# Supplemental Material 

# for <br> Asymptotic Size of Kleibergen's LM and Conditional LR Tests for Moment Condition Models 

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We let AG1 abbreviate the main paper Andrews and Guggenberger (2017) "Asymptotic Size of Kleibergen's LM and Conditional LR Tests for Moment Condition Models," Econometric Theory, forthcoming. References to Sections with Section numbers less than 9 refer to Sections of AG1. In consequence, the Section numbers in this Supplemental Material (SM) follow on from the main paper, starting with Section 9. Similarly, all theorems and lemmas with Section numbers less than 9 refer to results in AG1.

We let ACG abbreviate Andrews, Cheng, and Guggenberger (2009) "Generic Results for Establishing the Asymptotic Size of Confidence Sets and Tests," Cowles Foundation Discussion Paper No. 1813, Yale University. This SM makes use of some results in ACG.

We let AG2 abbreviate Andrews and Guggenberger (2014) "Identification- and SingularityRobust Inference for Moment Condition Models," Cowles Foundation Discussion Paper No. 1978, Yale University. AG2 utilizes some of the results in this SM.
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## 9 Outline

This SM provides proofs of the results stated in AG1. It also provides some complementary results to those in AG1.

Section 10 states some basic results that are used in all of the proofs. These results also are used in AG2 and should be useful for establishing the asymptotic sizes of other tests for moment condition models when strong identification is not assumed. Given the results in Section 10, Section 11 proves Theorem 4.1, Section 12 proves Theorem 6.1, and Section 13 proves Theorem 5.3.

Section 14 shows that the eigenvalue condition in $\mathcal{F}_{0}$, defined in (3.9), is not redundant in Theorems 4.1, 5.3, and 6.1.

Sections 15, 16, and 17 prove Lemma 10.2, Lemma 10.3, and Theorem 10.4, respectively, which appear in Section 10.

Section 18 proves that the conditions in (3.10) and (3.11) are sufficient for the second condition in $\mathcal{F}_{0 j}$.

Section 19 proves Theorem 5.1 and Lemma 5.2. Section 19 also determines the asymptotic size of Kleibergen's (2005) CLR test with Jacobian-variance weighting that employs the Robin and Smith (2000) rank statistic, defined in Section 5, for the general case of $p \geq 1$. When $p=1$, the asymptotic size of this test is correct. But, when $p \geq 2$, we cannot show that its asymptotic size is necessarily correct (because the sample moments and the rank statistic can be asymptotically dependent under some sequences of distributions). Section 19 provides some simulation results for this test.

Section 20 proves Theorem 7.1, which provides results for time series observations.
For notational simplicity, throughout the SM , we often suppress the argument $\theta_{0}$ for various quantities that depend on the null value $\theta_{0}$. Throughout the SM , the quantities $B_{F}, C_{F}$, and $\left(\tau_{1 F}, \ldots, \tau_{p F}\right)$ are defined using the general definitions given in 10.6 - 10.8 , rather than the definitions given in Section 3, which are a special case of the former definitions.

For notational simplicity, the proofs in Sections 1517 are for the sequence $\{n\}$, rather than a subsequence $\left\{w_{n}: n \geq 1\right\}$. The same proofs hold for any subsequence $\left\{w_{n}: n \geq 1\right\}$. The proofs in these three sections use the following simplified notation. Define

$$
\begin{align*}
D_{n} & :=E_{F_{n}} G_{i}, \Omega_{n}:=\Omega_{F_{n}}, B_{n}:=B_{F_{n}}, C_{n}:=C_{F_{n}}, B_{n}=\left(B_{n, q}, B_{n, p-q}\right), C_{n}=\left(C_{n, q}, C_{n, k-q}\right), \\
W_{n} & :=W_{F_{n}}, W_{2 n}:=W_{2 F_{n}}, U_{n}:=U_{F_{n}}, \text { and } U_{2 n}:=U_{2 F_{n}}, \tag{9.1}
\end{align*}
$$

where $q=q_{h}$ is defined in 10.16, $B_{n, q} \in R^{p \times q}, B_{n, p-q} \in R^{p \times(p-q)}, C_{n, q} \in R^{k \times q}$, and $C_{n, k-q} \in$

$$
R^{k \times(k-q)} \text {. Define }
$$

$$
\begin{align*}
\Upsilon_{n, q} & :=\operatorname{Diag}\left\{\tau_{1 F_{n}}, \ldots, \tau_{q F_{n}}\right\} \in R^{q \times q}, \Upsilon_{n, p-q}:=\operatorname{Diag}\left\{\tau_{(q+1) F_{n}}, \ldots, \tau_{p F_{n}}\right\} \in R^{(p-q) \times(p-q)}, \text { and } \\
\Upsilon_{n} & :=\left[\begin{array}{cc}
\Upsilon_{n, q} & 0^{q \times(p-q)} \\
0^{(p-q) \times q} & \Upsilon_{n, p-q} \\
0^{(k-p) \times q} & 0^{(k-p) \times(p-q)}
\end{array}\right] \in R^{k \times p} . \tag{9.2}
\end{align*}
$$

Note that $\Upsilon_{n}$ is the diagonal matrix of singular values of $W_{n} D_{n} U_{n}$, see (10.8).

## 10 Basic Framework and Results for the Proofs

### 10.1 Uniformity

The proofs of Theorems 4.1, 5.3, and 6.1 use Corollary 2.1(c) in ACG. The latter result provides general sufficient conditions for the correct asymptotic size and (uniform) asymptotic similarity of a sequence of tests.

We now state Corollary 2.1(c) of ACG. Let $\left\{\phi_{n}: n \geq 1\right\}$ be a sequence of tests of some null hypothesis whose null distributions are indexed by a parameter $\lambda$ with parameter space $\Lambda$. Let $R P_{n}(\lambda)$ denote the null rejection probability of $\phi_{n}$ under $\lambda$. For a finite nonnegative integer $J$, let $\left\{h_{n}(\lambda)=\left(h_{1 n}(\lambda), \ldots, h_{J n}(\lambda)\right)^{\prime} \in R^{J}: n \geq 1\right\}$ be a sequence of functions on $\Lambda$. Define

$$
\begin{align*}
H:=\{ & \left\{\left(R \cup(R \cup\{ \pm \infty\})^{J}: h_{w_{n}}\left(\lambda_{w_{n}}\right) \rightarrow h \text { for some subsequence }\left\{w_{n}\right\}\right.\right. \\
& \text { of } \left.\{n\} \text { and some sequence }\left\{\lambda_{w_{n}} \in \Lambda: n \geq 1\right\}\right\} . \tag{10.1}
\end{align*}
$$

Assumption B*: For any subsequence $\left\{w_{n}\right\}$ of $\{n\}$ and any sequence $\left\{\lambda_{w_{n}} \in \Lambda: n \geq 1\right\}$ for which $h_{w_{n}}\left(\lambda_{w_{n}}\right) \rightarrow h \in H, R P_{w_{n}}\left(\lambda_{w_{n}}\right) \rightarrow \alpha$ for some $\alpha \in(0,1)$.

Proposition 10.1 (ACG, Corollary 2.1(c)) Under Assumption B*, the tests $\left\{\phi_{n}: n \geq 1\right\}$ have asymptotic size $\alpha$ and are asymptotically similar (in a uniform sense). That is, AsySz $:=\limsup _{n \rightarrow \infty}$ $\sup _{\lambda \in \Lambda} R P_{n}(\lambda)=\alpha$ and $\liminf _{n \rightarrow \infty} \inf _{\lambda \in \Lambda} R P_{n}(\lambda)=\limsup _{n \rightarrow \infty} \sup _{\lambda \in \Lambda} R P_{n}(\lambda)$.

Comments: (i) By Comment 4 to Theorem 2.1 of ACG, Proposition 10.1 provides asymptotic size and similarity results for nominal $1-\alpha$ confidence sets (CS's), rather than tests, by defining $\lambda$ as one would for a test, but having it depend also on the parameter that is restricted by the null hypothesis, by enlarging the parameter space $\Lambda$ correspondingly (so it includes all possible values of the parameter that is restricted by the null hypothesis), and by replacing (i) $\phi_{n}$ by a CS based on a sample of size $n$, (ii) $\alpha$ by $1-\alpha$, (iii) $R P_{n}(\lambda)$ by $C P_{n}(\lambda)$, where $C P_{n}(\lambda)$ denotes the
coverage probability of the CS under $\lambda$ when the sample size is $n$, and (iv) the first $\limsup _{n \rightarrow \infty} \sup _{\lambda \in \Lambda}$ that appears by $\liminf _{n \rightarrow \infty} \inf _{\lambda \in \Lambda}$. In the present case, where the null hypotheses are of the form $H_{0}: \theta=\theta_{0}$ for $\theta \in \Theta$, for CS's, $\theta_{0}$ is taken to be a subvector of $\lambda$ and $\Lambda$ is specified so that the value of this subvector ranges over $\Theta$.
(ii) In the application of Proposition 10.1 to prove Theorems 4.1 and 6.1, one takes $\Lambda$ to be a one-to-one transformation of $\mathcal{F}_{0}$ for tests, and one takes $\Lambda$ to be a one-to-one transformation of $\mathcal{F}_{\Theta, 0}$ for CS's. With these changes, the proofs for tests and CS's are the same. In consequence, we provide explicit proofs for tests only and obtain the proofs for CS's by analogous applications of Proposition 10.1. In the application of Proposition 10.1 to prove Theorem 5.3, the same is done but with $\mathcal{F}_{J V W, p=1}$ in place of $\mathcal{F}_{0}$.
(iii) We prove the test results in Theorems $4.1,5.3$, and 6.1 using Proposition 10.1 by verifying Assumption $\mathrm{B}^{*}$ for suitable choices of $\lambda$ and $h_{n}(\lambda)$.

### 10.2 Random Weight Matrices $\widehat{W}_{n}$ and $\widehat{U}_{n}$

We prove results for statistics that depend on random weight matrices $\widehat{W}_{n} \in R^{k \times k}$ and $\widehat{U}_{n} \in$ $R^{p \times p}$. In particular, we consider statistics of the form $\widehat{W}_{n} \widehat{D}_{n} \widehat{U}_{n}$ and functions of this statistic, where $\widehat{D}_{n}$ is defined in 3.2. The definitions of the random weight matrices $\widehat{W}_{n}$ and $\widehat{U}_{n}$ depend upon the statistic that is of interest. They are taken to be of the form

$$
\begin{equation*}
\widehat{W}_{n}:=W_{1}\left(\widehat{W}_{2 n}\right) \in R^{k \times k} \text { and } \widehat{U}_{n}:=U_{1}\left(\widehat{U}_{2 n}\right) \in R^{p \times p} \tag{10.2}
\end{equation*}
$$

where $\widehat{W}_{2 n}$ and $\widehat{U}_{2 n}$ are random finite-dimensional quantities, such as matrices, and $W_{1}(\cdot)$ and $U_{1}(\cdot)$ are nonrandom functions that are assumed below to be continuous on certain sets. The estimators $\widehat{W}_{2 n}$ and $\widehat{U}_{2 n}$ have corresponding population quantities $W_{2 F}$ and $U_{2 F}$, respectively. For examples, see Examples 1-3 immediately below. Thus, the population quantities corresponding to $\widehat{W}_{n}$ and $\widehat{U}_{n}$ are

$$
\begin{equation*}
W_{F}:=W_{1}\left(W_{2 F}\right) \text { and } U_{F}:=U_{1}\left(U_{2 F}\right), \tag{10.3}
\end{equation*}
$$

respectively.
Example 1: With Kleibergen's (2005) LM test and the CLR test with moment-variance weighting, which are considered in Sections 4 and 6, respectively, we take

$$
\begin{equation*}
\widehat{W}_{n}=\widehat{\Omega}_{n}^{-1 / 2} \text { and } \widehat{U}_{n}=I_{p} . \tag{10.4}
\end{equation*}
$$

In this case, the functions $W_{1}(\cdot)$ and $U_{1}(\cdot)$ are the identity functions, and the corresponding popu-
lation quantities are $W_{F}=W_{2 F}=\Omega_{F}^{-1 / 2}$, where $\Omega_{F}:=E_{F} g_{i} g_{i}^{\prime}$, see 3.6, and $U_{F}=U_{2 F}=I_{p}$.
Example 2: For a CLR test based on an equally-weighted statistic other than $\widehat{\Omega}_{n}^{-1 / 2} \widehat{D}_{n}$, such as $\widetilde{W}_{n} \widehat{D}_{n}$, as in Comment (ii) to Theorem 6.1. one defines a pd matrix $\widetilde{W}_{n}$ as desired and one takes $\widehat{W}_{n}=\widetilde{W}_{n}$ and $\widehat{U}_{n}=U_{F}=U_{2 F}=I_{p}$.

Example 3: With Kleibergen's (2005) CLR test with Jacobian-variance weighting and $p=1$, which is considered in Section 5, we determine the asymptotic distribution of the rank statistic in 5.10 by taking $\widehat{W}_{n}=\widetilde{V}_{D n}^{-1 / 2}$ and $\widehat{U}_{n}=I_{p}$. In this case, the functions $W_{1}(\cdot)$ and $U_{1}(\cdot)$ are as in Example 1, and the corresponding population quantities are $W_{F}=W_{2 F}=\left(\operatorname{Var}_{F}\left(\operatorname{vec}\left(G_{i}\right)\right)-\right.$ $\left.\Gamma_{F}^{v e c}\left(G_{i}\right) \Omega_{F}^{-1} \Gamma_{F}^{v e c}\left(G_{i}\right) \prime\right)^{-1 / 2}=\left(\Psi_{F}^{v e c}\left(G_{i}\right)-E_{F} G_{i} E_{F} G_{i}^{\prime}\right)^{-1 / 2}$, and $U_{F}=U_{2 F}=I_{p}$. For this test, we need the asymptotic distribution of the LM statistic. In consequence, for this test, we also establish some asymptotic results with $\widehat{W}_{n}$ and $\widehat{U}_{n}$ defined as in Example 1.

Examples 4 \& 5: The results of this section are used in AG2 when the asymptotic sizes of two new SR-CQLR tests are determined. For the SR-CQLR tests, $\widehat{W}_{n}=\widehat{\Omega}_{n}^{-1 / 2}$ and it is convenient to take $W_{1}(\cdot)=(\cdot)^{-1 / 2}$ and $\widehat{W}_{2 n}=\widehat{\Omega}_{n}$, and the matrix $\widehat{U}_{n}$ is a nonlinear transformation $U_{1}(\cdot)$ of a matrix estimator, which is different for the two tests. For brevity, we do not define the nonlinear transformation or the two matrix estimators here.

We provide results for distributions $F$ in the following set of null distributions:

$$
\begin{equation*}
\mathcal{F}_{W U}:=\left\{F \in \mathcal{F}: \lambda_{\min }\left(W_{F}\right) \geq \delta_{W U}, \lambda_{\min }\left(U_{F}\right) \geq \delta_{W U},\left\|W_{F}\right\| \leq M_{W U}, \text { and }\left\|U_{F}\right\| \leq M_{W U}\right\} \tag{10.5}
\end{equation*}
$$

for some constants $\delta_{W U}>0$ and $M_{W U}<\infty$, where $\mathcal{F}$ is defined in (3.3). The set $\mathcal{F}_{W U} \cap \mathcal{F}_{0}$ is used to establish results for Kleibergen's LM and the CLR test with moment-variance weighting, considered in Section 6, using the fact that $\mathcal{F}_{0}=\mathcal{F}_{W U} \cap \mathcal{F}_{0}$ for $\delta_{W U}>0$ sufficiently small and $M_{W U}<\infty$ sufficiently large. This holds because for all $F \in \mathcal{F}_{0}, \lambda_{\min }\left(W_{F}\right)=\lambda_{\min }\left(\Omega_{F}^{-1 / 2}\right)=$ $\lambda_{\max }^{-1 / 2}\left(\Omega_{F}\right) \geq\left\|\Omega_{F}\right\|^{-1 / 2} \geq M_{*}^{-1 / 2}$ for some $M_{*}<\infty$ (because $\left\|\Omega_{F}\right\|=\left\|E_{F} g_{i} g_{i}^{\prime}\right\| \leq M_{*}$ for some $M_{*}<\infty$ by the moment conditions in $\mathcal{F}$ ), $\left\|W_{F}\right\|=\left\|\Omega_{F}^{-1 / 2}\right\| \leq \lambda_{\min }^{-1 / 2}\left(\Omega_{F}\right) \leq \delta^{-1 / 2}$ (using the $\lambda_{\min }\left(E_{F} g_{i} g_{i}^{\prime}\right) \geq \delta$ condition in $\left.\mathcal{F}\right)$, where $\delta>0, \lambda_{\min }\left(U_{F}\right)=\lambda_{\min }\left(I_{p}\right)=1$, and $\left\|U_{F}\right\|=\left\|I_{p}\right\|=p$.

### 10.3 Reparametrization

To apply Proposition 10.1, we reparametrize the null distribution $F$ to a vector $\lambda$. The vector $\lambda$ is chosen such that for a subvector of $\lambda$ convergence of a drifting subsequence of the subvector (after suitable renormalization) yields convergence in distribution of the test statistic and convergence in distribution of the critical value in the case of the CLR tests.

To be consistent with the use of general weight matrices $\widehat{W}_{n}$ and $\widehat{U}_{n}$ in this section, we provide more general definitions of $\tau_{j F}, B_{F}$, and $C_{F}$ here than are given in Section 3. These general definitions reduce to the definitions given in Section 3 when $W_{F}=\Omega_{F}^{-1 / 2}$ and $U_{F}=I_{p}$.

The vector $\lambda$ depends on the following quantities. Let
$B_{F}$ denote a $p \times p$ orthogonal matrix of eigenvectors of $U_{F}^{\prime}\left(E_{F} G_{i}\right)^{\prime} W_{F}^{\prime} W_{F}\left(E_{F} G_{i}\right) U_{F}$
ordered so that the corresponding eigenvalues $\left(\kappa_{1 F}, \ldots, \kappa_{p F}\right)$ are nonincreasing. The matrix $B_{F}$ is such that the columns of $W_{F}\left(E_{F} G_{i}\right) U_{F} B_{F}$ are orthogonal. Let

$$
\begin{equation*}
C_{F} \text { denote a } k \times k \text { orthogonal matrix of eigenvectors of } W_{F}\left(E_{F} G_{i}\right) U_{F} U_{F}^{\prime}\left(E_{F} G_{i}\right)^{\prime} W_{F}^{\prime} \tag{10.7}
\end{equation*}
$$

ordered so that the corresponding eigenvalues are $\left(\kappa_{1 F}, \ldots, \kappa_{p F}, 0, \ldots, 0\right) \in R^{k}$. The matrices $B_{F}$ and $C_{F}$ are not uniquely defined. We let $B_{F}$ denote one choice of the matrix of eigenvectors of $U_{F}^{\prime}\left(E_{F} G_{i}\right)^{\prime} W_{F}^{\prime} W_{F}\left(E_{F} G_{i}\right) U_{F}$ and analogously for $C_{F}$. Let

$$
\begin{equation*}
\left(\tau_{1 F}, \ldots, \tau_{p F}\right) \text { denote the } p \text { singular values of } W_{F}\left(E_{F} G_{i}\right) U_{F} \text {, } \tag{10.8}
\end{equation*}
$$

which are nonnegative, ordered so that $\tau_{j F}$ is nonincreasing. (Some of these singular values may be zero.) As is well-known, the squares of the $p$ singular values of a $k \times p$ matrix $A$ with $k \geq p$ equal the $p$ eigenvalues of $A^{\prime} A$ and the largest $p$ eigenvalues of $A A^{\prime}$. In consequence, $\kappa_{j F}=\tau_{j F}^{2}$ for $j=1, \ldots, p$.

Define the elements of $\lambda$ to be

$$
\begin{align*}
\lambda_{1, F} & :=\left(\tau_{1 F}, \ldots, \tau_{p F}\right)^{\prime} \in R^{p}, \\
\lambda_{2, F} & :=B_{F} \in R^{p \times p}, \\
\lambda_{3, F} & :=C_{F} \in R^{k \times k}, \\
\lambda_{4, F} & :=\left(E_{F} G_{i 1}, \ldots, E_{F} G_{i p}\right) \in R^{k \times p}, \\
\lambda_{5, F} & :=E_{F}\binom{g_{i}}{\operatorname{vec}\left(G_{i}\right)}\binom{g_{i}}{\operatorname{vec}\left(G_{i}\right)}^{\prime} \in R^{(p+1) k \times(p+1) k}, \\
\lambda_{6, F} & =\left(\lambda_{6,1 F}, \ldots, \lambda_{6,(p-1) F}\right)^{\prime}:=\left(\frac{\tau_{2 F}}{\tau_{1 F}}, \ldots, \frac{\tau_{p F}}{\tau_{(p-1) F}}\right)^{\prime} \in R^{p-1}, \text { where } 0 / 0:=0, \\
\lambda_{7, F} & :=W_{2 F}, \\
\lambda_{8, F} & :=U_{2 F}, \\
\lambda_{9, F} & :=F, \text { and } \\
\lambda & =\lambda_{F}:=\left(\lambda_{1, F}, \ldots, \lambda_{9, F}\right) . \tag{10.9}
\end{align*}
$$

For simplicity, when writing $\lambda=\left(\lambda_{1, F}, \ldots, \lambda_{9, F}\right)$, we allow the elements to be scalars, vectors, matrices, and distributions and likewise in similar expressions. If $p=1$, no vector $\lambda_{6, F}$ appears in $\lambda$ because $\lambda_{1, F}$ only contains a single element. The vector $\lambda_{6, F}$ is only used in the proofs for CLR tests. It could be deleted when considering only an LM test. The dimensions of $W_{2 F}$ and $U_{2 F}$ depend on the choices of $\widehat{W}_{n}=W_{1}\left(\widehat{W}_{2 n}\right)$ and $\widehat{U}_{n}=U_{1}\left(\widehat{U}_{2 n}\right)$. We let $\lambda_{5, g F}$ denote the upper left $k \times k$ submatrix of $\lambda_{5, F}$. Thus, $\lambda_{5, g F}=E_{F} g_{i} g_{i}^{\prime}=\Omega_{F}$.

We consider the parameter space $\Lambda_{0}$ for $\lambda$, which corresponds to $\mathcal{F}_{W U} \cap \mathcal{F}_{0}$, where $\mathcal{F}_{W U}$ and $\mathcal{F}_{0}$ are defined in (10.5) and (3.9), respectively. The parameter space $\Lambda_{0}$ and the function $h_{n}(\lambda)$ are defined by

$$
\begin{align*}
\Lambda_{0} & :=\left\{\lambda: \lambda=\left(\lambda_{1, F}, \ldots, \lambda_{9, F}\right) \text { for some } F \in \mathcal{F}_{W U} \cap \mathcal{F}_{0}\right\} \text { and } \\
h_{n}(\lambda) & :=\left(n^{1 / 2} \lambda_{1, F}, \lambda_{2, F}, \lambda_{3, F}, \lambda_{4, F}, \lambda_{5, F}, \lambda_{6, F}, \lambda_{7, F}, \lambda_{8, F}\right) . \tag{10.10}
\end{align*}
$$

By the definition of $\mathcal{F}, \Lambda_{0}$ indexes distributions that satisfy the null hypothesis $H_{0}: \theta=\theta_{0}$. The dimension $J$ of $h_{n}(\lambda)$ equals the number of elements in $\left(\lambda_{1, F}, \ldots, \lambda_{8, F}\right)$. Redundant elements in $\left(\lambda_{1, F}, \ldots, \lambda_{8, F}\right)$, such as the redundant off-diagonal elements of the symmetric matrix $\lambda_{5, F}$, are not needed, but do not cause any problem. Note that two parameter spaces denoted by $\Lambda_{1}$ and $\Lambda_{2}$, which are larger than $\Lambda_{0}$, are considered for the two SR-CQLR tests analyzed in AG2. (We also use $\Lambda_{2}$ in this paper, see (10.11) below.)

We define $\lambda$ and $h_{n}(\lambda)$ as in 10.9 and 10.10 because, as shown below, the asymptotic distributions of the test statistics under a sequence $\left\{F_{n}: n \geq 1\right\}$ for which $h_{n}\left(\lambda_{F_{n}}\right) \rightarrow h \in H$ depend on the behavior of $\lim n^{1 / 2} \lambda_{1, F_{n}}$, as well as $\lim \lambda_{m, F_{n}}$ for $m=2, \ldots, 8$. For example, the LM statistic in 4.2 depends on $\widehat{\Omega}_{n}^{-1 / 2} \widehat{D}_{n}$, or equivalently, on $n^{1 / 2} \widehat{\Omega}_{n}^{-1 / 2} \widehat{D}_{n} B_{F_{n}} S_{n}$ (because projections are invariant to rescaling and right-hand side (rhs) transformations by nonsingular matrices), where $S_{n}$ is a pd diagonal matrix that is designed to make this quantity $O_{p}(1)$ and not $o_{p}(1)$. We show that this quantity is asymptotically equivalent to $n^{1 / 2} \Omega_{F_{n}}^{-1 / 2} \widehat{D}_{n} B_{F_{n}} S_{n}$. In turn, the latter quantity depends on $n^{1 / 2} \Omega_{F_{n}}^{-1 / 2} \widehat{G}_{n} B_{F_{n}}=n^{1 / 2} \Omega_{F_{n}}^{-1 / 2}\left(\widehat{G}_{n} B_{F_{n}}-E_{F_{n}} G_{i} B_{F_{n}}\right)+n^{1 / 2} \Omega_{F_{n}}^{-1 / 2} E_{F_{n}} G_{i} B_{F_{n}}$. The quantity $\operatorname{vec}\left(n^{1 / 2} \Omega_{F_{n}}^{-1 / 2}\left(\widehat{G}_{n} B_{F_{n}}-E_{F_{n}} G_{i} B_{F_{n}}\right)\right)$ has a nondegenerate asymptotic normal distribution by the central limit theorem (CLT), using the behavior of $\lim \lambda_{s, F_{n}}$ for $s=2,4,5$, the fact that $B_{F_{n}}$ is an orthogonal matrix, and the restriction in $\mathcal{F}_{0}$. Hence, the asymptotic behavior of $\operatorname{vec}\left(n^{1 / 2} \Omega_{F_{n}}^{-1 / 2} \widehat{G}_{n} B_{F_{n}}\right)$ depends on that of $n^{1 / 2} \Omega_{F_{n}}^{-1 / 2} E_{F_{n}} G_{i} B_{F_{n}}$. Using the SVD of $\Omega_{F_{n}}^{-1 / 2} E_{F_{n}} G_{i}$, the latter is shown below to equal $\lambda_{3, F_{n}} \operatorname{Diag}\left\{n^{1 / 2} \lambda_{1, F_{n}}\right\}$, where $\operatorname{Diag}\left\{n^{1 / 2} \lambda_{1, F_{n}}\right\}$ denotes the $k \times p$ matrix with $n^{1 / 2} \lambda_{1, F_{n}}$ on the main diagonal and zeros elsewhere.

In Example 1 of Section 10.2 applied to the linear model 2.2), we have $W_{F}=\Omega_{F}^{-1 / 2}$ and $\tau_{j F}$ is the $j$ th singular value of $-\Omega_{F}^{-1 / 2} E_{F} Z_{i} Y_{2 i}^{\prime}=-\Omega_{F}^{-1 / 2} E_{F} Z_{i} Z_{i}^{\prime} \pi$, where $\Omega_{F}=E_{F} u_{i}^{2} Z_{i} Z_{i}^{\prime}$ for $j=1, \ldots, p$. As is well known, if $\pi$ is close to zero, weak instrument problems occur. But, as we show, matrices $\pi$ that are close to being singular, without their columns being close to zero, also lead to weak IV problems. This is captured in the present set-up by $\tau_{p F}$ being close to zero in the sense that $\lim n^{1 / 2} \tau_{p F_{n}}<\infty$. If this occurs, then weak identification problems arise.

For notational convenience,
$\left\{\lambda_{n, h}: n \geq 1\right\}$ denotes a sequence $\left\{\lambda_{n} \in \Lambda_{2}: n \geq 1\right\}$ for which $h_{n}\left(\lambda_{n}\right) \rightarrow h \in H$, where

$$
\begin{equation*}
\Lambda_{2}:=\left\{\lambda: \lambda=\left(\lambda_{1, F}, \ldots, \lambda_{9, F}\right) \text { for some } F \in \mathcal{F}_{W U}\right\} \tag{10.11}
\end{equation*}
$$

and $H$ is defined in (10.1) with $\Lambda$ replaced by $\Lambda_{2}$. Analogously, for any subsequence $\left\{w_{n}: n \geq 1\right\}$, $\left\{\lambda_{w_{n}, h}: n \geq 1\right\}$ denotes a sequence $\left\{\lambda_{w_{n}} \in \Lambda_{2}: n \geq 1\right\}$ for which $h_{w_{n}}\left(\lambda_{w_{n}}\right) \rightarrow h \in H$. By definition, $\Lambda_{0} \subset \Lambda_{2}$. We use the parameter space $\Lambda_{2}$ in many places in the paper, rather than $\Lambda_{0}$, for two reasons. First, this makes it clear where the conditions specified in $\mathcal{F}_{0}$ (and $\Lambda_{0}$ ) are really needed. Second, some of the results given here are used in AG2, which does not employ the smaller set $\Lambda_{0}$, but does use $\Lambda_{2}$. By the definitions of $\Lambda_{2}$ and $\mathcal{F}_{W U},\left\{\lambda_{n, h}: n \geq 1\right\}$ is a sequence of distributions that satisfies the null hypothesis $H_{0}: \theta=\theta_{0}$.

We decompose $h$ (defined by (10.1), 10.9), and 10.10) analogously to the decomposition of the first eight components of $\lambda: h=\left(h_{1}, \ldots, h_{8}\right)$, where $\lambda_{m, F}$ and $h_{m}$ have the same dimensions
for $m=1, \ldots, 8$. We further decompose the vector $h_{1}$ as $h_{1}=\left(h_{1,1}, \ldots, h_{1, p}\right)^{\prime}$, where the elements of $h_{1}$ could equal $\infty$. We decompose $h_{6}$ as $h_{6}=\left(h_{6,1}, \ldots, h_{6, p-1}\right)^{\prime}$. In addition, we let $h_{5, g}$ denote the upper left $k \times k$ submatrix of $h_{5}$. In consequence, under a sequence $\left\{\lambda_{n, h}: n \geq 1\right\}$, we have

$$
\begin{align*}
n^{1 / 2} \tau_{j F_{n}} & \rightarrow h_{1, j} \geq 0 \forall j \leq p, \lambda_{m, F_{n}} \rightarrow h_{m} \forall m=2, \ldots, 8, \\
\lambda_{5, g F_{n}} & =\Omega_{F_{n}}=E_{F_{n}} g_{i} g_{i}^{\prime} \rightarrow h_{5, g}, \text { and } \lambda_{6, j F_{n}} \rightarrow h_{6, j} \forall j=1, \ldots, p-1 . \tag{10.12}
\end{align*}
$$

By the conditions in $\mathcal{F}$, defined in (3.3), $h_{5, g}$ is pd.
The smallest and largest singular values of $W_{F}\left(E_{F} G_{i}\right) U_{F}$ (i.e., $\tau_{p F}$ and $\tau_{1 F}$ ) can be related to those of $E_{F} G_{i}$ (i.e., $s_{p F}$ and $s_{1 F}$ ) for $F \in \mathcal{F}_{W U}$ via

$$
\begin{equation*}
c_{1} s_{j F} \leq \tau_{j F} \leq c_{2} s_{j F} \text { for } j=1 \text { and } j=p \text { for some constants } 0<c_{1}<c_{2}<\infty \tag{10.13}
\end{equation*}
$$

that do not depend on $F$. As shown below, the parameter $\theta$ is strongly or semi-strongly identified under $\left\{\lambda_{n, h}: n \geq 1\right\}$ if $\lim n^{1 / 2} \tau_{p F_{n}}=\infty$. In consequence of 10.13), this holds iff $\lim n^{1 / 2} s_{p F_{n}}=$ $\infty$. The parameters are weakly identified in the standard sense if $\lim n^{1 / 2} \tau_{j F_{n}}<\infty \forall j \leq p$ or, equivalently, if $\lim n^{1 / 2} \tau_{1 F_{n}}<\infty$, which holds by (10.13) iff $\lim n^{1 / 2} s_{1 F_{n}}<\infty$. The parameters are weakly identified in the non-standard sense if $\lim n^{1 / 2} \tau_{1 F_{n}}=\infty$ and $\lim n^{1 / 2} \tau_{p F_{n}}<\infty$, which holds by 10.13 iff $\lim n^{1 / 2} s_{1 F_{n}}=\infty$ and $\lim n^{1 / 2} s_{p F_{n}}<\infty$.

The proof of 10.13 ) is as follows. For notational simplicity, we drop the subscript $F$ in some of the calculations. We have

$$
\begin{align*}
& \lambda_{\min }\left(U^{\prime} E G_{i}^{\prime} W^{\prime} W E G_{i} U\right) \\
= & \min _{\lambda:\|\lambda\|=1}(U \lambda /\|U \lambda\|)^{\prime} E G_{i}^{\prime} W^{\prime} W E G_{i}(U \lambda /\|U \lambda\|) \cdot\|U \lambda\|^{2} \\
\leq & \min _{\lambda:\|\lambda\|=1} \lambda^{\prime} E G_{i}^{\prime} W^{\prime} W E G_{i} \lambda \cdot \lambda_{\max }\left(U^{\prime} U\right) \\
= & \min _{\lambda:\|\lambda\|=1}\left(E G_{i} \lambda /\left\|E G_{i} \lambda\right\|\right)^{\prime} W^{\prime} W\left(E G_{i} \lambda /\left\|E G_{i} \lambda\right\|\right) \cdot\left\|E G_{i} \lambda\right\| \|^{2} \cdot \lambda_{\max }\left(U^{\prime} U\right) \\
\leq & \lambda_{\max }\left(W^{\prime} W\right) \lambda_{\min }\left(E G_{i}^{\prime} E G_{i}\right) \lambda_{\max }\left(U^{\prime} U\right) \\
\leq & c_{2}^{2} \lambda_{\min }\left(E G_{i}^{\prime} E G_{i}\right), \text { where } \\
c_{2}:= & \sup _{F \in \mathcal{F}_{W U}}\left[\lambda_{\max }\left(W_{F}^{\prime} W_{F}\right) \lambda_{\max }\left(U_{F}^{\prime} U_{F}\right)\right]^{1 / 2}<\infty \tag{10.14}
\end{align*}
$$

and the last inequality holds by the conditions in $\mathcal{F}_{W U}$ (defined in 10.5). Because the smallest eigenvalues of $U^{\prime} E G_{i}^{\prime} W^{\prime} W E G_{i} U$ and $E G_{i}^{\prime} E G_{i}$ equal the squares of the smallest singular values of $W E G_{i} U$ and $E G_{i}$, respectively, 10.14 establishes the second inequality in (10.13) for $j=p$. Analogous calculations establish the lower bound in (10.14) for $j=p$ and the bounds for $j=1$
by replacing min and $\leq$ by max and $\geq$, respectively, in the appropriate places and taking $c_{1}:=$ $\inf _{F \in \mathcal{F}_{W U}}\left[\lambda_{\min }\left(W_{F}^{\prime} W_{F}\right) \lambda_{\min }\left(U_{F}^{\prime} U_{F}\right)\right]^{1 / 2}>0$.

### 10.4 Assumption WU

We assume that the random weight matrices $\widehat{W}_{n}=W_{1}\left(\widehat{W}_{2 n}\right)$ and $\widehat{U}_{n}=U_{1}\left(\widehat{U}_{2 n}\right)$ defined in (10.2) satisfy the following assumption that depends on a suitably chosen parameter space $\Lambda_{*}$ $\left(\subset \Lambda_{2}\right)$, such as $\Lambda_{2}, \Lambda_{0}$, or $\Lambda_{1}$.

Assumption WU for the parameter space $\Lambda_{*} \subset \Lambda_{2}$ : Under all subsequences $\left\{w_{n}\right\}$ and all sequences $\left\{\lambda_{w_{n}, h}: n \geq 1\right\}$ with $\lambda_{w_{n}, h} \in \Lambda_{*}$,
(a) $\widehat{W}_{2 w_{n}} \rightarrow_{p} h_{7}\left(:=\lim W_{2 F_{w_{n}}}\right)$,
(b) $\widehat{U}_{2 w_{n}} \rightarrow_{p} h_{8}\left(:=\lim U_{2 F_{w_{n}}}\right)$, and
(c) $W_{1}(\cdot)$ is a continuous function at $h_{7}$ on some set $\mathcal{W}_{2}$ that contains $\left\{\lambda_{7, F}\left(=W_{2 F}\right): \lambda=\right.$ $\left.\left(\lambda_{1, F}, \ldots, \lambda_{9, F}\right) \in \Lambda_{*}\right\}$ and contains $\widehat{W}_{2 w_{n}} \mathrm{wp} \rightarrow 1$ and $U_{1}(\cdot)$ is a continuous function at $h_{8}$ on some set $\mathcal{U}_{2}$ that contains $\left\{\lambda_{8, F}\left(=U_{2 F}\right): \lambda=\left(\lambda_{1, F}, \ldots, \lambda_{9, F}\right) \in \Lambda_{*}\right\}$ and contains $\widehat{U}_{2 w_{n}} \mathrm{wp} \rightarrow 1$.

In Assumption WU and elsewhere below, "all sequences $\left\{\lambda_{w_{n}, h}: n \geq 1\right\}$ " means "all sequences $\left\{\lambda_{w_{n}, h}: n \geq 1\right\}$ for any $h \in H "$ and likewise with $n$ in place of $w_{n}$. Note that, by definition, a sequence $\left\{\lambda_{w_{n}, h}: n \geq 1\right\}$ determines a sequence of distributions $\left\{F_{w_{n}}: n \geq 1\right\}$, see 10.9.

Assumption WU for the parameter space $\Lambda_{0}$ is verified in Comment (ii) to Theorem 12.1 given below for the CLR test with moment-variance weighting, which is considered in Section 6. It also holds for Kleibergen's LM test (for the same parameter space $\Lambda_{0}$ ) by the same argument (because $\widehat{W}_{2 n}, \widehat{U}_{2 n}, W_{1}(\cdot)$, and $U_{1}(\cdot)$ are the same for these two tests, see 10.4$)$.

### 10.5 Basic Results

For any square-integrable random vector $a_{i}$ and $F, F_{n} \in \mathcal{F}$, define

$$
\begin{equation*}
\Phi_{F}^{a_{i}}:=\operatorname{Var}_{F}\left(a_{i}-\left(E_{F} a_{\ell} g_{\ell}^{\prime}\right) \Omega_{F}^{-1} g_{i}\right) \text { and } \Phi_{h}^{a_{i}}:=\lim \Phi_{F_{w_{n}}}^{a_{i}} \tag{10.15}
\end{equation*}
$$

whenever the limit exists, where the distributions $\left\{F_{w_{n}}: n \geq 1\right\}$ correspond to $\left\{\lambda_{w_{n}, h}: n \geq 1\right\}$ for any subsequence $\left\{w_{n}: n \geq 1\right\}$. Note that $\Phi_{F}^{a_{i}}=\Psi_{F}^{a_{i}}-E_{F} a_{i} E_{F} a_{i}^{\prime}$ (because $\Psi_{F}^{a_{i}}=E_{F} b_{i} b_{i}^{\prime}$ for $b_{i}=a_{i}-\left(E_{F} a_{\ell} g_{\ell}^{\prime}\right) \Omega_{F}^{-1} g_{i}$ and $\left.E_{F} g_{i}=0^{k}\right)$.

A basic result that is used in the proofs of results for all of the tests considered in this paper and AG2 is the following.

Lemma 10.2 Under all sequences $\left\{\lambda_{n, h}: n \geq 1\right\}$,

$$
n^{1 / 2}\binom{\widehat{g}_{n}}{\operatorname{vec}\left(\widehat{D}_{n}-E_{F_{n}} G_{i}\right)} \rightarrow_{d}\binom{\bar{g}_{h}}{\operatorname{vec}\left(\bar{D}_{h}\right)} \sim N\left(0^{(p+1) k},\left(\begin{array}{cc}
h_{5, g} & 0^{k \times p k} \\
0^{p k \times k} & \Phi_{h}^{v e c}\left(G_{i}\right)
\end{array}\right)\right) .
$$

Under all subsequences $\left\{w_{n}\right\}$ and all sequences $\left\{\lambda_{w_{n}, h}: n \geq 1\right\}$, the same result holds with $n$ replaced with $w_{n}$.

Comments: (i) The variance matrix $\Phi_{h}^{v e c\left(G_{i}\right)}$ depends on $h$ only through $h_{4}$ and $h_{5}$. The assumptions allow $\Phi_{h}^{v e c\left(G_{i}\right)}$ to be singular.
(ii) Suppose one eliminates the $\lambda_{\min }\left(E_{F} g_{i} g_{i}^{\prime}\right) \geq \delta$ condition in $\mathcal{F}$ and one defines $\widehat{D}_{n}$ in 3.2) with $\widehat{\Omega}_{n}$ replaced by an eigenvalue-adjusted matrix, denoted by $\widehat{\Omega}_{n}^{\varepsilon}$, which is constructed to have its smallest eigenvalue greater than or equal to $\varepsilon>0$ multiplied by its largest eigenvalue, see AG2 for the details of such a construction. In this case, the result of Lemma 10.2 still holds and all of the other asymptotic results following from Lemma 10.2 still hold, except the independence of $\bar{g}_{h}$ and $\bar{D}_{h}$. However, this independence is key because it is used in the conditioning argument that establishes the correct asymptotic size of all of the tests that are shown to have correct asymptotic size. Without it, these tests do not necessarily have correct asymptotic size. In consequence, we define $\widehat{D}_{n}$ in (3.2) using $\widehat{\Omega}_{n}$, not $\widehat{\Omega}_{n}^{\varepsilon}$.

