptotic normality and independence is proved.

Now, consider $\operatorname{LM}\left(\hat{\beta}, \pi_{0}\right)$ written as in (13) and condition it on $\hat{\beta}$. We can see that $\left(Z^{\prime} Z\right)^{-1 / 2} Z^{\prime}\left(\hat{\beta} \hat{U}+V_{0}\right)$ is normal $k \times 1$ - vector conditionally on $\hat{\beta}$, and

$$
\left(1+\hat{\beta}^{2}\right) I_{k}-\frac{\hat{\beta}^{2}\left(Z^{\prime} Z\right)^{1 / 2} \pi_{0} \pi_{0}^{\prime}\left(Z^{\prime} Z\right)^{1 / 2}}{\pi_{0}^{\prime} Z^{\prime} Z \pi_{0}}
$$

is its conditional covariance matrix. Thus, conditionally on $\hat{\beta}$ statistic $L M\left(\hat{\beta}, \pi_{0}\right)$ has $\chi_{k}^{2}$ asymptotic distribution.

Notation. Whenever a function is given with no argument, it means it is evaluated at the true $\theta_{0}$. Whenever we evaluate a function at a point other than $\theta_{0}$, we write it explicitly. For martingales $S$ and $A$ and for information quantities $J$ and $I$ indexes 1 and 2 stand for the sub-matrices corresponding to $\alpha_{1}$ and $\alpha_{2}$ respectively, and $\beta$ for the sub-matrix corresponding to $\beta$.

For the functions $\ell, m, n$ and $r$ only, the subscript 1 stands for the partial derivative with respect to $\alpha_{1}$, subscript 2 stands for the partial derivative with respect to $\alpha_{2}$, and subscript $\beta$ denotes the partial derivative with respect to $\beta$. For the second derivatives the order of subscripts determine the dimension of the corresponding matrix, for example, $\ell_{1,2}=\frac{\partial^{2}}{\partial \alpha_{1}^{\prime} \partial \alpha_{2}}$ is $k_{1} \times k_{2}$ matrix. In general the third derivative is a third order tensor, the forth derivative is a fourth order tensor, and any statement about convergence of tensors means element-by-element convergence unless otherwise specified. $N_{1}$ is the set
of indexes of $\alpha_{1}$, so $\alpha_{(n)}$ is an element of $\alpha_{1}$ if $n \in N_{1} . N_{2}$ is the set of indexes for $\alpha_{2}$, and $N_{\beta}$ is the set of indexes for $\beta$. For example, $\ell_{1 \beta(n)}=\frac{\partial}{\partial \alpha_{(n)}} \ell_{1 \beta}$ is a matrix of size $k_{1} \times k_{\beta}$ of third derivatives, and if $n \in N_{2}$ then $\ell_{1 \beta(n)}$ is sub-matrix of tensor $\ell_{1 \beta 2}$. Likewise $A_{\beta 1}$ is $k_{\beta} \times k_{1}$ matrix, while $A_{\beta(n)}$ is $k_{\beta} \times 1$ sub-matrix of $A_{\beta 1}$ if $n \in N_{1}$.

Consider the following normalization: $K_{1, T}=\frac{1}{\sqrt{T}} I d_{k_{1}} ; K_{2, T}=I d_{k_{2}} ; K_{\beta, T}=I d_{k_{\beta}} ;$ $K_{\alpha, T}=\left(\begin{array}{cc}K_{1, T} & 0 \\ 0 & K_{2, T}\end{array}\right)$ and $K_{T}=\left(\begin{array}{cc}K_{\alpha, T} & 0 \\ 0 & K_{\beta, T}\end{array}\right)$.