The reason that independence does not necessarily hold when $\widehat{D}_{n}$ is defined using $\widehat{\Omega}_{n}^{\varepsilon}$, rather than $\widehat{\Omega}_{n}$, is that the covariance term $E_{F_{n}}\left[G_{i j}-E_{F_{n}} G_{i j}-\left(E_{F_{n}} G_{\ell j} g_{\ell}^{\prime}\right)\left(\Omega_{F_{n}}^{\varepsilon}\right)^{-1} g_{i}\right] g_{i}^{\prime}$ typically does not equal $0^{k \times k}$ when $\Omega_{F_{n}}^{\varepsilon} \neq \Omega_{F_{n}}$, whereas $E_{F_{n}}\left[G_{i j}-E_{F_{n}} G_{i j}-\left(E_{F_{n}} G_{\ell j} g_{\ell}^{\prime}\right) \Omega_{F_{n}}^{-1} g_{i}\right] g_{i}^{\prime}$ necessarily equals $0^{k \times k}$, see the proof of Lemma 10.2 in Section 15 below for more details.
(iii) The proofs of Lemma 10.2 and other results in this section are given in Sections 1517 below.

The following is a key definition. Consider a sequence $\left\{\lambda_{n, h}: n \geq 1\right\}$. Let $q=q_{h}(\in\{0, \ldots, p\})$ be such that

$$
\begin{equation*}
h_{1, j}=\infty \text { for } 1 \leq j \leq q_{h} \text { and } h_{1, j}<\infty \text { for } q_{h}+1 \leq j \leq p, \tag{10.16}
\end{equation*}
$$

where $h_{1, j}:=\lim n^{1 / 2} \tau_{j F_{n}} \geq 0$ for $j=1, \ldots, p$ by 10.12 and the distributions $\left\{F_{n}: n \geq 1\right\}$ correspond to $\left\{\lambda_{n, h}: n \geq 1\right\}$ defined in 10.11. Such a $q$ exists because $\left\{h_{1, j}: j \leq p\right\}$ are nonincreasing in $j$ (since $\left\{\tau_{j F}: j \leq p\right\}$ are the ordered singular values of $W_{F}\left(E_{F} G_{i}\right) U_{F}$, as defined in (10.8). As defined, $q$ is the number of singular values of $W_{F_{n}}\left(E_{F_{n}} G_{i}\right) U_{F_{n}}$ that diverge to infinity when multiplied by $n^{1 / 2}$. Roughly speaking, $q$ is the number of parameters, or one-to-one transformations of the parameters, that are strongly or semi-strongly identified.

The following quantities appear in Lemma 10.3 below, which gives the asymptotic distribution
of $\widehat{D}_{n}$ after suitable rotations and rescaling, but without the recentering (by subtracting $E_{F_{n}} G_{i}$ ) that appears in Lemma 10.2. We partition $h_{2}$ and $h_{3}$ and define $\bar{\Delta}_{h}$ as follows:

$$
\begin{align*}
h_{2} & =\left(h_{2, q}, h_{2, p-q}\right), h_{3}=\left(h_{3, q}, h_{3, k-q}\right), h_{1, p-q}^{\diamond}:=\left[\begin{array}{c}
0^{q \times(p-q)} \\
\operatorname{Diag}\left\{h_{1, q+1}, \ldots, h_{1, p}\right\} \\
0^{(k-p) \times(p-q)}
\end{array}\right] \in R^{k \times(p-q)}, \\
\bar{\Delta}_{h} & =\left(\bar{\Delta}_{h, q}, \bar{\Delta}_{h, p-q}\right) \in R^{k \times p}, \bar{\Delta}_{h, q}:=h_{3, q}, \bar{\Delta}_{h, p-q}:=h_{3} h_{1, p-q}^{\diamond}+h_{71} \bar{D}_{h} h_{81} h_{2, p-q}, \\
h_{71} & :=W_{1}\left(h_{7}\right), \text { and } h_{81}:=U_{1}\left(h_{8}\right), \tag{10.17}
\end{align*}
$$

where $h_{2, q} \in R^{p \times q}, h_{2, p-q} \in R^{p \times(p-q)}, h_{3, q} \in R^{k \times q}, h_{3, k-q} \in R^{k \times(k-q)}, \bar{\Delta}_{h, q} \in R^{k \times q}, \bar{\Delta}_{h, p-q} \in$ $R^{k \times(p-q)}, h_{71} \in R^{k \times k}, h_{81} \in R^{p \times p}$, and $\bar{D}_{h}$ is defined in Lemma 10.2. For simplicity, there is some abuse of notation here, e.g., $h_{2, q}$ and $h_{2, p-q}$ denote different matrices even if $p-q$ happens to equal q. Note that when Assumption WU holds $h_{71}=\lim W_{F_{n}}=\lim W_{1}\left(W_{2 F_{n}}\right)$ and $h_{81}=\lim U_{F_{n}}=$ $\lim U_{1}\left(U_{2 F_{n}}\right)$ under $\left\{\lambda_{n, h}: n \geq 1\right\}$.

The case where $q=p$ (i.e., $n^{1 / 2} \tau_{j F_{n}} \rightarrow \infty$ for all $j \leq p$ ) is the strong or semi-strong identification case. In this case, no $h_{2, p-q}, h_{1, p-q}^{\diamond}$, and $\bar{\Delta}_{h, p-q}$ matrices appear in 10.17, $\bar{\Delta}_{h}=h_{3, q}=h_{3, p}$, and $\bar{\Delta}_{h}$ is non-random. In consequence, the limit in distribution (or probability) of the normalized matrix $n^{1 / 2} W_{F_{n}} \widehat{D}_{n} U_{F_{n}} T_{n}$, where $T_{n} \in R^{p \times p}$ is defined below, is non-random, see Lemma 10.3 below. When $q<p$, identification is weak and the limit of this matrix is random.

Now we provide some motivation for Lemma 10.3, which is stated below. To show that the LM statistic has a $\chi_{p}^{2}$ asymptotic distribution we need to determine the asymptotic behavior of $\widehat{D}_{n}$ without the recentering by $E_{F_{n}} G_{i}$ that occurs in Lemma 10.2. In addition, to determine the asymptotic distribution of the $r k_{n}$ statistic in (6.2), we need to determine the asymptotic distribution of $W_{F_{n}} \widehat{D}_{n} U_{F_{n}}$ without recentering by $E_{F_{n}} G_{i}$. (Furthermore, to determine the asymptotic distributions of the two SR-CQLR test statistics and conditional critical values considered in AG2, we need to determine the asymptotic distribution of $W_{F_{n}} \widehat{D}_{n} U_{F_{n}}$ without recentering by $E_{F_{n}} G_{i}$.) To do so, we post-multiply $W_{F_{n}} \widehat{D}_{n} U_{F_{n}}$ first by $B_{F_{n}}$ and then by a nonrandom diagonal matrix $S_{n} \in R^{p \times p}$ (which may depend on $F_{n}$ and $h$ ). The matrix $S_{n}$ rescales the columns of $W_{F_{n}} \widehat{D}_{n} U_{F_{n}} B_{F_{n}}$ to ensure that $n^{1 / 2} W_{F_{n}} \widehat{D}_{n} U_{F_{n}} B_{F_{n}} S_{n}$ converges in distribution to a (possibly) random matrix that is finite a.s. and not almost surely zero. For $F \in \mathcal{F}_{W U} \cap \mathcal{F}_{0}$, it ensures that the (possibly) random limit matrix has full column rank with probability one. For example, in the case of the LM statistic, these transformations are applied with $W_{F_{n}}=\Omega_{F_{n}}^{-1 / 2}$ and $U_{F_{n}}=I_{p}$.

For the LM statistic and the CLR statistics that employ it, we need the full column rank property of the limit random matrix in order to apply the continuous mapping theorem (CMT).

For the LM statistic, the full rank property ensures that the quantity $\widehat{D}_{n}^{\prime} \widehat{\Omega}_{n}^{-1} \widehat{D}_{n}$ (whose inverse appears in the expression for $L M_{n}$, see (4.2), is nonsingular asymptotically with probability one after $\widehat{D}_{n}$ has been transformed and rescaled to yield $n^{1 / 2} \Omega_{F_{n}}^{-1 / 2} \widehat{D}_{n} B_{F_{n}} S_{n}$. Note that $P_{\widehat{\Omega}_{n}^{-1 / 2}} \widehat{D}_{n}$, which appears in the definition of $L M_{n}$ in (4.2), can be written as

$$
\begin{align*}
P_{\widehat{\Omega}_{n}^{-1 / 2}} \widehat{D}_{n}:= & \widehat{\Omega}_{n}^{-1 / 2} \widehat{D}_{n}\left(\widehat{D}_{n}^{\prime} \widehat{\Omega}_{n}^{-1} \widehat{D}_{n}\right)^{-1} \widehat{D}_{n}^{\prime} \widehat{\Omega}_{n}^{-1 / 2} \\
= & \left(\widehat{\Omega}_{n}^{-1 / 2} \Omega_{n}^{1 / 2}\right)\left(n^{1 / 2} \Omega_{n}^{-1 / 2} \widehat{D}_{n} T_{n}\right)\left[\left(n^{1 / 2} \Omega_{n}^{-1 / 2} \widehat{D}_{n} T_{n}\right)^{\prime}\left(\widehat{\Omega}_{n}^{-1 / 2} \Omega_{n}^{1 / 2}\right)^{\prime}\left(\widehat{\Omega}_{n}^{-1 / 2} \Omega_{n}^{1 / 2}\right)\right. \\
& \left.\times\left(n^{1 / 2} \Omega_{n}^{-1 / 2} \widehat{D}_{n} T_{n}\right)\right]^{-1}\left(n^{1 / 2} \Omega_{n}^{-1 / 2} \widehat{D}_{n} T_{n}\right)^{\prime}\left(\Omega_{n}^{1 / 2} \widehat{\Omega}_{n}^{-1 / 2}\right), \text { where } \\
T_{n}:= & B_{F_{n}} S_{n} \in R^{p \times p} \text { and } \Omega_{n}:=\Omega_{F_{n}}\left(=E_{F_{n}} g_{i} g_{i}^{\prime}\right), \tag{10.18}
\end{align*}
$$

provided $T_{n}$ has full rank and $\Omega_{n}$ is pd. In consequence, these transformations do not affect the value or distribution of the LM statistic.

Note that the two SR-CQLR test statistics considered in AG2 do not depend on an LM statistic and do not require the asymptotic distribution of $n^{1 / 2} W_{F_{n}} \widehat{D}_{n} U_{F_{n}} B_{F_{n}} S_{n}$ to have full column rank a.s.

Define

$$
\begin{equation*}
S_{n}:=\operatorname{Diag}\left\{\left(n^{1 / 2} \tau_{1 F_{n}}\right)^{-1}, \ldots,\left(n^{1 / 2} \tau_{q F_{n}}\right)^{-1}, 1, \ldots, 1\right\} \in R^{p \times p} \tag{10.19}
\end{equation*}
$$

where $q=q_{h}$ is defined in (10.16). Note that $\tau_{j F_{n}}>0$ for $n$ large for $j \leq q$ and, hence, $S_{n}$ is well defined for $n$ large, because $n^{1 / 2} \tau_{j F_{n}} \rightarrow \infty$ for all $j \leq q$.

The proof of Theorem 11.1 for the LM test, the proofs of Theorems 10.4 and 12.1 for the CLR test with moment-variance weighting, and the proofs for the two SR-CQLR tests in AG2 use the following lemma. The $p \times p$ matrix $T_{n}$ is defined in 10.18).

Lemma 10.3 Suppose Assumption WU holds for some non-empty parameter space $\Lambda_{*} \subset \Lambda_{2}$. Under all sequences $\left\{\lambda_{n, h}: n \geq 1\right\}$ with $\lambda_{n, h} \in \Lambda_{*}$,

$$
n^{1 / 2}\left(\widehat{g}_{n}, \widehat{D}_{n}-E_{F_{n}} G_{i}, W_{F_{n}} \widehat{D}_{n} U_{F_{n}} T_{n}\right) \rightarrow_{d}\left(\bar{g}_{h}, \bar{D}_{h}, \bar{\Delta}_{h}\right)
$$

where (a) $\left(\bar{g}_{h}, \bar{D}_{h}\right)$ are defined in Lemma 10.2, (b) $\bar{\Delta}_{h}$ is the nonrandom function of $h$ and $\bar{D}_{h}$ defined in 10.17), (c) ( $\bar{D}_{h}, \bar{\Delta}_{h}$ ) and $\bar{g}_{h}$ are independent, (d) if Assumption WU holds with $\Lambda_{*}=\Lambda_{0}$, $W_{F}=\Omega_{F}^{-1 / 2}$, and $U_{F}=I_{p}$, then $\bar{\Delta}_{h}$ has full column rank $p$ with probability one, and (e) under all subsequences $\left\{w_{n}\right\}$ and all sequences $\left\{\lambda_{w_{n}, h}: n \geq 1\right\}$ with $\lambda_{w_{n}, h} \in \Lambda_{*}$, the convergence result above and the results of parts (a)-(d) hold with $n$ replaced with $w_{n}$.

Comments: (i) Lemma 10.3 (c)-(d) are key properties of the asymptotic distribution of $n^{1 / 2}\left(\widehat{g}_{n}\right.$,
$\left.W_{F_{n}} \widehat{D}_{n} U_{F_{n}} T_{n}\right)$ that lead to the LM statistic having a $\chi_{p}^{2}$ asymptotic distribution and the CLR test with moment-variance weighting having correct asymptotic size. Lemma 10.3 (c) is a key property that leads to the correct asymptotic size of the two SR-CQLR tests in AG2. Lemma 10.3(d) is not needed for these tests because they do not rely on an LM statistic.
(ii) The conditions in $\mathcal{F}_{0}$ are used in the proofs to obtain the result of Lemma 10.3(d) and are not used elsewhere in the proofs, except where Lemma 10.3 (d) is used.

The following theorems are used only for the CLR tests. For the proof of Theorem4.1 concerning Kleibergen's (2005) LM test, one can go from here to Section 11 .

Let

$$
\begin{equation*}
\widehat{\kappa}_{j n} \text { denote the } j \text { th eigenvalue of } n \widehat{U}_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} \widehat{U}_{n}, \forall j=1, \ldots, p, \tag{10.20}
\end{equation*}
$$

ordered to be nonincreasing in $j$. By definition, $\lambda_{\min }\left(n \widehat{U}_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} \widehat{U}_{n}\right)=\widehat{\kappa}_{p n}$. Also, the $j$ th singular value of $n^{1 / 2} \widehat{W}_{n} \widehat{D}_{n} \widehat{U}_{n}$ equals $\widehat{\kappa}_{j n}^{1 / 2}$.

Theorem 10.4 Suppose Assumption WU holds for some non-empty parameter space $\Lambda_{*} \subset \Lambda_{2}$. Under all sequences $\left\{\lambda_{n, h}: n \geq 1\right\}$ with $\lambda_{n, h} \in \Lambda_{*}$,
(a) $\widehat{\kappa}_{p n} \rightarrow_{p} \infty$ if $q=p$,
(b) $\widehat{\kappa}_{p n} \rightarrow_{d} \lambda_{\min }\left(\bar{\Delta}_{h, p-q}^{\prime} h_{3, k-q} h_{3, k-q}^{\prime} \bar{\Delta}_{h, p-q}\right)$ if $q<p$,
(c) $\widehat{\kappa}_{j n} \rightarrow_{p} \infty$ for all $j \leq q$,
(d) the (ordered) vector of the smallest $p-q$ eigenvalues of $n \widehat{U}_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} \widehat{U}_{n}$, i.e., $\left(\widehat{\kappa}_{(q+1) n}, \ldots\right.$, $\left.\widehat{\kappa}_{p n}\right)^{\prime}$, converges in distribution to the (ordered) $p-q$ vector of the eigenvalues of $\bar{\Delta}_{h, p-q}^{\prime} h_{3, k-q} h_{3, k-q}^{\prime}$ $\times \bar{\Delta}_{h, p-q} \in R^{(p-q) \times(p-q)}$,
(e) the convergence in parts (a)-(d) holds jointly with the convergence in Lemma 10.3, and
(f) under all subsequences $\left\{w_{n}\right\}$ and all sequences $\left\{\lambda_{w_{n}, h}: n \geq 1\right\}$ with $\lambda_{w_{n}, h} \in \Lambda_{*}$, the results in parts (a)-(e) hold with $n$ replaced with $w_{n}$.

Comments: (i) The statistic $\widehat{\kappa}_{p n}=\lambda_{\text {min }}\left(n \widehat{U}_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} \widehat{U}_{n}\right)$ in Theorem 10.4 (a) and (b) is a Robin and Smith (2000)-type rank statistic.
(ii) Theorem 10.4 (a) and (b) is used to determine the asymptotic behavior of the statistic $r k_{n}$ defined in 6.2) (which is employed by the CLR test with moment-variance weighting that is considered in Section 6). More specifically, Theorem 10.4(a) and (b) is used to verify Assumption $R$ in Section 12 below.
(iii) Theorem 10.4 (c) and (d) is used to determine the asymptotic behavior of the critical value functions for the two SR-CQLR tests considered in AG2 (with $\widehat{W}_{n}$ and $\widehat{U}_{n}$ defined suitably). Because Theorem 10.4 (c) and (d) are immediate by-products of the proofs of Theorem 10.4(a) and (b), they are stated and proved here, rather than in AG2.
(iv) The statement of Theorem 3 in Kleibergen (2005) is difficult to interpret because the expression given for the conditional asymptotic distribution of the CLR statistic involves Kleibergen's (2005) statistic $\mathrm{rk}\left(\theta_{0}\right)$, which is a finite-sample object. Based on Theorem 10.4, (12.7) below provides the asymptotic distribution of a class of CLR statistics in terms of an asymptotic version of the rank statistic employed, which is necessary for a precise statement of the asymptotic distribution. The class of CLR statistics considered are those defined in (5.1) and based on the rank statistic in Theorem 10.4 for some choices of $\widehat{W}_{n}$ and $\widehat{U}_{n}$, which is a Robin and Smith (2000)-type rank statistic. In particular, taking $\widehat{W}_{n}=\widehat{\Omega}_{n}^{-1 / 2}$ and $\widehat{U}_{n}=I_{p}$ gives the rank statistic defined in (6.2).

## 11 Asymptotic Size of the Nonlinear LM Test

In this section, we prove Theorem 4.1 for the LM test.
We state a theorem that verifies Assumption B* of ACG (stated in Section 10) for the LM test. The following theorem applies with $\widehat{W}_{n}=\widehat{\Omega}_{n}^{-1 / 2}, W_{F}=\Omega_{F}^{-1 / 2}$, and $\widehat{U}_{n}=U_{F}=I_{p}$. (These definitions affect the definition of $\lambda_{n, h}$, which appears in the theorem).

Theorem 11.1 The asymptotic null rejection probabilities of the nominal size $\alpha \in(0,1)$ LM test equal $\alpha$ under all subsequences $\left\{w_{n}\right\}$ and all sequences $\left\{\lambda_{w_{n}, h}: n \geq 1\right\}$ with $\lambda_{w_{n}, h} \in \Lambda_{0} \forall n \geq 1$.

Comments: (i) The requirement that $\lambda_{w_{n}, h} \in \Lambda_{0}$ (defined in 10.10) implies that the parameter space for $F$ is $\mathcal{F}_{0}$ (defined in (3.9) for the results given in Theorems 4.1 and 11.1 (because the restrictions in $\mathcal{F}_{W U}$ are not binding, see the discussion in the paragraph containing (10.5).
(ii) Proposition 10.1 and Theorem 11.1 prove Theorem 4.1 for the LM test. The proof of Theorem 4.1 for the LM CS is analogous, see Comments (i) and (ii) to Proposition 10.1.

For notational simplicity, we prove Theorem 11.1 for the sequence $\{n\}$, rather than a subsequence $\left\{w_{n}: n \geq 1\right\}$. We note here that the same proof holds for any subsequence $\left\{w_{n}: n \geq 1\right\}$.

Proof of Theorem 11.1. Let $\Omega_{n}:=\Omega_{F_{n}}$. We derive the limiting distribution of the statistic $L M_{n}$ using the CMT applied to $\Omega_{n}^{-1 / 2} n^{1 / 2} \widehat{g}_{n}, \widehat{\Omega}_{n}^{-1 / 2} \Omega_{n}^{1 / 2}$, and $n^{1 / 2} \Omega_{n}^{-1 / 2} \widehat{D}_{n} T_{n}$, where the latter two quantities appear in the expression on the rhs of 10.18. Note that $\widehat{\Omega}_{n} \rightarrow_{p} h_{5, g}$ by the WLLN, $\Omega_{n} \rightarrow$ $h_{5, g}$, and $h_{5, g}$ is pd. Thus, $\widehat{\Omega}_{n}^{-1 / 2} \Omega_{n}^{1 / 2} \rightarrow_{p} I_{k}$. By Lemma 10.3 applied with $W_{F}=\Omega_{F}^{-1 / 2}$ and $U_{F}=I_{p}$ (which results from taking $\widehat{W}_{n}=\widehat{\Omega}_{n}^{-1 / 2}$ and $\widehat{U}_{n}=I_{p}$ ), we get $\left(\Omega_{n}^{-1 / 2} n^{1 / 2} \widehat{g}_{n}, n^{1 / 2} \Omega_{n}^{-1 / 2} \widehat{D}_{n} T_{n}\right) \rightarrow_{d}$ $\left(h_{5, g}^{-1 / 2} \bar{g}_{h}, \bar{\Delta}_{h}\right)$. For the CMT to apply, it is enough to show that the function $f: R^{k \times p} \rightarrow R^{k \times k}$ defined by $f(D):=D\left(D^{\prime} D\right)^{-1} D^{\prime}$ for $D \in R^{k \times p}$ is continuous on a set $C \subset R^{k \times p}$ with $P\left(\bar{\Delta}_{h} \in\right.$ $C)=1$. This holds because the function $f_{2}(D, L):=L D\left((L D)^{\prime}(L D)\right)^{-1} D^{\prime} L^{\prime}$ for a nonsingular
$k \times k$ matrix $L$ is continuous at $\left(D, I_{k}\right)$ if $f(D)$ is continuous at $D$. Note that $f$ is continuous at each $D$ that has full column rank. And, by Lemma 10.3(d), $\bar{\Delta}_{h}$ has full column rank a.s. because $\lambda_{n, h} \in \Lambda_{0}, F_{n} \in \mathcal{F}_{0}, W_{F}=\Omega_{F}^{-1 / 2}$, and $U_{F}=I_{p}$. Hence, $f$ is continuous a.s. By $\widehat{\Omega}_{n}^{-1 / 2} \Omega_{n}^{1 / 2} \rightarrow_{p} I_{k}$, the convergence result in Lemma 10.3 , and the CMT, we have

$$
\begin{equation*}
P_{D_{n}^{\diamond}} \widehat{\Omega}_{n}^{-1 / 2} n^{1 / 2} \widehat{g}_{n}=D_{n}^{\diamond}\left(D_{n}^{\diamond \prime} D_{n}^{\diamond}\right)^{-1} D_{n}^{\diamond} \widehat{\Omega}_{n}^{-1 / 2} n^{1 / 2} \widehat{g}_{n} \rightarrow_{d} \bar{v}_{h}:=P_{\bar{\Delta}_{h}} h_{5, g}^{-1 / 2} \bar{g}_{h}, \tag{11.1}
\end{equation*}
$$

where $D_{n}^{\diamond}:=\left(\widehat{\Omega}_{n}^{-1 / 2} \Omega_{n}^{1 / 2}\right) n^{1 / 2} \Omega_{n}^{-1 / 2} \widehat{D}_{n} T_{n}$.
Conditional on $\bar{\Delta}_{h}, \bar{v}_{h}^{\prime} \bar{v}_{h}$ is distributed as $\chi_{p}^{2}$ because (i) $\bar{\Delta}_{h}$ and $\bar{g}_{h}$ are independent by property (c) in Lemma 10.3 , (ii) $h_{5, g}^{-1 / 2} \bar{g}_{h}$ is conditionally distributed as $N\left(0^{k}, I_{k}\right)$ by $\bar{g}_{h} \sim N\left(0^{k}, h_{5, g}\right)$ and (i), and (iii) $P_{\bar{\Delta}_{h}}$ is fixed given $\bar{\Delta}_{h}$ and projects onto a space of dimension $p$ a.s. by property (d) in Lemma 10.3. Because the $\chi_{p}^{2}$ distribution does not depend on $\bar{\Delta}_{h}, \bar{v}_{h}^{\prime} \bar{v}_{h}$ is unconditionally distributed as $\chi_{p}^{2}$ as well. In consequence, using the CMT again, we have

$$
\begin{equation*}
L M_{n} \rightarrow_{d} \overline{L M}_{h}:=\bar{v}_{h}^{\prime} \bar{v}_{h} \sim \chi_{p}^{2} . \tag{11.2}
\end{equation*}
$$

Given this result and the use of the $\chi_{p, 1-\alpha}^{2}$ critical value by the LM test, we obtain the conclusion of Theorem 11.1 for the LM test: $\lim P_{F_{n}}\left(L M_{n}>\chi_{p, 1-\alpha}^{2}\right)=\alpha$.

## 12 Asymptotic Size of the CLR Test with Moment-Variance Weighting

In this section, we prove Theorem 6.1, which concerns the CLR test (and CS) with momentvariance weighting based on the Robin-Smith rank statistic. In fact, for the CLR test defined by (5.1)- (5.2), we prove a stronger result than that given in Theorem 6.1. We establish Theorem 6.1 for a CLR test that is based on any rank statistic $r k_{n}$ that satisfies a high-level assumption, denoted Assumption R, not just the rank statistic $r k_{n}\left(\theta_{0}\right)$ defined in (6.2). Then, we verify Assumption R for the moment-variance-weighted Robin-Smith rank statistic $r k_{n}\left(\theta_{0}\right)$ in (6.2). Note that Assumption R does not hold for the rank statistic in (5.5) when $p \geq 2$.

Section 19.5 below provides additional asymptotic size results for equally-weighted CLR tests (and CS's), which are CLR tests that are based on $r k_{n}$ statistics that depend on $\widehat{D}_{n}$ only through $\widetilde{W}_{n} \widehat{D}_{n}$ for some $k \times k$ weighting matrix $\widetilde{W}_{n}$. These results show that equally-weighted CLR tests (and CS's) based on the Robin and Smith (2000) rank statistic with a general weight matrix $\widetilde{W}_{n}$ $\left(\in R^{k \times k}\right)$ have correct asymptotic size under suitable conditions on $\widetilde{W}_{n}$. One can view these results as verifying Assumption R for a broad class of $r k_{n}$ statistics. In contrast, the results in the present
section establish the correct asymptotic size of CLR tests (and CS's) under the high-level condition Assumption R and for the Robin and Smith (2000) rank statistic when $\widetilde{W}_{n}$ is the moment-variance weighting matrix $\widehat{\Omega}_{n}^{-1 / 2}$, see Comment (ii) to Theorem 12.1 below.

The high-level condition on the rank statistic $r k_{n}$ is the following.
Assumption R: For any subsequence $\left\{w_{n}\right\}$ and any sequence $\left\{\lambda_{w_{n}, h}: n \geq 1\right\}$ with $\lambda_{w_{n}, h} \in \Lambda_{0}$ $\forall n \geq 1$ either (a) $r k_{w_{n}} \rightarrow r_{h}=\infty$ or (b) $r k_{w_{n}} \rightarrow_{d} r_{h}\left(\bar{D}_{h}\right)$ for some nonrandom function $r_{h}$ : $R^{k \times p} \rightarrow R$, where $\bar{D}_{h}$ is defined in Lemma 10.2, and the convergence is joint with that in Lemma 10.2

In Assumption R, by $r k_{w_{n}} \rightarrow_{p} \infty$, we mean that for every $K<\infty$ we have $P_{\theta_{0}, \lambda_{w_{n}}}\left(r k_{w_{n}}>\right.$ $K) \rightarrow 1$, where $P_{\theta_{0}, \lambda_{w_{n}}}(\cdot)$ denotes probability under $\lambda_{w_{n}}$ when the true parameter vector equals $\theta_{0}$.

The following theorem applies when the LM statistic is defined as in 4.2) with projection onto $\widehat{\Omega}_{n}^{-1 / 2} \widehat{D}_{n}$. In consequence, the quantities in (10.2) in the present case are $\widehat{W}_{n}=\widehat{\Omega}_{n}^{-1 / 2}, W_{F}=\Omega_{F}^{-1 / 2}$, and $\widehat{U}_{n}=U_{F}=I_{p}$. (These definitions affect the definition of $\lambda_{n, h}$, which appears in the theorem).

Theorem 12.1 For any statistic $r k_{n}$ that satisfies Assumption R , the asymptotic null rejection probabilities of the nominal size $\alpha \in(0,1)$ CLR test defined in 4.2-5.2) based on rk $k_{n}$ equal $\alpha$ under all subsequences $\left\{w_{n}\right\}$ and all sequences $\left\{\lambda_{w_{n}, h}: n \geq 1\right\}$ with $\lambda_{w_{n}, h} \in \Lambda_{0} \forall n \geq 1$.

Comments: (i) Theorem 12.1 and Proposition 10.1 imply that a nominal size $\alpha$ CLR test based on any rank statistic that satisfies Assumption R has asymptotic size $\alpha$ and is asymptotically similar. Analogous CS results (to the test results stated in Theorem 12.1) hold for a parameter space $\Lambda_{\Theta, 0}$ that is a reparametrization of $\mathcal{F}_{\Theta, 0}$ and is defined as $\Lambda_{0}$ is defined, but with the adjustments outlined in Comments (i) and (ii) to Proposition 10.1.
(ii) Theorems 10.4 and 12.1 and Proposition 10.1 establish the test results of Theorem 6.1 . This holds because Theorem 10.4 (a), (b), (e), and (f) with $\widehat{W}_{n}=\widehat{\Omega}_{n}^{-1 / 2}$ and $\widehat{U}_{n}=I_{p}$ imply that Assumption R holds for the CLR test with moment-variance weighting, that is considered in Section 6, which uses the Robin and Smith (2000) $r k_{n}$ statistic defined in 6.2). (In the present context, Theorem 10.4 requires that Assumption WU holds for the parameter space $\Lambda_{0}$. It holds with $\widehat{W}_{n}=\widehat{W}_{2 n}, W_{1}(w)=w$ for $w \in R^{k \times k}, \mathcal{W}_{2}=R^{k \times k}, \widehat{U}_{n}=\widehat{U}_{2 n}, U_{1}(u)=u$ for $u \in R^{p \times p}$, and $\mathcal{U}_{2}=R^{p \times p}$, because $\widehat{W}_{n}=\widehat{\Omega}_{n}^{-1 / 2} \rightarrow_{p} h_{5, g}^{-1 / 2}$ under all sequences $\left\{\lambda_{n, h}: n \geq 1\right\}$ with $\lambda_{n, h} \in \Lambda_{0}$ and $\widehat{U}_{n}=I_{p}$ for all $n \geq 1$.) In particular, Assumption R holds with $r_{h}=\infty$ if $q=p$ and with $r_{h}\left(\bar{D}_{h}\right)$ equal to the smallest eigenvalue of $\bar{\Delta}_{h, p-q}^{\prime} h_{3, k-q} h_{3, k-q}^{\prime} \bar{\Delta}_{h, p-q}$ if $q<p$ (where $\bar{\Delta}_{h, p-q}$ and $h_{3, k-q}$ are defined in 10.17 based on $W_{F}=\Omega_{F}^{-1 / 2}$ and $U_{F}=I_{p}$. The CS results of Theorem 6.1 hold by Theorem 10.4. Comment (i) to Theorem 12.1, and Comment (i) to Proposition 10.1.
(iii) Theorem 5.1 shows that Assumption R does not hold in general for rank statistics based on $\widetilde{V}_{D n}$ and $\hat{D}_{n}^{\dagger}$, defined in (5.3)-5.4, when $p \geq 2$. The reason is that for some sequences of distributions the asymptotic distribution of $\widehat{D}_{n}^{\dagger}$ and, hence, the rank statistic $r k_{n}$ depends on $\bar{D}_{h}$ and $\bar{M}_{h}^{\dagger} \neq 0^{k \times p}$, not just on $\bar{D}_{h}$ alone.

For notational simplicity, the following proof is for the sequence $\{n\}$, rather than a subsequence $\left\{w_{n}: n \geq 1\right\}$. The same proof holds for any subsequence $\left\{w_{n}: n \geq 1\right\}$.

Proof of Theorem 12.1. Let

$$
\begin{equation*}
J_{n}:=n \widehat{g}_{n}^{\prime} \widehat{\Omega}_{n}^{-1 / 2} M_{\widehat{\Omega}_{n}^{-1 / 2} \widehat{D}_{n}} \widehat{\Omega}_{n}^{-1 / 2} \widehat{g}_{n} . \tag{12.1}
\end{equation*}
$$

It follows from (4.2) that

$$
\begin{equation*}
A R_{n}=L M_{n}+J_{n} . \tag{12.2}
\end{equation*}
$$

We now distinguish two cases. First, suppose Assumption R(a) holds: $r k_{n} \rightarrow_{p} \infty$. By (12.2) and some algebra, we have $\left(A R_{n}-r k_{n}\right)^{2}+4 L M_{n} \cdot r k_{n}=\left(L M_{n}-J_{n}+r k_{n}\right)^{2}+4 L M_{n} \cdot J_{n}$. Therefore,

$$
\begin{equation*}
C L R_{n}=\frac{1}{2}\left(L M_{n}+J_{n}-r k_{n}+\sqrt{\left(L M_{n}-J_{n}+r k_{n}\right)^{2}+4 L M_{n} \cdot J_{n}}\right) . \tag{12.3}
\end{equation*}
$$

Using a mean-value expansion of the square-root expression in 12.3) about $\left(L M_{n}-J_{n}+r k_{n}\right)^{2}$, we have

$$
\begin{equation*}
\sqrt{\left(L M_{n}-J_{n}+r k_{n}\right)^{2}+4 L M_{n} \cdot J_{n}}=L M_{n}-J_{n}+r k_{n}+\left(2 \sqrt{\zeta_{n}}\right)^{-1} 4 L M_{n} \cdot J_{n} \tag{12.4}
\end{equation*}
$$

for an intermediate value $\zeta_{n}$ between $\left(L M_{n}-J_{n}+r k_{n}\right)^{2}$ and $\left(L M_{n}-J_{n}+r k_{n}\right)^{2}+4 L M_{n} \cdot J_{n}$. It follows that $C L R_{n}=L M_{n}+o_{p}(1) \rightarrow_{d} \chi_{p}^{2}$ using $\sqrt{11.2}$ ) and $\left(\sqrt{\zeta_{n}}\right)^{-1}=o_{p}(1)$ (which holds because $r k_{n} \rightarrow_{p} \infty, L M_{n}=O_{p}(1)$, and $J_{n}=O_{p}(1)$ by 12.6 below). Analogously, it can be shown that the critical value $c\left(1-\alpha, r k_{n}\right)$, defined above 5.2$)$, of the CLR test converges in probability to $\chi_{p, 1-\alpha}^{2}$. The result of Theorem 12.1 then follows by the definition of convergence in distribution.

Second, suppose Assumption R(b) holds. Then, using Lemma 10.2, we have $\left(n^{1 / 2} \widehat{g}_{n}, n^{1 / 2}\left(\widehat{D}_{n}-\right.\right.$ $\left.\left.E_{F_{n}} G_{i}\right), r k_{n}\right) \rightarrow_{d}\left(\bar{g}_{h}, \bar{D}_{h}, r_{h}\left(\bar{D}_{h}\right)\right)$. By the proof of Lemma 10.3 applied with $W_{F}=\Omega_{F}^{-1 / 2}$ and $U_{F}=I_{p}$ (which correspond to $\widehat{W}_{n}=\widehat{\Omega}_{n}^{-1 / 2}$ and $\widehat{U}_{n}=I_{p}$ ), using the former result in place of $\left(n^{1 / 2} \widehat{g}_{n}, n^{1 / 2}\left(\widehat{D}_{n}-E_{F_{n}} G_{i}\right)\right) \rightarrow_{d}\left(\bar{g}_{h}, \bar{D}_{h}\right)$ gives

$$
\begin{equation*}
\left(n^{1 / 2} \widehat{g}_{n}, n^{1 / 2}\left(\widehat{D}_{n}-E_{F_{n}} G_{i}\right), n^{1 / 2} \Omega_{n}^{-1 / 2} \widehat{D}_{n} T_{n}, r k_{n}\right) \rightarrow_{d}\left(\bar{g}_{h}, \bar{D}_{h}, \bar{\Delta}_{h}, r_{h}\left(\bar{D}_{h}\right)\right), \tag{12.5}
\end{equation*}
$$

where $\Omega_{n}:=\Omega_{F_{n}},\left(\bar{D}_{h}, \bar{\Delta}_{h}\right)$ and $\bar{g}_{h}$ are independent, and $\bar{\Delta}_{h}$ has full column rank $p$ with probability
one by Lemma 10.3 (d) (because we are considering sequences $\left\{\lambda_{w_{n}, h}: n \geq 1\right\}$ with $\lambda_{w_{n}, h} \in \Lambda_{0}$ $\forall n \geq 1, W_{F}=\Omega_{F}^{-1 / 2}$, and $U_{F}=I_{p}$ ). In addition, $\widehat{\Omega}_{n} \rightarrow_{p} h_{5, g}, h_{5, g}$ is pd, and $M_{\widehat{\Omega}_{n}^{-1 / 2} \widehat{D}_{n}}=$ $M_{\left(\widehat{\Omega}_{n}^{-1 / 2} \Omega_{n}^{1 / 2}\right) n^{1 / 2} \Omega_{n}^{-1 / 2} \widehat{D}_{n} T_{n}}$ because $T_{n}$ (defined in 10.18 ) and $\Omega_{n}^{-1 / 2}$ are nonsingular. These results and the CMT imply that

$$
\begin{equation*}
J_{n} \rightarrow_{d} \bar{J}_{h}:=\bar{g}_{h}^{\prime} h_{5, g}^{-1 / 2} M_{\bar{\Delta}_{h}} h_{5, g}^{-1 / 2} \bar{g}_{h} . \tag{12.6}
\end{equation*}
$$

The convergence results in 11.2 and 12.6 and $r k_{n} \rightarrow_{d} r_{h}\left(\bar{D}_{h}\right)$ hold jointly by 12.5 and the definitions of $L M_{n}$ and $J_{n}$ in (4.2) and (12.1).

Note that $\overline{L M}_{h}=\bar{g}_{h}^{\prime} h_{5, g}^{-1 / 2} P_{\bar{\Delta}_{h}} h_{5, g}^{-1 / 2} \bar{g}_{h}$ by 11.1 and 11.2 . Conditional on $\bar{\Delta}_{h}, P_{\bar{\Delta}_{h}} h_{5, g}^{-1 / 2} \bar{g}_{h}$ and $M_{\bar{\Delta}_{h}} h_{5, g}^{-1 / 2} \bar{g}_{h}$ have a joint normal distribution with zero covariance (because $\operatorname{Var}\left(h_{5, g}^{-1 / 2} \bar{g}_{h}\right)=I_{k}$ and $P_{\bar{\Delta}_{h}} M_{\bar{\Delta}_{h}}=0^{k \times k}$ ) and, hence, are independent. The same holds true conditional on $\bar{D}_{h}$, because $\bar{\Delta}_{h}$ is a nonrandom function of $\bar{D}_{h}$ and $\bar{D}_{h}$ is independent of $\bar{g}_{h}$. In consequence, conditional on $\bar{D}_{h}, \overline{L M}_{h}$ and $\bar{J}_{h}$ are independent and distributed as $\chi_{p}^{2}$ and $\chi_{k-p}^{2}$, respectively.

Using the convergence results in (12.5) and 12.6), the definition of $C L R_{n}$ in with $A R_{n}=$ $L M_{n}+J_{n}$ substituted in, and the CMT, we obtain

$$
\begin{equation*}
C L R_{n} \rightarrow_{d} \overline{C L R}_{h}:=\frac{1}{2}\left(\overline{L M}_{h}+\bar{J}_{h}-\bar{r}_{h}+\sqrt{\left(\overline{L M}_{h}+\bar{J}_{h}-\bar{r}_{h}\right)^{2}+4 \overline{L M} \bar{r}_{h}}\right) \tag{12.7}
\end{equation*}
$$

where $\bar{r}_{h}:=r_{h}\left(\bar{D}_{h}\right)$.
The function $c(1-\alpha, r)$ (defined in (5.2) is continuous in $r$ on $R_{+}$by the absolute continuity of the distributions of $\chi_{p}^{2}$ and $\chi_{k-p}^{2}$, which appear in $\operatorname{clr}(r)$ (also defined in 5.2 ), and the continuity of $\operatorname{clr}(r)$ in $r$ a.s. This, $r k_{n} \rightarrow_{d} \bar{r}_{h}$, and (12.7) yield

$$
\begin{equation*}
C L R_{n}-c\left(1-\alpha, r k_{n}\right) \rightarrow_{d} \overline{C L R}_{h}-c\left(1-\alpha, \bar{r}_{h}\right) \tag{12.8}
\end{equation*}
$$

Therefore, by the definition of convergence in distribution, we have

$$
\begin{equation*}
P_{\theta_{0}, \lambda_{n}}\left(C L R_{n}>c\left(1-\alpha, r k_{n}\right)\right) \rightarrow P\left(\overline{C L R}_{h}>c\left(1-\alpha, \bar{r}_{h}\right)\right) \tag{12.9}
\end{equation*}
$$

provided $P\left(\overline{C L R}_{h}=c\left(1-\alpha, \bar{r}_{h}\right)\right)=0$, which holds because $P\left(\overline{C L R}_{h}=c\left(1-\alpha, \bar{r}_{h}\right) \mid \bar{D}_{h}\right)=0$ a.s. The latter holds because conditional on $\bar{D}_{h}, \overline{C L R}_{h}$ is absolutely continuous (by 12.7) since $\overline{L M}_{h}$ and $\bar{J}_{h}$ are independent and distributed as $\chi_{p}^{2}$ and $\chi_{k-p}^{2}$ and $\bar{r}_{h}$ is a nonrandom function of $\bar{D}_{h}$ ) and $c\left(1-\alpha, \bar{r}_{h}\right)$ is a constant.