Lemma 6 Under the conditions of Theorem 3 the following statements hold:
(a) $K_{T} J_{T} K_{T} \rightarrow^{p} \Sigma$, where $\Sigma$ is non-degenerate finite matrix;
(b) $J_{\alpha \alpha, T}^{-1} I_{\alpha \alpha, T} \rightarrow^{p} I d_{k_{\alpha}}$;
(c) $\frac{1}{T^{3 / 2}} \ell_{111} \rightarrow^{p} 0 ; \frac{1}{T} \ell_{112} \rightarrow^{p} 0 ; \frac{1}{\sqrt{T}} \ell_{122} \rightarrow^{p} 0 ; \ell_{222} \rightarrow^{p} 0$;
(d) $\frac{1}{T} \ell_{11 \beta} \rightarrow^{p} 0 ; \frac{1}{T} \ell_{112 \beta} \rightarrow^{p} 0 ; \frac{1}{\sqrt{T}} \ell_{122 \beta} \rightarrow^{p} 0 ; \ell_{222 \beta} \rightarrow^{p} 0$;
(e) $\ell_{22(j)}-\left[A_{2(j)}, S_{2}\right]-\left[S_{2}, A_{2(j)}\right] \rightarrow^{p} 0$ for any $j \in N_{\beta}$;
(f) $\frac{1}{\sqrt{T}}\left(\ell_{\beta 1(j)}-\left[A_{\beta(j)}, S_{1}\right]\right) \rightarrow^{p} 0$ for any $j \in N_{2}$;
(g) $\frac{A_{\beta 1}}{\sqrt{T}} \rightarrow^{p} 0 ;\left[\frac{A_{\beta(j)}}{\sqrt{T}}, \frac{S_{1}}{\sqrt{T}}\right] \rightarrow^{p} 0$ and $\left[\frac{A_{\beta(j)}}{\sqrt{T}}, S_{2}\right] \rightarrow^{p} 0$ for any $j \in N_{1}$;

Proof of Lemma 6. Let us denote $H_{1}=\sum_{t=1}^{T}\left(H\left(x_{t}\right)-\dot{A}\right)$ (which is a $p \times 1$ vector), $H_{2}=\sum_{t=1}^{T}\left(H\left(x_{t}\right)-\dot{A}\right)\left(H\left(x_{t}\right)-\dot{A}\right)^{\prime}($ a $p \times p$ matrix $)$. According to the conditions of Theorem 3, $\frac{1}{T} H_{1} \rightarrow^{p} 0$ and $\frac{1}{T} H_{2} \rightarrow^{p}-\ddot{A}$.
(a) One can check that

$$
S_{T}=\left(\begin{array}{c}
\left(m_{1}+\frac{1}{\sqrt{T}} n_{1} \alpha_{2}+\frac{1}{\sqrt{T}} r_{1}\right)^{\prime} H_{1} \\
\frac{n^{\prime}}{\sqrt{T}} H_{1} \\
\frac{\left(n_{\beta} \alpha_{2}+r_{\beta}\right)^{\prime}}{\sqrt{T}} H_{1}
\end{array}\right)=\frac{\partial \eta^{\prime}}{\partial \theta} H_{1},
$$

where $\frac{\partial \eta}{\partial \theta}=\left(\left(m_{1}+\frac{1}{\sqrt{T}} n_{1} \alpha_{2}+\frac{1}{\sqrt{T}} r_{1}\right), \frac{n}{\sqrt{T}}, \frac{\left(n_{\beta} \alpha_{2}+r_{\beta}\right)}{\sqrt{T}}\right)$ is $p \times k$ matrix. It is easy to show that $J_{T}=\frac{\partial \eta^{\prime}}{\partial \theta} H_{2} \frac{\partial \eta}{\partial \theta}$. Using the normalization $K_{T}$ we have:

$$
K_{T} J_{T} K_{T} \rightarrow^{p}-\left(\begin{array}{c}
m_{1}^{\prime} \\
n^{\prime} \\
\left(n_{\beta} \alpha_{2}+r_{\beta}\right)^{\prime}
\end{array}\right) \ddot{A}\left(m_{1}, n, n_{\beta} \alpha_{2}+r_{\beta}\right)=\Sigma .
$$

Due to the rank assumption, $\Sigma$ is positive-definite.
(b) We calculate $I_{\alpha \alpha, T}$ :

$$
I_{\alpha \alpha, T}={ }^{a s y}\left(\begin{array}{cc}
m_{11} H_{1}-T m_{1}^{\prime} \ddot{A} m_{1} & \frac{n_{1}^{\prime}}{\sqrt{T}} H_{1}-T \frac{n}{\sqrt{T}} \ddot{A} m_{1} \\
\frac{n_{1}^{\prime}}{\sqrt{T}} H_{1}-T \frac{n}{\sqrt{T}} \ddot{A} m_{1} & -T \frac{n^{\prime}}{\sqrt{T}} \ddot{A} \frac{n}{\sqrt{T}}
\end{array}\right)
$$