From above, conditional on $\bar{D}_{h}, \overline{L M}_{h}$ and $\bar{J}_{h}$ are independent and distributed as $\chi_{p}^{2}$ and $\chi_{k-p}^{2}$, respectively, and $\bar{r}_{h}$ is a constant. Thus, conditional on $\bar{D}_{h}, \overline{C L R}_{h}$ and $\operatorname{clr}\left(\bar{r}_{h}\right)$ have the same distribution. By definition, $c\left(1-\alpha, \bar{r}_{h}\right)$ is the $1-\alpha$ quantile of the absolutely continuous random
variable $\operatorname{clr}\left(\bar{r}_{h}\right)$ for any constant $\bar{r}_{h}$. Hence,

$$
\begin{equation*}
P\left(\overline{C L R}_{h}>c\left(1-\alpha, \bar{r}_{h}\right) \mid \bar{D}_{h}\right)=\alpha \text { a.s. } \tag{12.10}
\end{equation*}
$$

Because the left-hand side conditional probability equals $\alpha$ a.s. and $\alpha$ does not depend on $\bar{D}_{h}$, the unconditional probability $P\left(\overline{C L R}_{h}>c\left(1-\alpha, \bar{r}_{h}\right)\right)$ equals $\alpha$ as well. Combined with 12.9), this gives the desired result.

## 13 Asymptotic Size of the CLR Test with Jacobian-Variance Weighting when $\mathrm{p}=1$

In this section, we prove the test results of Theorem 5.3, which concerns Kleibergen's CLR test (and CS) with Jacobian-variance weighting when $p=1$. The CS results of Theorem 5.3 hold by an analogous argument, see Comments (i) and (ii) to Proposition 10.1 .

Proof of Theorem 5.3. We prove the test results of Theorem 5.3 using Proposition 10.1 and results (or variants of results) in Lemma 10.3 and Theorems 10.4, 11.1, and 12.1. The proof is made more complicated by the fact that we need to use two different definitions of $\widehat{W}_{n}$. To obtain the asymptotic distribution of the LM statistic (which is a component of the CLR statistic), we need to take $\widehat{W}_{n}=\widehat{\Omega}_{n}^{-1 / 2}$ and $\widehat{U}_{n}=1$, because the LM statistic (defined in 4.2) depends on $\widehat{\Omega}_{n}^{-1 / 2} \widehat{D}_{n}$. But, to obtain the asymptotic distribution of the rank statistic $r k_{n}:=n \widehat{D}_{n}{ }^{\prime} \widetilde{V}_{D n}^{-1} \widehat{D}_{n}$ (defined in (5.10), we need to take $\widehat{W}_{n}=\widetilde{V}_{D n}^{-1 / 2}$ and $\widehat{U}_{n}=1$, because $r k_{n}$ depends on $\widetilde{V}_{D n}^{-1 / 2} \widehat{D}_{n}$.

For notational simplicity, we establish results below for sequences $\{n\}$, rather than subsequences $\left\{w_{n}\right\}$ of $\{n\}$. Subsequence results hold by replacing $n$ by $w_{n}$ in the proofs.

We proceed as follows. First, we apply Lemma 10.3 exactly as in the proof of Theorem 11.1 with $\widehat{W}_{n}=\widehat{\Omega}_{n}^{-1 / 2}, \widehat{U}_{n}=1, W_{F}=\Omega_{F}^{-1 / 2}$, and $U_{F}=1$. This yields $n^{1 / 2}\left(\widehat{g}_{n}, \widehat{D}_{n}-E_{F_{n}} G_{i}, W_{F_{n}} \widehat{D}_{n} U_{F_{n}} T_{n}\right) \rightarrow_{d}$ $\left(\bar{g}_{h}, \bar{D}_{h}, \bar{\Delta}_{h}\right)$ for sequences $\left\{\lambda_{n, h}: n \geq 1\right\}$ that correspond to distributions $F$ in $\mathcal{F}_{W U} \cap \mathcal{F}_{0}$ based on these definitions of $W_{F}$ and $U_{F}$. As discussed in the paragraph containing 10.5, $\mathcal{F}_{0}=\mathcal{F}_{W U} \cap \mathcal{F}_{0}$ for $\delta_{W U}$ sufficiently small and $M_{W U}$ sufficiently large. We employ constants $\delta_{W U}$ and $M_{W U}$ for which this holds. The joint convergence result above yields the asymptotic distributions of the $A R_{n}, L M_{n}$, and $J_{n}$ statistics via the calculations in (11.1), (11.2), (12.1), (12.2), and (12.6).

Next, we take $\widehat{W}_{n}=\widetilde{V}_{D n}^{-1 / 2}, \widehat{U}_{n}=1, W_{F}=W_{2 F}=\left(\operatorname{Var}_{F}\left(G_{i}\right)-\Gamma_{F}^{G_{i}} \Omega_{F}^{-1} \Gamma_{F}^{G_{i} \prime}\right)^{-1 / 2}$, where $\Gamma_{F}^{G_{i}}$ and $\Omega_{F}$ are defined in (3.6), $W_{1}(\cdot)$ equals the identity function on $\mathcal{W}_{2}:=R^{k \times k}, U_{F}=U_{2 F}=1$, and $U_{1}(\cdot)$ equals the identity function on $\mathcal{U}_{2}:=R$. We consider distributions in $\mathcal{F}_{J V W, p=1}$ (which is a subset of $\mathcal{F}_{0}$ when $\delta_{3}=\delta_{2}$ by the paragraph following (5.9). We obtain the asymptotic distribution of $r k_{n}$
under the corresponding sequences $\left\{\lambda_{n, h}: n \geq 1\right\}$ (which differ from the sequences $\left\{\lambda_{n, h}: n \geq 1\right\}$ in the previous paragraph due to the difference between the two definitions of $W_{F}$ ). More specifically, we verify the convergence results in Assumption R for $r k_{n}:=n \widehat{D}_{n}^{\prime} \widetilde{V}_{D n}^{-1} \widehat{D}_{n}$ (defined in 5.10) for the $\left\{\lambda_{n, h}: n \geq 1\right\}$ sequences of this paragraph. The result of Theorem 10.4(a), (b), (e), and (f) verifies the convergence results in Assumption $R$ for sequences $\left\{\lambda_{n, h}: n \geq 1\right\}$ for which $F_{n} \in \mathcal{F}_{J V W, p=1}$ $\forall n \geq 1$ provided Assumption WU holds for such sequences with $\widehat{W}_{2 n}=\widehat{W}_{n}=\widetilde{V}_{D n}^{-1 / 2}, W_{1}(\cdot)$ equal to the identity function, $\widehat{U}_{2 n}=\widehat{U}_{n}=1, U_{1}(\cdot)$ equal to the identity function, and the parameter space $\Lambda_{*}$ being equal to $\Lambda_{J V W, p=1}:=\left\{\lambda: \lambda=\left(\lambda_{1, F}, \ldots, \lambda_{9, F}\right)\right.$ for some $\left.F \in \mathcal{F}_{W U} \cap \mathcal{F}_{J V W, p=1}\right\}$. Here $\mathcal{F}_{W U}$ is defined in 10.5 with $W_{F}=\left(\operatorname{Var}_{F}\left(G_{i}\right)-\Gamma_{F}^{G_{i}} \Omega_{F}^{-1} \Gamma_{F}^{G_{i} \prime}\right)^{-1 / 2}$ and $U_{F}=1$. Note that $\mathcal{F}_{J V W, p=1}=\mathcal{F}_{W U} \cap \mathcal{F}_{J V W, p=1}$ for $\delta_{W U}>0$ sufficiently small and $M_{W U}<\infty$ sufficiently large (and we employ constants $\delta_{W U}$ and $M_{W U}$ that satisfy these conditions). This holds because for all $F \in \mathcal{F}_{J V W, p=1}, \lambda_{\min }\left(W_{F}\right)=\lambda_{\min }\left(\left(\operatorname{Var}_{F}\left(G_{i}\right)-\Gamma_{F}^{G_{i}} \Omega_{F}^{-1} \Gamma_{F}^{G_{i} \prime}\right)^{-1 / 2}\right)=\lambda_{\max }^{-1 / 2}\left(\operatorname{Var}_{F}\left(G_{i}\right)-\Gamma_{F}^{G_{i}} \Omega_{F}^{-1} \Gamma_{F}^{G_{i}}{ }^{\prime}\right)$ $\geq \lambda_{\max }^{-1 / 2}\left(E_{F} G_{i} G_{i}^{\prime}\right) \geq M_{+}^{-1 / 2}$ for some $M_{+}<\infty$ (because $E_{F} G_{i} G_{i}^{\prime}-\left(\operatorname{Var}_{F}\left(G_{i}\right)-\Gamma_{F}^{G_{i}} \Omega_{F}^{-1} \Gamma_{F}^{G_{i}{ }^{\prime}}\right)=$ $E_{F} G_{i} E_{F} G_{i}^{\prime}+\Gamma_{F}^{G_{i}} \Omega_{F}^{-1} \Gamma_{F}^{G_{i}}$ is psd and $\left\|E_{F} G_{i} G_{i}^{\prime}\right\| \leq M_{+}$for some $M_{+}<\infty$ by the moment conditions in $\mathcal{F}),\left\|W_{F}\right\|=\left\|\left(\operatorname{Var}_{F}\left(G_{i}\right)-\Gamma_{F}^{G_{i}} \Omega_{F}^{-1} \Gamma_{F}^{G_{i}}\right)^{-1 / 2}\right\| \leq \lambda_{\min }^{-1 / 2}\left(\operatorname{Var}_{F}\left(G_{i}\right)-\Gamma_{F}^{G_{i}} \Omega_{F}^{-1} \Gamma_{F}^{G_{i}}\right) \leq \delta_{3}^{-1 / 2}$ (using the condition in $\mathcal{F}_{J V W, p=1}$ and the fact that $\operatorname{Var}_{F}\left(G_{i}\right)-\Gamma_{F}^{G_{i}} \Omega_{F}^{-1} \Gamma_{F}^{G_{i}}=\Psi_{F}^{G_{i}}-E_{F} G_{i} E_{F} G_{i}^{\prime}$ using the definition of $\Psi_{F}^{G_{i}}$ in 3.6 ), where $\delta_{3}>0$, and $\left\|U_{F}\right\|=\lambda_{\min }\left(U_{F}\right)=1$.

Assumption $\mathrm{WU}(\mathrm{b})$ holds automatically with $h_{8}=1$ because $\widehat{U}_{2 n}:=1$. The requirement of Assumption $\mathrm{WU}(\mathrm{c})$ that $W_{1}(\cdot)$ is continuous at $h_{7}$ and $U_{1}(\cdot)$ is continuous at $h_{8}$ also holds automatically because $W_{1}(\cdot)$ and $U_{1}(\cdot)$ are identity functions.

Assumption WU(a) for the parameter space $\Lambda_{J V W, p=1}$ requires that $\widehat{W}_{2 n} \rightarrow_{p} h_{7}\left(:=\lim W_{2 F_{n}}\right)$. For sequences $\left\{\lambda_{n, h}: n \geq 1\right\}$, we have

$$
\begin{align*}
\widetilde{V}_{D n} & :=n^{-1} \sum_{i=1}^{n}\left(G_{i}-\widehat{G}_{n}\right)\left(G_{i}-\widehat{G}_{n}\right)^{\prime}-\widehat{\Gamma}_{n} \widehat{\Omega}_{n}^{-1} \widehat{\Gamma}_{n}^{\prime} \\
& =E_{F_{n}}\left(G_{i}-E_{F_{n}} G_{i}\right)\left(G_{i}-E_{F_{n}} G_{i}\right)^{\prime}-\Gamma_{F_{n}}^{G_{i}} \Omega_{F_{n}}^{-1} \Gamma_{F_{n}}^{G_{i}}+o_{p}(1) \\
& =W_{2 F_{n}}^{-2}+o_{p}(1) \\
& \rightarrow p h_{7}^{-2}, \tag{13.1}
\end{align*}
$$

where the first equality holds by (5.3), the second equality holds by the WLLN's applied multiple times and Slutsky's Theorem using the conditions in $\mathcal{F}$, the third equality holds by the definition of $W_{2 F}$, and the convergence holds because $W_{2 F_{n}}=\lambda_{7, F_{n}} \rightarrow h_{7}$ by the definition of the sequence $\left\{\lambda_{n, h}: n \geq 1\right\}$ and $h_{7}$ is pd (since $h_{7}=\lim W_{2 F_{n}}$ and the eigenvalues of $W_{2 F}^{-2}$ are bounded above for $F \in \mathcal{F}$ ). Equation 13.1 and Slutsky's Theorem give $\widetilde{V}_{D n}^{-1 / 2} \rightarrow_{p} h_{7}$ because $h_{7}^{-2}$ is pd using
the condition in $\mathcal{F}_{J V W, p=1}$ that $\lambda_{\min }\left(\Psi_{F}^{G_{i}}-E_{F} G_{i} E_{F} G_{i}^{\prime}\right) \geq \delta$. In consequence, Assumption WU(a) holds.

This completes the verification of Assumption WU for the parameter space $\Lambda_{J V W, p=1}$ and, in consequence, the verification of the convergence results of Assumption R for $r k_{n}$ for sequences $\left\{\lambda_{n, h}: n \geq 1\right\}$ defined in the fourth paragraph of this proof.

Now we consider sequences $\left\{\lambda_{n, h}: n \geq 1\right\}$ that satisfy the conditions on $\left\{\lambda_{n, h}: n \geq 1\right\}$ given in both the third and fourth paragraphs of this proof. These sequences correspond to distributions $F$ in $\mathcal{F}_{J V W, p=1}$. These sequences satisfy the convergence conditions in (8.11) using the definitions in (8.9) and (8.10) with $\tau_{j F}, B_{F}, C_{F}$, and $W_{2 F}$ defined based on $W_{F}=\Omega_{F}^{-1 / 2}$ and with these quantities based on $W_{F}=\left(\operatorname{Var}_{F}\left(G_{i}\right)-\Gamma_{F}^{G_{i}} \Omega_{F}^{-1} \Gamma_{F}^{G_{i}}\right)^{-1 / 2}$. In consequence, for these sequences of distributions $\left\{\lambda_{n, h}: n \geq 1\right\}$, the results above establish the asymptotic distributions of the $A R_{n}, L M_{n}, J_{n}$, and $r k_{n}$ statistics and the convergence is joint because all of the convergence results are based on the underlying CLT result in Lemma 10.2. Given this joint convergence, by the same arguments as given in the proof of Theorem 12.1, we obtain that the CLR test with Jacobian-variance weighting has asymptotic null rejection probabilities equal to $\alpha$ under all such sequences $\left\{\lambda_{n, h}: n \geq 1\right\}$ (and all subsequences of such sequences).

Finally, we apply Proposition 8.1 with $\lambda$ and $h_{n}(\theta)$ given by the concatenation of the $\lambda$ vectors and $h_{n}(\lambda)$ functions used in the third and fourth paragraphs above and with $\Lambda$ given by the product space of the $\Lambda$ spaces used in these paragraphs. (Redundant elements of $\lambda$ and $h_{n}(\lambda)$ do not cause any problems.) The result of the previous paragraph verifies Assumption $\mathrm{B}^{*}$ for this choice $\lambda$, $h_{n}(\lambda)$, and $\Lambda$. In consequence, Proposition 8.1 implies that the Jacobian-variance weighted CLR test has correct asymptotic size and is asymptotically similar when $p=1$.

## 14 The Eigenvalue Condition in $\mathcal{F}_{0}$

In this section, we show that the restriction $\lambda_{p-j}\left(\Psi_{j F}(\xi)\right) \geq \delta_{1}>0$ in $\mathcal{F}_{0 j}$, defined in 3.9, is not redundant. If this restriction is weakened to $\lambda_{p-j}\left(\Psi_{j F}(\xi)\right)>0$, we show that, for some models, some sequences of distributions, and some (consistent) choices of variance and covariance estimators, the LM statistic in (4.2) has a $\chi_{k}^{2}$ asymptotic distribution. This leads to over-rejection of the null when the standard $\chi_{p}^{2}$ critical value is used and the parameters are over-identified (i.e., $k>p)$. On the other hand, we show that the LM statistic equals zero a.s. for some models and some distributions $F$ if the condition $\lambda_{p-j}\left(\Psi_{j F}(\xi)\right) \geq \delta_{1}>0$ is removed entirely. This implies that the LM test also under-rejects the null hypothesis and is nonsimilar in both finite samples and asymptotically for some $F$.

All of the CLR tests considered in Sections 5 and 6, except that of Smith (2007), are functions of the LM statistic in (4.2) (and other statistics). In consequence, the aberrant behavior of the LM statistic and test demonstrated in this section, when the restriction $\lambda_{p-j}\left(\Psi_{j F}(\xi)\right) \geq \delta_{1}>0$ in $\mathcal{F}_{0}$ is weakened or eliminated, carries over to the CLR statistics and tests in Sections 5 and 6 . Smith's (2007) CLR test is a function of the LM statistic in (4.2) but with $\widehat{\Omega}_{n}^{-1 / 2} \widehat{D}_{n}$ replaced by $\widehat{D}_{n}^{\dagger}$.

### 14.1 Eigenvalue Condition Counter-Examples

For simplicity, we consider the case $p=1$ in this section. As above, the null hypothesis is $H_{0}: \theta=\theta_{0}$.

Lemma 14.1 (a) Suppose $\mathcal{F}_{0}$ is defined with the condition $\lambda_{p-j}\left(\Psi_{j F}(\xi)\right)>0$ in place of $\lambda_{p-j}\left(\Psi_{j F}(\xi)\right) \geq \delta_{1}>0$ in $\mathcal{F}_{0 j}$ for all $j \in\{0, \ldots, p\}$, where $p=1$. Suppose $\widehat{\Omega}_{n}(\theta)$ is defined in (3.1) and $\widehat{\Gamma}_{1 n}(\theta)=n^{-1} \sum_{i=1}^{n} G_{i}(\theta) g_{i}(\theta)^{\prime}$ (which differs from its definition in 3.2). Then, there exist moment functions $g\left(W_{i}, \theta\right)$ and a sequence of null distributions $\left\{F_{n} \in \mathcal{F}_{0}: n \geq 1\right\}$ for which $\widehat{\Omega}_{n}=\widehat{\Omega}_{n}\left(\theta_{0}\right)$ and $\widehat{\Gamma}_{1 n}=\widehat{\Gamma}_{1 n}\left(\theta_{0}\right)$ are well-behaved (in the sense that $\widehat{\Omega}_{n}-E_{F_{n}} g_{i} g_{i}^{\prime} \rightarrow_{p} 0^{k \times k}$ and $\left.\widehat{\Gamma}_{1 n}-E_{F_{n}} G_{i} g_{i}^{\prime} \rightarrow_{p} 0^{k \times k}\right)$ and $L M_{n}\left(\theta_{0}\right)=A R_{n}\left(\theta_{0}\right)+o_{p}(1) \rightarrow_{d} \chi_{k}^{2}$.
(b) Suppose $\mathcal{F}_{0}$ is defined with the condition $\lambda_{p-j}\left(\Psi_{j F}(\xi)\right) \geq \delta_{1}>0$ deleted in $\mathcal{F}_{0 j}$ for all $j \in\{0, \ldots, p\}$, where $p=1$. Suppose $\widehat{\Omega}_{n}(\theta)$ and $\widehat{\Gamma}_{1 n}(\theta)$ are defined in 3.1 and 3.2, respectively. Then, there exists moment functions and a null distribution $F \in \mathcal{F}_{0}$ for which $L M_{n}\left(\theta_{0}\right)=0$ a.s. for all $n \geq 1$.

Comments: (i) The model we use to prove Lemma 14.1(a) is the linear IV regression model with one endogenous rhs variable and (for simplicity) no exogenous variables. Specifically, the model is

$$
\begin{equation*}
y_{1 i}=y_{2 i} \theta+u_{i} \text { and } y_{2 i}=Z_{i}^{\prime} \pi+v_{2 i}, \tag{14.1}
\end{equation*}
$$

where $y_{1 i}, \theta, y_{2 i}, v_{2 i} \in R, Z_{i}, \pi \in R^{k}, v_{2 i}=\rho u_{i}+\delta \xi_{i}$ for some random variable $\xi_{i}, \delta=\left(1-\rho^{2}\right)^{1 / 2}$, and the observations are i.i.d. across $i$ for any given $n$. The parameter space $\mathcal{F}^{*}$ for the distribution $F$ of the random vector $W_{i}=\left(y_{1 i}, y_{2 i}, Z_{i}^{\prime}\right)^{\prime}$ is

$$
\begin{align*}
\mathcal{F}^{*}:= & \left\{F: 14.1 \text { holds with } \theta=\theta_{0}, \pi=\pi_{F} \in R^{k}, \rho=\rho_{F} \in(-1,1),\right. \\
& Z_{i}, u_{i}, \text { and } \xi_{i} \text { are mutually independent, } E_{F} u_{i}=E_{F} \xi_{i}=0, \\
& \left.E_{F} u_{i}^{2}=E_{F} \xi_{i}^{2}=1, E_{F}\left\|\left(u_{i}, \xi_{i}, Z_{i}^{\prime} Z_{i}\right)\right\|^{2+\gamma} \leq M, \text { and } \lambda_{\min }\left(E_{F} Z_{i} Z_{i}^{\prime}\right) \geq \delta\right\} \tag{14.2}
\end{align*}
$$

for some $\gamma, \delta>0$ and $M<\infty$. As defined, $\rho$ is the correlation between $u_{i}$ and $v_{2 i}$.

The moment functions are $g\left(W_{i}, \theta\right)=Z_{i}\left(y_{1 i}-y_{2 i} \theta\right)$. When the null value $\theta_{0}$ is the true value, this gives $g_{i}=g_{i}\left(\theta_{0}\right)=Z_{i} u_{i}$ and $G_{i}=G_{i}\left(\theta_{0}\right)=-Z_{i} y_{2 i}$. The set $\mathcal{F}^{*}$ is a subset of $\mathcal{F}_{0}$ when the latter is defined with the condition $\lambda_{p-j}\left(\Psi_{j F}(\xi)\right)>0$ in place of $\lambda_{p-j}\left(\Psi_{j F}(\xi)\right) \geq \delta_{1}>0$. This holds because (i) for all $F \in \mathcal{F}^{*}, \lambda_{\min }\left(\Psi_{F}^{v e c\left(G_{i}\right)}\right)>0$ (by the argument in the paragraph that contains (3.12) because $\lambda_{\min }\left(E_{F} Z_{i} Z_{i}^{\prime}\right)>0$ and $\lambda_{\min }\left(E_{F} \varepsilon_{i} \varepsilon_{i}^{\prime}\right)>0$, where $\varepsilon_{i}=\left(u_{i},-\rho u_{i}-\delta \xi_{i}\right)^{\prime}$ for $\rho \in$ $(-1,1))$, (ii) $\lambda_{\min }\left(E_{F} g_{i} g_{i}^{\prime}\right)=E_{F} u_{i}^{2} \lambda_{\min }\left(E_{F} Z_{i} Z_{i}^{\prime}\right) \geq \delta>0$, and (iii) $\lambda_{p-j}\left(\Psi_{F}^{C_{F, k-j}^{\prime} \Omega_{F}^{-1 / 2} G_{i} B_{F, p-j} \xi}\right) \geq$ $\lambda_{\min }\left(\Psi_{F}^{v e c\left(G_{i}\right)}\right) M^{-2 /(2+\gamma)}$ for all $\xi \in R^{p-j}$ with $\|\xi\|=1$ and all $j \in\{0, \ldots, p\}$ (by the results and arguments in the paragraphs that contain (18.1)-(18.3), which verify that condition (iv), stated in 3.10, is a sufficient condition for the $\lambda_{p-j}(\cdot)$ condition in $\left.\mathcal{F}_{0 j}\right)$. The quantity $\lambda_{\min }\left(\Psi_{F}^{v e c\left(G_{i}\right)}\right)$ is arbitrarily close to zero for $\rho$ arbitrarily close to one.

We consider a sequence of distributions $\left\{F_{n} \in \mathcal{F}^{*}: n \geq 1\right\}$ for which $\pi_{F_{n}}=0^{k}$ for all $n \geq 1, \rho_{n}$ $\left(=\rho_{F_{n}}\right) \rightarrow 1$, and $E_{F_{n}} Z_{i} Z_{i}^{\prime}$ does not depend on $n$. For these distributions,

$$
\begin{equation*}
G_{i}=-\rho_{n} g_{i}+\delta_{n} G_{i}^{*}, \text { where } G_{i}^{*}:=-Z_{i} \xi_{i} \text { and } \delta_{n}:=\left(1-\rho_{n}^{2}\right)^{1 / 2} . \tag{14.3}
\end{equation*}
$$

In this case, the IV's are irrelevant and the degree of endogeneity is close to perfect for $n$ large.
(ii) The model we consider in Lemma 14.1(b) is the same as that in part (a) except that $\mathcal{F}^{*}$ allows for $\rho=\rho_{F} \in(-1,1]$ and we consider a single distribution $F$ with $\pi=0^{k}$ and $\rho=1$, rather than a drifting sequence of distributions. For this distribution, $\lambda_{\min }\left(\Psi_{F}^{v e c}\left(G_{i}\right)\right)=0$.
(iii) The intuition for the results in Lemma 14.1 (a) and (b) is as follows. As 14.3) shows, $G_{i}$ is close to being proportional to $g_{i}$ when $\pi_{F_{n}}=0^{k}$ and $\rho_{n}$ is close to one. And, when $\pi_{F_{n}}=0^{k}$ and $\rho_{n}=1$, they are exactly proportional. By averaging over $i=1, \ldots, n$ and by taking expectations, the same properties are seen to hold for $\widehat{G}_{n}$ and $\widehat{g}_{n}$ and their population counterparts. In consequence, $\widehat{D}_{n}\left(:=\widehat{G}_{n}-\widehat{\Gamma}_{n} \widehat{\Omega}_{n}^{-1} \widehat{g}_{n}\right.$ when $\left.p=1\right)$ is close to $0^{k}$ (because it is a sample version of the $L^{2}(F)$ projection of $G_{i}$ on $g_{i}$ ) and the same is true of the population counterpart of $\widehat{D}_{n}$ (because it is the $L^{2}(F)$ projection of $G_{i}$ on $\left.g_{i}\right)$. The latter implies that the direction of the $k$-vector $\widehat{D}_{n}$ is primarily random. In consequence, this direction turns out to be sensitive to the specification of the sample matrices $\widehat{\Gamma}_{n}$ and $\widehat{\Omega}_{n}$ even within the class of consistent estimators of their population counterparts.

One consistent choice of $\widehat{\Gamma}_{n}$ and $\widehat{\Omega}_{n}$ (used in Lemma 14.1 (a)) yields $\widehat{D}_{n}$ to be very close to being proportional to $\widehat{g}_{n}$. In this case, the projection of $\widehat{\Omega}_{n}^{-1 / 2} \widehat{g}_{n}$ onto $\widehat{\Omega}_{n}^{-1 / 2} \widehat{D}_{n}$ is asymptotically equivalent to $\widehat{\Omega}_{n}^{-1 / 2} \widehat{g}_{n}$ itself. The LM statistic is a quadratic form in this projection $k$-vector (i.e., $P_{\widehat{\Omega}_{n}^{-1 / 2} \widehat{D}_{n}} \widehat{\Omega}_{n}^{-1 / 2} \widehat{g}_{n}$ ) multiplied by $n$. Hence, it behaves asymptotically like a quadratic form in $\widehat{\Omega}_{n}^{-1 / 2} \widehat{g}_{n}$ multiplied by $n$, which is just the AR statistic. This explains the result in Lemma 14.1 (a).

On the other hand, when $\rho_{n}=1$ (which implies that $\widehat{G}_{n}=-\widehat{g}_{n}$ by 14.3), another consistent choice of $\widehat{\Gamma}_{n}$ and $\widehat{\Omega}_{n}$ (used in Lemma 14.1(b)) yields $\widehat{D}_{n}=0^{k}$ a.s. In this case, the projection of $\widehat{\Omega}_{n}^{-1 / 2} \widehat{g}_{n}$ onto $\widehat{\Omega}_{n}^{-1 / 2} \widehat{D}_{n}$ equals $0^{k}$ a.s. Hence, the LM statistic (which is a quadratic form in this projection times $n$ ) equals zero a.s. This explains the result in Lemma 14.1(b).
(iv) The result of Lemma 14.1 (a) also holds for the model described in Comment (ii). Hence, drifting sequences of distributions are not required to show the result of Lemma 14.1 (a) if one removes the condition $\lambda_{p-j}\left(\Psi_{j F}(\xi)\right) \geq \delta_{1}>0$ entirely from $\mathcal{F}_{0 j}$. Furthermore, the result of Lemma 14.1 (a) can be extended to cover weak IV cases (in which $\pi=\pi_{n} \neq 0^{k}$, but $\pi_{n} \rightarrow 0^{k}$ sufficiently quickly as $n \rightarrow \infty$ ), rather than the irrelevant IV case (in which $\pi=0^{k}$ ).
(v) In the extreme case of the model, where $\rho=1$ and $\pi=0$, the endogenous variables $y_{1 i}$ and $y_{2 i}$ are identical, which is similar to perfect multicollinearity in linear regression. However, the result of Lemma 14.1(a) does not require either $\rho$ to be exactly equal to one or $\pi$ to be exactly equal to zero.
(vi) Finite sample simulations corroborate the asymptotic result given in Lemma 14.1(a). For the model and LM test described in Comment (i) with $k=5, \pi=0^{k}, \rho=1, Z_{i} \sim N\left(0^{5}, I_{5}\right)$, $\left(u_{i}, \xi_{i}\right) \sim N\left(0^{2}, I_{2}\right)$, and $Z_{i}$ independent of $\left(u_{i}, \xi_{i}\right)$, the null rejection rate of the nominal $5 \% \mathrm{LM}$ test is $59.4 \%$ when $n=200$ and $57.6 \%$ when $n=1000$. However, when $\rho$ deviates from 1 even by a small amount, the magnitude of over-rejection drops very quickly. The null rejection rate of this nominal $5 \%$ LM test is $10.1 \%$ when $\rho=0.99$ and $n=200$ and $12.9 \%$ when $\rho=0.998$ and $n=1000$. (These simulation results are based on 50,000 simulation repetitions.)
(vii) The conditions of Lemma 14.1(a) and (b) are consistent with those of Theorem 1 of Kleibergen (2005). This implies that the $\chi_{p}^{2}$ asymptotic distribution of the LM statistic obtained in the latter only holds under additional conditions, such as those in $\mathcal{F}_{0}$.

### 14.2 Proof of Lemma 14.1

Proof of Lemma 14.1. To prove part (a), we use the model defined in (14.1)-(14.3). We have

$$
\begin{align*}
& \widehat{G}_{n}=-\rho_{n} \widehat{g}_{n}+\delta_{n} \widehat{G}_{n}^{*}, \text { where } \widehat{G}_{n}^{*}:=n^{-1} \sum_{i=1}^{n} G_{i}^{*}, \text { and } \\
& \widehat{\Gamma}_{1 n}=n^{-1} \sum_{i=1}^{n} G_{i} g_{i}^{\prime}=n^{-1} \sum_{i=1}^{n}\left(-\rho_{n} g_{i}+\delta_{n} G_{i}^{*}\right) g_{i}^{\prime}=-\rho_{n} \widehat{\Omega}_{n}-\rho_{n} \widehat{g}_{n} \widehat{g}_{n}^{\prime}+\delta_{n} \widehat{\Gamma}_{1 n}^{*}, \text { where } \\
& \widehat{\Gamma}_{1 n}^{*}:=n^{-1} \sum_{i=1}^{n} G_{i}^{*} g_{i}^{\prime} . \tag{14.4}
\end{align*}
$$

We choose $\left\{\rho_{n}: n \geq 1\right\}$ to converge to one sufficiently fast that $n \delta_{n} \rightarrow 0$, where $\delta_{n}=\left(1-\rho_{n}^{2}\right)^{1 / 2}$
by 14.3. For example, we can take $\rho_{n}=\left(1-n^{-3}\right)^{1 / 2}$. Using the results above, we obtain

$$
\begin{align*}
\widehat{D}_{n} & =\widehat{G}_{n}-\widehat{\Gamma}_{1 n} \widehat{\Omega}_{n}^{-1} \widehat{g}_{n} \\
& =-\rho_{n} \widehat{g}_{n}+\delta_{n} \widehat{G}_{n}^{*}-\left[-\rho_{n} \widehat{\Omega}_{n}-\rho_{n} \widehat{g}_{n} \widehat{g}_{n}^{\prime}+\delta_{n} \widehat{\Gamma}_{1 n}^{*}\right] \widehat{\Omega}_{n}^{-1} \widehat{g}_{n} \\
& =\rho_{n}\left(\widehat{g}_{n}^{\prime} \widehat{\Omega}_{n}^{-1} \widehat{g}_{n}\right) \widehat{g}_{n}+\delta_{n}\left(\widehat{G}_{n}^{*}-\widehat{\Gamma}_{1 n}^{*} \widehat{\Omega}_{n}^{-1} \widehat{g}_{n}\right) . \tag{14.5}
\end{align*}
$$

This gives

$$
\begin{align*}
\widetilde{g}_{n} & :=\widehat{g}_{n}+n \delta_{n} \zeta_{n}=\widehat{D}_{n} /\left(\rho_{n} \widehat{g}_{n}^{\prime} \widehat{\Omega}_{n}^{-1} \widehat{g}_{n}\right), \text { where } \\
\zeta_{n} & :=\left(\widehat{G}_{n}^{*}-\widehat{\Gamma}_{1 n}^{*} \widehat{\Omega}_{n}^{-1} \widehat{g}_{n}\right) /\left(\rho_{n} n \widehat{g}_{n}^{\prime} \widehat{\Omega}_{n}^{-1} \widehat{g}_{n}\right)=O_{p}\left(n^{-1 / 2}\right) \text { and } \widetilde{g}_{n}=\widehat{g}_{n}+o_{p}\left(n^{-1 / 2}\right), \tag{14.6}
\end{align*}
$$

where $\zeta_{n}=O_{p}\left(n^{-1 / 2}\right)$ because $\rho_{n} \rightarrow 1, \widehat{G}_{n}^{*}=O_{p}\left(n^{-1 / 2}\right)$ by the CLT since $E_{F_{n}} G_{i}^{*}=-E_{F_{n}} Z_{i}$. $E_{F_{n}} \xi_{i}=0^{k}, \widehat{\Gamma}_{1 n}^{*} \widehat{\Omega}_{n}^{-1}=O_{p}(1)$ by the WLLN applied twice and $\lambda_{\min }\left(E_{F_{n}} g_{i} g_{i}^{\prime}\right)=\lambda_{\min }\left(E_{F_{n}} Z_{i} Z_{i}^{\prime}\right) \geq$ $\delta>0, \widehat{g}_{n}=O_{p}\left(n^{-1 / 2}\right)$ by the CLT, and $\left(n \widehat{g}_{n}^{\prime} \widehat{\Omega}_{n}^{-1} \widehat{g}_{n}\right)^{-1}=O_{p}(1)$, which holds by the CMT because $A R_{n}=n \widehat{g}_{n}^{\prime} \widehat{\Omega}_{n}^{-1} \widehat{g}_{n} \rightarrow_{d} \chi_{k}^{2}$ (by the CLT, WLLN, and CMT) and $\chi_{k}^{2}>0$ a.s., and lastly the result for $\widetilde{g}_{n}$ in the second line of 14.6 holds by $\zeta_{n}=O_{p}\left(n^{-1 / 2}\right)$ and $n \delta_{n}=o(1)$.

Projections are invariant to nonzero scalar multiplications of the matrix that defines the projection. That is, $P_{A}=P_{c A}$ for any matrix $A$ and any scalar $c \neq 0$. We have $\rho_{n} \widehat{g}_{n}^{\prime} \widehat{\Omega}_{n}^{-1} \widehat{g}_{n} \neq 0 \mathrm{wp} \rightarrow 1$ because $\left(n \widehat{g}_{n}^{\prime} \widehat{\Omega}_{n}^{-1} \widehat{g}_{n}\right)^{-1}=O_{p}(1)$ and $\rho_{n} \rightarrow 1$. So, the LM statistic is unchanged wp $\rightarrow 1$ when $\widehat{D}_{n}$ is replaced by $\widehat{D}_{n} /\left(\rho_{n} \widehat{g}_{n}^{\prime} \widehat{\Omega}_{n}^{-1} \widehat{g}_{n}\right)=\widetilde{g}_{n}=\widehat{g}_{n}+o_{p}\left(n^{-1 / 2}\right)$ using 14.6. Thus, we have

$$
\begin{align*}
L M_{n} & :=n \widehat{g}_{n}^{\prime} \widehat{\Omega}_{n}^{-1 / 2} P_{\widehat{\Omega}_{n}^{-1 / 2}} \widehat{D}_{n} \widehat{\Omega}_{n}^{-1 / 2} \widehat{g}_{n} \\
& =n \widehat{g}_{n}^{\prime} \widehat{\Omega}_{n}^{-1 / 2} P_{\widehat{\Omega}_{n}^{-1 / 2} \widehat{g}_{n}} \widehat{\Omega}_{n}^{-1 / 2} \widehat{g}_{n}+o_{p}(1) \\
& =n \widehat{g}_{n}^{\prime} \widehat{\Omega}_{n}^{-1} \widetilde{g}_{n}\left(\widetilde{g}_{n}^{\prime} \widehat{\Omega}_{n}^{-1} \widetilde{g}_{n}\right)^{-1} \widetilde{g}_{n}^{\prime} \widehat{\Omega}_{n}^{-1} \widehat{g}_{n}+o_{p}(1) \\
& =n \widehat{g}_{n}^{\prime} \widehat{\Omega}_{n}^{-1} \widehat{g}_{n}+o_{p}(1)=A R_{n}+o_{p}(1) \rightarrow_{d} \chi_{k}^{2}, \tag{14.7}
\end{align*}
$$

which completes the proof of part (a).
Next, we prove part (b). In this case, we use the model in (14.1)-14.3) with $\rho_{n}=1$ and $\delta_{n}=0$ for all $n \geq 1$. In consequence, $G_{i}=-g_{i}$ and $\widehat{G}_{n}=-\widehat{g}_{n}$. Given the definitions of $\widehat{\Omega}_{n}$ and $\widehat{\Gamma}_{1 n}$ in (3.1) and (3.2), this yields

$$
\begin{align*}
\widehat{\Gamma}_{1 n} & =n^{-1} \sum_{i=1}^{n} G_{i} g_{i}^{\prime}-\widehat{G}_{n} \widehat{g}_{n}^{\prime}=-n^{-1} \sum_{i=1}^{n} g_{i} g_{i}^{\prime}+\widehat{g}_{n} \widehat{g}_{n}^{\prime}=-\widehat{\Omega}_{n}, \\
\widehat{D}_{n} & =\widehat{G}_{n}-\widehat{\Gamma}_{1 n} \widehat{\Omega}_{n}^{-1} \widehat{g}_{n}=0^{k}, \text { and } \\
L M_{n} & :=n \widehat{g}_{n}^{\prime} \widehat{\Omega}_{n}^{-1 / 2} P_{\widehat{\Omega}_{n}^{-1 / 2} \widehat{D}_{n}} \widehat{\Omega}_{n}^{-1 / 2} \widehat{g}_{n}=n \widehat{g}_{n}^{\prime} \widehat{\Omega}_{n}^{-1 / 2} P_{0^{k}} \widehat{\Omega}_{n}^{-1 / 2} \widehat{g}_{n}=0 \tag{14.8}
\end{align*}
$$

for all $n \geq 1$, where the projection matrix, $P_{0^{k}}$, onto $0^{k}$ equals $0^{k \times k}$.

## 15 Proof of Lemma 10.2

Lemma 10.2 of AG1. Under all sequences $\left\{\lambda_{n, h}: n \geq 1\right\}$,

$$
n^{1 / 2}\binom{\widehat{g}_{n}}{\operatorname{vec}\left(\widehat{D}_{n}-E_{F_{n}} G_{i}\right)} \rightarrow_{d}\binom{\bar{g}_{h}}{\operatorname{vec}\left(\bar{D}_{h}\right)} \sim N\left(0^{(p+1) k},\left(\begin{array}{cc}
h_{5, g} & 0^{k \times p k} \\
0^{p k \times k} & \Phi_{h}^{v e c\left(G_{i}\right)}
\end{array}\right)\right) .
$$

Under all subsequences $\left\{w_{n}\right\}$ and all sequences $\left\{\lambda_{w_{n}, h}: n \geq 1\right\}$, the same result holds with $n$ replaced with $w_{n}$.

Proof of Lemma 10.2, We have

$$
\begin{align*}
n^{1 / 2} \operatorname{vec}\left(\widehat{D}_{n}-D_{n}\right) & =n^{-1 / 2} \sum_{i=1}^{n} \operatorname{vec}\left(G_{i}-D_{n}\right)-\left(\begin{array}{c}
\widehat{\Gamma}_{1 n} \\
\vdots \\
\widehat{\Gamma}_{p n}
\end{array}\right) \widehat{\Omega}_{n}^{-1} n^{1 / 2} \widehat{g}_{n}  \tag{15.1}\\
& =n^{-1 / 2} \sum_{i=1}^{n}\left[\operatorname{vec}\left(G_{i}-D_{n}\right)-\left(\begin{array}{c}
E_{F_{n}} G_{\ell 1} g_{\ell}^{\prime} \\
\vdots \\
E_{F_{n}} G_{\ell p} g_{\ell}^{\prime}
\end{array}\right) \Omega_{F_{n}}^{-1} g_{i}\right]+o_{p}(1),
\end{align*}
$$

where the second equality holds by (i) the weak law of large numbers (WLLN) applied to $n^{-1} \sum_{\ell=1}^{n}$ $G_{\ell j} g_{\ell}^{\prime}$ for $j=1, \ldots, p, n^{-1} \sum_{\ell=1}^{n} \operatorname{vec}\left(G_{\ell}\right)$, and $n^{-1} \sum_{\ell=1}^{n} g_{\ell} g_{\ell}^{\prime}$, (ii) $E_{F_{n}} g_{i}=0^{k}$, (iii) $h_{5, g}=\lim \Omega_{F_{n}}$ is pd , and (iv) the CLT, which implies that $n^{1 / 2} \widehat{g}_{n}=O_{p}(1)$.