Now it is straightforward to show that $K_{\alpha, T} I_{\alpha \alpha, T} K_{\alpha, T}$ converges to the same limit as $K_{\alpha, T} J_{\alpha \alpha, T} K_{\alpha, T}$. This means that $J_{\alpha \alpha, T}^{-1} I_{\alpha \alpha, T} \rightarrow^{p} I d_{k_{\alpha}}$.
(c) Can be proved by tedious differentiation. Below we drop all terms that are of obviously smaller order:

$$
\begin{gathered}
\frac{1}{T^{3 / 2}} \ell_{111}=\frac{1}{T^{3 / 2}}\left(H_{1} m_{111}-3 T \ddot{A} m_{1} m_{11}-T \dddot{A} m_{1}^{3}\right)+o_{p}(1) \rightarrow^{p} 0 \\
\frac{1}{T} \ell_{112}=\frac{1}{T}\left(H_{1} \frac{n_{11}}{\sqrt{T}}-T \ddot{A} m_{11} \frac{n}{\sqrt{T}}-2 T \ddot{A} m_{1} \frac{n_{1}}{\sqrt{T}}-T \dddot{A} m_{1}^{2} \frac{n}{\sqrt{T}}\right) \rightarrow^{p} 0 \\
\frac{1}{\sqrt{T}} \ell_{122}=\frac{1}{\sqrt{T}}\left(T \dddot{A} m_{1}\left(\frac{n}{\sqrt{T}}\right)^{2}-2 T \ddot{A} \frac{n}{\sqrt{T}} \frac{n_{1}}{\sqrt{T}}\right) \rightarrow 0 \\
\ell_{222}=-T \dddot{A}\left(\frac{n}{\sqrt{T}}\right)^{3} \rightarrow 0
\end{gathered}
$$

The last two statements employ that $\alpha_{2}$ enters linearly, so any time we differentiate with respect to $\alpha_{2}$ a term including $\frac{n}{\sqrt{T}}$ appears.
(d) As above, the idea here is that since $\eta_{T}$ is linear in $\alpha_{2}$, each additional derivative with respect to $\alpha_{2}$ generates $\frac{n}{\sqrt{T}}$.

$$
\frac{1}{T} \ell_{11 \beta}=\frac{1}{T}\left(H_{1} \frac{n_{11 \beta} \alpha_{2}+r_{11 \beta}}{\sqrt{T}}-T \ddot{A} m_{11} \frac{n_{\beta} \alpha_{2}+r_{\beta}}{\sqrt{T}}-2 T \ddot{A} m_{1} \frac{n_{1 \beta} \alpha_{2}+r_{1 \beta}}{\sqrt{T}}-T \dddot{A} m_{1}^{2} \frac{n_{\beta} \alpha_{2}+r_{\beta}}{\sqrt{T}}\right) \rightarrow 0
$$

$$
\begin{gathered}
\frac{1}{\sqrt{T}} \ell_{122 \beta}=\frac{1}{\sqrt{T}}\left(-T \dddot{A} m_{1 \beta}\left(\frac{n}{\sqrt{T}}\right)^{2}-2 T \dddot{A} m_{1} \frac{n}{\sqrt{T}} \frac{n_{\beta}}{\sqrt{T}}-T \dddot{A} m_{1}\left(\frac{n}{\sqrt{T}}\right)^{2} \frac{n_{\beta} \alpha_{2}+r_{\beta}}{\sqrt{T}}-\right. \\
\left.-2 T \ddot{A} \frac{n_{\beta}}{\sqrt{T}} \frac{n_{1}}{\sqrt{T}}-2 T \ddot{A} \frac{n}{\sqrt{T}} \frac{n_{1 \beta}}{\sqrt{T}}-2 T \dddot{A} \frac{n}{\sqrt{T}} \frac{n_{1}}{\sqrt{T}} \frac{n_{\beta} \alpha_{2}+r_{\beta}}{\sqrt{T}}\right) \rightarrow 0 \\
\ell_{222 \beta}=-3 T \dddot{A}\left(\frac{n}{\sqrt{T}}\right)^{2} \frac{n_{\beta}}{\sqrt{T}}-T \dddot{A}\left(\frac{n}{\sqrt{T}}\right)^{3} \frac{n_{\beta} \alpha_{2}+r_{\beta}}{\sqrt{T}} \rightarrow 0
\end{gathered}
$$