Using (15.1), the convergence result of Lemma 10.2 holds (with $n$ in place of $w_{n}$ ) by the Lyapunov triangular-array multivariate CLT using the moment restrictions in $\mathcal{F}$. The limiting covariance matrix between $n^{1 / 2} \operatorname{vec}\left(\widehat{D}_{n}-D_{n}\right)$ and $n^{1 / 2} \widehat{g}_{n}$ in Lemma 10.2 is a zero matrix because

$$
\begin{equation*}
E_{F_{n}}\left[G_{i j}-D_{n j}-\left(E_{F_{n}} G_{\ell j} g_{\ell}^{\prime}\right) \Omega_{F_{n}}^{-1} g_{i}\right] g_{i}^{\prime}=0^{k \times k}, \tag{15.2}
\end{equation*}
$$

where $D_{n j}$ denotes the $j$ th column of $D_{n}$, using $E_{F_{n}} g_{i}=0^{k}$ for $j=1, \ldots, p$. By the CLT, the limiting variance matrix of $n^{1 / 2} \operatorname{vec}\left(\widehat{D}_{n}-D_{n}\right)$ in Lemma 10.2 equals

$$
\begin{equation*}
\lim \operatorname{Var}_{F_{n}}\left(\operatorname{vec}\left(G_{i}\right)-\left(E_{F_{n}} \operatorname{vec}\left(G_{\ell}\right) g_{\ell}^{\prime}\right) \Omega_{F_{n}}^{-1} g_{i}\right)=\lim \Phi_{F_{n}}^{\operatorname{vec}\left(G_{i}\right)}=\Phi_{h}^{v e c}\left(G_{i}\right), \tag{15.3}
\end{equation*}
$$

see 10.15 , and the limit exists because (i) the components of $\Phi_{F_{n}}^{v e c\left(G_{i}\right)}$ are comprised of $\lambda_{4, F_{n}}$ and submatrices of $\lambda_{5, F_{n}}$ and (ii) $\lambda_{s, F_{n}} \rightarrow h_{s}$ for $s=4,5$. By the CLT, the limiting variance matrix of
$n^{1 / 2} \widehat{g}_{n}$ equals $\lim E_{F_{n}} g_{i} g_{i}^{\prime}=h_{5, g}$.

## 16 Proof of Lemma 10.3

Lemma 10.3 of AG1. Suppose Assumption WU holds for some non-empty parameter space $\Lambda_{*} \subset \Lambda_{2}$. Under all sequences $\left\{\lambda_{n, h}: n \geq 1\right\}$ with $\lambda_{n, h} \in \Lambda_{*}$,

$$
n^{1 / 2}\left(\widehat{g}_{n}, \widehat{D}_{n}-E_{F_{n}} G_{i}, W_{F_{n}} \widehat{D}_{n} U_{F_{n}} T_{n}\right) \rightarrow_{d}\left(\bar{g}_{h}, \bar{D}_{h}, \bar{\Delta}_{h}\right),
$$

where (a) ( $\bar{g}_{h}, \bar{D}_{h}$ ) are defined in Lemma 10.2 , (b) $\bar{\Delta}_{h}$ is the nonrandom function of $h$ and $\bar{D}_{h}$ defined in 10.17), (c) ( $\bar{D}_{h}, \bar{\Delta}_{h}$ ) and $\bar{g}_{h}$ are independent, (d) if Assumption WU holds with $\Lambda_{*}=\Lambda_{0}$, $W_{F}=\Omega_{F}^{-1 / 2}$, and $U_{F}=I_{p}$, then $\bar{\Delta}_{h}$ has full column rank $p$ with probability one and (e) under all subsequences $\left\{w_{n}\right\}$ and all sequences $\left\{\lambda_{w_{n}, h}: n \geq 1\right\}$ with $\lambda_{w_{n}, h} \in \Lambda_{*}$, the convergence result above and the results of parts (a)-(d) hold with $n$ replaced with $w_{n}$.

The proof of part (d) of Lemma 10.3 uses the following two lemmas and corollary.
Lemma 16.1 Suppose $\Delta \in R^{k \times p}$ has a multivariate normal distribution (with possibly singular variance matrix), $k \geq p$, and the variance matrix of $\Delta \xi \in R^{k}$ has rank at least $p$ for all nonrandom vectors $\xi \in R^{p}$ with $\|\xi\|=1$. Then, $P(\Delta$ has full column rank $p)=1$.

Comments: (i) Let Condition $\Delta$ denote the condition of the lemma on the variance of $\Delta \xi$. A sufficient condition for Condition $\Delta$ is that $\operatorname{vec}(\Delta)$ has a pd variance matrix (because $\Delta \xi=$ $\left.\left(\xi^{\prime} \otimes I_{k}\right) v e c(\Delta)\right)$. The converse is not true. This is proved in Comment (iii) below.
(ii) A weaker sufficient condition for Condition $\Delta$ is that the variance matrix of $\Delta \xi \in R^{k}$ has rank $k$ for all constant vectors $\xi \in R^{p}$ with $\|\xi\|=1$. The latter condition holds iff $\operatorname{Var}\left(\zeta^{\prime} \operatorname{vec}(\Delta)\right)>0$ for all $\zeta \in R^{p k}$ of the form $\zeta=\xi \otimes \mu$ for some $\xi \in R^{p}$ and $\mu \in R^{k}$ with $\|\xi\|=1$ and $\|\mu\|=1$ (because $\left(\xi^{\prime} \otimes \mu^{\prime}\right) \operatorname{vec}(\Delta)=\operatorname{vec}\left(\mu^{\prime} \Delta \xi\right)=\mu^{\prime} \Delta \xi$ ). In contrast, $\operatorname{vec}(\Delta)$ has a pd variance matrix iff $\operatorname{Var}\left(\zeta^{\prime} \operatorname{vec}(\Delta)\right)>0$ for all $\zeta \in R^{p k}$ with $\|\zeta\|=1$.
(iii) For example, the following matrix $\Delta$ satisfies the sufficient condition given in Comment (ii) for Condition $\Delta$ (and hence Condition $\Delta$ holds), but not the sufficient condition given in Comment (i). Let $Z_{j}$ for $j=1,2,3$ be independent standard normal random variables. Define

$$
\Delta=\left(\begin{array}{ll}
Z_{1} & Z_{2}  \tag{16.1}\\
Z_{3} & Z_{1}
\end{array}\right)
$$

Obviously, $\operatorname{Var}(\operatorname{vec}(\Delta))$ is not pd. On the other hand, writing $\xi=\left(\xi_{1}, \xi_{2}\right)^{\prime}$ and $\mu=\left(\mu_{1}, \mu_{2}\right)^{\prime}$, we
have

$$
\begin{align*}
\operatorname{Var}\left(\mu^{\prime} \Delta \xi\right) & =\operatorname{Var}\left(\mu_{1}\left[Z_{1} \xi_{1}+Z_{2} \xi_{2}\right]+\mu_{2}\left[Z_{3} \xi_{1}+Z_{1} \xi_{2}\right]\right) \\
& =\operatorname{Var}\left(\left(\mu_{1} \xi_{1}+\mu_{2} \xi_{2}\right) Z_{1}+\mu_{1} \xi_{2} Z_{2}+\mu_{2} \xi_{1} Z_{3}\right) \\
& =\left(\mu_{1} \xi_{1}+\mu_{2} \xi_{2}\right)^{2}+\left(\mu_{1} \xi_{2}\right)^{2}+\left(\mu_{2} \xi_{1}\right)^{2} \tag{16.2}
\end{align*}
$$

Now, $\left(\mu_{1} \xi_{2}\right)^{2}=0$ implies $\mu_{1}=0$ or $\xi_{2}=0$ and $\left(\mu_{2} \xi_{1}\right)^{2}=0$ implies $\mu_{2}=0$ or $\xi_{1}=0$. In addition, $\mu_{1}=0$ implies $\mu_{2} \neq 0, \xi_{2}=0$ implies $\xi_{1} \neq 0$, etc. So, the two cases where $\left(\mu_{1} \xi_{2}\right)^{2}=\left(\mu_{2} \xi_{1}\right)^{2}=0$ are: $\left(\mu_{1}, \xi_{1}\right)=(0,0)$ and $\left(\mu_{2}, \xi_{2}\right)=(0,0)$. But, $\left(\mu_{1}, \xi_{1}\right)=(0,0)$ implies $\left(\mu_{1} \xi_{1}+\mu_{2} \xi_{2}\right)^{2}=\left(\mu_{2} \xi_{2}\right)^{2}>0$ and $\left(\mu_{2}, \xi_{2}\right)=(0,0)$ implies $\left(\mu_{1} \xi_{1}+\mu_{2} \xi_{2}\right)^{2}=\left(\mu_{1} \xi_{1}\right)^{2}>0$. Hence, $\operatorname{Var}\left(\mu^{\prime} \Delta \xi\right)>0$ for all $\mu$ and $\xi$ with $\|\mu\|=\|\xi\|=1, \operatorname{Var}(\Delta \xi)$ is pd for all $\xi \in R^{2}$ with $\|\xi\|^{2}=1$, and the sufficient condition given in Comment (ii) for Condition $\Delta$ holds.
(iv) Condition $\Delta$ allows for redundant rows in $\Delta$, which corresponds to redundant moment conditions in the application of Lemma 16.1. Suppose a matrix $\Delta$ satisfies Condition $\Delta$. Then, one adds one or more rows to $\Delta$, which consist of one or more of the existing rows of $\Delta$ or some linear combinations of them. (In fact, the added rows can be arbitrary provided the resulting matrix has a multivariate normal distribution.) Call the new matrix $\Delta_{+}$. The matrix $\Delta_{+}$also satisfies Condition $\Delta$ (because the rank of the variance of $\Delta_{+} \xi$ is at least as large as the rank of the variance of $\Delta \xi$, which is $p$ ).

Corollary 16.2 Suppose $\Delta_{q_{*}} \in R^{k \times q_{*}}$ is a nonrandom matrix with full column rank $q_{*}$ and $\Delta_{p-q_{*}} \in$ $R^{k \times\left(p-q_{*}\right)}$ has a multivariate normal distribution (with possibly singular variance matrix) and $k \geq p$. Let $M \in R^{k \times k}$ be a nonsingular matrix such that $M \Delta_{q_{*}}=\left(e_{1}, \ldots, e_{q_{*}}\right)$, where $e_{l}$ denotes the $l$-th coordinate vector in $R^{k}$. Decompose $M=\left(M_{1}^{\prime}, M_{2}^{\prime}\right)^{\prime}$ with $M_{1} \in R^{q_{*} \times k}$ and $M_{2} \in R^{\left(k-q_{*}\right) \times k}$. Suppose the variance matrix of $M_{2} \Delta_{p-q_{*}} \xi_{2} \in R^{k-q_{*}}$ has rank at least $p-q_{*}$ for all nonrandom vectors $\xi_{2} \in R^{p-q_{*}}$ with $\left\|\xi_{2}\right\|=1$. Then, for $\Delta=\left(\Delta_{q_{*}}, \Delta_{p-q_{*}}\right) \in R^{k \times p}$, we have $P(\Delta$ has full column rank $p)=1$.

Comment: Corollary 16.2 follows from Lemma 16.1 by the following argument. We have

$$
M \Delta=\left(\begin{array}{cc}
M_{1} \Delta_{q_{*}} & M_{1} \Delta_{p-q_{*}}  \tag{16.3}\\
M_{2} \Delta_{q_{*}} & M_{2} \Delta_{p-q_{*}}
\end{array}\right)=\left(\begin{array}{cc}
I_{q_{*}} & M_{1} \Delta_{p-q_{*}} \\
0^{\left(k-q_{*}\right) \times q_{*}} & M_{2} \Delta_{p-q_{*}}
\end{array}\right) .
$$

The matrix $\Delta$ has full column rank $p$ iff $M \Delta$ has full column rank $p$ iff $M_{2} \Delta_{p-q_{*}}$ has full column rank $p-q_{*}$. The Corollary now follows from Lemma 16.1 applied with $\Delta, k, p$, and $\xi$ replaced by $M_{2} \Delta_{p-q_{*}}, k-q_{*}, p-q_{*}$, and $\xi_{2}$, respectively.

The following lemma is a special case of Cauchy's interlacing eigenvalues result, e.g., see Hwang (2004). As above, for a symmetric matrix $A$, let $\lambda_{1}(A) \geq \lambda_{2}(A) \geq \ldots$ denote the eigenvalues of $A$. Let $A_{-r}$ denote a principal submatrix of $A$ of order $r \geq 1$. That is, $A_{-r}$ denotes $A$ with some choice of $r$ rows and the same $r$ columns deleted.

Proposition 16.3 Let $A$ by a symmetric $k \times k$ matrix. Then, $\lambda_{k}(A) \leq \lambda_{k-1}\left(A_{-1}\right) \leq \lambda_{k-1}(A) \leq$ $\ldots \leq \lambda_{2}(A) \leq \lambda_{1}\left(A_{-1}\right) \leq \lambda_{1}(A)$.

The following is a straightforward corollary of Proposition 16.3 .

Corollary 16.4 Let $A$ by a symmetric $k \times k$ matrix and let $r \in\{1, \ldots, k-1\}$. Then, (a) $\lambda_{m}(A) \geq$ $\lambda_{m}\left(A_{-r}\right)$ for $m=1, \ldots, k-r$ and (b) $\lambda_{m}(A) \leq \lambda_{m-r}\left(A_{-r}\right)$ for $m=r+1, \ldots, k$.

Proof of Lemma 10.3. First, we prove the convergence result in Lemma 10.3. The singular value decomposition of $W_{n} D_{n} U_{n}$ is

$$
\begin{equation*}
W_{n} D_{n} U_{n}=C_{n} \Upsilon_{n} B_{n}^{\prime}, \tag{16.4}
\end{equation*}
$$

because $B_{n}$ is a matrix of eigenvectors of $U_{n}^{\prime} D_{n}^{\prime} W_{n}^{\prime} W_{n} D_{n} U_{n}, C_{n}$ is a matrix of eigenvectors of $W_{n} D_{n} U_{n} U_{n}^{\prime} D_{n}^{\prime} W_{n}^{\prime}$, and $\Upsilon_{n}$ is the $k \times p$ matrix with the singular values $\left\{\tau_{j F_{n}}: j \leq p\right\}$ of $W_{n} D_{n} U_{n}$ on the diagonal (ordered so that $\tau_{j F_{n}} \geq 0$ is nonincreasing in $j$ ).

Using (16.4), we get

$$
\begin{equation*}
W_{n} D_{n} U_{n} B_{n, q} \Upsilon_{n, q}^{-1}=C_{n} \Upsilon_{n} B_{n}^{\prime} B_{n, q} \Upsilon_{n, q}^{-1}=C_{n} \Upsilon_{n}\binom{I_{q}}{0^{(p-q) \times q}} \Upsilon_{n, q}^{-1}=C_{n}\binom{I_{q}}{0^{(k-q) \times q}}=C_{n, q}, \tag{16.5}
\end{equation*}
$$

where the second equality uses $B_{n}^{\prime} B_{n}=I_{p}$. Hence, we obtain

$$
\begin{align*}
W_{n} \widehat{D}_{n} U_{n} B_{n, q} \Upsilon_{n, q}^{-1} & =W_{n} D_{n} U_{n} B_{n, q} \Upsilon_{n, q}^{-1}+W_{n} n^{1 / 2}\left(\widehat{D}_{n}-D_{n}\right) U_{n} B_{n, q}\left(n^{1 / 2} \Upsilon_{n, q}\right)^{-1} \\
& =C_{n, q}+o_{p}(1) \rightarrow_{p} h_{3, q}=\bar{\Delta}_{h, q}, \tag{16.6}
\end{align*}
$$

where the second equality uses $n^{1 / 2} \tau_{j F_{n}} \rightarrow \infty$ for all $j \leq q$ (by the definition of $q$ in 10.16), $W_{n}=O(1)$ (by the condition $\left\|W_{F}\right\| \leq M_{1}<\infty \forall F \in \mathcal{F}_{W U}$, see 10.5 ), $n^{1 / 2}\left(\widehat{D}_{n}-D_{n}\right)=O_{p}(1)$ (by Lemma 10.2), $U_{n}=O(1)$ (by the condition $\left\|U_{F}\right\| \leq M_{1}<\infty \forall F \in \mathcal{F}_{W U}$, see 10.5), and $B_{n, q} \rightarrow h_{2, q}$ with $\left\|v e c\left(h_{2, q}\right)\right\|<\infty$ (by (10.12) using the definitions in 10.17) and (9.1)). The convergence in (16.6) holds by (10.12), 10.17), and (9.1), and the last equality in (16.6) holds by the definition of $\bar{\Delta}_{h, q}$ in 10.17 .

Using (16.4) again, we have

$$
\begin{align*}
& n^{1 / 2} W_{n} D_{n} U_{n} B_{n, p-q}=n^{1 / 2} C_{n} \Upsilon_{n} B_{n}^{\prime} B_{n, p-q}=n^{1 / 2} C_{n} \Upsilon_{n}\binom{0^{q \times(p-q)}}{I_{p-q}} \\
& =C_{n}\left(\begin{array}{c}
0^{q \times(p-q)} \\
n^{1 / 2} \Upsilon_{n, p-q} \\
0^{(k-p) \times(p-q)}
\end{array}\right) \rightarrow h_{3}\left(\begin{array}{c}
0^{q \times(p-q)} \\
\operatorname{Diag}\left\{h_{1, q+1}, \ldots, h_{1, p}\right\} \\
0^{(k-p) \times(p-q)}
\end{array}\right)=h_{3} h_{1, p-q}^{\diamond}, \tag{16.7}
\end{align*}
$$

where the second equality uses $B_{n}^{\prime} B_{n}=I_{p}$, the convergence holds by 10.12 using the definitions in (10.17) and (9.2), and the last equality holds by the definition of $h_{1, p-q}^{\circ}$ in 10.17).

Using 16.7) and Lemma 10.2, we get

$$
\begin{align*}
n^{1 / 2} W_{n} \widehat{D}_{n} U_{n} B_{n, p-q} & =n^{1 / 2} W_{n} D_{n} U_{n} B_{n, p-q}+W_{n} n^{1 / 2}\left(\widehat{D}_{n}-D_{n}\right) U_{n} B_{n, p-q} \\
& \rightarrow{ }_{d} h_{3} h_{1, p-q}^{\diamond}+h_{71} \bar{D}_{h} h_{81} h_{2, p-q}=\bar{\Delta}_{h, p-q}, \tag{16.8}
\end{align*}
$$

where $B_{n, p-q} \rightarrow h_{2, p-q}, W_{n} \rightarrow h_{71}$, and $U_{n} \rightarrow h_{81}$ by (10.3), 10.12, 10.17), and Assumption WU using the definitions in (9.1) and the last equality holds by the definition of $\bar{\Delta}_{h, p-q}$ in 10.17).

Equations 16.6 and 16.8 combine to prove

$$
\begin{align*}
n^{1 / 2} W_{n} \widehat{D}_{n} U_{n} T_{n} & =n^{1 / 2} W_{n} \widehat{D}_{n} U_{n} B_{n} S_{n}=\left(W_{n} \widehat{D}_{n} U_{n} B_{n, q} \Upsilon_{n, q}^{-1}, n^{1 / 2} W_{n} \widehat{D}_{n} U_{n} B_{n, p-q}\right) \\
& \rightarrow_{d}\left(\bar{\Delta}_{h, q}, \bar{\Delta}_{h, p-q}\right)=\bar{\Delta}_{h} \tag{16.9}
\end{align*}
$$

using the definition of $S_{n}$ in 10.19 . The convergence is joint with that in Lemma 10.2 because it just relies on the convergence of $n^{1 / 2}\left(\widehat{D}_{n}-D_{n}\right)$, which is part of the former. This establishes the convergence result of Lemma 10.3 .

Properties (a) and (b) in Lemma 10.3 hold by definition. Property (c) in Lemma 10.3 holds by Lemma 10.2 and property (b) in Lemma 10.3 .

Now, we prove property (d). We have

$$
\begin{equation*}
h_{2, p-q}^{\prime} h_{2, p-q}=\lim B_{n, p-q}^{\prime} B_{n, p-q}=I_{p-q} \text { and } h_{3, q}^{\prime} h_{3, q}=\lim C_{n, q}^{\prime} C_{n, q}=I_{q} \tag{16.10}
\end{equation*}
$$

because $B_{n}$ and $C_{n}$ are orthogonal matrices by 10.6 and 10.7. Hence, if $q=p$, then $\bar{\Delta}_{h}=$ $\bar{\Delta}_{h, q}=h_{3, q}, \bar{\Delta}_{h}^{\prime} \bar{\Delta}_{h}=I_{p}$, and $\bar{\Delta}_{h}$ has full column rank.

Hence, it suffices to consider the case where $q<p$ and $\lambda_{n, h} \in \Lambda_{0} \forall n \geq 1$, which is assumed in part (d). We prove part (d) for this case by applying Corollary 16.2 with $q_{*}=q, \Delta_{q_{*}}=\bar{\Delta}_{h, q}\left(=h_{3, q}\right)$,
$\Delta_{p-q_{*}}=\bar{\Delta}_{h, p-q}, M=h_{3}^{\prime}, M_{1}=h_{3, q}^{\prime}, M_{2}=h_{3, k-q}^{\prime}, \xi_{2} \in R^{p-q}$, and $\Delta=\bar{\Delta}_{h}$. Corollary 16.2 gives the desired result that $P\left(\bar{\Delta}_{h}\right.$ has full column rank $\left.p\right)=1$. The condition in Corollary 16.2 that " $M \Delta_{q_{*}}=\left(e_{1}, \ldots, e_{q_{*}}\right)$ " holds in this case because $h_{3}^{\prime} \bar{\Delta}_{h, q}=h_{3}^{\prime} h_{3, q}=\left(e_{1}, \ldots, e_{q}\right)$. The condition in Corollary 16.2 that "the variance matrix of $M_{2} \Delta_{p-q_{*}} \xi_{2} \in R^{k-q_{*}}$ has rank at least $p-q_{*}$ for all nonrandom vectors $\xi_{2} \in R^{p-q_{*}}$ with $\left\|\xi_{2}\right\|=1$ " in this case becomes "the variance matrix of $h_{3, k-q}^{\prime} \bar{\Delta}_{h, p-q} \xi_{2} \in R^{k-q}$ has rank at least $p-q$ for all nonrandom vectors $\xi_{2} \in R^{p-q}$ with $\left\|\xi_{2}\right\|=1$." It remains to establish the latter property, which is equivalent to

$$
\begin{equation*}
\lambda_{p-q}\left(\operatorname{Var}\left(h_{3, k-q}^{\prime} \bar{\Delta}_{h, p-q} \xi_{2}\right)\right)>0 \forall \xi_{2} \in R^{p-q} \text { with }\left\|\xi_{2}\right\|=1 . \tag{16.11}
\end{equation*}
$$

We have

$$
\begin{align*}
\operatorname{Var}\left(h_{3, k-q}^{\prime} \bar{\Delta}_{h, p-q} \xi_{2}\right) & =\operatorname{Var}\left(h_{3, k-q}^{\prime} h_{5, g}^{-1 / 2} \bar{D}_{h} h_{2, p-q} \xi_{2}\right) \\
& =\left(\left(h_{2, p-q} \xi_{2}\right)^{\prime} \otimes\left(h_{3, k-q}^{\prime} h_{5, g}^{-1 / 2}\right)\right) \operatorname{Var}\left(\operatorname{vec}\left(\bar{D}_{h}\right)\right)\left(\left(h_{2, p-q} \xi_{2}\right) \otimes\left(h_{3, k-q}^{\prime} h_{5, g}^{-1 / 2}\right)^{\prime}\right) \\
& =\left(\left(h_{2, p-q} \xi_{2}\right)^{\prime} \otimes\left(h_{3, k-q}^{\prime} h_{5, g}^{-1 / 2}\right)\right) \Phi_{h}^{\operatorname{vec}\left(G_{i}\right)}\left(\left(h_{2, p-q} \xi_{2}\right) \otimes\left(h_{3, k-q}^{\prime} h_{5, g}^{-1 / 2}\right)^{\prime}\right) \\
& =\Phi_{h}^{h_{3, k-q}^{\prime} h_{5, g}^{-1 / 2} G_{i} h_{2, p-q} \xi_{2}}, \tag{16.12}
\end{align*}
$$

where the first equality holds by the definition of $\bar{\Delta}_{h, p-q}$ in 10.17 and the fact that $h_{71}=h_{5, g}^{-1 / 2}$ and $h_{81}=I_{p}$ by the conditions in part (d) of Lemma 10.3 , the second and fourth equalities use the general formula $\operatorname{vec}(A B C)=\left(C^{\prime} \otimes A\right) \operatorname{vec}(B)$, the third equality holds because $\operatorname{vec}\left(\bar{D}_{h}\right) \sim$ $N\left(0^{p k}, \Phi_{h}^{v e c\left(G_{i}\right)}\right)$ by Lemma 10.2 , and the fourth equality uses the definition of the variance matrix $\Phi_{h}^{a_{i}}$ in 10.15 for an arbitrary random vector $a_{i}$.

Next, we show that $\Phi_{h}^{h_{3, k-q}^{\prime} h_{5, g}^{-1 / 2} G_{i} h_{2, p-q} \xi_{2}}$ equals the expected outer-product matrix $\lim \Psi_{F_{n}}^{C_{n, k-q}^{\prime} \Omega_{n}^{-1 / 2} G_{i} B_{n, p-q} \xi_{2}}$ :

$$
\begin{align*}
& \Phi_{h}^{h_{3, k-q}^{\prime} h_{5, g}^{-1 / 2} G_{i} h_{2, p-q} \xi_{2}} \\
= & \left(\left(h_{2, p-q} \xi_{2}\right)^{\prime} \otimes\left(h_{3, k-q}^{\prime} h_{5, g}^{-1 / 2}\right)\right) \Phi_{h}^{v e c\left(G_{i}\right)}\left(\left(h_{2, p-q} \xi_{2}\right) \otimes\left(h_{3, k-q}^{\prime} h_{5, g}^{-1 / 2}\right)^{\prime}\right) \\
= & \lim \left(\left(B_{n, p-q} \xi_{2}\right)^{\prime} \otimes\left(C_{n, k-q}^{\prime} \Omega_{n}^{-1 / 2}\right)\right) \Phi_{F_{n}}^{v e c\left(G_{i}\right)}\left(\left(B_{n, p-q} \xi_{2}\right) \otimes\left(C_{n, k-q}^{\prime} \Omega_{n}^{-1 / 2}\right)^{\prime}\right) \\
= & \lim \left(\left(B_{n, p-q} \xi_{2}\right)^{\prime} \otimes\left(C_{n, k-q}^{\prime} \Omega_{n}^{-1 / 2}\right)\right) \Psi_{F_{n}}^{v e c\left(G_{i}\right)}\left(\left(B_{n, p-q} \xi_{2}\right) \otimes\left(C_{n, k-q}^{\prime} \Omega_{n}^{-1 / 2}\right)^{\prime}\right) \\
& \quad-\lim \left(\left(B_{n, p-q} \xi_{2}\right)^{\prime} \otimes\left(C_{n, k-q}^{\prime} \Omega_{n}^{-1 / 2}\right)\right) E_{F_{n}} v e c\left(G_{i}\right) \cdot E_{F_{n}} v e c\left(G_{i}\right)^{\prime}\left(\left(B_{n, p-q} \xi_{2}\right) \otimes\left(C_{n, k-q}^{\prime} \Omega_{n}^{-1 / 2}\right)^{\prime}\right) \\
= & \lim \left(\left(B_{n, p-q} \xi_{2}\right)^{\prime} \otimes\left(C_{n, k-q}^{\prime} \Omega_{n}^{-1 / 2}\right)\right) \Psi_{F_{n}}^{v e c\left(G_{i}\right)}\left(\left(B_{n, p-q} \xi_{2}\right) \otimes\left(C_{n, k-q}^{\prime} \Omega_{n}^{-1 / 2}\right)^{\prime}\right) \\
& \quad-\lim E_{F_{n}} v e c\left(C_{n, k-q}^{\prime} \Omega_{n}^{-1 / 2} G_{i} B_{n, p-q} \xi_{2}\right) \cdot E_{F_{n}} v e c\left(C_{n, k-q}^{\prime} \Omega_{n}^{-1 / 2} G_{i} B_{n, p-q} \xi_{2}\right)^{\prime} \\
= & \lim \Psi_{F_{n}}^{C_{n, k-q}^{\prime} \Omega_{n}^{-1 / 2} G_{i} B_{n, p-q} \xi_{2}}, \tag{16.13}
\end{align*}
$$

where the general formula $\operatorname{vec}(A B C)=\left(C^{\prime} \otimes A\right) \operatorname{vec}(B)$ is used multiple times, the limits exist by the conditions imposed on the sequence $\left\{\lambda_{n, h}: n \geq 1\right\}$, the second equality uses $B_{n, p-j} \rightarrow h_{2, p-j}$, $C_{n, k-q} \rightarrow h_{3, k-q}$, and $\Omega_{n}^{-1 / 2} \rightarrow h_{5, g}^{-1 / 2}$, the third equality uses the definitions of $\Psi_{F}^{a_{i}}$ and $\Phi_{F}^{a_{i}}$ given in 3.6 and 10.15, respectively, and the last equality uses $E_{F_{n}} v e c\left(C_{n, k-q}^{\prime} \Omega_{n}^{-1 / 2} G_{i} B_{n, p-q}\right)=$ $\operatorname{vec}\left(C_{n, k-q}^{\prime} \Omega_{n}^{-1 / 2} D_{n} B_{n, p-q}\right)=O\left(n^{-1 / 2}\right)$ by 16.7 with $W_{n}=\Omega_{n}^{-1 / 2}$.

We can write $\lim \Psi_{F_{n}}^{v e c}\left(C_{n}^{\prime} \Omega_{n}^{-1 / 2} G_{i} B_{n}\right)$ as the limit of a subsequence $\left\{n_{m}: m \geq 1\right\}$ of matrices $\Psi_{F_{n_{m}}}^{v e c\left(C_{n_{m}}^{\prime} \Omega_{n_{m}}^{-1 / 2} G_{i} B_{n_{m}}\right)}$ for which $F_{n_{m}} \in \mathcal{F}_{0 j}$ for all $m \geq 1$ for some $j=0, \ldots, q$. It cannot be the case that $j>q$, because if $j>q$, then we obtain a contradiction because $n_{m}^{1 / 2} \tau_{j F_{n_{m}}} \rightarrow \infty$ as $m \rightarrow \infty$ by the first condition of $\mathcal{F}_{0 j}$ and $n_{m}^{1 / 2} \tau_{j F_{n_{m}}} \nrightarrow \infty$ as $m \rightarrow \infty$ by the definition of $q$ in 10.16.

Now, we fix an arbitrary $j \in\{0, \ldots, q\}$. The continuity of the $\lambda_{p-j}(\cdot)$ function and the $\lambda_{p-j}(\cdot)$ condition in $\mathcal{F}_{0 j}$ imply that, for all $\xi \in R^{p-j}$ with $\|\xi\|=1$,

$$
\begin{equation*}
\lambda_{p-j}\left(\lim \Psi_{F_{n_{m}}}^{C_{m, j-j}^{\prime} \Omega_{n}^{-1 / 2} G_{i} B_{n_{m}, p-j} \xi}\right)=\lim \lambda_{p-j}\left(\Psi_{F_{n_{m}}}^{C_{n_{m}, k-j}^{\prime} \Omega_{n m}^{-1 / 2} G_{i} B_{n_{m}, p-j} \xi}\right)>0 . \tag{16.14}
\end{equation*}
$$

For all $\xi_{2} \in R^{p-q}$ with $\left\|\xi_{2}\right\|=1$, let $\xi=\left(0^{q-j}, \xi_{2}^{\prime}\right)^{\prime} \in R^{p-j}$. Then, $B_{n_{m}, p-j} \xi=B_{n_{m}, p-q} \xi_{2}$ and, by (16.14),

$$
\begin{equation*}
\lambda_{p-j}\left(\lim \Psi_{F_{n}}^{C_{n_{m}, k-j}^{\prime} \Omega_{n_{m}}^{-1 / 2} G_{i} B_{n_{m}, p-q} \xi_{2}}\right)>0 \forall \xi_{2} \in R^{p-q} \text { with }\left\|\xi_{2}\right\|=1 . \tag{16.15}
\end{equation*}
$$

Next, we apply Corollary 16.4 (b) with $A=\lim \Psi_{F_{n_{m}}}^{C_{n_{m}, k-j}^{\prime} \Omega_{n_{m}}^{-1 / 2} G_{i} B_{n_{m}, p-q} \xi_{2}}$ and $A_{-(q-j)}=\lim$ $\Psi_{F_{n_{m}}}^{C_{n_{m}, k-q}^{\prime} \Omega_{n}^{-1 / 2} G_{i} B_{n_{m}, p-q} \xi_{2}}, m=p-j, r=q-j$, where $A_{-(q-j)}$ equals $A$ with its first $q-j$ rows and columns deleted in the present case and $p>q$ implies that $m=p-j \geq 1$ for all $j=0, \ldots, q$. Corollary 16.4 and 16.15 give

$$
\begin{equation*}
\lambda_{p-q}\left(\lim \Psi_{F_{n_{m}}}^{C_{n}^{\prime}, k-q} \Omega_{n_{m}}^{-1 / 2} G_{i} B_{n_{m}, p-q} \xi_{2}\right)>0 \forall \xi_{2} \in R^{p-q} \text { with }\left\|\xi_{2}\right\|=1 . \tag{16.16}
\end{equation*}
$$

Equations (16.12, 16.13), and 16.16) combine to establish 16.11) and the proof of part (d) is complete.

Part (e) of the Lemma holds by replacing $n$ by the subsequence value $w_{n}$ throughout the arguments given above.

Proof of Lemma 16.1. It suffices to show that $P\left(\Delta \xi=0^{k}\right.$ for some $\xi \in R^{p}$ with $\left.\|\xi\|=1\right)=0$.
For any constant $\gamma>0$, there exists a constant $K_{\gamma}<\infty$ such that $P\left(\|v e c(\Delta)\|>K_{\gamma}\right) \leq \gamma$.
Given $\varepsilon>0$, let $\left\{B\left(\xi_{s}, \varepsilon\right): s=1, \ldots, N_{\varepsilon}\right\}$ be a finite cover of $\left\{\xi \in R^{p}:\|\xi\|=1\right\}$, where $\left\|\xi_{s}\right\|=1$ and $B\left(\xi_{s}, \varepsilon\right)$ is a ball in $R^{p}$ centered at $\xi_{s}$ of radius $\varepsilon$. It is possible to choose $\left\{\xi_{s}: s=1, \ldots, N_{\varepsilon}\right\}$ such that the number, $N_{\varepsilon}$, of balls in the cover is of order $\varepsilon^{-p+1}$. That is, $N_{\varepsilon} \leq C_{1} \varepsilon^{-p+1}$ for some
constant $C_{1}<\infty$.
Let $\Delta_{r}$ denote the $r$ th row of $\Delta$ for $r=1, \ldots, k$ written as a column vector. If $\xi \in B\left(\xi_{s}, \varepsilon\right)$, we have

$$
\begin{equation*}
\left\|\Delta \xi-\Delta \xi_{s}\right\|=\left(\sum_{r=1}^{k}\left(\Delta_{r}^{\prime}\left(\xi-\xi_{s}\right)\right)^{2}\right)^{1 / 2} \leq\left(\sum_{r=1}^{k}\left\|\Delta_{r}\right\|^{2}\left\|\xi-\xi_{s}\right\|^{2}\right)^{1 / 2}=\varepsilon\|\operatorname{vec}(\Delta)\| \tag{16.17}
\end{equation*}
$$

where the inequality holds by the Cauchy-Bunyakovsky-Schwarz inequality. If $\xi \in B\left(\xi_{s}, \varepsilon\right)$ and $\Delta \xi=0^{k}$, this gives

$$
\begin{equation*}
\left\|\Delta \xi_{s}\right\| \leq \varepsilon\|\operatorname{vec}(\Delta)\| \tag{16.18}
\end{equation*}
$$

Suppose $Z_{*} \in R^{p}$ has a multivariate normal distribution with pd variance matrix. Then, for any $\varepsilon>0$,

$$
\begin{equation*}
P\left(\left\|Z_{*}\right\| \leq \varepsilon\right)=\int_{\{\|z\| \leq \varepsilon\}} f_{Z_{*}}(z) d z \leq \sup _{z \in R^{k}} f_{Z_{*}}(z) \int_{\{\|z\| \leq \varepsilon\}} d z \leq C_{2} \varepsilon^{p} \tag{16.19}
\end{equation*}
$$

for some constant $C_{2}<\infty$, where $f_{Z_{*}}(z)$ denotes the density of $Z_{*}$ with respect to Lebesgue measure, which exists because the variance matrix of $Z_{*}$ is pd , and the inequalities hold because the density of a multivariate normal is bounded and the volume of a sphere in $R^{p}$ of radius $\varepsilon$ is proportional to $\varepsilon^{p}$.

For any $\xi \in R^{p}$ with $\|\xi\|=1$, let $B_{\xi} \Lambda_{\xi} B_{\xi}^{\prime}$ be a spectral decomposition of $\operatorname{Var}(\Delta \xi)$, where $\Lambda_{\xi}$ is the diagonal $k \times k$ matrix with the eigenvalues of $\operatorname{Var}(\Delta \xi)$ on its diagonal in nonincreasing order and $B_{\xi}$ is an orthogonal $k \times k$ matrix whose columns are eigenvectors of $\operatorname{Var}(\Delta \xi)$ that correspond to the eigenvalues in $\Lambda_{\xi}$. By assumption, the rank of $\operatorname{Var}(\Delta \xi)$ is $p$ or larger. In consequence, the first $p$ diagonal elements of $\Lambda_{\xi}$ are positive. We have $\|\Delta \xi\|=\left\|B_{\xi}^{\prime} \Delta \xi\right\|$ and $\operatorname{Var}\left(B_{\xi}^{\prime} \Delta \xi\right)=$ $B_{\xi}^{\prime} \operatorname{Var}(\Delta \xi) B_{\xi}=\Lambda_{\xi}$. Let $\left(B_{\xi}^{\prime} \Delta \xi\right)_{p}$ denote the $p$ vector that contains the first $p$ elements of the $k$ vector $B_{\xi}^{\prime} \Delta \xi$. Let $\Lambda_{\xi p}$ denote the upper left $p \times p$ submatrix of $\Lambda_{\xi}$. We have $\operatorname{Var}\left(\left(B_{\xi}^{\prime} \Delta \xi\right)_{p}\right)=\Lambda_{\xi p}$ and $\Lambda_{\xi p}$ is pd (because the first $p$ diagonal elements of $\Lambda_{\xi}$ are positive).

Now, given any $\gamma>0$ and $\varepsilon>0$, we have

$$
\begin{align*}
& P\left(\Delta \xi=0^{k} \text { for some } \xi \in R^{p} \text { with }\|\xi\|=1\right) \\
= & P\left(\cup_{s=1}^{N_{\varepsilon}} \cup_{\xi \in B\left(\xi_{s}, \varepsilon\right):\|\xi\|=1}\left\{\Delta \xi=0^{k}\right\}\right) \\
\leq & P\left(\cup_{s=1}^{N_{\varepsilon}}\left\{\left\|\Delta \xi_{s}\right\| \leq \varepsilon \| \text { vec }(\Delta) \|\right\}\right) \\
\leq & P\left(\cup_{s=1}^{N_{\varepsilon}}\left\{\left\|\Delta \xi_{s}\right\| \leq \varepsilon \| \text { vec }(\Delta) \|\right\} \cap\left\{\| \text { vec }(\Delta) \| \leq K_{\gamma}\right\}\right)+P\left(\| \text { vec }(\Delta) \|>K_{\gamma}\right) \\
\leq & P\left(\cup_{s=1}^{N_{\varepsilon}}\left\{\left\|\Delta \xi_{s}\right\| \leq \varepsilon K_{\gamma}\right\}\right)+\gamma \\
\leq & \sum_{s=1}^{N_{\varepsilon}} P\left(\left\|\Delta \xi_{s}\right\| \leq \varepsilon K_{\gamma}\right)+\gamma \\
\leq & \sum_{s=1}^{N_{\varepsilon}} P\left(\left\|\left(B_{\xi_{s}}^{\prime} \Delta \xi_{s}\right)_{p}\right\| \leq \varepsilon K_{\gamma}\right)+\gamma \\
\leq & N_{\varepsilon} C_{2} K_{\gamma}^{p} \varepsilon^{p}+\gamma \\
\leq & C_{1} \varepsilon^{-p+1} C_{2} K_{\gamma}^{p} \varepsilon^{p}+\gamma \\
\rightarrow & \gamma \text { as } \varepsilon \rightarrow 0 \tag{16.20}
\end{align*}
$$

where the first inequality holds by 16.18 using $\xi \in B\left(\xi_{s}, \varepsilon\right)$, the third inequality uses the definition of $K_{\gamma}$, the third last inequality holds because $\left\|\left(B_{\xi_{s}}^{\prime} \Delta \xi_{s}\right)_{p}\right\| \leq\left\|B_{\xi_{s}}^{\prime} \Delta \xi_{s}\right\|=\left\|\Delta \xi_{s}\right\|$ using the definitions in the paragraph that follows the paragraph that contains 16.19), the second last inequality holds by 16.19 with $Z_{*}=\left(B_{\xi_{s}}^{\prime} \Delta \xi_{s}\right)_{p}$ and the fact that the variance matrix of $\left(B_{\xi_{s}}^{\prime} \Delta \xi_{s}\right)_{p}$ is pd by the argument given in the paragraph following (16.19), and the last inequality holds by the bound given above on $N_{\varepsilon}$.

Because $\gamma>0$ is arbitrary, 16.20 implies that $P\left(\Delta \xi=0^{k}\right.$ for some $\xi \in R^{p}$ with $\left.\|\xi\|=1\right)=0$, which completes the proof.

## 17 Proof of Theorem 10.4

Theorem 10.4 of AG1. Suppose Assumption WU holds for some non-empty parameter space $\Lambda_{*} \subset \Lambda_{2}$. Under all sequences $\left\{\lambda_{n, h}: n \geq 1\right\}$ with $\lambda_{n, h} \in \Lambda_{*}$,
(a) $\widehat{\kappa}_{p n} \rightarrow_{p} \infty$ if $q=p$,
(b) $\widehat{\kappa}_{p n} \rightarrow_{d} \lambda_{\min }\left(\bar{\Delta}_{h, p-q}^{\prime} h_{3, k-q} h_{3, k-q}^{\prime} \bar{\Delta}_{h, p-q}\right)$ if $q<p$,
(c) $\widehat{\kappa}_{j n} \rightarrow_{p} \infty$ for all $j \leq q$,
(d) the (ordered) vector of the smallest $p-q$ eigenvalues of $n \widehat{U}_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} \widehat{U}_{n}$, i.e., $\left(\widehat{\kappa}_{(q+1) n}, \ldots\right.$, $\left.\widehat{\kappa}_{p n}\right)^{\prime}$, converges in distribution to the (ordered) $p-q$ vector of the eigenvalues of $\bar{\Delta}_{h, p-q}^{\prime} h_{3, k-q} h_{3, k-q}^{\prime}$ $\times \bar{\Delta}_{h, p-q} \in R^{(p-q) \times(p-q)}$,
(e) the convergence in parts (a)-(d) holds jointly with the convergence in Lemma 10.3 , and
(f) under all subsequences $\left\{w_{n}\right\}$ and all sequences $\left\{\lambda_{w_{n}, h}: n \geq 1\right\}$ with $\lambda_{w_{n}, h} \in \Lambda_{*}$, the results in parts (a)-(e) hold with $n$ replaced with $w_{n}$.

The proof of Theorem 10.4 uses the following rate of convergence lemma. This lemma is a key technical contribution of the paper.

Lemma 17.1 Suppose Assumption WU holds for some non-empty parameter space $\Lambda_{*} \subset \Lambda_{2}$. Under all sequences $\left\{\lambda_{n, h}: n \geq 1\right\}$ with $\lambda_{n, h} \in \Lambda_{*}$ and for which $q$ defined in 10.16) satisfies $q \geq 1$, we have (a) $\widehat{\kappa}_{j n} \rightarrow_{p} \infty$ for $j=1, \ldots, q$ and (b) when $p>q, \widehat{\kappa}_{j n}=o_{p}\left(\left(n^{1 / 2} \tau_{\ell F_{n}}\right)^{2}\right)$ for all $\ell \leq q$ and $j=q+1, \ldots, p$. Under all subsequences $\left\{w_{n}\right\}$ and all sequences $\left\{\lambda_{w_{n}, h}: n \geq 1\right\}$ with $\lambda_{w_{n}, h} \in \Lambda_{*}$, the same result holds with $n$ replaced with $w_{n}$.

Proof of Lemma 17.1. By the definitions in 10.9) and 10.12), $h_{6, j}:=\lim \tau_{(j+1) F_{n}} / \tau_{j F_{n}}$ for $j=1, \ldots, p-1$. By the definition of $q$ in 10.16), $h_{6, q}=0$ if $q<p$. If $q=p, h_{6, q}$ is not defined by (10.9) and 10.12 and we define it here to equal zero. Because $\tau_{j F}$ is nonnegative and nonincreasing in $j, h_{6, j} \in[0,1]$. If $h_{6, j}>0$, then $\left\{\tau_{j F_{n}}: n \geq 1\right\}$ and $\left\{\tau_{(j+1) F_{n}}: n \geq 1\right\}$ are of the same order of magnitude, i.e., $0<\lim \tau_{(j+1) F_{n}} / \tau_{j F_{n}} \leq 1{ }^{11}$ We group the first $q$ singular values into groups that have the same order of magnitude within each group. Let $G_{h}(\in\{1, \ldots, q\})$ denote the number of groups. (We have $G_{h} \geq 1$ because $q \geq 1$ is assumed in the statement of the lemma.) Note that $G_{h}$ equals the number of values in $\left\{h_{6,1}, \ldots, h_{6, q}\right\}$ that equal zero. Let $r_{g}$ and $r_{g}^{\diamond}$ denote the indices of the first and last singular values, respectively, in the $g$ th group for $g=1, \ldots, G_{h}$. Thus, $r_{1}=1$, $r_{g}^{\diamond}=r_{g+1}-1$, where $r_{G_{h}+1}$ is defined to equal $q+1$, and $r_{G_{h}}^{\diamond}=q$. Note that $r_{g}$ and $r_{g}^{\diamond}$ depend on $h$. By definition, the singular values in the $g$ th group, which have the $g$ th largest order of magnitude, are $\left\{\tau_{r_{g} F_{n}}: n \geq 1\right\}, \ldots,\left\{\tau_{r_{g}^{\diamond} F_{n}}: n \geq 1\right\}$. By construction, $h_{6, j}>0$ for all $j \in\left\{r_{g}, \ldots, r_{g}^{\diamond}-1\right\}$ for $g=1, \ldots, G_{h}$. (The reason is: if $h_{6, j}$ is equal to zero for some $j \in\left\{r_{g}, \ldots, r_{g}^{\diamond}-1\right\}$, then $\left\{\tau_{r_{g}^{\diamond} F_{n}}: n \geq 1\right\}$ is of smaller order of magnitude than $\left\{\tau_{r_{g} F_{n}}: n \geq 1\right\}$, which contradicts the definition of $r_{g}^{\diamond}$.) Also by construction, $\lim \tau_{j^{\prime} F_{n}} / \tau_{j F_{n}}=0$ for any $\left(j, j^{\prime}\right)$ in groups $\left(g, g^{\prime}\right)$, respectively, with $g<g^{\prime}$. Note that when $p=1$ we have $G_{h}=1$ and $r_{1}=r_{1}^{\diamond}=1$.

The eigenvalues $\left\{\widehat{\kappa}_{j n}: j \leq p\right\}$ of $n \widehat{U}_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} \widehat{U}_{n}$ are solutions to the determinantal equation $\left|n \widehat{U}_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} \widehat{U}_{n}-\kappa I_{p}\right|=0$. Equivalently, by multiplying this equation by $\tau_{r_{1} F_{n}}^{-2} n^{-1}\left|B_{n}^{\prime} U_{n}^{\prime} \widehat{U}_{n}^{-1^{\prime}}\right|$ $\times\left|\widehat{U}_{n}^{-1} U_{n} B_{n}\right|$, they are solutions to

$$
\begin{equation*}
\left|\tau_{r_{1} F_{n}}^{-2} B_{n}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n}-\left(n^{1 / 2} \tau_{r_{1} F_{n}}\right)^{-2} \kappa B_{n}^{\prime} U_{n}^{\prime} \widehat{U}_{n}^{-1} \prime \widehat{U}_{n}^{-1} U_{n} B_{n}\right|=0 \tag{17.1}
\end{equation*}
$$

[^0]$\mathrm{wp} \rightarrow 1$, using $\left|A_{1} A_{2}\right|=\left|A_{1}\right| \cdot\left|A_{2}\right|$ for any conformable square matrices $A_{1}$ and $A_{2},\left|B_{n}\right|>0,\left|U_{n}\right|>0$ (by the conditions in $\mathcal{F}_{W U}$ in 10.5) because $\Lambda_{*} \subset \Lambda_{2}$ and $\Lambda_{2}$ only contains distributions in $\mathcal{F}_{W U}$ ), $\left|\widehat{U}_{n}^{-1}\right|>0 \mathrm{wp} \rightarrow 1$ (because $\widehat{U}_{n} \rightarrow_{p} h_{81}$ by 10.2, 10.12, 10.17, and Assumption WU(b) and (c) and $h_{81}$ is pd), and $\tau_{r_{1} F_{n}}>0$ for $n$ large (because $n^{1 / 2} \tau_{r_{1} F_{n}} \rightarrow \infty$ for $r_{1} \leq q$ ). (For simplicity, we omit the qualifier $\mathrm{wp} \rightarrow 1$ from some statements below.) Thus, $\left\{\left(n^{1 / 2} \tau_{r_{1} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}: j \leq p\right\}$ solve
\[

$$
\begin{align*}
& \left|\tau_{r_{1} F_{n}}^{-2} B_{n}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n}-\kappa\left(I_{p}+\widehat{A}_{n}\right)\right|=0 \text { or } \\
& \left|\left(I_{p}+\widehat{A}_{n}\right)^{-1} \tau_{r_{1} F_{n}}^{-2} B_{n}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n}-\kappa I_{p}\right|=0, \text { where } \\
& \widehat{A}_{n}=\left[\begin{array}{cc}
\widehat{A}_{1 n} & \widehat{A}_{2 n} \\
\widehat{A}_{2 n}^{\prime} & \widehat{A}_{3 n}
\end{array}\right]:=B_{n}^{\prime} U_{n}^{\prime} \widehat{U}_{n}^{-1} \widehat{U}_{n}^{-1} U_{n} B_{n}-I_{p} \tag{17.2}
\end{align*}
$$
\]

for $\widehat{A}_{1 n} \in R^{r_{1}^{\diamond} \times r_{1}^{\circ}}, \widehat{A}_{2 n} \in R^{r_{1}^{\diamond} \times\left(p-r_{1}^{\circ}\right)}$, and $\widehat{A}_{3 n} \in R^{\left(p-r_{1}^{\circ}\right) \times\left(p-r_{1}^{\circ}\right)}$ and the second line is obtained by multiplying the first line by $\left|\left(I_{p}+\widehat{A}_{n}\right)^{-1}\right|$.

We have

$$
\begin{align*}
& \tau_{r_{1} F_{n}}^{-1} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n} \\
= & \tau_{r_{1} F_{n}}^{-1}\left(\widehat{W}_{n} W_{n}^{-1}\right) W_{n} D_{n} U_{n} B_{n}-\left(n^{1 / 2} \tau_{r_{1} F_{n}}\right)^{-1} \widehat{W}_{n} n^{1 / 2}\left(\widehat{D}_{n}-D_{n}\right) U_{n} B_{n} \\
= & \tau_{r_{1} F_{n}}^{-1}\left(\widehat{W}_{n} W_{n}^{-1}\right) C_{n} \Upsilon_{n}+O_{p}\left(\left(n^{1 / 2} \tau_{r_{1} F_{n}}\right)^{-1}\right)  \tag{17.3}\\
= & \left(I_{k}+o_{p}(1)\right) C_{n}\left[\begin{array}{cc}
h_{6, r_{1}^{\diamond}}^{\diamond}+o(1) & 0^{r_{1}^{\diamond} \times\left(p-r_{1}^{\diamond}\right)} \\
0^{\left(p-r_{1}^{\diamond}\right) \times r_{1}^{\diamond}} & O\left(\tau_{r_{2} F_{n}} / \tau_{r_{1} F_{n}}\right)^{\left(p-r_{1}^{\diamond}\right) \times\left(p-r_{1}^{\diamond}\right)} \\
0^{(k-p) \times r_{1}^{\diamond}} & 0^{(k-p) \times\left(p-r_{1}^{\diamond}\right)}
\end{array}\right]+O_{p}\left(\left(n^{1 / 2} \tau_{r_{1} F_{n}}\right)^{-1}\right) \\
\rightarrow & p h_{3}\left[\begin{array}{cc}
h_{6, r_{1}^{\diamond}}^{\diamond} & 0^{0_{1}^{\diamond} \times\left(p-r_{1}^{\diamond}\right)} \\
0^{\left(k-r_{1}^{\diamond}\right) \times r_{1}^{\diamond}} & 0^{\left(k-r_{1}^{\diamond}\right) \times\left(p-r_{1}^{\diamond}\right)}
\end{array}\right], \text { where } h_{6, r_{1}^{\diamond}}^{\diamond}:=\operatorname{Diag}\left\{1, h_{6,1}, h_{6,1} h_{6,2}, \ldots, \prod_{\ell=1}^{r_{1}^{\diamond-1}} h_{6, \ell}\right\},
\end{align*}
$$

$h_{6, r_{1}^{\diamond}}^{\diamond} \in R^{r_{1}^{\diamond} \times r_{1}^{\diamond}}, h_{6, r_{1}^{\diamond}}^{\diamond}:=1$ when $r_{1}^{\diamond}=1, O\left(\tau_{r_{2} F_{n}} / \tau_{r_{1} F_{n}}\right)^{\left(p-r_{1}^{\diamond}\right) \times\left(p-r_{1}^{\diamond}\right)}$ denotes a diagonal $\left(p-r_{1}^{\diamond}\right) \times(p-$ $r_{1}^{\diamond}$ ) matrix whose diagonal elements are $O\left(\tau_{r_{2} F_{n}} / \tau_{r_{1} F_{n}}\right)$, the second equality uses 16.4, $\widehat{W}_{n} \rightarrow_{p} h_{71}$ (by Assumption WU(a) and (c)), $\left\|h_{71}\right\|=\left\|\lim W_{n}\right\|<\infty$ (by the conditions in $\mathcal{F}_{W U}$ defined in 10.5), $n^{1 / 2}\left(\widehat{D}_{n}-D_{n}\right)=O_{p}(1)$ (by Lemma 10.2), $U_{n}=O(1)$ (by the conditions in $\mathcal{F}_{W U}$ ), and $B_{n}=$ $O(1)$ (because $B_{n}$ is orthogonal), the third equality uses $\widehat{W}_{n} W_{n}^{-1} \rightarrow_{p} I_{k}$ (because $\widehat{W}_{n} \rightarrow_{p} h_{71}, h_{71}:=$ $\lim W_{n}$, and $h_{71}$ is pd by the conditions in $\left.\mathcal{F}_{W U}\right), \tau_{j F_{n}} / \tau_{r_{1} F_{n}}=\prod_{\ell=1}^{j-1}\left(\tau_{(\ell+1) F_{n}} / \tau_{\ell F_{n}}\right)=\prod_{\ell=1}^{j-1} h_{6, \ell}+o(1)$ for $j=2, \ldots, r_{1}^{\diamond}$, and $\tau_{j F_{n}} / \tau_{r_{1} F_{n}}=O\left(\tau_{r_{2} F_{n}} / \tau_{r_{1} F_{n}}\right)$ for $j=r_{2}, \ldots, p$ (because $\left\{\tau_{j F_{n}}: j \leq p\right\}$ are nonincreasing in $j$ ), and the convergence uses $C_{n} \rightarrow h_{3}, \tau_{r_{2} F_{n}} / \tau_{r_{1} F_{n}} \rightarrow 0$ (by the definition of $r_{2}$ ), and $n^{1 / 2} \tau_{r_{1} F_{n}} \rightarrow \infty$ (by 10.16 because $\left.r_{1} \leq q\right)$. Note that, for matrices that are written as $O(\cdot)$,
we sometimes provide the dimensions of the matrix as superscripts for clarity, and sometimes we do not provide the dimensions for simplicity.

Equation (17.3) yields

$$
\begin{align*}
\tau_{r_{1} F_{n}}^{-2} B_{n}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n} & \rightarrow p\left[\begin{array}{cc}
h_{6, r_{1}^{\circ}}^{\diamond} & 0^{r_{1}^{\circ} \times\left(p-r_{1}^{\circ}\right)} \\
0^{\left(k-r_{1}^{\circ}\right) \times r_{1}^{\circ}} & 0^{\left(k-r_{1}^{\circ}\right) \times\left(p-r_{1}^{\circ}\right)}
\end{array}\right]^{\prime} h_{3}^{\prime} h_{3}\left[\begin{array}{cc}
h_{6, r_{1}^{\circ}}^{\diamond} & 0^{r_{1}^{\circ} \times\left(p-r_{1}^{\circ}\right)} \\
0^{\left(k-r_{1}^{\circ}\right) \times r_{1}^{\circ}} & 0^{\left(k-r_{1}^{\circ}\right) \times\left(p-r_{1}^{\circ}\right)}
\end{array}\right] \\
& =\left[\begin{array}{cc}
h_{6, r_{1}^{\circ}}^{\diamond 2} & 0^{r_{1}^{\circ} \times\left(p-r_{1}^{\circ}\right)} \\
0^{\left(p-r_{1}^{\circ}\right) \times r_{1}^{\circ}} & 0^{\left(p-r_{1}^{\circ}\right) \times\left(p-r_{1}^{\circ}\right)}
\end{array}\right], \tag{17.4}
\end{align*}
$$

where the equality holds because $h_{3}^{\prime} h_{3}=\lim C_{n}^{\prime} C_{n}=I_{k}$ using 10.7.
In addition, we have

$$
\begin{equation*}
\widehat{A}_{n}:=B_{n}^{\prime} U_{n}^{\prime} \widehat{U}_{n}^{-1} / \widehat{U}_{n}^{-1} U_{n} B_{n}-I_{p} \rightarrow_{p} 0^{p \times p} \tag{17.5}
\end{equation*}
$$

using $\widehat{U}_{n}^{-1} U_{n} \rightarrow_{p} I_{p}$ (because $\widehat{U}_{n} \rightarrow_{p} h_{81}$ by Assumption $\mathrm{WU}(\mathrm{b})$ and (c), $h_{81}:=\lim U_{n}$, and $h_{81}$ is pd by the conditions in $\mathcal{F}_{W U}$ ), $B_{n} \rightarrow h_{2}$, and $h_{2}^{\prime} h_{2}=I_{p}$ (because $B_{n}$ is orthogonal for all $n \geq 1$ ).

The ordered vector of eigenvalues of a matrix is a continuous function of the matrix by Elsner's Theorem, see Stewart (2001, Thm. 3.1, pp. 37-38). Hence, by the second line of (17.2), 17.4, 17.5), and Slutsky's Theorem, the largest $r_{1}^{\diamond}$ eigenvalues of $\tau_{r_{1} F_{n}}^{-2} B_{n}^{\prime} \widehat{U}_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} \widehat{U}_{n} B_{n}$ (i.e., $\left\{\left(n^{1 / 2} \tau_{r_{1} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}: j \leq r_{1}^{\diamond}\right\}$ by the definition of $\left.\widehat{\kappa}_{j n}\right)$, satisfy

$$
\begin{align*}
& \left(\left(n^{1 / 2} \tau_{r_{1} F_{n}}\right)^{-2} \widehat{\kappa}_{1 n}, \ldots,\left(n^{1 / 2} \tau_{r_{1} F_{n}}\right)^{-2} \widehat{\kappa}_{r_{1}^{\circ} n}\right) \rightarrow_{p}\left(1, h_{6,1}^{2}, h_{6,1}^{2} h_{6,2}^{2}, \ldots, \prod_{\ell=1}^{r_{1}^{\diamond}-1} h_{6, \ell}^{2}\right) \text { and so } \\
& \widehat{\kappa}_{j n} \rightarrow p \infty \forall j=1, \ldots, r_{1}^{\diamond} \tag{17.6}
\end{align*}
$$

because $n^{1 / 2} \tau_{r_{1} F_{n}} \rightarrow \infty$ (by 10.16) since $\left.r_{1} \leq q\right)$ and $h_{6, \ell}>0$ for all $\ell \in\left\{1, \ldots, r_{1}^{\diamond}-1\right\}$ (as noted above). By the same argument, the smallest $p-r_{1}^{\diamond}$ eigenvalues of $\tau_{r_{1} F_{n}}^{-2} B_{n}^{\prime} \widehat{U}_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} \widehat{U}_{n} B_{n}$, i.e., $\left\{\left(n^{1 / 2} \tau_{r_{1} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}: j=r_{1}^{\diamond}+1, \ldots, p\right\}$, satisfy

$$
\begin{equation*}
\left(n^{1 / 2} \tau_{r_{1} F_{n}}\right)^{-2} \widehat{\kappa}_{j n} \rightarrow_{p} 0 \forall j=r_{1}^{\diamond}+1, \ldots, p \tag{17.7}
\end{equation*}
$$

If $G_{h}=1$, 17.6) proves part (a) of the lemma and 17.7) proves part (b) of the lemma (because in this case $r_{1}^{\diamond}=q$ and $\tau_{r_{1} F_{n}} / \tau_{\ell F_{n}}=O(1)$ for all $\ell \leq q$ by the definitions of $q$ and $\left.G_{h}\right)$. Hence, from here on, we assume that $G_{h} \geq 2$.

Next, define $B_{n, j_{1}, j_{2}}$ to be the $p \times\left(j_{2}-j_{1}\right)$ matrix that consists of the $j_{1}+1, \ldots, j_{2}$ columns of $B_{n}$ for $0 \leq j_{1}<j_{2} \leq p$. Note that the difference between the two subscripts $j_{1}$ and $j_{2}$ equals the number of columns of $B_{n, j_{1}, j_{2}}$, which is useful for keeping track of the dimensions of the $B_{n, j_{1}, j_{2}}$
matrices that appear below. By definition, $B_{n}=\left(B_{n, 0, r_{1}^{\diamond}}, B_{n, r_{1}^{\diamond}, p}\right)$.
By 17.3 (excluding the convergence part) applied once with $B_{n, r_{1}^{\diamond}, p}$ in place of $B_{n}$ as the farright multiplicand and applied a second time with $B_{n, 0, r_{1}^{\circ}}$ in place of $B_{n}$ as the far-right multiplicand, we have

$$
\left.\begin{array}{rl}
\varrho_{n}:= & \tau_{r_{1} F_{n}}^{-2} B_{n, 0, r_{1}^{\diamond}}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n, r_{1}^{\diamond}, p} \\
= & {\left[\begin{array}{c}
h_{6, r_{1}^{\diamond}}^{\diamond}+o(1) \\
0^{\left(k-r_{1}^{\diamond}\right) \times r_{1}^{\diamond}}
\end{array}\right]^{\prime} C_{n}^{\prime}\left(I_{k}+o_{p}(1)\right) C_{n}\left[\begin{array}{c}
0_{1}^{r_{1}^{\diamond} \times\left(p-r_{1}^{\diamond}\right)} \\
\\
\\
\\
+
\end{array} O_{p}\left(\left(n^{1 / 2} \tau_{r_{2} F_{n}}\right)^{-1}\right)\right.} \\
= & \left.\tau_{p}\left(\tau_{r_{1} F_{n} F_{n}}\right)^{\left(k-r_{1}^{\diamond}\right) \times\left(p-r_{1}^{\diamond}\right)}\right]
\end{array}\right]
$$

where the last equality holds because (i) $C_{n}^{\prime}\left(I_{k}+o_{p}(1)\right) C_{n}=I_{k}+o_{p}(1)$, (ii) when $I_{k}$ appears in place of $C_{n}^{\prime}\left(I_{k}+o_{p}(1)\right) C_{n}$, the first summand on the left-hand side (lhs) of the last equality equals $0^{r_{1}^{\diamond} \times\left(p-r_{1}^{\diamond}\right)}$, and (iii) when $o_{p}(1)$ appears in place of $C_{n}^{\prime}\left(I_{k}+o_{p}(1)\right) C_{n}$, the first summand on the lhs of the last equality equals an $r_{1}^{\diamond} \times\left(p-r_{1}^{\diamond}\right)$ matrix with elements that are $o_{p}\left(\tau_{r_{2} F_{n}} / \tau_{r_{1} F_{n}}\right)$.

Define

$$
\begin{align*}
& \widehat{\xi}_{1 n}(\kappa):=\tau_{r_{1} F_{n}}^{-2} B_{n, 0, r_{1}^{\diamond}}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n, 0, r_{1}^{\diamond}}-\kappa\left(I_{r_{1}^{\diamond}}+\widehat{A}_{1 n}\right) \in R^{r_{1}^{\diamond} \times r_{1}^{\diamond}}, \\
& \widehat{\xi}_{2 n}(\kappa):=\varrho_{n}-\kappa \widehat{A}_{2 n} \in R^{r \gtrdot} \times\left(p-r_{1}^{\diamond}\right)  \tag{17.9}\\
& \widehat{\xi}_{3 n}(\kappa):=\tau_{r_{1} F_{n}}^{-2} B_{n, r_{1}^{\diamond}, p}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n, r_{1}^{\diamond}, p}-\kappa\left(I_{p-r_{1}^{\diamond}}+\widehat{A}_{3 n}\right) \in R^{\left(p-r_{1}^{\diamond}\right) \times\left(p-r_{1}^{\diamond}\right)} .
\end{align*}
$$

As in the first line of $17.2,\left\{\left(n^{1 / 2} \tau_{r_{1} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}: j \leq p\right\}$ solve

$$
\begin{align*}
0= & \left|\tau_{r_{1} F_{n}}^{-2} B_{n}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n}-\kappa\left(I_{p}+\widehat{A}_{n}\right)\right| \\
= & \left|\left[\begin{array}{cc}
\widehat{\xi}_{1 n}(\kappa) & \widehat{\xi}_{2 n}(\kappa) \\
\widehat{\xi}_{2 n}(\kappa)^{\prime} & \widehat{\xi}_{3 n}(\kappa)
\end{array}\right]\right| \\
= & \left|\widehat{\xi}_{1 n}(\kappa)\right| \cdot\left|\widehat{\xi}_{3 n}(\kappa)-\widehat{\xi}_{2 n}(\kappa)^{\prime} \widehat{\xi}_{1 n}^{-1}(\kappa) \widehat{\xi}_{2 n}(\kappa)\right| \\
= & \left|\widehat{\xi}_{1 n}(\kappa)\right| \cdot \mid \tau_{r_{1} F_{n}}^{-2} B_{n, r_{1}^{\circ}, p}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n, r_{1}^{\diamond}, p}-\varrho_{n}^{\prime} \widehat{\xi}_{1 n}^{-1}(\kappa) \varrho_{n} \\
& -\kappa\left(I_{p-r_{1}^{\circ}}+\widehat{A}_{3 n}-\widehat{A}_{2 n}^{\prime} \widehat{\xi}_{1 n}^{-1}(\kappa) \varrho_{n}-\varrho_{n}^{\prime} \widehat{\xi}_{1 n}^{-1}(\kappa) \widehat{A}_{2 n}+\kappa \widehat{A}_{2 n}^{\prime} \widehat{\xi}_{1 n}^{-1}(\kappa) \widehat{A}_{2 n}\right) \mid, \tag{17.10}
\end{align*}
$$

where the third equality uses the standard formula for the determinant of a partitioned matrix (i.e., the determinant of $\xi=\left[\begin{array}{cc}\xi_{1} & \xi_{2} \\ \xi_{2}^{\prime} & \xi_{3}\end{array}\right]$ equals $|\xi|=\left|\xi_{1}\right| \cdot\left|\xi_{3}-\xi_{2}^{\prime} \xi_{1}^{-1} \xi_{2}\right|$ provided $\xi_{1}$ is nonsingular, e.g., see Rao (1973, p. 32)) and the result given in 17.11 below, which shows that $\widehat{\xi}_{1 n}(\kappa)$ is nonsingular
$\omega p \rightarrow 1$ for $\kappa$ equal to any solution $\left(n^{1 / 2} \tau_{r_{1} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}$ to the first equality in 17.10 for $j \leq p$, and the last equality holds by algebra.

Now we show that, for $j=r_{1}^{\diamond}+1, \ldots, p,\left(n^{1 / 2} \tau_{r_{1} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}$ cannot solve the determinantal equation $\left|\widehat{\xi}_{1 n}(\kappa)\right|=0, \mathrm{wp} \rightarrow 1$, where this determinant is the first multiplicand on the rhs of 17.10 . This implies that $\left\{\left(n^{1 / 2} \tau_{r_{1} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}: j=r_{1}^{\diamond}+1, \ldots, p\right\}$ must solve the determinantal equation based on the second multiplicand on the rhs of $17.10 \mathrm{wp} \rightarrow 1$. For $j=r_{1}^{\diamond}+1, \ldots, p$, we have

$$
\begin{align*}
\widetilde{\xi}_{j 1 n} & :=\widehat{\xi}_{1 n}\left(\left(n^{1 / 2} \tau_{r_{1} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}\right) \\
& =\tau_{r_{1} F_{n}}^{-2} B_{n, 0, r_{1}^{r}}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n, 0, r_{1}^{\diamond}}-\left(n^{1 / 2} \tau_{r_{1} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}\left(I_{r_{1}^{\diamond}}+\widehat{A}_{1 n}\right) \\
& =h_{6, r_{1}^{\diamond}}^{\diamond 2}+o_{p}(1)-o_{p}(1)\left(I_{r_{1}^{\diamond}}+o_{p}(1)\right) \\
& =h_{6, r_{1}^{\diamond}}^{\diamond 2}+o_{p}(1), \tag{17.11}
\end{align*}
$$

where the second last equality holds by $\sqrt{17.4}, \sqrt{17.5}$, and 17.7 . Equation 17.11 and $\lambda_{\min }\left(h_{6, r_{1}^{\circ}}^{\diamond 2}\right)>$ 0 (which follows from the definition of $h_{6, r_{1}^{\circ}}^{\diamond}$ in 17.3) and the fact that $h_{6, \ell}>0$ for all $\ell \in$ $\left.\left\{1, \ldots, r_{1}^{\diamond}-1\right\}\right)$ establish the result stated in the first sentence of this paragraph.

For $j=r_{1}^{\diamond}+1, \ldots, p$, plugging $\left(n^{1 / 2} \tau_{r_{1} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}$ into the second multiplicand on the rhs of 17.10 ) gives

$$
\begin{align*}
& 0=\mid \tau_{r_{1} F_{n}}^{-2} B_{n, r_{1}^{\circ}, p}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n, r_{1}^{\diamond}, p}+o_{p}\left(\left(\tau_{r_{2} F_{n}} / \tau_{r_{1} F_{n}}\right)^{2}\right)+O_{p}\left(\left(n^{1 / 2} \tau_{r_{1} F_{n}}\right)^{-2}\right) \\
&-\left(n^{1 / 2} \tau_{r_{1} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}\left(I_{p-r_{1}^{\diamond}}+\widehat{A}_{j 2 n}\right) \mid, \text { where }  \tag{17.12}\\
& \widehat{A}_{j 2 n}:=\widehat{A}_{3 n}-\widehat{A}_{2 n}^{\prime} \widetilde{\xi}_{j 1 n}^{-1} \varrho_{n}-\varrho_{n}^{\prime} \widetilde{\xi}_{j 1 n}^{-1} \widehat{A}_{2 n}+\left(n^{1 / 2} \tau_{r_{1} F_{n}}\right)^{-2} \widehat{\kappa}_{j n} \widehat{A}_{2 n}^{\prime} \widetilde{\xi}_{j 1 n}^{-1} \widehat{A}_{2 n} \in R^{\left(p-r_{1}^{\diamond}\right) \times\left(p-r_{1}^{\circ}\right)}
\end{align*}
$$

using (17.8) and 17.11. Multiplying 17.12 by $\tau_{r_{1} F_{n}}^{2} / \tau_{r_{2} F_{n}}^{2}$ gives

$$
\begin{equation*}
0=\left|\tau_{r_{2} F_{n}}^{-2} B_{n, r_{1}^{\diamond}, p}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n, r_{1}^{\diamond}, p}+o_{p}(1)-\left(n^{1 / 2} \tau_{r_{2} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}\left(I_{p-r_{1}^{\diamond}}+\widehat{A}_{j 2 n}\right)\right| \tag{17.13}
\end{equation*}
$$

using $O_{p}\left(\left(n^{1 / 2} \tau_{r_{2} F_{n}}\right)^{-2}\right)=o_{p}(1)$ (because $r_{2} \leq q$ by the definition of $r_{2}$ and $n^{1 / 2} \tau_{j F_{n}} \rightarrow \infty$ for all $j \leq q$ by the definition of $q$ in 10.16).

Thus, $\left\{\left(n^{1 / 2} \tau_{r_{2} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}: j=r_{1}^{\diamond}+1, \ldots, p\right\}$ solve

$$
\begin{equation*}
0=\left|\tau_{r_{2} F_{n}}^{-2} B_{n, r_{1}^{\diamond}, p}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n, r_{1}^{\diamond}, p}+o_{p}(1)-\kappa\left(I_{p-r_{1}^{\diamond}}+\widehat{A}_{j 2 n}\right)\right| . \tag{17.14}
\end{equation*}
$$

For $j=r_{1}^{\diamond}+1, \ldots, p$, we have

$$
\begin{equation*}
\widehat{A}_{j 2 n}=o_{p}(1) \tag{17.15}
\end{equation*}
$$

because $\widehat{A}_{2 n}=o_{p}(1)$ and $\widehat{A}_{3 n}=o_{p}(1)($ by 17.5$), \widetilde{\xi}_{j 1 n}^{-1}=O_{p}(1)($ by 17.11$), \varrho_{n}=o_{p}(1) \quad$ by 17.8 ) since $\tau_{r_{2} F_{n}} \leq \tau_{r_{1} F_{n}}$ and $\left.n^{1 / 2} \tau_{r_{1} F_{n}} \rightarrow \infty\right)$, and $\left(n^{1 / 2} \tau_{r_{1} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}=o_{p}(1)$ for $j=r_{1}^{\diamond}+1, \ldots, p$ (by (17.7)).

Now, we repeat the argument from $(17.2)$ to $(17.15$ with the expression in 17.14 replacing that in the first line of 17.2 , with 17.15 replacing 17.5 , and with $j=r_{2}^{\diamond}+1, \ldots, p, \widehat{A}_{j 2 n}, B_{n, p-r_{1}^{\diamond}}, \tau_{r_{2} F_{n}}$, $\tau_{r_{3} F_{n}}, r_{2}^{\diamond}-r_{1}^{\diamond}, p-r_{2}^{\diamond}$, and $h_{6, r_{2}^{\diamond}}^{\diamond}=\operatorname{Diag}\left\{1, h_{6, r_{1}^{\diamond}+1}, h_{6, r_{1}^{\diamond}+1} h_{6, r_{1}^{\diamond}+2}, \ldots, \prod_{\ell=r_{1}^{\diamond+1}}^{r_{2}^{\diamond}-1} h_{6, \ell}\right\} \in R^{\left(r_{2}^{\diamond}-r_{1}^{\diamond}\right) \times\left(r_{2}^{\diamond}-r_{1}^{\diamond}\right)}$ in place of $j=r_{1}^{\diamond}+1, \ldots, p, \widehat{A}_{n}, B_{n}, \tau_{r_{1} F_{n}}, \tau_{r_{2} F_{n}}, r_{1}^{\diamond}, p-r_{1}^{\diamond}$, and $h_{6, r_{1}^{\diamond}}^{\diamond}$, respectively. (The fact that $\widehat{A}_{j 2 n}$ depends on $j$, whereas $\widehat{A}_{n}$ does not, does not affect the argument.) In addition, $B_{n, 0, r_{1}^{\diamond}}$ and $B_{n, r_{1}^{\diamond}, p}$ in $17.8-17.10$ are replaced by the matrices $B_{n, r_{1}^{\diamond}, r_{2}^{\diamond}}$ and $B_{n, r_{2}^{\diamond}, p}$ (which consist of the $r_{1}^{\diamond}+1, \ldots, r_{2}^{\diamond}$ columns of $B_{n}$ and the last $p-r_{2}^{\diamond}$ columns of $B_{n}$, respectively.) This argument gives the analogues of 17.6 and 17.7 ), which are

$$
\begin{equation*}
\widehat{\kappa}_{j n} \rightarrow_{p} \infty \forall j=r_{2}, \ldots, r_{2}^{\diamond} \text { and }\left(n^{1 / 2} \tau_{r_{2} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}=o_{p}(1) \forall j=r_{2}^{\diamond}+1, \ldots, p \tag{17.16}
\end{equation*}
$$

In addition, the analogue of 17.14 is the same as 17.14 but with $\widehat{A}_{j 3 n}$ in place of $\widehat{A}_{j 2 n}$, where $\widehat{A}_{j 3 n}$ is defined just as $\widehat{A}_{j 2 n}$ is defined in 17.12 but with $\widehat{A}_{2 j 2 n}$ and $\widehat{A}_{3 j 2 n}$ in place of $\widehat{A}_{2 n}$ and $\widehat{A}_{3 n}$, respectively, where

$$
\widehat{A}_{j 2 n}=\left[\begin{array}{ll}
\widehat{A}_{1 j 2 n} & \widehat{A}_{2 j 2 n}  \tag{17.17}\\
\widehat{A}_{2 j 2 n}^{\prime} & \widehat{A}_{3 j 2 n}
\end{array}\right]
$$

for $\widehat{A}_{1 j 2 n} \in R^{r_{2}^{\diamond} \times r_{2}^{\diamond}}, \widehat{A}_{2 j 2 n} \in R^{r_{2}^{\diamond} \times\left(p-r_{1}^{\diamond}-r_{2}^{\diamond}\right)}$, and $\widehat{A}_{3 j 2 n} \in R^{\left(p-r_{1}^{\left.\diamond-r_{2}^{\diamond}\right) \times\left(p-r_{1}^{\diamond}-r_{2}^{\diamond}\right)} \text {. }\right.}$
Repeating the argument $G_{h}-2$ more times yields

$$
\begin{equation*}
\widehat{\kappa}_{j n} \rightarrow_{p} \infty \forall j=1, \ldots, r_{G_{h}}^{\diamond} \text { and }\left(n^{1 / 2} \tau_{r_{g} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}=o_{p}(1) \forall j=r_{g}^{\diamond}+1, \ldots, p, \forall g=1, \ldots, G_{h} . \tag{17.18}
\end{equation*}
$$

A formal proof of this "repetition of the argument $G_{h}-2$ more times" is given below using induction. Because $r_{G_{h}}^{\diamond}=q$, the first result in 17.18 proves part (a) of the lemma.

The second result in 17.18 with $g=G_{h}$ implies: for all $j=q+1, \ldots, p$,

$$
\begin{equation*}
\left(n^{1 / 2} \tau_{r_{G_{h}} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}=o_{p}(1) \tag{17.19}
\end{equation*}
$$

because $r_{G_{h}}^{\diamond}=q$. Either $r_{G_{h}}=r_{G_{h}}^{\diamond}=q$ or $r_{G_{h}}<r_{G_{h}}^{\diamond}=q$. In the former case, $\left(n^{1 / 2} \tau_{q F_{n}}\right)^{-2} \widehat{\kappa}_{j n}=$
$o_{p}(1)$ for $j=q+1, \ldots, p$ by 17.19$)$. In the latter case, we have

$$
\begin{equation*}
\lim \frac{\tau_{q F_{n}}}{\tau_{r_{G_{h}} F_{n}}}=\lim \frac{\tau_{r_{G_{h}}^{\diamond} F_{n}}^{\tau_{r_{G_{h}} F_{n}}}=\prod_{j=r_{G_{h}}}^{r_{G_{h}}^{\diamond}-1} h_{6, j}>0, ~, ~, ~ . ~}{} \tag{17.20}
\end{equation*}
$$

where the inequality holds because $h_{6, \ell}>0$ for all $\ell \in\left\{r_{G_{h}}, \ldots, r_{G_{h}}^{\diamond}-1\right\}$, as noted at the beginning of the proof. Hence, in this case too, $\left(n^{1 / 2} \tau_{q F_{n}}\right)^{-2} \widehat{\kappa}_{j n}=o_{p}(1)$ for $j=q+1, \ldots, p$ by 17.19 and 17.20). Because $\tau_{\ell F_{n}} \geq \tau_{q F_{n}}$ for all $\ell \leq q$, this establishes part (b) of the lemma.

Now we establish by induction the results given in 17.18 that are obtained heuristically by "repeating the argument $G_{h}-2$ more times." The induction proof shows that subtleties arise when establishing the asymptotic negligibility of certain terms.

Let $o_{g p}$ denote a symmetric $\left(p-r_{g-1}^{\diamond}\right) \times\left(p-r_{g-1}^{\diamond}\right)$ matrix whose $(\ell, m)$ element for $\ell, m=$ $1, \ldots, p-r_{g-1}^{\diamond}$ is $o_{p}\left(\tau_{\left(r_{g-1}^{\diamond}+\ell\right) F_{n}} \tau_{\left(r_{g-1}^{\diamond}+m\right) F_{n}} / \tau_{r_{g} F_{n}}^{2}\right)+O_{p}\left(\left(n^{1 / 2} \tau_{r_{g} F_{n}}\right)^{-1}\right)$. Note that $o_{g p}=o_{p}(1)$ because $r_{g-1}^{\diamond}+\ell \geq r_{g}$ for $\ell \geq 1$ (since $\tau_{j F_{n}}$ are nonincreasing in $j$ ) and $n^{1 / 2} \tau_{r_{g} F_{n}} \rightarrow \infty$ for $g=1, \ldots, G_{h}$.

We now show by induction over $g=1, \ldots, G_{h}$ that wp $\rightarrow 1\left\{\left(n^{1 / 2} \tau_{r_{g} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}: j=r_{g-1}^{\diamond}+1, \ldots, p\right\}$ solve

$$
\begin{equation*}
\left|\tau_{r_{g} F_{n}}^{-2} B_{n, r_{g-1}^{\diamond}, p}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n, r_{g-1}^{\diamond}, p}+o_{g p}-\kappa\left(I_{p-r_{g-1}^{\diamond}}+\widehat{A}_{j g n}\right)\right|=0 \tag{17.21}
\end{equation*}
$$

for some $\left(p-r_{g-1}^{\diamond}\right) \times\left(p-r_{g-1}^{\diamond}\right)$ symmetric matrices $\widehat{A}_{j g n}=o_{p}(1)$ and $o_{g p}$ (where the matrices that are $o_{g p}$ may depend on $j$ ).

The initiation step of the induction proof holds because 17.21) holds with $g=1$ by the first line of 17.2 with $\widehat{A}_{j g n}:=\widehat{A}_{n}$ and $o_{g p}=0$ for $g=1$ (and using the fact that, for $g=1, r_{g-1}^{\diamond}=r_{0}^{\diamond}:=0$ and $\left.B_{n, r_{g-1}^{\circ}, p}=B_{n, 0, p}=B_{n}\right)$.