The expression for $\frac{1}{T} \ell_{112 \beta}$ is quite long, and we omit it here.
(e) One can easily check that $\ell_{22}=-n^{\prime} \ddot{A} n$, so for any $j \in N_{2}$ we have

$$
\ell_{22(j)}=-n_{(j)}^{\prime} \ddot{A} n-n^{\prime} \ddot{A} n_{(j)}-n^{\prime} \dddot{A} A n \frac{n_{(j)} \alpha_{2}+r_{(j)}}{\sqrt{T}} \rightarrow-n_{(j)}^{\prime} \ddot{A} n-n^{\prime} \ddot{A} n_{(j)} .
$$

We also have

$$
A_{2(j)}=\frac{n_{(j)}^{\prime}}{\sqrt{T}} H_{1}+o_{p}(1)
$$

and as a result

$$
\left[A_{2(j)}, S_{2}\right]=\frac{n_{(j)}^{\prime}}{\sqrt{T}} H_{1} \frac{n}{\sqrt{T}}+o_{p}(1) \rightarrow^{p}-n_{(j)}^{\prime} \ddot{A} n
$$

where we have used assumptions (a)-(c) of Theorem 3 several times.
(f) Taking derivatives one can check that for any $j \in N_{2}$

$$
\frac{1}{\sqrt{T}} \ell_{1 \beta(j)} \rightarrow^{p}-m_{1}^{\prime} \ddot{A}\left(n e_{(j)}\right)_{\beta},
$$

where $e_{(j)}$ is a vector of size $k_{2} \times 1$ with $(j)$-th component 1 and zeros in all other entries.

$$
A_{\beta(j)}=H_{1}^{\prime} \frac{\left(n e_{(j)}\right)_{\beta}}{\sqrt{T}}+o_{p}(1)
$$

As a result,

$$
\left[\frac{S_{1}}{\sqrt{T}}, A_{\beta(j)}\right]=m_{1}^{\prime} \frac{1}{T} H_{2}\left(n e_{(j)}\right)_{\beta} \rightarrow^{p}-m_{1}^{\prime} \ddot{A}\left(n e_{(j)}\right)_{\beta} .
$$

(g) For $j \in N_{1}$

$$
\frac{A_{\beta(j)}}{\sqrt{T}}=\frac{1}{\sqrt{T}} H_{1}^{\prime} \frac{n_{\beta(j)}}{\sqrt{T}}+o_{p}(1) .
$$

As a result,

$$
\left[\frac{A_{\beta(j)}}{\sqrt{T}}, \frac{S_{1}}{\sqrt{T}}\right]=\frac{1}{\sqrt{T}} m_{1}^{\prime} \frac{H_{2}}{T} n_{\beta(j)}+o_{p}(1) \rightarrow^{p} 0
$$

and

$$
\left[\frac{A_{\beta(j)}}{\sqrt{T}}, S_{2}\right]=\frac{1}{\sqrt{T}} n^{\prime} \frac{H_{2}}{T} n_{\beta(j)}+o_{p}(1) \rightarrow^{p} 0 .
$$

## Proof of Theorem 3

Let us consider the restricted ML estimator for $\alpha$, namely $\hat{\alpha}$ solving $S_{\alpha, T}\left(\hat{\alpha}, \beta_{0}\right)=0$. Taylor expanding this expression around $\alpha_{0}$ gives us

$$
K_{\alpha, T} S_{\alpha}(\hat{\alpha}, \beta)=K_{\alpha, T} S_{\alpha}^{0}-K_{\alpha, T} I_{\alpha \alpha, T} K_{\alpha, T} K_{\alpha, T}^{-1}\left(\hat{\alpha}-\alpha_{0}\right)+R
$$

where

$$
R_{(i)}=\left\{\begin{array}{cl}
\frac{1}{\sqrt{T}}\left(\frac{1}{2}\left(\alpha_{1}^{*}-\alpha_{1}^{0}\right)^{\prime} \ell_{11(i)}\left(\alpha_{1}^{*}-\alpha_{1}^{0}\right)+\left(\alpha_{1}^{*}-\alpha_{1}^{0}\right)^{\prime} \ell_{12(i)}\left(\alpha_{2}^{*}-\alpha_{2}^{0}\right)\right) & \text { if } i \in I_{1} ; \\
\left(\left(\alpha_{1}^{*}-\alpha_{1}^{0}\right)^{\prime} \ell_{12(i)}\left(\alpha_{2}^{*}-\alpha_{2}^{0}\right)+\frac{1}{2}\left(\alpha_{2}^{*}-\alpha_{2}^{0}\right)^{\prime} \ell_{22(i)}\left(\alpha_{2}^{*}-\alpha_{2}^{0}\right)\right) & \text { if } i \in I_{2}
\end{array}\right.
$$

Lemma 6 (c) implies that $R \rightarrow^{p} 0$. Using also Lemma 6 (b) we get:

$$
K_{\alpha, T}^{-1}\left(\hat{\alpha}-\alpha_{0}\right)=K_{\alpha, T}^{-1}\left(I_{\alpha \alpha, T}\right)^{-1} S_{\alpha}+o_{p}(1)=K_{\alpha, T}^{-1}\left(J_{\alpha \alpha, T}\right)^{-1} S_{\alpha}+o_{p}(1)
$$

This means that the equation for the restricted MLE is asymptotically linear in $S_{\alpha}$, and thus that the asymptotic distribution of $\hat{\alpha}$ will be normal despite the fact that $\hat{\alpha}_{2}$ is not consistent.

Now let us consider the score at the restricted MLE

$$
\begin{align*}
S_{\beta, T}\left(\hat{\alpha}, \beta_{0}\right)=S_{\beta}-I_{\beta \alpha}\left(\hat{\alpha}-\alpha_{0}\right) & +\frac{1}{2} \sum_{i \in I_{2}} \ell_{\beta 2(i)}\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right)\left(\hat{\alpha}_{(i)}-\alpha_{(i)}^{0}\right)+ \\
& +\sum_{i \in I_{2}} \ell_{\beta 1(i)}\left(\hat{\alpha}_{1}-\alpha_{1}^{0}\right)\left(\hat{\alpha}_{(i)}-\alpha_{(i)}^{0}\right)+\tilde{R} \tag{33}
\end{align*}
$$

where for each $i \in I_{\beta}$

$$
\begin{aligned}
\tilde{R}_{(i)}=\frac{1}{2}\left(\alpha_{1}^{*}-\alpha_{1}^{0}\right)^{\prime} \ell_{11(i)}\left(\alpha_{1}^{*}-\alpha_{1}^{0}\right)+ & \frac{1}{2} \sum_{j \in I_{2}}\left(\alpha_{1}^{*}-\alpha_{1}^{0}\right)^{\prime} \ell_{11(i)(j)}\left(\alpha_{1}^{*}-\alpha_{1}^{0}\right)\left(\alpha_{(j)}^{*}-\alpha_{(j)}^{0}\right)^{2}+ \\
+ & \frac{1}{2} \sum_{j \in I_{2}}\left(\alpha_{1}^{*}-\alpha_{1}^{0}\right)^{\prime} \ell_{12(i)(j)}\left(\alpha_{2}^{*}-\alpha_{2}^{0}\right)\left(\alpha_{(j)}^{*}-\alpha_{(j)}^{0}\right)^{2}+ \\
& +\frac{1}{6} \sum_{j \in I_{2}}\left(\alpha_{2}^{*}-\alpha_{2}^{0}\right)^{\prime} \ell_{22(i)(j)}\left(\alpha_{2}^{*}-\alpha_{2}^{0}\right)\left(\alpha_{(j)}^{*}-\alpha_{(j)}^{0}\right) .
\end{aligned}
$$

Statement (d) of Lemma 6 gives us that $\tilde{R} \rightarrow^{p} 0$.
Now we plug statements (e) and (f) of Lemma 6 in to equation (33):

$$
\begin{aligned}
S_{\beta, T}\left(\hat{\alpha}, \beta_{0}\right)=S_{\beta}-I_{\beta \alpha}\left(\hat{\alpha}-\alpha_{0}\right)+ & \sum_{i \in I_{2}}\left[A_{\beta(i)}, S_{2}\right]\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right)\left(\hat{\alpha}_{(i)}-\alpha_{(i)}^{0}\right)+ \\
& +\sum_{i \in I_{2}}\left[A_{\beta(i)}, S_{1}\right]\left(\hat{\alpha}_{1}-\alpha_{1}^{0}\right)\left(\hat{\alpha}_{(i)}-\alpha_{(i)}^{0}\right) .
\end{aligned}
$$

By definition

$$
\begin{array}{r}
I_{\beta \alpha}\left(\hat{\alpha}-\alpha_{0}\right)=J_{\beta \alpha} J_{\alpha \alpha}^{-1} S_{\alpha}-A_{\beta 1}\left(\hat{\alpha}_{1}-\alpha_{1}^{0}\right)-A_{\beta 2}\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right)= \\
=J_{\beta \alpha} J_{\alpha \alpha}^{-1} S_{\alpha}-A_{\beta 1}\left(\hat{\alpha}_{1}-\alpha_{1}^{0}\right)-\sum_{i \in I_{2}} A_{\beta(i)}\left(\hat{\alpha}_{(i)}-\alpha_{(i)}^{0}\right)
\end{array}
$$