For the induction step of the proof, we assume that (17.21) holds for some $g \in\left\{1, \ldots, G_{h}-1\right\}$ and show that it then also holds for $g+1$. By an argument analogous to that in 17.3 , we have

$$
\begin{align*}
& \tau_{r_{g} F_{n}}^{-1} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n, r_{g-1}^{\diamond}, p}=\left(I_{k}+o_{p}(1)\right) C_{n}\left[\begin{array}{c}
0^{r_{g-1}^{\diamond} \times\left(p-r_{g-1}^{\diamond}\right)} \\
\operatorname{Diag}\left\{\tau_{r_{g} F_{n}}, \ldots, \tau_{p F_{n}}\right\} / \tau_{r_{g} F_{n}} \\
0^{(k-p) \times\left(p-r_{g-1}^{\diamond}\right)}
\end{array}\right]+O_{p}\left(\left(n^{1 / 2} \tau_{r_{g} F_{n}}\right)^{-1}\right) \\
& \rightarrow_{p} h_{3}\left(\left[\begin{array}{c}
0^{r_{g-1}^{\diamond} \times\left(r_{g}^{\diamond}-r_{g-1}^{\diamond}\right)} \\
h_{6, r_{g}^{\diamond}}^{\diamond} \\
0^{\left(k-r_{g}^{\diamond}\right) \times\left(r_{g}^{\diamond}-r_{g-1}^{\diamond}\right)}
\end{array}\right], 0^{k \times\left(p-r_{g}^{\diamond}\right)}\right), \text { where } h_{6, r_{g}^{\diamond}}^{\diamond}:=\operatorname{Diag}\left\{1, h_{6, r_{g}}, \ldots, \prod_{j=r_{g-1}^{\diamond}+1}^{r_{g}^{\diamond}-1} h_{6, j}\right\}, \tag{17.22}
\end{align*}
$$

$h_{6, r_{g}^{\diamond}}^{\diamond} \in R^{\left(r_{g}^{\diamond}-r_{g-1}^{\diamond}\right) \times\left(r_{g}^{\diamond}-r_{g-1}^{\diamond}\right)}$, and $h_{6, r_{g}^{\diamond}}^{\diamond}:=1$ when $r_{g}^{\diamond}=1$.

Equation 17.22) and $h_{3}^{\prime} h_{3}=\lim C_{n}^{\prime} C_{n}=I_{k}$ yield

$$
\tau_{r_{g} F_{n}}^{-2} B_{n, r_{g-1}^{\diamond}, p}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n, r_{g-1}^{\diamond}, p} \rightarrow_{p}\left[\begin{array}{cc}
h_{6, r_{g}^{\diamond}}^{\diamond 2} & 0^{\left(r_{g}^{\diamond}-r_{g-1}^{\diamond}\right) \times\left(p-r_{g}^{\diamond}\right)}  \tag{17.23}\\
0^{\left(p-r_{g}^{\diamond}\right) \times\left(r_{g}^{\diamond}-r_{g-1}^{\diamond}\right)} & 0^{\left(p-r_{g}^{\diamond}\right) \times\left(p-r_{g}^{\diamond}\right)}
\end{array}\right] .
$$

By 17.21 and $o_{g p}=o_{p}(1)$, we have $\mathrm{wp} \rightarrow 1\left\{\left(n^{1 / 2} \tau_{r_{g} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}: j=r_{g-1}^{\diamond}+1, \ldots, p\right\}$ solve $\left|\left(I_{p-r_{g-1}^{\diamond}}+\widehat{A}_{j g n}\right)^{-1} \tau_{r_{g} F_{n}}^{-2} B_{n, r_{g-1}^{\circ}, p}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n, r_{g-1}^{\diamond}, p}+o_{p}(1)-\kappa I_{p-r_{g-1}^{\diamond}}\right|=0$. Hence, by 17.23), $\widehat{A}_{j g n}=o_{p}(1)$ (which holds by the induction assumption), and the same argument as used to establish 17.6 and 17.7), we obtain

$$
\begin{equation*}
\widehat{\kappa}_{j n} \rightarrow_{p} \infty \forall j=r_{g-1}^{\diamond}+1, \ldots, r_{g}^{\diamond} \text { and }\left(n^{1 / 2} \tau_{r_{g} F_{n}}\right)^{-2} \widehat{\kappa}_{j n} \rightarrow_{p} 0 \forall j=r_{g}^{\diamond}+1, \ldots, p \tag{17.24}
\end{equation*}
$$

Let $o_{g p}^{*}$ denote an $\left(r_{g}^{\diamond}-r_{g-1}^{\diamond}\right) \times\left(p-r_{g}^{\diamond}\right)$ matrix whose elements in column $j$ for $j=1, \ldots, p-r_{g}^{\diamond}$ are $o_{p}\left(\tau_{\left(r_{g}^{\diamond}+j\right) F_{n}} / \tau_{r_{g} F_{n}}\right)+O_{p}\left(\left(n^{1 / 2} \tau_{r_{g} F_{n}}\right)^{-1}\right)$. Note that $o_{g p}^{*}=o_{p}(1)$.

By 17.22 applied once with $B_{n, r_{g}^{\diamond}, p}$ in place of $B_{n, r_{g-1}^{\diamond}, p}$ as the far-right multiplicand and applied a second time with $B_{n, r_{g-1}^{\diamond}, r_{g}^{\diamond}}$ in place of $B_{n, r_{g-1}^{\diamond}, p}$ as the far-right multiplicand, we have

$$
\begin{align*}
& \varrho_{g n} \\
:= & \tau_{r_{g} F_{n}}^{-2} B_{n, r_{g-1}^{\diamond}, r_{g}^{\diamond}}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n, r_{g}^{\diamond}, p} \\
= & {\left[\begin{array}{c}
0^{r_{g-1}^{\diamond} \times\left(r_{g}^{\diamond}-r_{g-1}^{\diamond}\right)} \\
\operatorname{Diag}\left\{\tau_{\left(r_{g-1}^{\diamond}+1\right) F_{n}}, \ldots, \tau_{r_{g}^{\diamond} F_{n}}\right\} / \tau_{r_{g} F_{n}} \\
0^{\left(k-r_{g}^{\diamond}\right) \times\left(r_{g}^{\diamond}-r_{g-1}^{\diamond}\right)}
\end{array}\right]^{\prime} C_{n}^{\prime}\left(I_{k}+o_{p}(1)\right) C_{n}\left[\begin{array}{c}
0^{r_{g}^{\diamond} \times\left(p-r_{g}^{\diamond}\right)} \\
\operatorname{Diag}\left\{\tau_{\left(r_{g}^{\circ}+1\right) F_{n}}, \ldots, \tau_{p F_{n}}\right\} / \tau_{r_{g} F_{n}} \\
0^{(k-p) \times\left(p-r_{g}^{\diamond}\right)}
\end{array}\right] } \\
& +O_{p}\left(\left(n^{1 / 2} \tau_{r_{g} F_{n}}\right)^{-1}\right) \\
= & o_{g p}^{*}, \tag{17.25}
\end{align*}
$$

where $\varrho_{g n} \in R^{\left(r_{g}^{\diamond}-r_{g-1}^{\diamond}\right) \times\left(p-r_{g}^{\diamond}\right)}, \operatorname{Diag}\left\{\tau_{\left(r_{g-1}^{\diamond}+1\right) F_{n}}, \ldots, \tau_{r_{g}^{\diamond} F_{n}}\right\} / \tau_{r_{g} F_{n}}=h_{6, r_{g}^{\diamond}}^{\diamond}+o(1)=O(1)$ and the last equality holds because (i) $C_{n}^{\prime}\left(I_{k}+o_{p}(1)\right) C_{n}=I_{k}+o_{p}(1)$, (ii) when $I_{k}$ appears in place of $C_{n}^{\prime}\left(I_{k}+o_{p}(1)\right) C_{n}$, then the contribution from the first summand on the lhs of the last equality in 17.25 equals $0^{\left(r_{g}^{\diamond}-r_{g-1}^{\diamond}\right) \times\left(p-r_{g}^{\diamond}\right)}$, and (iii) when $o_{p}(1)$ appears in place of $C_{n}^{\prime}\left(I_{k}+o_{p}(1)\right) C_{n}$, the contribution from the first summand on the lhs of the last inequality in 17.25 equals an $o_{g p}^{*}$ matrix.

We partition the $\left(p-r_{g-1}^{\diamond}\right) \times\left(p-r_{g-1}^{\diamond}\right)$ matrices $o_{g p}$ and $\widehat{A}_{j g n}$ as follows:

$$
o_{g p}=\left(\begin{array}{cc}
o_{1 g p} & o_{2 g p}  \tag{17.26}\\
o_{2 g p}^{\prime} & o_{3 g p}
\end{array}\right) \text { and } \widehat{A}_{j g n}=\left[\begin{array}{cc}
\widehat{A}_{1 j g n} & \widehat{A}_{2 j g n} \\
\widehat{A}_{2 j g n}^{\prime} & \widehat{A}_{3 j g n}
\end{array}\right],
$$

where $o_{1 g p}, \widehat{A}_{1 j g n} \in R^{\left(r_{g}^{\diamond}-r_{g-1}^{\diamond}\right) \times\left(r_{g}^{\diamond}-r_{g-1}^{\diamond}\right)}$, oo $o_{2 g p}, \widehat{A}_{2 j g n} \in R^{\left(r_{g}^{\diamond}-r_{g-1}^{\diamond}\right) \times\left(p-r_{g}^{\diamond}\right)}$, and $o_{3 g p}, \widehat{A}_{3 j g n}$ $\in R^{\left(p-r_{g}^{\diamond}\right) \times\left(p-r_{g}^{\diamond}\right)}$, for $j=r_{g-1}^{\diamond}+1, \ldots, p$ and $g=1, \ldots, G_{h}$. Define

$$
\begin{align*}
& \widehat{\xi}_{1 j g n}(\kappa):=\tau_{r g}^{-2} B_{n, r_{g-1}^{\diamond}, r_{g}^{\diamond}}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n, r_{g-1}^{\diamond}, r_{g}^{\diamond}}+o_{1 g p}-\kappa\left(I_{r_{g}^{\diamond}-r_{g-1}^{\diamond}}+\widehat{A}_{1 j g n}\right), \\
& \widehat{\xi}_{2 j g n}(\kappa):=\varrho_{g n}+o_{2 g p}-\kappa \widehat{A}_{2 j g n}, \text { and }  \tag{17.27}\\
& \widehat{\xi}_{3 j g n}(\kappa):=\tau_{r_{g} F_{n}}^{-2} B_{n, r_{g}^{\diamond}, p}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n, r_{g}^{\diamond}, p}+o_{3 g p}-\kappa\left(I_{p-r_{g}^{\diamond}}+\widehat{A}_{3 j g n}\right),
\end{align*}
$$

where $\widehat{\xi}_{1 j g n}(\kappa), \widehat{\xi}_{2 j g n}(\kappa)$, and $\widehat{\xi}_{3 j g n}(\kappa)$ have the same dimensions as $o_{1 g p}, o_{2 g p}$, and $o_{3 g p}$, respectively.
From 17.21, we have $\mathrm{wp} \rightarrow 1\left\{\left(n^{1 / 2} \tau_{r_{g} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}: j=r_{g-1}^{\diamond}+1, \ldots, p\right\}$ solve

$$
\begin{align*}
0= & \left|\tau_{r_{g} F_{n}}^{-2} B_{n, r_{g-1}^{\circ}, p}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n, r_{g-1}^{\diamond}, p}+o_{g p}-\kappa\left(I_{p-r_{g-1}^{\diamond}}+\widehat{A}_{j g n}\right)\right| \\
= & \left|\widehat{\xi}_{1 j g n}(\kappa)\right| \cdot\left|\widehat{\xi}_{3 j g n}(\kappa)-\widehat{\xi}_{2 j g n}(\kappa)^{\prime} \widehat{\xi}_{1 j g n}^{-1}(\kappa) \widehat{\xi}_{2 j g n}(\kappa)\right| \\
= & \left|\widehat{\xi}_{1 j g n}(\kappa)\right| \cdot \mid \tau_{r_{g} F_{n}}^{-2} B_{n, r_{g}^{\diamond}, p}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n, r_{g}^{\diamond}, p}+o_{3 g p}-\left(\varrho_{g n}+o_{2 g p}\right)^{\prime} \widehat{\xi}_{1 j g n}^{-1}(\kappa)\left(\varrho_{g n}+o_{2 g p}\right) \\
& \quad-\kappa\left[I_{p-r_{g}^{\diamond}}+\widehat{A}_{3 j g n}-\widehat{A}_{2 j g n}^{\prime} \widehat{\xi}_{1 j g n}^{-1}(\kappa)\left(\varrho_{g n}+o_{2 g p}\right)-\left(\varrho_{g n}+o_{2 g p}\right)^{\prime} \widehat{\xi}_{1 j g n}^{-1}(\kappa) \widehat{A}_{2 j g n}\right. \\
& \left.\quad+\kappa \widehat{A}_{2 j g n}^{\prime} \widehat{\xi}_{1 j g n}^{-1}(\kappa) \widehat{A}_{2 j g n}\right] \mid, \tag{17.28}
\end{align*}
$$

where the second equality holds by the same argument as for 17.10 and uses the result given in 17.29 below which shows that $\widehat{\xi}_{1 j g n}(\kappa)$ is nonsingular wp $\rightarrow 1$ when $\kappa$ equals $\left(n^{1 / 2} \tau_{r_{g} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}$ for $j=r_{g}^{\diamond}+1, \ldots, p$.

Now we show that, for $j=r_{g}^{\diamond}+1, \ldots, p,\left(n^{1 / 2} \tau_{r_{g} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}$ cannot solve the determinantal equation $\left|\widehat{\xi}_{1 j g n}(\kappa)\right|=0$ for $n$ large, where this determinant is the first multiplicand on the rhs of 17.28 and, hence, it must solve the determinantal equation based on the second multiplicand on the rhs of 17.28 ). For $j=r_{g}^{\diamond}+1, \ldots, p$, we have

$$
\begin{equation*}
\widetilde{\xi}_{1 j g n}:=\widehat{\xi}_{1 j g n}\left(\left(n^{1 / 2} \tau_{r_{g} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}\right)=h_{6, r_{g}^{\diamond}}^{\diamond 2}+o_{p}(1) \tag{17.29}
\end{equation*}
$$

by the same argument as in 17.11, using $o_{1 g p}=o_{p}(1)$ and $\widehat{A}_{1 j g n}=o_{p}(1)$ (which holds by the definition of $\widehat{A}_{1 j g n}$ following 17.21 ). Equation 17.29 and $\lambda_{\min }\left(h_{6, r_{g}^{\diamond}}^{\diamond 2}\right)>0$ establish the result stated in the first sentence of this paragraph.

For $j=r_{g}^{\diamond}+1, \ldots, p$, plugging $\left(n^{1 / 2} \tau_{r_{g} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}$ into the second multiplicand on the rhs of 17.28 )
gives

$$
\begin{align*}
0= & \mid \tau_{r_{g} F_{n}}^{-2} B_{n, r_{g}^{\circ}, p}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n, r_{g}^{\diamond}, p}+o_{3 g p}-\left(\varrho_{g n}+o_{2 g p}\right)^{\prime} \widetilde{\xi}_{1 j g n}^{-1}\left(\varrho_{g n}+o_{2 g p}\right) \\
& -\left(n^{1 / 2} \tau_{r_{g} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}\left(I_{p-r_{g}^{\diamond}}+\widehat{A}_{j(g+1) n}\right) \mid, \text { where } \\
\widehat{A}_{j(g+1) n}:= & \widehat{A}_{3 j g n}-\widehat{A}_{2 j g n}^{\prime} \widetilde{\xi}_{1 j g n}^{-1}\left(\varrho_{g n}+o_{2 g p}\right)-\left(\varrho_{g n}+o_{2 g p}\right)^{\prime} \widetilde{\xi}_{1 j g n}^{-1} \widehat{A}_{2 j g n} \\
& +\left(n^{1 / 2} \tau_{r_{g} F_{n}}\right)^{-2} \widehat{\kappa}_{j n} \widehat{A}_{2 j g n}^{\prime} \widetilde{\xi}_{1 j g n}^{-1} \widehat{A}_{2 j g n} \tag{17.30}
\end{align*}
$$

and $\widehat{A}_{j(g+1) n} \in R^{\left(p-r_{g}^{\circ}\right) \times\left(p-r_{g}^{\circ}\right)}$. The last two summands on the rhs of the first line of 17.30 satisfy

$$
\begin{align*}
& o_{3 g p}-\left(\varrho_{g n}+o_{2 g p}\right)^{\prime} \widetilde{\xi}_{1 j g n}^{-1}\left(\varrho_{g n}+o_{2 g p}\right)=o_{3 g p}-\left(o_{g p}^{*}+o_{2 g p}\right)^{\prime}\left(h_{6, r_{g}^{\circ}}^{\diamond-2}+o_{p}(1)\right)\left(o_{g p}^{*}+o_{2 g p}\right) \\
= & o_{3 g p}-o_{g p}^{* \prime} o_{g p}^{*}=\left(\tau_{r_{g+1} F_{n}}^{2} / \tau_{r_{g} F_{n}}^{2}\right) o_{(g+1) p}, \tag{17.31}
\end{align*}
$$

where (i) the first equality uses 17.25 and 17.29 , (ii) the second equality uses $o_{2 g p}=o_{g p}^{*}$ (which holds because the $(j, m)$ element of $o_{2 g p}$ for $j=1, \ldots, r_{g}^{\diamond}-r_{g-1}^{\diamond}$ and $m=1, \ldots, p-r_{g}^{\diamond}$ is $o_{p}\left(\tau_{\left(r_{g-1}^{\diamond}+j\right) F_{n}}\right.$ $\left.\times \tau_{\left(r_{g}^{\diamond}+m\right) F_{n}} / \tau_{r_{g} F_{n}}^{2}\right)+O_{p}\left(\left(n^{1 / 2} \tau_{r_{g} F_{n}}\right)^{-1}\right)=o_{p}\left(\tau_{\left(r_{g}^{\diamond}+m\right) F_{n}} / \tau_{r_{g} F_{n}}\right)+O_{p}\left(\left(n^{1 / 2} \tau_{r_{g} F_{n}}\right)^{-1}\right)$ since $r_{g-1}^{\diamond}+j \geq$ $r_{g}$ ) and $\left(h_{6, r_{g}^{\circ}}^{\diamond-2}+o_{p}(1)\right) o_{g p}^{*}=o_{g p}^{*}$ (which holds because $h_{6, r_{g}^{\diamond}}^{\diamond}$ is diagonal and $\lambda_{\min }\left(h_{6, r_{g}^{\diamond}}^{\diamond}\right)>0$ ), (iii) the last equality uses the fact that the $(j, m)$ element of $\left(\tau_{r_{g} F_{n}}^{2} / \tau_{r_{g+1} F_{n}}^{2}\right) o_{g p}^{* \prime} o_{g p}^{*}$ for $j, m=1, \ldots, p-r_{g}^{\diamond}$ is the sum of a term that is $o_{p}\left(\tau_{\left(r_{g}^{\diamond}+j\right) F_{n}} \tau_{\left(r_{g}^{\diamond}+m\right) F_{n}} / \tau_{r_{g} F_{n}}^{2}\right)\left(\tau_{r_{g} F_{n}}^{2} / \tau_{r_{g+1} F_{n}}^{2}\right)=o_{p}\left(\tau_{\left(r_{g}^{\diamond}+j\right) F_{n}} \tau_{\left(r_{g}^{\diamond}+m\right) F_{n}}\right.$ $\left./ \tau_{r_{g+1} F_{n}}^{2}\right)$ and a term that is $O_{p}\left(\left(n^{1 / 2} \tau_{r_{g} F_{n}}\right)^{-2}\right)\left(\tau_{r_{g} F_{n}}^{2} / \tau_{r_{g+1} F_{n}}^{2}\right)=O_{p}\left(\left(n^{1 / 2} \tau_{r_{g+1} F_{n}}\right)^{-2}\right)$ and, hence, $\left(\tau_{r_{g} F_{n}}^{2} / \tau_{r_{g+1} F_{n}}^{2}\right) o_{g p}^{* \prime} o_{g p}^{*}$ is $o_{(g+1) p}$ (using the definition of $o_{(g+1) p}$ ), and (iv) the last equality uses the fact that the $(j, m)$ element of $\left(\tau_{r_{g} F_{n}}^{2} / \tau_{r_{g+1} F_{n}}^{2}\right) o_{3 g p}$ for $j, m=1, \ldots, p-r_{g}^{\diamond}$ is $o_{p}\left(\tau_{\left(r_{g}^{\diamond}+j\right) F_{n}} \tau_{\left(r_{g}^{\diamond}+m\right) F_{n}}\right.$ $\left./ \tau_{r_{g} F_{n}}^{2}\right)\left(\tau_{r_{g} F_{n}}^{2} / \tau_{r_{g+1} F_{n}}^{2}\right)+O_{p}\left(\left(n^{1 / 2} \tau_{r_{g} F_{n}}\right)^{-1}\right)\left(\tau_{r_{g} F_{n}}^{2} / \tau_{r_{g+1} F_{n}}^{2}\right)=o_{p}\left(\tau_{\left(r_{g}^{\diamond}+j\right) F_{n}} \tau_{\left(r_{g}^{\diamond}+m\right) F_{n}} / \tau_{r_{g+1} F_{n}}^{2}\right)$ $+O_{p}\left(\left(n^{1 / 2} \tau_{r_{g+1} F_{n}}\right)^{-1}\right)\left(\tau_{r_{g} F_{n}} / \tau_{r_{g+1} F_{n}}\right)$, which again is the same order as the $(j, m)$ element of $o_{(g+1) p}$ (using $\tau_{r_{g} F_{n}} / \tau_{r_{g+1} F_{n}} \leq 1$ ).

The calculations in 17.31 are a key part of the induction proof. The definitions of the terms $o_{g p}$ and $o_{g p}^{*}$ (given preceding 17.21) and 17.25), respectively) are chosen so that the results in (17.31) hold.

For $j=r_{g}^{\diamond}+1, \ldots, p$, we have

$$
\begin{equation*}
\widehat{A}_{j(g+1) n}=o_{p}(1), \tag{17.32}
\end{equation*}
$$

because $\widehat{A}_{2 j g n}=o_{p}(1)$ and $\widehat{A}_{3 j g n}=o_{p}(1)$ by 17.21), $\widetilde{\xi}_{1 j g n}^{-1}=O_{p}(1)($ by 17.29$), \varrho_{g n}+o_{2 g p}=o_{p}(1)$ (by 17.25 ) since $o_{g p}^{*}=o_{p}(1)$ ), and $\left(n^{1 / 2} \tau_{r_{g} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}=o_{p}(1)$ (by 17.24 ).

Inserting 17.31 and 17.32 into 17.30 and multiplying by $\tau_{r_{g} F_{n}}^{2} / \tau_{r_{g+1} F_{n}}^{2}$ gives

$$
\begin{equation*}
0=\left|\tau_{r_{g+1} F_{n}}^{-2} B_{n, r_{g}^{\diamond}, p}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n, r_{g}^{\diamond}, p}+o_{(g+1) p}-\left(n^{1 / 2} \tau_{r_{g+1} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}\left(I_{p-r_{g}^{\diamond}}+\widehat{A}_{j(g+1) n}\right)\right| . \tag{17.33}
\end{equation*}
$$

Thus, $\mathrm{wp} \rightarrow 1,\left\{\left(n^{1 / 2} \tau_{r_{g+1} F_{n}}\right)^{-2} \widehat{\kappa}_{j n}: j=r_{g+1}, \ldots, p\right\}$ solve

$$
\begin{equation*}
0=\left|\tau_{r_{g+1} F_{n}}^{-2} B_{n, r_{g}^{\diamond}, p}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n, r_{g}^{\diamond}, p}+o_{(g+1) p}-\kappa\left(I_{p-r_{g}^{\diamond}}+\widehat{A}_{j(g+1) n}\right)\right| . \tag{17.34}
\end{equation*}
$$

This establishes the induction step and concludes the proof that 17.21) holds for all $g=1, \ldots, G_{h}$.
Finally, given that (17.21) holds for all $g=1, \ldots, G_{h}, 17.24$ gives the results stated in 17.18) and 17.18) gives the results stated in the Lemma by the argument in 17.18)-17.20).

Now we use the approach in Johansen (1991, pp. 1569-1571) and Robin and Smith (2000, pp. 172-173) to prove Theorem 10.4. In these papers, asymptotic results are established under a fixed true distribution under which certain population eigenvalues are either positive or zero. Here we need to deal with drifting sequences of distributions under which these population eigenvalues may be positive or zero for any given $n$, but the positive ones may drift to zero as $n \rightarrow \infty$, possibly at different rates. This complicates the proof. In particular, the rate of convergence result of Lemma 17.1 (b) is needed in the present context, but not in the fixed distribution scenario considered in Johansen (1991) and Robin and Smith (2000).

Proof of Theorem 10.4. Theorem 10.4 (a) and (c) follow immediately from Lemma 17.1 (a).
Next, we assume $q<p$ and we prove part (b). The eigenvalues $\left\{\widehat{\kappa}_{j n}: j \leq p\right\}$ of $n \widehat{U}_{n} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n}$ $\times \widehat{D}_{n} \widehat{U}_{n}$ are the ordered solutions to the determinantal equation $\left|n \widehat{U}_{n} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} \widehat{U}_{n}-\kappa I_{p}\right|=0$. Equivalently, with probability that goes to one ( $\mathrm{wp} \rightarrow 1$ ), they are the solutions to

$$
\begin{equation*}
\left|Q_{n}^{\diamond}(\kappa)\right|=0, \text { where } Q_{n}^{\diamond}(\kappa):=n S_{n} B_{n}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} U_{n} B_{n} S_{n}-\kappa S_{n}^{\prime} B_{n}^{\prime} U_{n}^{\prime} \widehat{U}_{n}^{-1} \widehat{U}_{n}^{-1} U_{n} B_{n} S_{n}, \tag{17.35}
\end{equation*}
$$

because $\left|S_{n}\right|>0,\left|B_{n}\right|>0,\left|U_{n}\right|>0$, and $\left|\widehat{U}_{n}\right|>0 \mathrm{wp} \rightarrow 1$. Thus, $\lambda_{\min }\left(n \widehat{U}_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} \widehat{U}_{n}\right)$ equals the smallest solution, $\widehat{\kappa}_{p n}$, to $\left|Q_{n}^{\diamond}(\kappa)\right|=0 \mathrm{wp} \rightarrow 1$. (For simplicity, we omit the qualifier $\mathrm{wp} \rightarrow 1$ that applies to several statements below.)

We write $Q_{n}^{\diamond}(\kappa)$ in partitioned form using

$$
\begin{align*}
B_{n} S_{n} & =\left(B_{n, q} S_{n, q}, B_{n, p-q}\right), \text { where } \\
S_{n, q} & :=\operatorname{Diag}\left\{\left(n^{1 / 2} \tau_{1 F_{n}}\right)^{-1}, \ldots,\left(n^{1 / 2} \tau_{q F_{n}}\right)^{-1}\right\} \in R^{q \times q} . \tag{17.36}
\end{align*}
$$

The convergence result of Lemma 10.3 for $n^{1 / 2} W_{n} \widehat{D}_{n} U_{n} T_{n}\left(=n^{1 / 2} W_{n} \widehat{D}_{n} U_{n} B_{n} S_{n}\right)$ can be written as

$$
\begin{equation*}
n^{1 / 2} W_{n} \widehat{D}_{n} U_{n} B_{n, q} S_{n, q} \rightarrow_{p} \bar{\Delta}_{h, q}:=h_{3, q} \text { and } n^{1 / 2} W_{n} \widehat{D}_{n} U_{n} B_{n, p-q} \rightarrow_{d} \bar{\Delta}_{h, p-q} \tag{17.37}
\end{equation*}
$$

where $\bar{\Delta}_{h, q}$ and $\bar{\Delta}_{h, p-q}$ are defined in 10.17 .
We have

$$
\begin{equation*}
\widehat{W}_{n} W_{n}^{-1} \rightarrow_{p} I_{k} \text { and } \widehat{U}_{n} U_{n}^{-1} \rightarrow_{p} I_{p} \tag{17.38}
\end{equation*}
$$

because $\widehat{W}_{n} \rightarrow_{p} h_{71}:=\lim W_{n}$ (by Assumption WU(a) and (c)), $\widehat{U}_{n} \rightarrow_{p} h_{81}:=\lim U_{n}$ (by Assumption $\mathrm{WU}(\mathrm{b})$ and (c)), and $h_{71}$ and $h_{81}$ are pd (by the conditions in $\mathcal{F}_{W U}$ ).

By (17.35)-17.38), we have

$$
\begin{align*}
Q_{n}^{\diamond}(\kappa)= & {\left[\begin{array}{cc}
I_{q}+o_{p}(1) & h_{3, q}^{\prime} n^{1 / 2} W_{n} \widehat{D}_{n} U_{n} B_{n, p-q}+o_{p}(1) \\
n^{1 / 2} B_{n, p-q}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} W_{n}^{\prime} h_{3, q}+o_{p}(1) & n^{1 / 2} B_{n, p-q}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} W_{n}^{\prime} W_{n} n^{1 / 2} \widehat{D}_{n} U_{n} B_{n, p-q}+o_{p}(1)
\end{array}\right] } \\
& -\kappa\left[\begin{array}{cc}
S_{n, q}^{2} & 0^{q \times(p-q)} \\
0^{(p-q) \times q} & I_{p-q}
\end{array}\right]-\kappa\left[\begin{array}{cc}
S_{n, q} A_{1 n} S_{n, q} & S_{n, q} A_{2 n} \\
A_{2 n}^{\prime} S_{n, q} & A_{3 n}
\end{array}\right], \text { where }  \tag{17.39}\\
\widehat{A}_{n}= & {\left[\begin{array}{cc}
A_{1 n} & A_{2 n} \\
A_{2 n}^{\prime} & A_{3 n}
\end{array}\right]:=B_{n}^{\prime} U_{n}^{\prime} \widehat{U}_{n}^{-1} \widehat{U}_{n}^{-1} U_{n} B_{n}-I_{p}=o_{p}(1) \text { for } A_{1 n} \in R^{q \times q}, A_{2 n} \in R^{q \times(p-q)}, }
\end{align*}
$$

and $A_{3 n} \in R^{(p-q) \times(p-q)}, \widehat{A}_{n}$ is defined in 17.39 just as in 17.5), and the first equality uses $\bar{\Delta}_{h, q}:=h_{3, q}$ and $\bar{\Delta}_{h, q}^{\prime} \bar{\Delta}_{h, q}=h_{3, q}^{\prime} h_{3, q}=\lim C_{n, q}^{\prime} C_{n, q}=I_{q}($ by 10.7, 10.9, 10.12, and 10.17) . Note that $A_{j n}$ and $\widehat{A}_{j n}$ (defined in 17.2 ) are not the same in general for $j=1,2,3$, because their dimensions differ. For example, $A_{1 n} \in R^{q \times q}$, whereas $\widehat{A}_{1 n} \in R^{r_{1}^{\diamond} \times r_{1}^{\circ}}$.

If $q=0(<p)$, then $B_{n}=B_{n, p-q}$ and

$$
\begin{align*}
& n B_{n}^{\prime} \widehat{U}_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} \widehat{U}_{n} B_{n} \\
= & n B_{n}^{\prime}\left(U_{n}^{-1} \widehat{U}_{n}\right)^{\prime} B_{n}^{-1 \prime} B_{n}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} W_{n}^{\prime}\left(\widehat{W}_{n} W_{n}^{-1}\right)^{\prime}\left(\widehat{W}_{n} W_{n}^{-1}\right)\left(W_{n} \widehat{D}_{n} U_{n} B_{n}\right) B_{n}^{-1}\left(U_{n}^{-1} \widehat{U}_{n}\right) B_{n} \\
\rightarrow & { }_{d} \bar{\Delta}_{h, p-q}^{\prime} \bar{\Delta}_{h, p-q}, \tag{17.40}
\end{align*}
$$

where the convergence holds by 17.37 and 17.38 and $\bar{\Delta}_{h, p-q}$ is defined as in 10.17 with $q=0$. The smallest eigenvalue of a matrix is a continuous function of the matrix (by Elsner's Theorem, see Stewart (2001, Thm. 3.1, pp. 37-38)). Hence, the smallest eigenvalue of $n B_{n}^{\prime} \widehat{U}_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime} \widehat{W}_{n} \widehat{D}_{n} \widehat{U}_{n} B_{n}$ converges in distribution to the smallest eigenvalue of $\bar{\Delta}_{h, p-q}^{\prime} h_{3, k-q} h_{3, k-q}^{\prime} \bar{\Delta}_{h, p-q}$ (using $h_{3, k-q} h_{3, k-q}^{\prime}=$ $h_{3} h_{3}^{\prime}=I_{k}$ when $q=0$ ), which proves part (b) of Theorem 10.4 when $q=0$.

In the remainder of the proof of part (b), we assume $1 \leq q<p$, which is the remaining case
to be considered in the proof of part (b). The formula for the determinant of a partitioned matrix and (17.39) give

$$
\begin{align*}
\left|Q_{n}^{\diamond}(\kappa)\right|= & \left|Q_{1 n}^{\diamond}(\kappa)\right| \cdot\left|Q_{2 n}^{\diamond}(\kappa)\right|, \text { where } \\
Q_{1 n}^{\diamond}(\kappa):= & I_{q}+o_{p}(1)-\kappa S_{n, q}^{2}-\kappa S_{n, q} A_{1 n} S_{n, q}, \\
Q_{2 n}^{\diamond}(\kappa):= & n^{1 / 2} B_{n, p-q}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} W_{n}^{\prime} W_{n} n^{1 / 2} \widehat{D}_{n} U_{n} B_{n, p-q}+o_{p}(1)-\kappa I_{p-q}-\kappa A_{3 n} \\
& -\left[n^{1 / 2} B_{n, p-q}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} W_{n}^{\prime} h_{3, q}+o_{p}(1)-\kappa A_{2 n}^{\prime} S_{n, q}\right]\left(I_{q}+o_{p}(1)-\kappa S_{n, q}^{2}-\kappa S_{n, q} A_{1 n} S_{n, q}\right)^{-1} \\
& \times\left[h_{3, q}^{\prime} n^{1 / 2} W_{n} \widehat{D}_{n} U_{n} B_{n, p-q}+o_{p}(1)-\kappa S_{n, q} A_{2 n}\right], \tag{17.41}
\end{align*}
$$

none of the $o_{p}(1)$ terms depend on $\kappa$, and the equation in the first line holds provided $Q_{1 n}^{\diamond}(\kappa)$ is nonsingular.

By Lemma 17.1(b) (which applies for $1 \leq q<p$ ), for $j=q+1, \ldots, p$, we have $\widehat{\kappa}_{j n} S_{n, q}^{2}=o_{p}(1)$ and $\widehat{\kappa}_{j n} S_{n, q} A_{1 n} S_{n, q}=o_{p}(1)$. Thus,

$$
\begin{equation*}
Q_{1 n}^{\diamond}\left(\widehat{\kappa}_{j n}\right)=I_{q}+o_{p}(1)-\widehat{\kappa}_{j n} S_{n, q}^{2}-\widehat{\kappa}_{j n} S_{n, q} A_{1 n} S_{n, q}=I_{q}+o_{p}(1) . \tag{17.42}
\end{equation*}
$$

By 17.35 and 17.41), $\left|Q_{n}^{\diamond}\left(\widehat{\kappa}_{j n}\right)\right|=\left|Q_{1 n}^{\diamond}\left(\widehat{\kappa}_{j n}\right)\right| \cdot\left|Q_{2 n}^{\diamond}\left(\widehat{\kappa}_{j n}\right)\right|=0$ for $j=1, \ldots, p$. By 17.42), $\left|Q_{1 n}^{\diamond}\left(\widehat{\kappa}_{j n}\right)\right| \neq 0$ for $j=q+1, \ldots, p \mathrm{wp} \rightarrow 1$. Hence, $\mathrm{wp} \rightarrow 1$,

$$
\begin{equation*}
\left|Q_{2 n}^{\diamond}\left(\widehat{\kappa}_{j n}\right)\right|=0 \text { for } j=q+1, \ldots, p \tag{17.43}
\end{equation*}
$$

Now we plug in $\widehat{\kappa}_{j n}$ for $j=q+1, \ldots, p$ into $Q_{2 n}^{\diamond}(\kappa)$ in (17.41) and use 17.42). We have

$$
\begin{align*}
Q_{2 n}^{\diamond}\left(\widehat{\kappa}_{j n}\right)= & n B_{n, p-q}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} W_{n}^{\prime} W_{n} \widehat{D}_{n} U_{n} B_{n, p-q}+o_{p}(1) \\
& -\left[n^{1 / 2} B_{n, p-q}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} W_{n}^{\prime} h_{3, q}+o_{p}(1)\right]\left(I_{q}+o_{p}(1)\right)\left[h_{3, q}^{\prime} n^{1 / 2} W_{n} \widehat{D}_{n} U_{n} B_{n, p-q}+o_{p}(1)\right] \\
& -\widehat{\kappa}_{j n}\left[I_{p-q}+A_{3 n}-\left(n^{1 / 2} B_{n, p-q}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} W_{n}^{\prime} h_{3, q}+o_{p}(1)\right)\left(I_{q}+o_{p}(1)\right) S_{n, q} A_{2 n}\right. \\
& -A_{2 n}^{\prime} S_{n, q}\left(I_{q}+o_{p}(1)\right)\left(h_{3, q}^{\prime} n^{1 / 2} W_{n} \widehat{D}_{n} U_{n} B_{n, p-q}+o_{p}(1)\right) \\
& \left.+\widehat{\kappa}_{j n} A_{2 n}^{\prime} S_{n, q}\left(I_{q}+o_{p}(1)\right) S_{n, q} A_{2 n}\right] . \tag{17.44}
\end{align*}
$$

The term in square brackets on the last three lines of 17.44 that multiplies $\widehat{\kappa}_{j n}$ equals

$$
\begin{equation*}
I_{p-q}+o_{p}(1), \tag{17.45}
\end{equation*}
$$

because $A_{3 n}=o_{p}(1)($ by 17.39$), n^{1 / 2} W_{n} \widehat{D}_{n} U_{n} B_{n, p-q}=O_{p}(1)$ (by 17.37) ), $S_{n, q}=o(1)$ (by the definitions of $q$ and $S_{n, q}$ in 10.16) and 17.36, respectively, and $\left.h_{1, j}:=\lim n^{1 / 2} \tau_{j F_{n}}\right), A_{2 n}=o_{p}(1)$
(by 17.39) , and $\widehat{\kappa}_{j n} A_{2 n}^{\prime} S_{n, q}\left(I_{q}+o_{p}(1)\right) S_{n, q} A_{2 n}=A_{2 n}^{\prime} \widehat{\kappa}_{j n} S_{n, q}^{2} A_{2 n}+A_{2 n}^{\prime} \widehat{\kappa}_{j n} S_{n, q} o_{p}(1) S_{n, q} A_{2 n}=o_{p}(1)$ (using $\widehat{\kappa}_{j n} S_{n, q}^{2}=o_{p}(1)$ and $A_{2 n}=o_{p}(1)$ ).

Equations (17.44) and 17.45) give

$$
\begin{align*}
Q_{2 n}^{\diamond}\left(\widehat{\kappa}_{j n}\right) & =n^{1 / 2} B_{n, p-q}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} W_{n}^{\prime}\left[I_{k}-h_{3, q} h_{3, q}^{\prime}\right] n^{1 / 2} W_{n} \widehat{D}_{n} U_{n} B_{n, p-q}+o_{p}(1)-\widehat{\kappa}_{j n}\left[I_{p-q}+o_{p}(1)\right] \\
& =n^{1 / 2} B_{n, p-q}^{\prime} U_{n}^{\prime} \widehat{D}_{n}^{\prime} W_{n}^{\prime} h_{3, k-q} h_{3, k-q}^{\prime} n^{1 / 2} W_{n} \widehat{D}_{n} U_{n} B_{n, p-q}+o_{p}(1)-\widehat{\kappa}_{j n}\left[I_{p-q}+o_{p}(1)\right] \\
& :=M_{n, p-q}-\widehat{\kappa}_{j n}\left[I_{p-q}+o_{p}(1)\right], \tag{17.46}
\end{align*}
$$

where the second equality uses $I_{k}=h_{3} h_{3}^{\prime}=h_{3, q} h_{3, q}^{\prime}+h_{3, k-q} h_{3, k-q}^{\prime}$ (because $h_{3}=\lim C_{n}$ is an orthogonal matrix) and the last line defines the $(p-q) \times(p-q)$ matrix $M_{n, p-q}$.

Equations (17.43) and (17.46) imply that $\left\{\widehat{\kappa}_{j n}: j=q+1, \ldots, p\right\}$ are the $p-q$ eigenvalues of the matrix

$$
\begin{equation*}
M_{n, p-q}^{\diamond}:=\left[I_{p-q}+o_{p}(1)\right]^{-1 / 2} M_{n, p-q}\left[I_{p-q}+o_{p}(1)\right]^{-1 / 2} \tag{17.47}
\end{equation*}
$$

by pre- and post-multiplying the quantities in 17.46 by the rhs quantity $\left[I_{p-q}+o_{p}(1)\right]^{-1 / 2}$ in (17.46). By (17.37),

$$
\begin{equation*}
M_{n, p-q}^{\diamond} \rightarrow{ }_{d} \bar{\Delta}_{h, p-q}^{\prime} h_{3, k-q} h_{3, k-q}^{\prime} \bar{\Delta}_{h, p-q} . \tag{17.48}
\end{equation*}
$$

The vector of (ordered) eigenvalues of a matrix is a continuous function of the matrix (by Elsner's Theorem, see Stewart (2001, Thm. 3.1, pp. 37-38)). By (17.48), the matrix $M_{n, p-q}^{\diamond}$ converges in distribution. In consequence, by the CMT, the vector of eigenvalues of $M_{n, p-q}^{\diamond}$, viz., $\left\{\widehat{\kappa}_{j n}: j=q+1, \ldots, p\right\}$, converges in distribution to the vector of eigenvalues of the limit matrix $\bar{\Delta}_{h, p-q}^{\prime} h_{3, k-q} h_{3, k-q}^{\prime} \bar{\Delta}_{h, p-q}$, which proves part (d) of Theorem 10.4 . In addition, $\lambda_{\min }\left(n \widehat{U}_{n}^{\prime} \widehat{D}_{n}^{\prime} \widehat{W}_{n}^{\prime}\right.$ $\times \widehat{W}_{n} \widehat{D}_{n} \widehat{U}_{n}$ ), which equals the smallest eigenvalue, $\widehat{\kappa}_{p n}$, converges in distribution to the smallest eigenvalue of $\bar{\Delta}_{h, p-q}^{\prime} h_{3, k-q} h_{3, k-q}^{\prime} \bar{\Delta}_{h, p-q}$, which completes the proof of part (b) of Theorem 10.4 .