So,

$$
\begin{array}{r}
S_{\beta, T}\left(\hat{\alpha}, \beta_{0}\right)=S_{\beta}-J_{\beta \alpha} J_{\alpha, \alpha}^{-1} S_{\alpha}-A_{\beta 1}\left(\hat{\alpha}_{1}-\alpha_{1}^{0}\right)- \\
-\sum_{i \in I_{2}}\left(A_{\beta(i)}-\left[A_{\beta(i)}, S_{1}\right]\left(\hat{\alpha}_{1}-\alpha_{1}^{0}\right)-\left[A_{\beta(i)}, S_{2}\right]\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right)\right)\left(\hat{\alpha}_{(i)}-\alpha_{(i)}^{0}\right)= \\
=S_{\beta}-J_{\beta \alpha} J_{\alpha, \alpha}^{-1} S_{\alpha}-\sum_{i \in I_{2}}\left(A_{\beta(i)}-\left[A_{\beta(i)}, S_{\alpha}\right]\left(\hat{\alpha}-\alpha^{0}\right)\right)\left(\hat{\alpha}_{(i)}-\alpha_{(i)}^{0}\right)= \\
=S_{\beta}^{\perp}-A_{\beta 2}^{\perp}\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right)
\end{array}
$$

where $S_{\beta}^{\perp}$ and $A_{\beta 2}^{\perp}$ are asymptotically normal and independent from $\frac{1}{\sqrt{T}} S_{1}$ and $S_{2}$. In the derivation above we used statement (g) of Lemma 6.

More explicitly,

$$
S_{\beta}^{\perp}=S_{\beta}^{0}-J_{\alpha \beta}^{0} J_{\alpha, \alpha}^{-1} S_{\alpha}
$$

and for any $i \in I_{2}$

$$
A_{\beta(i)}^{\perp}=A_{\beta(i)}-\left[A_{\beta(i)}, S_{\alpha}\right]\left(\hat{\alpha}-\alpha^{0}\right)=A_{\beta(i)}-\left[A_{\beta(i)}, S_{\alpha}\right] J_{\alpha \alpha}^{-1} S_{\alpha} .
$$

To summarize, we showed that

$$
\begin{equation*}
S_{\beta, T}\left(\hat{\alpha}, \beta_{0}\right)=S_{\beta}^{\perp}-A_{\beta 2}^{\perp}\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right) . \tag{34}
\end{equation*}
$$

This statistic is asymptotically normal conditional on $K_{\alpha, T} S_{\alpha, T}$ (and thus conditional on $\left.K_{\alpha, T}^{-1}\left(\hat{\alpha}-\alpha_{0}\right)\right)$.

Now, we consider the denominator. Let us consider $s_{\beta, t}(\theta)=\frac{n_{\beta}\left(\alpha_{1}, \beta\right) \alpha_{2}+r_{\beta}\left(\alpha_{1}, \beta\right)}{\sqrt{T}}\left(H\left(x_{t}\right)-\right.$ $\left.\dot{A}\left(\eta_{t}(\theta)\right)\right)$. We may notice that

$$
\begin{equation*}
s_{\beta, t}\left(\hat{\alpha}, \beta_{0}\right)=s_{\beta, t}\left(\theta_{0}\right)+\frac{\partial}{\partial \alpha_{2}} s_{\beta, t}\left(\theta_{0}\right)\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right)+O_{p}\left(\frac{1}{T}\right) \tag{35}
\end{equation*}
$$

This can be seen from observing that $\frac{\partial}{\partial \alpha_{1}} s_{\beta, t}\left(\theta_{0}\right)=O_{p}\left(\frac{1}{\sqrt{T}}\right), \frac{\partial^{2}}{\partial \alpha_{2}^{2}} s_{\beta, t}=O_{p}\left(\frac{1}{T}\right),\left(\hat{\alpha}_{1}-\alpha_{1}^{0}\right)=$
$O_{p}\left(\frac{1}{\sqrt{T}}\right)$, etc.
Similarly,

$$
s_{1, t}\left(\hat{\alpha}, \beta_{0}\right)=s_{1, t}\left(\theta_{0}\right)+O_{p}\left(\frac{1}{\sqrt{T}}\right)
$$

and

$$
s_{2, t}\left(\hat{\alpha}, \beta_{0}\right)=s_{2, t}\left(\theta_{0}\right)+O_{p}\left(\frac{1}{T}\right) .
$$

Joining the last two together we get:

$$
\begin{equation*}
K_{\alpha, T} s_{\alpha, t}\left(\hat{\alpha}, \beta_{0}\right)=K_{\alpha, T} s_{\alpha, t}\left(\theta_{0}\right)+O_{p}\left(\frac{1}{T}\right) \tag{36}
\end{equation*}
$$

Based on the definition of martingale $A_{T}$ we can re-write (35) as

$$
\begin{equation*}
s_{\beta, t}\left(\hat{\alpha}, \beta_{0}\right)=s_{\beta, t}\left(\theta_{0}\right)+a_{\beta 2, t}\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right)+s_{\beta, t}\left(\theta_{0}\right) s_{2, t}\left(\theta_{0}\right)^{\prime}\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right)+O_{p}\left(\frac{1}{T}\right) . \tag{37}
\end{equation*}
$$

From (37) and Lemma 4

$$
J_{\beta \beta}\left(\hat{\alpha}, \beta_{0}\right)=J_{\beta \beta}+\left[S_{\beta}, A_{\beta 2}\right]\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right)+\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right)^{\prime}\left[A_{2 \beta}, S_{\beta}\right]+\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right)^{\prime}\left[A_{2 \beta}, A_{\beta, 2}\right]\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right)+o_{p}(1) .
$$

Joining together (37) and (36) we get

$$
K_{\alpha, T} J_{\alpha \beta}\left(\hat{\alpha}, \beta_{0}\right)=K_{\alpha, T} J_{\alpha \beta}\left(\theta_{0}\right)+K_{\alpha, T}\left[S_{\alpha}, A_{\beta 2}\right]\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right)+o_{p}(1) .
$$

Finally (36) implies that

$$
K_{\alpha, T} J_{\alpha \alpha}\left(\hat{\alpha}, \beta_{0}\right) K_{\alpha, T}=K_{\alpha, T} J_{\alpha \alpha}\left(\theta_{0}\right) K_{\alpha, T}+o_{p}(1) .
$$

Putting together the last three statements we have

$$
\begin{array}{r}
\left.\left(J_{\beta \beta}-J_{\beta \alpha} J_{\alpha \alpha}^{-1} J_{\alpha \beta}\right)\right|_{\theta=\left(\hat{\alpha}, \beta_{0}\right)}=J_{\beta \beta}-J_{\beta \alpha} J_{\alpha \alpha}^{-1} J_{\alpha \beta}+ \\
+\left[S_{\beta}, A_{\beta 2}\right]\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right)+\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right)^{\prime}\left[A_{2 \beta}, S_{\beta}\right]+\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right)^{\prime}\left[A_{2 \beta}, A_{\beta, 2}\right]\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right)- \\
-J_{\beta \alpha} J_{\alpha \alpha}^{-1}\left[S_{\alpha}, A_{\beta 2}\right]\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right)-\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right)^{\prime}\left[A_{\beta 2}, S_{\alpha}\right] J_{\alpha \alpha}^{-1} J_{\alpha \beta}- \\
-\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right)^{\prime}\left[A_{\beta 2}, S_{\alpha}\right] J_{\alpha \alpha}^{-1}\left[S_{\alpha}, A_{\beta 2}\right]\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right)+o_{p}(1) .
\end{array}
$$

All functions on the right side are evaluated at $\theta_{0}$. One can see that the quantity on the
right side of the last equation is the conditional variance of the score from equation (34) given $\left(\hat{\alpha}_{2}-\alpha_{2}^{0}\right)$.