The convergence in parts (a)-(d) of Theorem 10.4 is joint with that in Lemma 10.3 because it just relies on the convergence in distribution of $n^{1 / 2} W_{n} \widehat{D}_{n} U_{n} T_{n}$, which is part of the former. This establishes part (e) of Theorem 10.4

Part (f) of Theorem 10.4 holds by the same proof as used for parts (a)-(e) with $n$ replaced by $w_{n}$.

## 18 Proofs of Sufficiency of Several Conditions for the $\boldsymbol{\lambda}_{\mathbf{p}-\mathbf{j}}(\cdot)$ Condition in $\mathcal{F}_{\mathbf{0 j}}$

In this section, we show that the conditions in (3.10) and (3.11) are sufficient for the second condition in $\mathcal{F}_{0 j}$, which is $\lambda_{p-j}\left(\Psi_{F}^{C_{F, k-j}^{\prime} \Omega_{F}^{-1 / 2} G_{i} B_{F, p-j} \xi}\right) \geq \delta_{1} \forall \xi \in R^{p-j}$ with $\|\xi\|=1$.

Condition (i) in (3.10) is sufficient by the following argument:

$$
\begin{align*}
& \lambda_{p-j}\left(\Psi_{F}^{C_{F, k-j}^{\prime} \Omega_{F}^{-1 / 2} G_{i} B_{F, p-j} \xi}\right) \\
& \geq \lambda_{p-j}\left(\Psi_{F}^{\bar{C}_{F, p-j}^{\prime} \Omega_{F}^{-1 / 2} G_{i} B_{F, p-j} \xi}\right) \\
&= \lambda_{\min }\left(\left(\xi^{\prime} \otimes I_{p-j}\right) \Psi_{F}^{v e c}\left(\bar{C}_{F, p-j}^{\prime} \Omega_{F}^{-1 / 2} G_{i} B_{F, p-j}\right)\right. \\
&\left.\left(\xi \otimes I_{p-j}\right)\right) \\
&= \min _{\lambda \in R^{p-j}:\|\lambda\|=1}\left(\frac{\left(\xi \otimes I_{p-j}\right) \lambda}{\left\|\left(\xi \otimes I_{p-j}\right) \lambda\right\|}\right)^{\prime} \Psi_{F}^{v e c\left(\bar{C}_{F, p-j}^{\prime} \Omega_{F}^{-1 / 2} G_{i} B_{F, p-j}\right)} \frac{\left(\xi \otimes I_{p-j}\right) \lambda}{\left\|\left(\xi \otimes I_{p-j}\right) \lambda\right\|} \times\left\|\left(\xi \otimes I_{p-j}\right) \lambda\right\|^{2} \\
& \geq \min _{\eta \in R^{(p-j)^{2}}:\|\eta\|=1} \eta^{\prime} \Psi_{F}^{v e c\left(\bar{C}_{F, p-j}^{\prime} \Omega_{F}^{-1 / 2} G_{i} B_{F, p-j}\right)} \eta \times \min _{\lambda \in R^{p-j}:\|\lambda\|=1}\left\|\left(\xi \otimes I_{p-j}\right) \lambda\right\|^{2}  \tag{18.1}\\
&= \lambda_{\min }\left(\Psi_{F}^{v e c\left(\bar{C}_{F, p-j}^{\prime} \Omega_{F}^{-1 / 2} G_{i} B_{F, p-j}\right)}\right),
\end{align*}
$$

where the first inequality holds by Corollary 16.4(a) with $m=p-j$ and $r=k-p$ (because $\Psi_{F}^{\bar{C}_{F, p-j}^{\prime} \Omega_{F}^{-1 / 2} G_{i} B_{F, p-j} \xi}$ is a submatrix of $\Psi_{F}^{C_{F, k-j}^{\prime} \Omega_{F}^{-1 / 2} G_{i} B_{F, p-j} \xi}$, since $\Psi_{F}^{C_{F, k-j}^{\prime} \Omega_{F}^{-1 / 2} G_{i} B_{F, p-j} \xi}=$ $C_{F, k-j}^{\prime} \Psi_{F}^{\Omega_{F}^{-1 / 2} G_{i} B_{F, p-j} \xi} C_{F, k-j}$, likewise with $C_{F, k-j}^{\prime}$ replaced by $\bar{C}_{F, p-j}^{\prime}$, and by definition the rows of $\bar{C}_{F, p-j}^{\prime}$ are a collection of $p-j$ rows of $C_{F, k-j}^{\prime}$ ), the first equality holds because the ( $p-j$ )-th largest eigenvalue of a $(p-j) \times(p-j)$ matrix equals its minimum eigenvalue and by the general formula $\operatorname{vec}(A B C)=\left(C^{\prime} \otimes A\right) \operatorname{vec}(B)$, and the last equality holds because $\left\|\left(\xi \otimes I_{p-j}\right) \lambda\right\|^{2}=\lambda^{\prime}\left(\xi^{\prime} \xi \otimes I_{p-j}\right) \lambda=$ $\lambda^{\prime} \lambda=1$ using $\|\xi\|=\|\lambda\|=1$.

Condition (ii) in (3.10) is sufficient by sufficient condition (i) in (3.10) and the following:

$$
\begin{align*}
& \lambda_{\min }\left(\Psi_{F}^{v e c}\left(\bar{C}_{F, p-j}^{\prime} \Omega_{F}^{-1 / 2} G_{i} B_{F, p-j}\right)\right. \\
= & \min _{\eta \in R^{(p-j)^{2}}:\|\eta\|=1}\left(\frac{\left(I_{p-j} \otimes \bar{C}_{F, p-j}\right) \eta}{\left\|\left(I_{p-j} \otimes \bar{C}_{F, p-j}\right) \eta\right\|}\right)^{\prime} \Psi_{F}^{v e c\left(\Omega_{F}^{-1 / 2} G_{i} B_{F, p-j}\right)} \frac{\left(I_{p-j} \otimes \bar{C}_{F, p-j}\right) \eta}{\left\|\left(I_{p-j} \otimes \bar{C}_{F, p-j}\right) \eta\right\|} \\
& \times\left\|\left(I_{p-j} \otimes \bar{C}_{F, p-j}\right) \eta\right\|^{2} \\
\geq & \min _{\zeta \in R^{(p-j) k}:\|\zeta\|=1} \zeta^{\prime} \Psi_{F}^{v e c\left(\Omega_{F}^{-1 / 2} G_{i} B_{F, p-j}\right)} \zeta \times \min _{\eta \in R^{(p-j)^{2}}:\|\eta\|=1}\left\|\left(I_{p-j} \otimes \bar{C}_{F, p-j}\right) \eta\right\|^{2} \\
= & \lambda_{\min }\left(\Psi_{F}^{v e c\left(\Omega_{F}^{-1 / 2} G_{i} B_{F, p-j}\right)}\right), \tag{18.2}
\end{align*}
$$

where the last equality uses $\left\|\left(I_{p-j} \otimes \bar{C}_{F, p-j}\right) \eta\right\|^{2}=\eta^{\prime}\left(I_{p-j} \otimes \bar{C}_{F, p-j}^{\prime} \bar{C}_{F, p-j}\right) \eta=1$ because the rows of $\bar{C}_{F, p-j}^{\prime}$ are orthonormal and $\|\eta\|=1$.

Condition (iii) in (3.10) is sufficient by sufficient condition (ii) in (3.10) and a similar argument to that given in 18.2 using the fact that $\min _{\psi \in R^{p k}:\|\psi\|=1}\left\|\left(B_{F, p-j}^{\prime} \otimes I_{k}\right) \psi\right\|^{2}=1$ because the columns of $B_{F, p-j}$ are orthonormal.

Condition (iv) in (3.10) is sufficient by sufficient condition (iii) in (3.10) and a similar argument to that given in 18.2 using $\min _{\phi \in R^{p k}:\|\phi\|=1}\left\|\left(I_{p} \otimes \Omega_{F}^{-1 / 2}\right) \phi\right\|^{2} \geq M^{-2 /(2+\gamma)}$ for $M$ as in the definition of $\mathcal{F}$ in place of $\min _{\eta \in R^{(p-j)^{2}}:\|\eta\|=1}\left\|\left(I_{p-j} \otimes \bar{C}_{F, p-j}\right) \eta\right\|^{2}=1$. The latter inequality holds by the following calculations:

$$
\begin{align*}
\phi^{\prime}\left(I_{p} \otimes \Omega_{F}^{-1}\right) \phi & =\sum_{j=1}^{p}\left(\phi_{j} /\left\|\phi_{j}\right\|\right)^{\prime} \Omega_{F}^{-1}\left(\phi_{j} /\left\|\phi_{j}\right\|\right) \times\left\|\phi_{j}\right\|^{2} \\
& \geq \sum_{j=1}^{p} \lambda_{\min }\left(\Omega_{F}^{-1}\right) \times\left\|\phi_{j}\right\|^{2}=1 / \lambda_{\max }\left(\Omega_{F}\right) \geq M^{-2 /(2+\gamma)}, \tag{18.3}
\end{align*}
$$

where $\phi=\left(\phi_{1}^{\prime}, \ldots, \phi_{p}^{\prime}\right)^{\prime}$ for $\phi_{j} \in R^{k} \forall j \leq p$, the sums are over $j$ for which $\phi_{j} \neq 0^{k}$, the second equality uses $\|\phi\|=1$, and the last inequality holds because $\lambda_{\max }\left(\Omega_{F}\right)=\max _{\lambda \in R^{k}:\|\lambda\|=1} E_{F}\left(\lambda^{\prime} g_{i}\right)^{2} \leq$ $E_{F}\left\|g_{i}\right\|^{2}=\left(\left(E_{F}\left\|g_{i}\right\|^{2}\right)^{1 / 2}\right)^{2} \leq\left(\left(E_{F}\left\|g_{i}\right\|^{2+\gamma}\right)^{1 /(2+\gamma)}\right)^{2} \leq M^{2 /(2+\gamma)}$ by successively applying the Cauchy-Bunyakovsky-Schwarz inequality, Lyapunov's inequality, and the moment bound $E_{F}\left\|g_{i}\right\|^{2+\gamma}$ $\leq M$ in $\mathcal{F}$.

Conditions (v) and (vi) in (3.10) are sufficient by the following argument. Write

$$
\begin{equation*}
\Psi_{F}^{v e c\left(G_{i}\right)}=\left(M_{F}, I_{p k}\right) \Sigma_{F}^{f_{i}}\left(M_{F}, I_{p k}\right)^{\prime}, \text { where } M_{F}=-\left(E_{F} v e c\left(G_{i}\right) g_{i}^{\prime}\right)\left(E_{F} g_{i} g_{i}^{\prime}\right)^{-1} \in R^{p k \times k} \tag{18.4}
\end{equation*}
$$

We have

$$
\begin{align*}
\lambda_{\min }\left(\Psi_{F}^{v e c\left(G_{i}\right)}\right) & =\min _{\lambda \in R^{p k}:\|\lambda\|=1} \lambda^{\prime}\left(M_{F}, I_{p k}\right) \Sigma_{F}^{f_{i}}\left(M_{F}, I_{p k}\right)^{\prime} \lambda \\
& =\min _{\lambda \in R^{p k}:\|\lambda\|=1}\left(\frac{\left(M_{F}, I_{p k}\right)^{\prime} \lambda}{\left\|\left(M_{F}, I_{p k}\right)^{\prime} \lambda\right\|}\right)^{\prime} \Sigma_{F}^{f_{i}}\left(\frac{\left(M_{F}, I_{p k}\right)^{\prime} \lambda}{\left\|\left(M_{F}, I_{p k}\right)^{\prime} \lambda\right\|}\right) \times\left\|\left(M_{F}, I_{p k}\right)^{\prime} \lambda\right\|^{2} \\
& \geq \min _{\mu \in R^{(p+1) k}:\|\mu\|=1} \mu^{\prime} \Sigma_{F}^{f_{i}} \mu \\
& =\lambda_{\min }\left(\Sigma_{F}^{f_{i}}\right), \tag{18.5}
\end{align*}
$$

where the inequality uses $\left\|\left(M_{F}, I_{p k}\right)^{\prime} \lambda\right\|^{2}=\lambda^{\prime} \lambda+\lambda^{\prime} M_{F}^{\prime} M_{F} \lambda \geq 1$ for $\lambda \in R^{p k}$ with $\|\lambda\|=1$. This shows that condition (v) is sufficient for sufficient condition (iv) in 3.10. Since $\Sigma_{F}^{f_{i}}=$ $\operatorname{Var}_{F}\left(f_{i}\right)+E_{F} f_{i} E_{F} f_{i}^{\prime}$, condition (vi) is sufficient for sufficient condition (v) in (3.10).

The condition in (3.11) is sufficient by the following argument:

$$
\begin{equation*}
\lambda_{p-j}\left(\Psi_{F}^{C_{F, k-j}^{\prime} \Omega_{F}^{-1 / 2} G_{i} B_{F, p-j} \xi}\right) \geq \lambda_{p}\left(\Psi_{F}^{C_{F}^{\prime} \Omega_{F}^{-1 / 2} G_{i} B_{F, p-j} \xi}\right)=\lambda_{p}\left(\Psi_{F}^{\Omega_{F}^{-1 / 2} G_{i} B_{F, p-j} \xi}\right) \tag{18.6}
\end{equation*}
$$

where the first inequality holds by Corollary 16.4 (b) with $m=p$ and $r=j$ and the equality holds because $\Psi_{F}^{C_{F}^{\prime} \Omega_{F}^{-1 / 2} G_{i} B_{F, p-j} \xi}=C_{F}^{\prime} \Psi_{F}^{\Omega_{F}^{-1 / 2} G_{i} B_{F, p-j} \xi} C_{F}$ and $C_{F}$ is orthogonal.

## 19 Asymptotic Size of Kleibergen's CLR Test with JacobianVariance Weighting and the Proof of Theorem 5.1

In this section, we establish the asymptotic size of Kleibergen's CLR test with Jacobian-variance weighting when the Robin and Smith (2000) rank statistic (defined in (5.5) is employed. This rank statistic depends on a variance matrix estimator $\widetilde{V}_{D n}$. See Section 5 for the definition of the test. We provide a formula for the asymptotic size of the test that depends on the specifics of the moment conditions considered and does not necessarily equal its nominal size $\alpha$. First, in Section 19.1, we provide an example that illustrates the results in Theorem 5.1 and Comment (v) to Theorem 5.1 . In Section 19.2, we establish the asymptotic size of the test based on $\widetilde{V}_{D n}$ defined as in 5.3. In Section 19.3, we report some simulation results for a linear instrumental variable (IV) model with two rhs endogenous variables. In Section 19.4, we establish the asymptotic size of Kleibergen's CLR test with Jacobian-variance weighting under a general assumption that allows for other definitions of $\widetilde{V}_{D n}$.

In Section 19.5, we show that equally-weighted versions of Kleibergen's CLR test have correct asymptotic size when the Robin and Smith (2000) rank statistic is employed and a general equalweighting matrix $\widetilde{W}_{n}$ is employed. This result extends the result given in Theorem 6.1 in Section 6. which applies to the specific case where $\widetilde{W}_{n}=\widehat{\Omega}_{n}^{-1 / 2}$, as in 6.2. The results of Section 19.5 are a relatively simple by-product of the results in Section 19.4 .

Proofs of the results stated in this section are given in Section 19.6 .
Lemma 5.2 is proved in Section 19.7 .
Theorem 5.1 follows from Lemma 19.2 and Theorem 19.3 , which are stated in Section 19.4 .
As stated in footnote 4 in Section 2 of AG1, "under sequences $F_{n}$ such that $n^{1 / 2} E_{F_{n}} G\left(W_{i}, \theta\right)$ converges to a finite matrix, $n^{1 / 2} \widehat{D}_{n}(\theta)$ and $n^{1 / 2} \widehat{g}_{n}(\theta)\left(=n^{-1 / 2} \sum_{i=1}^{n} g\left(W_{i}, \theta\right)\right)$ are asymptotically independent (see Lemmas 10.2 and 10.3 in Section 10 in this SM). Therefore, if $r\left(\widehat{V}_{n}, n^{1 / 2} \widehat{D}_{n}(\theta)\right)$ is a continuous function of $n^{1 / 2} \widehat{D}_{n}(\theta)$ and a weighting matrix $\widehat{V}_{n}$ (that converges in probability to a positive definite matrix), then by the continuous mapping theorem (CMT),
$n^{1 / 2} \widehat{g}_{n}(\theta)$ and $r\left(\widehat{V}_{n}, n^{1 / 2} \widehat{D}_{n}(\theta)\right)$ are also asymptotically independent."
Footnote 4 of AG1 also states "however, under sequences for which a component of $n^{1 / 2} E_{F_{n}} G\left(W_{i}, \theta\right)$ diverges to plus or minus infinity, the CMT cannot be applied because $n^{1 / 2} \widehat{D}_{n}(\theta)$ does not converge in distribution, but rather, some component of it diverges to plus or minus infinity in probability (see Lemma 10.3 in Section 10 in this $S M$ when $h_{1, j}=\infty$ for some $\left.j \leq p\right)$. In this case, $r\left(\widehat{V}_{n}, n^{1 / 2} \widehat{D}_{n}(\theta)\right)$ may not have an asymptotic distribution, and if it does, $r\left(\widehat{V}_{n}, n^{1 / 2} \widehat{D}_{n}(\theta)\right)$ and $n^{1 / 2} \widehat{g}_{n}(\theta)$ are not necessarily asymptotically independent."

The following is a simple example of the latter situation when $p=2$. Let $r\left(\widehat{V}_{n}, n^{1 / 2} \widehat{D}_{n}(\theta)\right)=$ $\widehat{V}_{12 n}\left\|n^{1 / 2} \widehat{D}_{1 n}(\theta)\right\|$, where $\widehat{V}_{12 n}$ is the $(1,2)$ component of $\widehat{V}_{n}$ and $\widehat{D}_{1 n}(\theta)$ is the first column of $\widehat{D}_{n}(\theta)$. Assume $\widehat{V}_{n}-V \rightarrow_{p} 0$ for some matrix $V$ and $n^{1 / 2}\left(\widehat{V}_{n}-V\right) \rightarrow_{d} \xi$, where $\xi$ is a mean zero normal random matrix. Assume that under $F_{n}$ the first column $E_{F_{n}} G_{1}\left(W_{i}, \theta\right)$ of $E_{F_{n}} G\left(W_{i}, \theta\right)$ is a fixed nonzero vector, $G_{1}^{e}$ say. Assume that the $(1,2)$ element of $V$, denoted by $V_{12}$, equals zero under $F_{n}$. Then, $\widehat{D}_{1 n}(\theta) \rightarrow_{p} G_{1}^{e}$ (see Lemma 10.2 in Section 10 in this SM ) and $\widehat{V}_{12 n}\left\|n^{1 / 2} \widehat{D}_{1 n}(\theta)\right\|=$ $n^{1 / 2}\left(\widehat{V}_{12 n}-V_{12}\right)\left\|\widehat{D}_{1 n}(\theta)\right\| \rightarrow_{d} \xi_{12}\left\|G_{1}^{e}\right\|$. But, in general there is no reason why $\xi_{12}$ and the random limit of $n^{1 / 2} \widehat{g}_{n}(\theta)$ are independent. For simplicity, the previous example is somewhat contrived, because rank statistics typically are not of the form $\widehat{V}_{12 n}\left\|n^{1 / 2} \widehat{D}_{1 n}(\theta)\right\|$. But, components of rank statistics may be of this form. Section 19.1, which follows, provides a more concrete example of this type of situation.

### 19.1 An Example

Here we provide an example that illustrates the result of Theorem 5.1. In this example, the true distribution $F$ does not depend on $n$. Suppose $p=2, E_{F} G_{i}=\left(1^{k}, 0^{k}\right)$, where $c^{k}=(c, \ldots, c)^{\prime} \in R^{k}$ for $c=0,1, n^{1 / 2}\left(\widehat{D}_{n}-E_{F} G_{i}\right) \rightarrow{ }_{d} \bar{D}_{h}$ under $F$ for some random matrix $\bar{D}_{h}=\left(\bar{D}_{1 h}, \bar{D}_{2 h}\right) \in R^{k \times 2}$. Suppose for $\widetilde{M}_{n}=\widetilde{V}_{D n}^{-1 / 2}$ and $M_{F}=I_{2 k}$, we have $n^{1 / 2}\left(\widetilde{M}_{n}-M_{F}\right) \rightarrow_{d} \bar{M}_{h}$ under $F$ for some random matrix $\bar{M}_{h} \in R^{2 k \times 2 k}$. (The convergence results $n^{1 / 2}\left(\widehat{D}_{n}-E_{F} G_{i}\right) \rightarrow_{d} \bar{D}_{h}$ and $n^{1 / 2}\left(\widetilde{M}_{n}-M_{F}\right) \rightarrow_{d}$ $\bar{M}_{h}$ are established in Lemmas 10.2 and 19.2 , respectively, in Section 10 and Section 19.4 in this SM under general conditions.) We have

$$
\begin{equation*}
\widehat{D}_{n}^{\dagger}=\operatorname{vec}_{k, p}^{-1}\left(\widetilde{V}_{D n}^{-1 / 2} \operatorname{vec}\left(\widehat{D}_{n}\right)\right)=\left(\widetilde{M}_{11 n} \widehat{D}_{1 n}+\widetilde{M}_{12 n} \widehat{D}_{2 n}, \widetilde{M}_{21 n} \widehat{D}_{1 n}+\widetilde{M}_{22 n} \widehat{D}_{2 n}\right) \tag{19.1}
\end{equation*}
$$

where $\widehat{D}_{n}=\left(\widehat{D}_{1 n}, \widehat{D}_{2 n}\right), \widetilde{M}_{j \ell n}$ for $j, \ell=1,2$ are the four $k \times k$ submatrices of $\widetilde{M}_{n}$, and likewise for $M_{j \ell F}$ for $j, \ell=1,2$. Let $\bar{M}_{j \ell h}$ for $j, \ell=1,2$ denote the four $k \times k$ submatrices of $\bar{M}_{h}$. We let
$T_{n}^{\dagger}=\operatorname{Diag}\left\{n^{-1 / 2}, 1\right\}$. Then, we have

$$
\begin{align*}
n^{1 / 2} \widehat{D}_{n}^{\dagger} T_{n}^{\dagger} & =\left(\widetilde{M}_{11 n} \widehat{D}_{1 n}+\widetilde{M}_{12 n} \widehat{D}_{2 n}, n^{1 / 2} \widetilde{M}_{21 n} \widehat{D}_{1 n}+\widetilde{M}_{22 n} n^{1 / 2} \widehat{D}_{2 n}\right) \\
& \rightarrow{ }_{d}\left(I_{k} 1^{k}+0^{k \times k} 0^{k}, \bar{M}_{21 h} 1^{k}+I_{k} \bar{D}_{2 h}\right)=\left(1^{k}, \bar{M}_{21 h} 1^{k}+\bar{D}_{2 h}\right), \tag{19.2}
\end{align*}
$$

where the convergence uses $n^{1 / 2} \widetilde{M}_{21 n} \rightarrow{ }_{d} \bar{M}_{21 h}$ (because $M_{21 F}=0^{k \times k}$ ) and $n^{1 / 2} \widehat{D}_{2 n} \rightarrow{ }_{d} \bar{D}_{2 h}$ (because $E_{F} G_{i 2}=0^{k}$ ). Equation (19.2) shows that the asymptotic distribution of $n^{1 / 2} \hat{D}_{n}^{\dagger} T_{n}^{\dagger}$ depends on the randomness of the variance estimator $\widetilde{V}_{D n}$ through $\bar{M}_{21 h}$.

It may appear that this example is quite special and the asymptotic behavior in (19.2) only arises in special circumstances, because $E_{F} G_{i}=\left(1^{k}, 0^{k}\right), M_{21 F}=0^{k \times k}$, and $M_{F}=I_{2 k}$ in this example. But this is not true. The asymptotic behavior in 19.2 arises quite generally, as shown in Theorem 5.1, whenever $p \geq 2{ }^{12}$

If one replaces $\widetilde{V}_{D n}^{-1 / 2}$ by its probability limit, $M_{F}$, in the definition of $\widehat{D}_{n}^{\dagger}$, then the calculations in 19.2 hold but with $n^{1 / 2} \widetilde{M}_{21 n}$ replaced by $n^{1 / 2} M_{21 F}=0^{k \times k}$ in the first line and, hence, $\bar{M}_{21 h}$ replaced by $0^{k \times k}$ in the second line. Hence, in this case, the asymptotic distribution only depends on $\bar{D}_{h}$. Hence, Comment (iv) to Theorem 5.1 holds in this example.

Suppose one defines $\widehat{D}_{n}^{\dagger}$ by $\widetilde{W}_{n} \widehat{D}_{n}$ as in Comment (v) to Theorem 5.1. This yields equal weighting of each column of $\widehat{D}_{n}$. This is equivalent to replacing $\widetilde{V}_{D n}^{-1 / 2}$ by $I_{2} \otimes \widetilde{W}_{n}$ in the definition of $\widehat{D}_{n}^{\dagger}$ in 19.1. In this case, the off-diagonal $k \times k$ blocks of $I_{2} \otimes \widetilde{W}_{n}$ are $0^{k \times k}$ and, hence, $\widetilde{M}_{21 n}$ in the first line of 19.2 equals $0^{k \times k}$, which implies that $\bar{M}_{21 h}=0^{k \times k}$ in the second line of 19.2 . Thus, the asymptotic distribution of $\widehat{D}_{n}^{\dagger}$ does not depend on the asymptotic distribution of the (normalized) weight matrix estimator $\widetilde{W}_{n}$. It only depends on the probability limit of $\widetilde{W}_{n}$, as stated in Comment (v) to Theorem 5.1.

### 19.2 Asymptotic Size of Kleibergen's CLR Test with Jacobian-Variance Weighting

In this subsection, we determine the asymptotic size of Kleibergen's CLR test when $\widehat{D}_{n}$ is weighted by $\widetilde{V}_{D n}$, defined in (5.3), which yields what we call Jacobian-variance weighting, and the Robin and Smith (2000) rank statistic is employed. This rank statistic is defined in (5.5) with $\theta=\theta_{0}$. For convenience, we restate the definition here:

$$
\begin{equation*}
r k_{n}=r k_{n}^{\dagger}:=\lambda_{\min }\left(n\left(\widehat{D}_{n}^{\dagger}\right)^{\prime} \widehat{D}_{n}^{\dagger}\right), \text { where } \widehat{D}_{n}^{\dagger}:=\operatorname{vec}_{k, p}^{-1}\left(\widetilde{V}_{D n}^{-1 / 2} \operatorname{vec}\left(\widehat{D}_{n}\right)\right) \tag{19.3}
\end{equation*}
$$

[^1](so $\hat{D}_{n}^{\dagger}$ is as in 5.4 with $\theta=\theta_{0}$ ). As in Section 5, the function $v e c_{k, p}^{-1}(\cdot)$ is the inverse of the $v e c(\cdot)$ function for $k \times p$ matrices. Thus, the domain of $v e c_{k, p}^{-1}(\cdot)$ consists of $k p$-vectors and its range consists of $k \times p$ matrices. Let
\[

$$
\begin{equation*}
\widehat{\kappa}_{j n}^{\dagger} \text { denote the } j \text { th eigenvalue of } n\left(\widehat{D}_{n}^{\dagger}\right)^{\prime} \widehat{D}_{n}^{\dagger} \text {, for } j=1, \ldots, p \tag{19.4}
\end{equation*}
$$

\]

ordered to be nonincreasing in $j$. By definition, $\lambda_{\min }\left(n\left(\widehat{D}_{n}^{\dagger}\right)^{\prime} \widehat{D}_{n}^{\dagger}\right)=\widehat{\kappa}_{p n}^{\dagger}$. Also, the $j$ th singular value of $n^{1 / 2} \widehat{D}_{n}^{\dagger}$ equals $\left(\widehat{\kappa}_{j n}^{\dagger}\right)^{1 / 2}$.

Define the parameter space $\mathcal{F}_{K C L R}$ for the distribution $F$ by

$$
\begin{equation*}
\mathcal{F}_{K C L R}:=\left\{F \in \mathcal{F}: \lambda_{\min }\left(\operatorname{Var}_{F}\left(\left(g_{i}^{\prime}, \operatorname{vec}\left(G_{i}\right)^{\prime}\right)^{\prime}\right)\right) \geq \delta_{2}, E_{F}\left\|\left(g_{i}^{\prime}, \operatorname{vec}\left(G_{i}\right)^{\prime}\right)^{\prime}\right\|^{4+\gamma} \leq M\right\}, \tag{19.5}
\end{equation*}
$$

where $\delta_{2}>0$ and $\gamma>0$ and $M<\infty$ are as in the definition of $\mathcal{F}$ in (3.3). Note that $\mathcal{F}_{K C L R} \subset \mathcal{F}_{0}$ when $\delta_{1}$ in $\mathcal{F}_{0}$ satisfies $\delta_{1} \leq M^{-2 /(2+\gamma)} \delta_{2}$, by condition (vi) in (3.10). Let vech $(\cdot)$ denote the half vectorization operator that vectorizes the nonredundant elements in the columns of a symmetric matrix (that is, the elements on or below the main diagonal). The moment condition in $\mathcal{F}_{K C L R}$ is imposed because the asymptotic distribution of the rank statistic $r k_{n}^{\dagger}$ depends on a triangular array CLT for $\operatorname{vech}\left(f_{i}^{*} f_{i}^{* \prime}\right)$, which employs $4+\gamma$ moments for $f_{i}^{*}$, where $f_{i}^{*}:=\left(g_{i}^{\prime}, \operatorname{vec}\left(G_{i}-E_{F_{n}} G_{i}\right)^{\prime}\right)^{\prime}$ as in 5.6. The $\lambda_{\min }(\cdot)$ condition in $\mathcal{F}_{K C L R}$ ensures that $\widetilde{V}_{D n}$ is positive definite $\mathrm{wp} \rightarrow 1$, which is needed because $\widetilde{V}_{D n}$ enters the rank statistic $r k_{n}^{\dagger}$ via $\widetilde{V}_{D n}^{-1 / 2}$, see 19.3.

For a fixed distribution $F, \widetilde{V}_{D n}$ estimates $\Phi_{F}^{v e c\left(G_{i}\right)}$ defined in 10.15 , where $\Phi_{F}^{v e c\left(G_{i}\right)}$ is pd by its definition in 10.15 and the $\lambda_{\min }(\cdot)$ condition in $\mathcal{F}_{K C L R}$. More specifically, $\Phi_{F}^{v e c\left(G_{i}\right)}$ is pd because by $10.15 \Phi_{F}^{\operatorname{vec}\left(G_{i}\right)}:=\operatorname{Var}_{F}\left(\operatorname{vec}\left(G_{i}\right)-\left(E_{F} \operatorname{vec}\left(G_{\ell}\right) g_{\ell}^{\prime}\right) \Omega_{F}^{-1} g_{i}\right)=\left(-\left(E_{F} \operatorname{vec}\left(G_{\ell}\right) g_{\ell}^{\prime}\right) \Omega_{F}^{-1}, I_{p k}\right)$ $\operatorname{Var}_{F}\left(\left(g_{i}^{\prime}, \operatorname{vec}\left(G_{i}\right)^{\prime}\right)^{\prime}\right)\left(-\left(E_{F} \operatorname{vec}\left(G_{\ell}\right) g_{\ell}^{\prime}\right) \Omega_{F}^{-1}, I_{p k}\right)^{\prime}$, where $\left(-\left(E_{F} v e c\left(G_{\ell}\right) g_{\ell}^{\prime}\right) \Omega_{F}^{-1}, I_{p k}\right) \in R^{p k \times(p+1) k}$ has full row rank $p k$ and $\operatorname{Var}_{F}\left(\left(g_{i}^{\prime}, \operatorname{vec}\left(G_{i}\right)^{\prime}\right)^{\prime}\right)$ is pd by the $\lambda_{\min }(\cdot)$ condition in $\mathcal{F}_{K C L R}$. Let

$$
\begin{align*}
& M_{F}=\left[\begin{array}{ccc}
M_{11 F} & \cdots & M_{1 p F} \\
\vdots & \ddots & \vdots \\
M_{p 1 F} & \cdots & M_{p p F}
\end{array}\right]:=\left(\Phi_{F}^{v e c\left(G_{i}\right)}\right)^{-1 / 2} \text { and }  \tag{19.6}\\
& D_{F}^{\dagger}:=\sum_{j=1}^{p}\left(M_{1 j F} E_{F} G_{i j}, \ldots, M_{p j F} E_{F} G_{i j}\right) \in R^{k \times p}, \text { where } G_{i}=\left(G_{i 1}, \ldots, G_{i p}\right) \in R^{k \times p} .
\end{align*}
$$

Let $\left(\tau_{1 F}^{\dagger}, \ldots, \tau_{p F}^{\dagger}\right)$ denote the singular values of $D_{F}^{\dagger}$. Define

$$
\begin{align*}
& B_{F}^{\dagger} \in R^{p \times p} \text { to be an orthogonal matrix of eigenvectors of } D_{F}^{\dagger \prime} D_{F}^{\dagger} \text { and } \\
& C_{F}^{\dagger} \in R^{k \times k} \text { to be an orthogonal matrix of eigenvectors of } D_{F}^{\dagger} D_{F}^{\dagger \prime} \tag{19.7}
\end{align*}
$$

ordered so that the corresponding eigenvalues $\left(\kappa_{1 F}^{\dagger}, \ldots, \kappa_{p F}^{\dagger}\right)$ and $\left(\kappa_{1 F}^{\dagger}, \ldots, \kappa_{p F}^{\dagger}, 0, \ldots, 0\right) \in R^{k}$, respectively, are nonincreasing. We have $\kappa_{j F}^{\dagger}=\left(\tau_{j F}^{\dagger}\right)^{2}$ for $j=1, \ldots, p$. Note that 19.7 gives definitions of $B_{F}$ and $C_{F}$ that are similar to the definitions in 10.6 and 10.7, but differ because $D_{F}^{\dagger}$ replaces $W_{F}\left(E_{F} G_{i}\right) U_{F}$ in the definitions.

Define $\left(\lambda_{1, F}, \ldots, \lambda_{9, F}\right)$ as in 10.9 with $\lambda_{7, F}=W_{F}=\Omega_{F}^{-1 / 2}, \lambda_{8, F}=I_{p}$, and $W_{1}(\cdot)$ and $U_{1}(\cdot)$ equal to identity functions. Define

$$
\begin{equation*}
\lambda_{10, F}=\operatorname{Var}_{F}\binom{f_{i}^{*}}{\operatorname{vech}\left(f_{i}^{*} f_{i}^{* \prime}\right)} \in R^{d^{*} \times d^{*}}, \tag{19.8}
\end{equation*}
$$

where $d^{*}:=(p+1) k+(p+1) k((p+1) k+1) / 2$. Define $\left(\lambda_{1, F}^{\dagger}, \lambda_{2, F}^{\dagger}, \lambda_{3, F}^{\dagger}, \lambda_{6, F}^{\dagger}\right)$ as $\left(\lambda_{1, F}, \lambda_{2, F}, \lambda_{3, F}, \lambda_{6, F}\right)$ are defined in 10.9 but with $\left\{\tau_{j F}^{\dagger}: j \leq p\right\}, B_{F}^{\dagger}$, and $C_{F}^{\dagger}$ in place of $\left\{\tau_{j F}: j \leq p\right\}, B_{F}$, and $C_{F}$, respectively.

Define

$$
\begin{align*}
\lambda & =\lambda_{F}:=\left(\lambda_{1, F}, \ldots, \lambda_{10, F}, \lambda_{1, F}^{\dagger}, \lambda_{2, F}^{\dagger}, \lambda_{3, F}^{\dagger}, \lambda_{6, F}^{\dagger}\right),  \tag{19.9}\\
\Lambda_{K C L R} & :=\left\{\lambda: \lambda=\left(\lambda_{1, F}, \ldots, \lambda_{10, F}, \lambda_{1, F}^{\dagger}, \lambda_{2, F}^{\dagger}, \lambda_{3, F}^{\dagger}, \lambda_{6, F}^{\dagger}\right) \text { for some } F \in \mathcal{F}_{K C L R}\right\}, \text { and } \\
h_{n}(\lambda) & :=\left(n^{1 / 2} \lambda_{1, F}, \lambda_{2, F}, \lambda_{3, F}, \lambda_{4, F}, \lambda_{5, F}, \lambda_{6, F}, \lambda_{7, F}, \lambda_{8, F}, \lambda_{10, F}, n^{1 / 2} \lambda_{1, F}^{\dagger}, \lambda_{2, F}^{\dagger}, \lambda_{3, F}^{\dagger}, \lambda_{6, F}^{\dagger}\right) .
\end{align*}
$$

Let $\left\{\lambda_{n, h} \in \Lambda_{K C L R}: n \geq 1\right\}$ denote a sequence $\left\{\lambda_{n} \in \Lambda_{K C L R}: n \geq 1\right\}$ for which $h_{n}\left(\lambda_{n}\right) \rightarrow h \in H$, for $H$ as in 10.1). The asymptotic variance of $n^{1 / 2} \operatorname{vec}\left(\widehat{D}_{n}-E_{F_{n}} G_{i}\right)$ is $\Phi_{h}^{v e c}\left(G_{i}\right)$ under $\left\{\lambda_{n, h} \in\right.$ $\left.\Lambda_{K C L R}: n \geq 1\right\}$ by Lemma 10.2 .

Define $h_{1, j}$ for $j \leq p$ and $h_{s}$ for $s=2, \ldots, 8$ as in 10.12, $q=q_{h}$ as in 10.16, $, h_{2, q}, h_{2, p-q}, h_{3, q}$, $h_{3, p-q}$, and $h_{1, p-q}^{\diamond}$ as in 10.17, and $\Upsilon_{n}, \Upsilon_{n, q}$, and $\Upsilon_{n, p-q}$ as in 9.2 . Note that $h_{7}=h_{5, g}^{-1 / 2}$ and $h_{8}=I_{p}$ due to the definitions of $\lambda_{7, F}$ and $\lambda_{8, F}$ given above, where $h_{5, g}\left(=\lim E_{F_{n}} g_{i} g_{i}^{\prime}\right)$ denotes the upper left $k \times k$ submatrix of $h_{5}$, as in Section 10 .

For a sequence $\left\{\lambda_{n, h} \in \Lambda_{K C L R}: n \geq 1\right\}$, we have

$$
h_{10}=\left[\begin{array}{cc}
h_{10, f^{*}} & h_{10, f^{*} f^{* 2}}  \tag{19.10}\\
h_{10, f^{* 2} f^{*}} & h_{10, f^{* 2} f^{* 2}}
\end{array}\right]:=\lim \operatorname{Var}_{F_{n}}\binom{f_{i}^{*}}{\operatorname{vech}\left(f_{i}^{*} f_{i}^{* \prime}\right)} \in R^{d^{*} \times d^{*}} .
$$

Note that $h_{10, f^{*}} \in R^{(p+1) k \times(p+1) k}$ is pd by the definition of $\mathcal{F}_{K C L R}$ in 19.5.
With $\tau_{j F}^{\dagger}, B_{F}^{\dagger}$, and $C_{F}^{\dagger}$ in place of $\tau_{j F}, B_{F}$, and $C_{F}$, respectively, define $h_{1, j}^{\dagger}$ for $j \leq p$ and $h_{s}^{\dagger}$ for $s=2,3,6$ as in 10.12 as analogues to the quantities without the $\dagger$ superscript, define $q^{\dagger}=q_{h}^{\dagger}$ as in 10.16, define $h_{2, q^{\dagger}}^{\dagger}, h_{2, p-q^{\dagger}}^{\dagger}, h_{3, q^{\dagger}}^{\dagger}, h_{3, k-q^{\dagger}}^{\dagger}$, and $h_{1, p-q^{\dagger}}^{\dagger \diamond}$ as in 10.17, and define $\Upsilon_{n}^{\dagger}, \Upsilon_{n, q^{\dagger}}^{\dagger}$, and $\Upsilon_{n, p-q^{\dagger}}^{\dagger}$ as in 9.2 . The quantity $q^{\dagger}$ determines the asymptotic behavior of $r k_{n}^{\dagger}$. By definition, $q^{\dagger}$ is the largest value $j(\leq p)$ for which $\lim n^{1 / 2} \tau_{j F_{n}}^{\dagger}=\infty$ under $\left\{\lambda_{n, h} \in \Lambda_{K C L R}: n \geq 1\right\}$. It is shown below that if $q^{\dagger}=p$, then $r k_{n}^{\dagger} \rightarrow_{p} \infty$, whereas if $q^{\dagger}<p$, then $r k_{n}^{\dagger}$ converges in distribution to a nondegenerate random variable, see Lemma 19.4 .