## 10 Figures and Tables

Figure 1: Histogram of parameter estimates from jointly estimating $\alpha, \phi_{x}, \phi_{\pi}, \rho, \delta, \kappa, \sigma_{a}, \sigma_{u}, \sigma$


Table 1: Size of Classical Tests for $\theta_{0}$ in Simulation

| Test Satatistic | Size of $5 \%$ Test | Size of $10 \%$ Test |
| :---: | :---: | :---: |
| $L R\left(\theta_{0}\right)$ | $3.20 \%$ | $7.05 \%$ |
| Wald $\left(I\left(\theta_{0}\right)\right)$ | $65.45 \%$ | $67.20 \%$ |
| Wald $(I(\hat{\theta}))$ | $63.05 \%$ | $64.30 \%$ |
| Wald $\left(J\left(\theta_{0}\right)\right)$ | $68.05 \%$ | $70.80 \%$ |
| Wald $(J(\hat{\theta}))$ | $68.15 \%$ | $71.00 \%$ |
| $L M\left(\theta_{0}\right)$ | $6.55 \%$ | $8.60 \%$ |

Figure 2: CDF of simulated LM statistic (using J) compared to $\chi_{9}^{2}$


Table 2: Simulated Test Size for $\left(\phi_{x}, \phi_{\pi}, \kappa, \sigma_{a}, \sigma_{u}, \sigma\right)$

| Test Statistic | $5 \%$ | $10 \%$ |
| :---: | :---: | :---: |
| $\alpha M\left(\hat{\alpha}, \beta_{0}\right)$ | $7.99 \%$ | $15.28 \%$ |
| $\min _{\alpha} L M\left(\alpha, \beta_{0}\right)$ | $6.41 \%$ | $12.99 \%$ |

Table 3: Simulated Test Size for $\left(\alpha, \rho, \delta, \sigma_{a}, \sigma_{u}, \sigma\right)$

| Test Statistic | $5 \%$ | $10 \%$ |
| :---: | :---: | :---: |
| $\widehat{L M}\left(\hat{\alpha}, \beta_{0}\right)$ | $8.95 \%$ | $15.40 \%$ |
| $\min _{\alpha} \widehat{L M}\left(\alpha, \beta_{0}\right)$ | $7.50 \%$ | $13.30 \%$ |

Table 4: Simulated Test Size for one dimensional hypotheses for $\left(\phi_{x}, \phi_{\pi}\right)$

| Parameter | Test Statistic | $5 \%$ | $10 \%$ |
| :---: | :---: | :---: | :---: |
| $\phi_{x}$ | $\widehat{L M}\left(\hat{\alpha}, \beta_{0}\right)$ | $8.90 \%$ | $16.30 \%$ |
|  | $\min _{\alpha} \widehat{L M}\left(\alpha, \beta_{0}\right)$ | $3.80 \%$ | $8.00 \%$ |
| $\phi_{\pi}$ | $\widehat{L M}\left(\hat{\alpha}, \beta_{0}\right)$ | $9.90 \%$ | $18.90 \%$ |
|  | $\min _{\alpha} \widehat{L M}\left(\alpha, \beta_{0}\right)$ | $3.30 \%$ | $7.80 \%$ |

Table 5: Size of $5 \%$ and $10 \%$ Tests based on $\widetilde{L M}\left(\pi_{0}, \hat{\beta}_{M L}\right)$ for Nonlinear IV Model

| Sample Size | Rejection rate for $5 \%$ test | Rejection rate for $10 \%$ test |
| :---: | :---: | :---: |
| 100 | $6.49 \%$ | $12.70 \%$ |
| 10000 | $5.70 \%$ | $11.34 \%$ |

Figure 3: CDF of $\widetilde{L M}\left(\pi_{0}, \hat{\beta}_{M L}\right)$ for Nonlinear weak IV, $c=.01, \mathrm{~T}=100$


Table 6: Size of $5 \%$ and $10 \%$ Asymptotic Tests Based on $\widetilde{L M}\left(\beta_{1}^{0}, \hat{\beta}_{2}, \hat{\pi}_{1}, \hat{\pi}_{2}\right)$

| Sample Size | Rejection rate for $5 \%$ test | Rejection rate for $10 \%$ test |
| :---: | :---: | :---: |
| 100 | $1.81 \%$ | $4.89 \%$ |
| 10000 | $1.37 \%$ | $4.04 \%$ |

Figure 4: CDF of $\widetilde{L M}\left(\beta_{1}^{0}, \hat{\beta}_{2}, \hat{\pi}\right)$ for $T=100$
Empirical CDF



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[^1]:    ${ }^{4}$ Our model is based on slides from Lawrence Christiano on estimation of the Clarida Gali and Gertler (1999) model, which we have altered by setting the tax rate equal to zero to simplify the model.

[^2]:    ${ }^{5}$ We conducted extensive simulations, only some of which are presented here. Additional results are available from the authors by request.

[^3]:    ${ }^{6}$ We produced simulations for different values of $\pi$ and $\Sigma$, which are available from the authors by request.