By the CLT, for any sequence $\left\{\lambda_{n, h} \in \Lambda_{K C L R}: n \geq 1\right\}$,

$$
\begin{align*}
& n^{-1 / 2} \sum_{i=1}^{n}\binom{f_{i}^{*}}{\operatorname{vech}\left(f_{i}^{*} f_{i}^{* \prime}-E_{F_{n}} f_{i}^{*} f_{i}^{* \prime}\right)} \rightarrow_{d} \bar{L}_{h} \sim N\left(0^{d^{*}}, h_{10}\right), \text { where } \\
& \bar{L}_{h}=\left(\bar{L}_{h, 1}^{\prime}, \bar{L}_{h, 2}^{\prime}, \bar{L}_{h, 3}^{\prime}\right)^{\prime} \text { for } \bar{L}_{h, 1} \in R^{k}, \bar{L}_{h, 2} \in R^{k p}, \text { and } \bar{L}_{h, 3} \in R^{(p+1) k((p+1) k+1) / 2}( \tag{19.11}
\end{align*}
$$

and the CLT holds using the moment conditions in $\mathcal{F}_{K C L R}$. Note that by the definitions of $h_{4}:=$ $\lim E_{F_{n}} G_{i}$ and $h_{5}:=\lim E_{F_{n}}\left(g_{i}^{\prime}, \operatorname{vec}\left(G_{i}\right)^{\prime}\right)^{\prime}\left(g_{i}^{\prime}, \operatorname{vec}\left(G_{i}\right)^{\prime}\right)$, we have

$$
h_{10, f^{*}}=\left[\begin{array}{cc}
h_{5, g} & h_{5, g G}  \tag{19.12}\\
h_{5, G g} & h_{5, G}-\operatorname{vec}\left(h_{4}\right) \operatorname{vec}\left(h_{4}\right)^{\prime}
\end{array}\right], \text { where } h_{5}=\left[\begin{array}{cc}
h_{5, g} & h_{5, g G} \\
h_{5, G g} & h_{5, G}
\end{array}\right]
$$

for $h_{5, g} \in R^{k \times k}, h_{5, G g} \in R^{k p \times k}$, and $h_{5, G} \in R^{k p \times k p}$.
We now provide new, but distributionally equivalent, definitions of $\bar{g}_{h}$ and $\bar{D}_{h}$ :

$$
\begin{equation*}
\bar{g}_{h}:=\bar{L}_{h, 1} \text { and } \operatorname{vec}\left(\bar{D}_{h}\right):=\bar{L}_{h, 2}-h_{5, G g} h_{5, g}^{-1} \bar{L}_{h, 1} . \tag{19.13}
\end{equation*}
$$

These definitions are distributionally equivalent to the previous definitions of $\bar{g}_{h}$ and $\bar{D}_{h}$ given in Lemma 10.2 , because by either set of definitions $\bar{g}_{h}$ and $\operatorname{vec}\left(\bar{D}_{h}\right)$ are independent mean zero random vectors with variance matrices $h_{5, g}$ and $\Phi_{h}^{v e c\left(G_{i}\right)}\left(=h_{5, G}-\operatorname{vec}\left(h_{4}\right) v e c\left(h_{4}\right)^{\prime}-h_{5, G g} h_{5, g}^{-1} h_{5, G g}^{\prime}\right)$, respectively, where $\Phi_{h}^{\operatorname{vec}\left(G_{i}\right)}$ is defined in 10.15 and is pd (because $\Phi_{h}^{v e c\left(G_{i}\right)}=\lim \Phi_{F_{n}}^{v e c\left(G_{i}\right)}$ and $\lambda_{\min }\left(\Phi_{F_{n}}^{v e c}\left(G_{i}\right)\right.$ is bounded away from zero by its definition in 10.15 and the $\lambda_{\min }(\cdot)$ condition in $\left.\mathcal{F}_{\text {KCLR }}\right)$.

Define

$$
\bar{D}_{h}^{\dagger}:=\sum_{j=1}^{p}\left(M_{1 j h} \bar{D}_{j h}, \ldots, M_{p j h} \bar{D}_{j h}\right) \in R^{k \times p}, \text { where }\left[\begin{array}{ccc}
M_{11 h} & \cdots & M_{1 p h}  \tag{19.14}\\
\vdots & \ddots & \vdots \\
M_{p 1 h} & \cdots & M_{p p h}
\end{array}\right]:=\left(\Phi_{h}^{v e c\left(G_{i}\right)}\right)^{-1 / 2},
$$

$\bar{D}_{h}=\left(\bar{D}_{1 h}, \ldots, \bar{D}_{p h}\right)$, and $\bar{D}_{h}$ is defined in 19.13. Define

$$
\begin{align*}
\bar{\Delta}_{h}^{\dagger} & =\left(\bar{\Delta}_{h, q^{\dagger}}^{\dagger}, \bar{\Delta}_{h, p-q^{\dagger}}^{\dagger}\right) \in R^{k \times p}, \bar{\Delta}_{h, q^{\dagger}}^{\dagger}:=h_{3, q^{\dagger}}^{\dagger} \in R^{k \times q^{\dagger}}, \text { and } \\
\bar{\Delta}_{h, p-q^{\dagger}}^{\dagger} & :=h_{3}^{\dagger} h_{1, p-q^{\dagger}}^{\dagger \diamond}+\bar{D}_{h}^{\dagger} h_{2, p-q^{\dagger}}^{\dagger} \in R^{k \times\left(p-q^{\dagger}\right)} \tag{19.15}
\end{align*}
$$

Let $a(\cdot)$ be the function from $R^{d^{*}}$ to $R^{k p(k p+1) / 2}$ that maps

$$
\begin{align*}
& n^{-1} \sum_{i=1}^{n}\binom{f_{i}^{*}}{\operatorname{vech}\left(f_{i}^{*} f_{i}^{* \prime}\right)} \text { into }  \tag{19.16}\\
& A_{n}:=\operatorname{vech}\left(\left(n^{-1} \sum_{i=1}^{n} \operatorname{vec}\left(G_{i}-E_{F_{n}} G_{i}\right) \operatorname{vec}\left(G_{i}-E_{F_{n}} G_{i}\right)^{\prime}-\widetilde{\Gamma}_{n} \widetilde{\Omega}_{n}^{-1} \widetilde{\Gamma}_{n}^{\prime}\right)^{-1 / 2}\right), \text { where } \\
& \widetilde{\Omega}_{n}:=n^{-1} \sum_{i=1}^{n} g_{i} g_{i}^{\prime} \in R^{k \times k} \text { and } \widetilde{\Gamma}_{n}:=n^{-1} \sum_{i=1}^{n} \operatorname{vec}\left(G_{i}-E_{F_{n}} G_{i}\right) g_{i}^{\prime} \in R^{p k \times k}
\end{align*}
$$

Note that $a(\cdot)$ does not depend on the $n^{-1} \sum_{i=1}^{n} f_{i}^{*}$ part of its argument. Also, $a(\cdot)$ is well defined and continuously partially differentiable at any value of its argument for which $n^{-1} \sum_{i=1}^{n} f_{i}^{*} f_{i}^{* \prime}$ is pd. (The function $a(\cdot)$ is well defined in this case because $n^{-1} \sum_{i=1}^{n} \operatorname{vec}\left(G_{i}-E_{F_{n}} G_{i}\right) v e c\left(G_{i}-\right.$ $\left.E_{F_{n}} G_{i}\right)^{\prime}-\widetilde{\Gamma}_{n} \widetilde{\Omega}_{n}^{-1} \widetilde{\Gamma}_{n}^{\prime}=\left(-\widetilde{\Gamma}_{n} \widetilde{\Omega}_{n}^{-1}, I_{p k}\right) n^{-1} \sum_{i=1}^{n} f_{i}^{*} f_{i}^{* \prime}\left(-\widetilde{\Gamma}_{n} \widetilde{\Omega}_{n}^{-1}, I_{p k}\right)^{\prime}$ and $\left(-\widetilde{\Gamma}_{n} \widetilde{\Omega}_{n}^{-1}, I_{p k}\right) \in R^{p k \times(p+1) k}$ has full row rank $p k$.) We define $\bar{A}_{h}$ as follows:

$$
\begin{align*}
& \bar{A}_{h} \text { denotes the }(k p)(k p+1) / 2 \times d^{*} \text { matrix of partial derivatives of } a(\cdot) \\
& \text { evaluated at }\left(0^{(p+1) k \prime}, \operatorname{vech}\left(h_{10, f^{*}}\right)^{\prime}\right)^{\prime} \tag{19.17}
\end{align*}
$$

where the latter vector is the limit of the mean vector of $\left(f_{i}^{* \prime}, v e c h\left(f_{i}^{*} f_{i}^{* \prime}\right)^{\prime}\right)^{\prime}$ under $\left\{\lambda_{n, h} \in \Lambda_{K C L R}\right.$ : $n \geq 1\}$.

Define

$$
\begin{equation*}
\bar{M}_{h}:=\operatorname{vech}_{k p, k p}^{-1}\left(\bar{A}_{h} \bar{L}_{h}\right) \in R^{k p \times k p} \tag{19.18}
\end{equation*}
$$

where $v e c h_{k p, k p}^{-1}(\cdot)$ denotes the inverse of the $\operatorname{vech}(\cdot)$ operator applied to symmetric $k p \times k p$ matrices.

Define

$$
\begin{align*}
& \bar{M}_{h}^{\dagger}:=\left(\bar{M}_{h, q^{\dagger}}^{\dagger}, \bar{M}_{h, p-q^{\dagger}}^{\dagger}\right):=\left(0^{k \times q^{\dagger}}, \bar{M}_{h, p-q^{\dagger}}^{\dagger}\right) \in R^{k \times p} \text {, where }  \tag{19.19}\\
& \bar{M}_{h, p-q^{\dagger}}^{\dagger}:=\sum_{j=1}^{p}\left(\bar{M}_{1 j h} h_{4, j}, \ldots, \bar{M}_{p j h} h_{4, j}\right) h_{2, p-q^{\dagger}}^{\dagger} \in R^{k \times\left(p-q^{\dagger}\right)}, \bar{M}_{h}=\left[\begin{array}{ccc}
\bar{M}_{11 h} & \cdots & \bar{M}_{1 p h} \\
\vdots & \ddots & \vdots \\
\bar{M}_{p 1 h} & \cdots & \bar{M}_{p p h}
\end{array}\right],
\end{align*}
$$

and $h_{4}=\left(h_{4,1}, \ldots, h_{4, p}\right) \in R^{k \times p}$.
Below (in Lemma 19.4, we show that the asymptotic distribution of $r k_{n}^{\dagger}$ under sequences $\left\{\lambda_{n, h} \in \Lambda_{K C L R}: n \geq 1\right\}$ with $q^{\dagger}<p$ is given by

$$
\begin{equation*}
r_{h}\left(\bar{D}_{h}, \bar{M}_{h}\right):=\lambda_{\min }\left(\left(\bar{\Delta}_{h, p-q^{\dagger}}^{\dagger}+\bar{M}_{h, p-q^{\dagger}}^{\dagger}\right)^{\prime} h_{3, k-q^{\dagger}}^{\dagger} h_{3, k-q^{\dagger}}^{\dagger \prime}\left(\bar{\Delta}_{h, p-q^{\dagger}}^{\dagger}+\bar{M}_{h, p-q^{\dagger}}^{\dagger}\right)\right), \tag{19.20}
\end{equation*}
$$

where $\bar{\Delta}_{h, p-q^{\dagger}}^{\dagger}$ is a nonrandom function of $\bar{D}_{h}$ by 19.14 and 19.15 and $\bar{M}_{h, p-q^{\dagger}}^{\dagger}$ is a nonrandom function of $\bar{M}_{h}$ by 19.19. For sequences $\left\{\lambda_{n, h} \in \Lambda_{K C L R}: n \geq 1\right\}$ with $q^{\dagger}=p$, we show that $r k_{n} \rightarrow p \bar{r}_{h}:=\infty$.

We define $\bar{\Delta}_{h}$, as in 10.17, as follows:

$$
\begin{align*}
& \bar{\Delta}_{h}=\left(\bar{\Delta}_{h, q}, \bar{\Delta}_{h, p-q}\right) \in R^{k \times p}, \bar{\Delta}_{h, q}:=h_{3, q}, \text { and } \bar{\Delta}_{h, p-q}:=h_{3} h_{1, p-q}^{\diamond}+h_{7} \bar{D}_{h} h_{8} h_{2, p-q}, \text { where } \\
& h_{2}=\left(h_{2, q}, h_{2, p-q}\right), h_{3}=\left(h_{3, q}, h_{3, k-q}\right), h_{1, p-q}^{\diamond}:=\left[\begin{array}{c}
0^{q \times(p-q)} \\
\operatorname{Diag}\left\{h_{1, q+1}, \ldots, h_{1, p}\right\} \\
0^{(k-p) \times(p-q)}
\end{array}\right] \in R^{k \times(p-q)} . \tag{19.21}
\end{align*}
$$

In the present case, $h_{7}=h_{5, g}^{-1 / 2}$ and $h_{8}=I_{p}$ because the $C L R_{n}$ statistic depends on $\widehat{D}_{n}$ through $\widehat{\Omega}_{n}^{-1 / 2} \widehat{D}_{n}$, which appears in the $L M_{n}$ statistic. (The $C L R_{n}$ statistic also depends on $\widehat{D}_{n}$ through the rank statistic.) This means that Assumption WU for the parameter space $\Lambda_{K C L R}$ (defined in Section 10.4 holds with $\widehat{W}_{n}=\widehat{\Omega}_{n}^{-1 / 2}, \widehat{U}_{n}=I_{p}, h_{7}=h_{5, g}^{-1 / 2}$, and $h_{8}=I_{p}$. Thus, the distribution of $\bar{\Delta}_{h}$ depends on $\bar{D}_{h}, q$, and $h_{s}$ for $s=1,2,3,5$.

Below (in Lemma 19.5), we show that the asymptotic distribution of the $C L R_{n}$ statistic under sequences $\left\{\lambda_{n, h} \in \Lambda_{K C L R}: n \geq 1\right\}$ with $q^{\dagger}<p$ is given by

$$
\begin{align*}
\overline{C L R}_{h} & :=\frac{1}{2}\left(\overline{L M}_{h}+\bar{J}_{h}-\bar{r}_{h}+\sqrt{\left(\overline{L M}_{h}+\bar{J}_{h}-\bar{r}_{h}\right)^{2}+4 \overline{L M} \bar{r}_{h}}\right), \text { where } \\
\overline{L M}_{h} & :=\bar{v}_{h}^{\prime} \bar{v}_{h} \sim \chi_{p}^{2}, \bar{v}_{h}:=P_{\bar{\Delta}_{h}} h_{5, g}^{-1 / 2} \bar{g}_{h}, \bar{J}_{h}:=\bar{g}_{h}^{\prime} h_{5, g}^{-1 / 2} M_{\bar{\Delta}_{h}} h_{5, g}^{-1 / 2} \bar{g}_{h} \sim \chi_{k-p}^{2}, \text { and } \\
\bar{r}_{h} & :=r_{h}\left(\bar{D}_{h}, \bar{M}_{h}\right) . \tag{19.22}
\end{align*}
$$

The quantities ( $\bar{g}_{h}, \bar{D}_{h}, \bar{M}_{h}$ ) are specified in 19.13) and 19.18 (and ( $\bar{g}_{h}, \bar{D}_{h}$ ) are the same as in Lemma 10.2. The definitions of $\bar{v}_{h}, \overline{L M}_{h}, \bar{J}_{h}$, and $\overline{C L R}_{h}$ in 19.22 are the same as in 11.1, (11.2), (12.6), and 12.7), respectively.

Conditional on $\bar{D}_{h}, \overline{L M}_{h}$ and $\bar{J}_{h}$ are independent and distributed as $\chi_{p}^{2}$ and $\chi_{k-p}^{2}$, respectively (see the paragraph following 12.6). For sequences $\left\{\lambda_{n, h} \in \Lambda_{K C L R}: n \geq 1\right\}$ with $q^{\dagger}=p$, we show that the asymptotic distribution of the $C L R_{n}$ statistic is $\overline{C L R}_{h}:=\overline{L M}_{h}:=\bar{v}_{h}^{\prime} \bar{v}_{h} \sim \chi_{p}^{2}$, where $\bar{v}_{h}:=P_{\bar{\Delta}_{h}} h_{5, g}^{-1 / 2} \bar{g}_{h}$.

The critical value function $c(1-\alpha, r)$ is defined in (5.2) for $0 \leq r<\infty$. For $r=\infty$, we define $c(1-\alpha, r)$ to be the $1-\alpha$ quantile of the $\chi_{p}^{2}$ distribution.

Now we state the asymptotic size of Kleibergen's CLR test based on Robin and Smith (2000) statistic with $\widetilde{V}_{D n}$ defined in 5.3.

Theorem 19.1 Let the parameter space for $F$ be $\mathcal{F}_{K C L R}$. Suppose the variance matrix estimator $\widetilde{V}_{D n}$ employed by the rank statistic rkn (defined in 19.3) is defined by (5.3). Then, the asymptotic size of Kleibergen's CLR test based on the rank statistic r $k_{n}^{\dagger}$ is

$$
\text { AsySz }=\max \left\{\alpha, \sup _{h \in H} P\left(\overline{C L R}_{h}>c\left(1-\alpha, \bar{r}_{h}\right)\right)\right\}
$$

provided $P\left(\overline{C L R}_{h}=c\left(1-\alpha, \bar{r}_{h}\right)\right)=0$ for all $h \in H$.
Comments: (i) The proviso in Theorem 19.1 is a continuity condition on the distribution function of $\overline{C L R}_{h}-c\left(1-\alpha, \bar{r}_{h}\right)$ at zero. If the proviso in Theorem 19.1 does not hold, then the following weaker conclusion holds:

$$
\begin{equation*}
A s y S z \tag{19.23}
\end{equation*}
$$

$$
\in\left[\max \left\{\alpha, \sup _{h \in H} P\left(\overline{C L R}_{h}>c\left(1-\alpha, \bar{r}_{h}\right)\right)\right\}, \max \left\{\alpha, \sup _{h \in H} \lim _{x \uparrow 0} P\left(\overline{C L R}_{h}>c\left(1-\alpha, \bar{r}_{h}\right)+x\right)\right\}\right] .
$$

(ii) Conditional on ( $\bar{D}_{h}, \bar{M}_{h}$ ), $\bar{g}_{h}$ has a multivariate normal distribution a.s. (because ( $\bar{g}_{h}, \bar{D}_{h}$, $\bar{M}_{h}$ ) has a multivariate normal distribution unconditionally). Note that $\bar{g}_{h}$ is independent of $\bar{D}_{h}$. The proviso in Theorem 19.1 holds whenever $\bar{g}_{h}$ has a non-zero variance matrix conditional on $\left(\bar{D}_{h}, \bar{M}_{h}\right)$ a.s. for all $h \in H$. This holds because (a) $P\left(\overline{C L R}_{h}=c\left(1-\alpha, \bar{r}_{h}\right)\right)=E_{\left(\bar{D}_{h}, \bar{M}_{h}\right)} P\left(\overline{C L R}_{h}=\right.$ $\left.c\left(1-\alpha, \bar{r}_{h}\right) \mid \bar{D}_{h}, \bar{M}_{h}\right)$ by the law of iterated expectations, (b) some calculations show that $\overline{C L R}_{h}=$ $c\left(1-\alpha, \bar{r}_{h}\right)$ iff $\left(\bar{r}_{h}+c\right) \overline{L M}_{h}=-c \bar{J}_{h}+c^{2}+c \bar{r}_{h}$ iff $\bar{X}_{h}^{\prime} \bar{X}_{h}=c^{2}+c \bar{r}_{h}$, where $c:=c\left(1-\alpha, \bar{r}_{h}\right)$ and $\bar{X}_{h}:=\left(\left(\bar{r}_{h}+c\right)^{1 / 2}\left(P_{\bar{\Delta}_{h}} h_{5, g}^{-1 / 2} \bar{g}_{h}\right)^{\prime}, c^{1 / 2}\left(M_{\bar{\Delta}_{h}} h_{5, g}^{-1 / 2} \bar{g}_{h}\right)^{\prime}\right)^{\prime}$ using 19.22, (c) $P_{\bar{\Delta}_{h}}+M_{\bar{\Delta}_{h}}=I_{k}$ and $P_{\bar{\Delta}_{h}} M_{\bar{\Delta}_{h}}=0^{k \times k}$, and (d) conditional on ( $\bar{D}_{h}, \bar{M}_{h}$ ), $\bar{r}_{h}, c$, and $\bar{\Delta}_{h}$ are constants.
(iii) When $p=1$, the formula for AsySz in Theorem 19.1 reduces to $\alpha$ and the proviso holds automatically. That is, Kleibergen's CLR test has correct asymptotic size when $p=1$. This holds because when $p=1$ the quantity $\bar{M}_{h}^{\dagger}$ in 19.19) equals $0^{k \times p}$ by Comment (ii) to Theorem 19.3 below. This implies that $r_{h}\left(\bar{D}_{h}, \bar{M}_{h}\right)$ in 19.20) does not depend on $\bar{M}_{h}$. Given this, the proof that $P\left(\overline{C L R}_{h}>c\left(1-\alpha, \bar{r}_{h}\right)=\alpha\right.$ for all $h \in H$ and that the proviso holds is the same as in 12.9)-12.10) in the proof of Theorem 12.1.
(iv) Theorem 19.1 is proved by showing that it is a special case of Theorem 19.6 below, which is similar but applies not to $\widetilde{V}_{D n}$ defined in $\sqrt{5.3}$ ), but to an arbitrary estimator $\widetilde{V}_{D n}$ (of the asymptotic variance $\Phi_{h}^{\operatorname{vec}\left(G_{i}\right)}$ of $\left.n^{1 / 2} \operatorname{vec}\left(\widehat{D}_{n}-E_{F_{n}} G_{i}\right)\right)$ that satisfies an Assumption VD (which is stated below). Lemma 19.2 below shows that the estimator $\widetilde{V}_{D n}$ defined in 5.3) satisfies Assumption VD.
(v) A CS version of Theorem 19.1 holds with the parameter space $\mathcal{F}_{\Theta, K C L R}$ in place of $\mathcal{F}_{K C L R}$, where $\mathcal{F}_{\Theta, K C L R}:=\left\{\left(F, \theta_{0}\right): F \in \mathcal{F}_{K C L R}\left(\theta_{0}\right), \theta_{0} \in \Theta\right\}$ and $\mathcal{F}_{K C L R}\left(\theta_{0}\right)$ is the set $\mathcal{F}_{K C L R}$ defined in (19.5) with its dependence on $\theta_{0}$ made explicit. The proof of this CS result is as outlined in the Comment to Proposition 10.1. For the CS result, the $h$ index and its parameter space $H$ are as defined above, but $h$ also includes $\theta_{0}$ as a subvector, and $H$ allows this subvector to range over $\Theta$.

### 19.3 Simulation Results

In this section, for a particular linear IV regression model, we simulate (i) correlations between $\bar{M}_{h, p-q^{\dagger}}^{\dagger}$ (defined in 19.19) and $\bar{g}_{h}$ and (ii) some asymptotic null rejection probabilities (NRP's) of Kleibergen's CLR test that uses Jacobian-variance weighting and employs the Robin and Smith (2000) rank statistic. The model has $p=2$ rhs endogenous variables and $k=15$ IV's. The model is

$$
\begin{equation*}
y_{1 i}=Y_{2 i}^{\prime} \theta_{0}+u_{i} \text { and } Y_{2 i}=\pi^{\prime} Z_{i}+V_{2 i}, \tag{19.24}
\end{equation*}
$$

where $y_{1 i}, u_{i} \in R, Y_{2 i}, V_{2 i} \in R^{2}, \theta_{0} \in R^{2}, Z_{i}=\left(Z_{i 1}, \ldots, Z_{i k}\right)^{\prime} \in R^{k}$, and $\pi \in R^{k \times 2}$. We take $Z_{i j} \sim \chi_{1}^{2}-1$ i.i.d. for $j=1, \ldots, k, u_{i} \sim\left\|Z_{i}\right\| \widetilde{u}_{i},\left(\widetilde{u}_{i}, V_{2 i}^{\prime}\right)^{\prime} \sim N\left(0, \Sigma_{\rho}\right),\left(\widetilde{u}_{i}, V_{2 i}^{\prime}\right)^{\prime}$ independent of $Z_{i}$, and $\Sigma_{\rho} \in R^{3 \times 3}$ with diagonal elements 1 and off-diagonal elements $\rho$. This data generating process (DGP) involves an asymmetric distribution for $Z_{i j}$ and conditional heteroskedasticity in $u_{i}$. We take $\pi=\pi_{n}=\left(e_{1}, e_{2} c n^{-1 / 2}\right)$, where $e_{j} \in R^{k}$ denote the $j$ th coordinate vector for $j=1,2$. We consider integer values of the constant $c$ in $[0,30], \rho=.5, \theta_{0}=(0,0)^{\prime}$, and nominal size $5 \%$ for the tests. We also experimented with additional DGPs for $\left(u_{i}, V_{2 i}^{\prime}, Z_{i}^{\prime}\right)^{\prime}$ and $k \in\{5,10\}$ and nominal size of $1 \%$ but no important additional insights were gained from these simulations.

In this model, we have $g_{i}=Z_{i} u_{i}$ and $G_{i}=-Z_{i} Y_{2 i}^{\prime}$. Furthermore, $h_{1,1}=\infty$ and $h_{1,2}$ is a finite nonnegative number that depends on $c$. The quantities $h_{1, j}^{\dagger}$ for $j=1,2$ (defined just below
(19.10) are not available in closed form, so we simulate them using a very large value of $n$, viz., $n=2,000,000$. We use $4,000,000$ simulation repetitions to compute the correlations between the $j$ th elements of $\bar{M}_{h, p-q^{\dagger}}^{\dagger}$ and $\bar{g}_{h}$ for $j=1, \ldots, k$ and the asymptotic NRP's of the CLR test. To conserve space we do not report the correlations between the $j$ th and $k$ th elements of these vectors for $j \neq k$. The data-dependent critical values for the test are computed using a look-up table that gives the critical values for each fixed value $r$ of the rank statistic in a grid from 0 to 10,000 with a step size of $.005, .05$, and 1 for $r \in[0,100],[100,1000]$, and [1000, 10000], respectively. These critical values are computed using $4,000,000$ simulation repetitions. Note that for $p=2$, the dimension $d^{*}:=(p+1) k+(p+1) k((p+1) k+1) / 2$ in 19.8 equals 135,495 , and 1080 , for $k=5,10,15$, respectively, and simulation with 4 million repetitions becomes computationally involved for large $k$.
(i) The simulations provide evidence for the findings given in Theorem 5.1 that $\bar{M}_{h, p-q^{\dagger}}^{\dagger}$ (the second column of $\bar{M}_{h}^{\dagger} \in R^{k \times 2}$ ) and $\bar{g}_{h}$ are correlated asymptotically in some models under some sequences of distributions. For example, when $k=15$ the simulated correlations between the $j$ th elements of $\bar{M}_{h, p-q^{\dagger}}^{\dagger}$ and $\bar{g}_{h}$ for $j=1,8,15$ take on the values $.32, .11$, and -.06 , respectively, for all values $c \in[0,30]$. In consequence, it is not possible to show the Jacobian-variance weighted CLR test has correct asymptotic size via a conditioning argument that relies on the independence of $\bar{\Delta}_{h, p-q^{\dagger}}^{\dagger}+\bar{M}_{h, p-q^{\dagger}}^{\dagger}$ and $\bar{g}_{h}$.
(ii) Next, we report the asymptotic NRP results for Kleibergen's CLR test that uses Jacobianvariance weighting and the Robin and Smith (2000) rank statistic. The asymptotic NRP's are found to be between $4.99 \%$ and $5.11 \%$ for the values of $c$ considered. These values are close to the nominal size of $5.00 \%$. Whether the difference is due to simulation noise or not is not clear. The simulation standard error based on the formula $100 *(\alpha(1-\alpha) / \text { reps })^{1 / 2}$, where reps $=4,000,000$ is the number of simulation repetitions, is .01 . However, this formula does not take into account simulation error from the computation of the critical values and from error in approximation of $h_{1, j}^{\dagger}$. For comparison, we also simulated the asymptotic NRP of the LM test (that has asymptotic size equal to nominal size) and find them to be between $5.01 \%$ and $5.02 \%$ for the values of $c$ considered.

We conclude that, for the model and error distribution considered, the asymptotic NRP's of Kleibergen's CLR test with Jacobian-variance weighting is quite close to its nominal size. This occurs even though there are non-negligible correlations between $\bar{M}_{h, p-q^{\dagger}}^{\dagger}$ and $\bar{g}_{h}$. Whether this occurs for all parameters and distributions in the linear IV model, and whether it occurs in other moment condition model, is an open question. It appears to be a question that can only be answered on a case by case basis.

### 19.4 Asymptotic Size of Kleibergen's CLR Test for General $\widetilde{\mathrm{V}}_{\text {Dn }}$ Estimators

In this section, we determine the asymptotic size of Kleibergen's CLR test (defined in Section 5) using the Robin and Smith (2000) rank statistic based on a general "Jacobian-variance" estimator $\widetilde{V}_{D n}\left(=\widetilde{V}_{D n}\left(\theta_{0}\right)\right)$ that satisfies the following Assumption VD.

The first two results of this section, viz., Lemma 19.2 and Theorem 19.3 , combine to establish Theorem 5.1, see Comment (i) to Theorem 19.3. The first and last results of this section, viz., Lemma 19.2 and Theorem 19.6, combine to prove Theorem 19.1 .

The proofs of the results in this section are given in Section 19.6.
Assumption VD: For any sequence $\left\{\lambda_{n, h} \in \Lambda_{K C L R}: n \geq 1\right\}$, the estimator $\widetilde{V}_{D n}$ is such that $n^{1 / 2}\left(\widetilde{M}_{n}-M_{F_{n}}\right) \rightarrow{ }_{d} \bar{M}_{h}$ for some random matrix $\bar{M}_{h} \in R^{k p \times k p}$ (where $\widetilde{M}_{n}=\widetilde{V}_{D n}^{-1 / 2}$ and $M_{F_{n}}$ is defined in 19.6), the convergence is joint with

$$
n^{1 / 2}\binom{\widehat{g}_{n}}{\operatorname{vec}\left(\widehat{D}_{n}-E_{F_{n}} G_{i}\right)} \rightarrow_{d}\binom{\bar{g}_{h}}{\operatorname{vec}\left(\bar{D}_{h}\right)} \sim N\left(0^{(p+1) k},\left(\begin{array}{cc}
h_{5, g} & 0^{k \times p k}  \tag{19.25}\\
0^{p k \times k} & \Phi_{h}^{v e c\left(G_{i}\right)}
\end{array}\right)\right)
$$

and ( $\bar{g}_{h}, \bar{D}_{h}, \bar{M}_{h}$ ) has a mean zero multivariate normal distribution with pd variance matrix. The same condition holds for any subsequence $\left\{w_{n}\right\}$ and any sequence $\left\{\lambda_{w_{n}, h} \in \Lambda_{K C L R}: n \geq 1\right\}$ with $w_{n}$ in place of $n$ throughout.

Note that the convergence in 19.25 holds by Lemma 10.2 .
The following lemma verifies Assumption VD for the estimator $\widetilde{V}_{D n}$ defined in 5.3 .
Lemma 19.2 The estimator $\widetilde{V}_{D n}$ defined in 5.3. satisfies Assumption VD. Specifically, $n^{1 / 2}\left(\widehat{g}_{n}, \widehat{D}_{n}-E_{F_{n}} G_{i}, \widetilde{M}_{n}-M_{F_{n}}\right) \rightarrow_{d}\left(\bar{g}_{h}, \bar{D}_{h}, \bar{M}_{h}\right)$, where $\widetilde{M}_{n}:=\widetilde{V}_{D n}^{-1 / 2}, M_{F_{n}}:=\left(\Phi_{F_{n}}^{v e c\left(G_{i}\right)}\right)^{-1 / 2}$, and ( $\bar{g}_{h}, \bar{D}_{h}, \bar{M}_{h}$ ) has a mean zero multivariate normal distribution defined by 19.11 and 19.13)(19.18) with pd variance matrix.

Comment: As stated in the paragraph containing 19.21, $\widehat{D}_{n}$ is defined in Lemma 19.2 and Theorem 19.3 below with $\widehat{W}_{n}=\widehat{\Omega}_{n}^{-1 / 2}$ and $\widehat{U}_{n}=I_{p}$.

Define

$$
\begin{equation*}
S_{n}^{\dagger}:=\operatorname{Diag}\left\{\left(n^{1 / 2} \tau_{1 F_{n}}^{\dagger}\right)^{-1}, \ldots,\left(n^{1 / 2} \tau_{q F_{n}}^{\dagger}\right)^{-1}, 1, \ldots, 1\right\} \in R^{p \times p} \text { and } T_{n}^{\dagger}:=B_{n}^{\dagger} S_{n}^{\dagger}, \tag{19.26}
\end{equation*}
$$

where $B_{n}^{\dagger}$ is defined in (19.7).
The asymptotic distribution of $n^{1 / 2} \widehat{D}_{n}^{\dagger} T_{n}^{\dagger}$ is given in the following theorem.

Theorem 19.3 Suppose Assumption VD holds. For all sequences $\left\{\lambda_{n, h} \in \Lambda_{K C L R}: n \geq 1\right\}$, $n^{1 / 2}\left(\widehat{g}_{n}, \widehat{D}_{n}-E_{F_{n}} G_{i}, \widehat{D}_{n}^{\dagger} T_{n}^{\dagger}\right) \rightarrow_{d}\left(\bar{g}_{h}, \bar{D}_{h}, \bar{\Delta}_{h}^{\dagger}+\bar{M}_{h}^{\dagger}\right)$, where $\bar{\Delta}_{h}^{\dagger}$ is a nonrandom affine function of $\bar{D}_{h}$ defined in 19.14 and 19.15, $\bar{M}_{h}^{\dagger}$ is a nonrandom linear (i.e., affine and homogeneous of degree one) function of $\bar{M}_{h}$ defined in 19.19), $\left(\bar{g}_{h}, \bar{D}_{h}, \bar{M}_{h}\right)$ has a mean zero multivariate normal distribution, and $\bar{g}_{h}$ and $\bar{D}_{h}$ are independent. Under all subsequences $\left\{w_{n}\right\}$ and all sequences $\left\{\lambda_{w_{n}, h} \in \Lambda_{K C L R}: n \geq 1\right\}$, the same result holds with $n$ replaced with $w_{n}$.
Comments: (i) Note that the random variables $\left(\bar{g}_{h}, \bar{\Delta}_{h}^{\dagger}, \bar{M}_{h}^{\dagger}\right)$ in Theorem 5.1 have a multivariate normal distribution whose mean and variance matrix depend on $\lim \operatorname{Var}_{F_{n}}\left(\left(f_{i}^{* \prime}, \operatorname{vec}\left(f_{i}^{*} f_{i}^{* \prime}\right)^{\prime}\right)\right.$ and on the limits of certain functions of $E_{F_{n}} G_{i}$ by 19.11)-19.19). This, Lemma 19.2, and Theorem 19.3 combine to prove Theorem 5.1 of AG1.
(ii) From 19.19, $\bar{M}_{h}^{\dagger}=0^{k \times p}$ if $p=1$ (because $q^{\dagger}=0$ implies $q=0$ which, in turn, implies $h_{4}=0^{k}$ and $q^{\dagger}=1$ implies $\bar{M}_{h, p-q^{\dagger}}^{\dagger}$ has no columns). For $p \geq 2, \bar{M}_{h}^{\dagger}=0^{k \times p}$ if $p=q^{\dagger}$ (because $\bar{M}_{h, p-q^{\dagger}}^{\dagger}$ has no columns) or if $h_{4, j}=0^{k}$ for all $j \leq p$. The former holds if the singular values $\left(\tau_{1 F_{n}}, \ldots, \tau_{p F_{n}}\right)$ of $D_{F_{n}}^{\dagger}$ satisfy $n^{1 / 2} \tau_{j F_{n}} \rightarrow \infty$ for all $j \leq p$ (i.e., all parameters are strongly or semi-strongly identified). The latter occurs if $E_{F_{n}} G_{i} \rightarrow 0^{k \times p}$ (i.e., all parameters are either weakly identified in the standard sense or semi-strongly identified). These two condition fail to hold when one or more parameters are strongly identified and one or more parameters are weakly identified or jointly weakly identified.
(iii) For example, when $p=2$ the conditions in Comment (ii) (under which $\bar{M}_{h}^{\dagger}=0^{k \times p}$ ) fail to hold if $E_{F_{n}} G_{i 1} \neq 0^{k}$ does not depend on $n$ and $n^{1 / 2} E_{F_{n}} G_{i 2} \rightarrow c$ for some $c \in R^{k}$.

The following lemma establishes the asymptotic distribution of $r k_{n}^{\dagger}$.
Lemma 19.4 Let the parameter space for $F$ be $\mathcal{F}_{\text {KCLR }}$. Suppose the variance matrix estimator $\widetilde{V}_{D n}$ employed by the rank statistic $r k_{n}^{\dagger}$ (defined in 19.3) satisfies Assumption VD. Then, under all sequences $\left\{\lambda_{n, h} \in \Lambda_{K C L R}: n \geq 1\right\}$,
(a) $r k_{n}^{\dagger}:=\widehat{\kappa}_{p n}^{\dagger} \rightarrow_{p} \infty$ if $q^{\dagger}=p$,
(b) $r k_{n}^{\dagger}:=\widehat{\kappa}_{p n}^{\dagger} \rightarrow{ }_{d} r_{h}\left(\bar{D}_{h}, \bar{M}_{h}\right)$ if $q^{\dagger}<p$, where $r_{h}\left(\bar{D}_{h}, \bar{M}_{h}\right)$ is defined in 19.20 using 19.19) with $\bar{M}_{h}$ defined in Assumption VD (rather than in 19.18),
(c) $\widehat{\kappa}_{j n}^{\dagger} \rightarrow_{p} \infty$ for all $j \leq q^{\dagger}$,
(d) the (ordered) vector of the smallest $p-q^{\dagger}$ singular values of $n^{1 / 2} \widehat{D}_{n}^{\dagger}$, i.e., $\left(\left(\widehat{\kappa}_{\left(q^{\dagger}+1\right) n}^{\dagger}\right)^{1 / 2}, \ldots\right.$, $\left.\left(\widehat{\kappa}_{p n}^{\dagger}\right)^{1 / 2}\right)^{\prime}$, converges in distribution to the (ordered) $p-q^{\dagger}$ vector of the singular values of

[^2]$h_{3, k-q^{\dagger}}^{\dagger \prime}\left(\bar{\Delta}_{h, p-q^{\dagger}}^{\dagger}+\bar{M}_{h, p-q^{\dagger}}^{\dagger}\right) \in R^{\left(k-q^{\dagger}\right) \times\left(p-q^{\dagger}\right)}$, where $\bar{M}_{h, p-q^{\dagger}}^{\dagger}$ is defined in 19.19 with $\bar{M}_{h}$ defined in Assumption VD (rather than in 19.18)),
(e) the convergence in parts (a)-(d) holds jointly with the convergence in Theorem 19.3, and
(f) under all subsequences $\left\{w_{n}\right\}$ and all sequences $\left\{\lambda_{w_{n}, h} \in \Lambda_{K C L R}: n \geq 1\right\}$, parts (a)-(e) hold with $n$ replaced with $w_{n}$.

The following lemma gives the joint asymptotic distribution of $C L R_{n}$ and $r k_{n}^{\dagger}$ and the asymptotic null rejection probabilities of Kleibergen's CLR test.

Lemma 19.5 Let the parameter space for $F$ be $\mathcal{F}_{K C L R}$. Suppose the variance matrix estimator $\widetilde{V}_{D n}$ employed by the rank statistic $r k_{n}^{\dagger}$ (defined in 19.3) satisfies Assumption VD. Then, under all sequences $\left\{\lambda_{n, h} \in \Lambda_{K C L R}: n \geq 1\right\}$,
(a) $C L R_{n}=L M_{n}+o_{p}(1) \rightarrow_{d} \chi_{p}^{2}$ and $r k_{n}^{\dagger} \rightarrow_{p} \infty$ if $q^{\dagger}=p$,
(b) $\lim _{n \rightarrow \infty} P\left(C L R_{n}>c\left(1-\alpha, r k_{n}^{\dagger}\right)\right)=\alpha$ if $q^{\dagger}=p$,
(c) $\left(C L R_{n}, r k_{n}^{\dagger}\right) \rightarrow_{d}\left(\overline{C L R}_{h}, \bar{r}_{h}\right)$ if $q^{\dagger}<p$, and
(d) $\lim _{n \rightarrow \infty} P\left(C L R_{n}>c\left(1-\alpha, r k_{n}^{\dagger}\right)\right)=P\left(\overline{C L R}_{h}>c\left(1-\alpha, \bar{r}_{h}\right)\right)$ if $q^{\dagger}<p$, provided $P\left(\overline{C L R}_{h}=c\left(1-\alpha, \bar{r}_{h}\right)\right)=0$.

Under all subsequences $\left\{w_{n}\right\}$ and all sequences $\left\{\lambda_{w_{n}, h} \in \Lambda_{K C L R} \geq 1\right\}$, parts (a)-(d) hold with $n$ replaced with $w_{n}$.

Comments: (i) The CLR critical value function $c(1-\alpha, r)$ is the $1-\alpha$ quantile of $c l r(r)$. By definition,

$$
\begin{equation*}
c l r(r):=\frac{1}{2}\left(\chi_{p}^{2}+\chi_{k-p}^{2}-r+\sqrt{\left(\chi_{p}^{2}+\chi_{k-p}^{2}-r\right)^{2}+4 \chi_{p}^{2} r}\right), \tag{19.27}
\end{equation*}
$$

where the chi-square random variables $\chi_{p}^{2}$ and $\chi_{k-p}^{2}$ are independent. If $\bar{r}_{h}:=r_{h}\left(\bar{D}_{h}, \bar{M}_{h}\right)$ does not depend on $\bar{M}_{h}$, then, conditional on $\bar{D}_{h}, \bar{r}_{h}$ is a constant and $\overline{L M}_{h}$ and $\bar{J}_{h}$ are independent and distributed as $\chi_{p}^{2}$ and $\chi_{k-p}^{2}$ (see the paragraph following 12.6). In this case, even when $q^{\dagger}=p$,

$$
\begin{equation*}
P\left(\overline{C L R}_{h}>c\left(1-\alpha, \bar{r}_{h}\right)\right)=E_{\bar{D}_{h}} P\left(\overline{C L R}_{h}>c\left(1-\alpha, \bar{r}_{h}\right) \mid \bar{D}_{h}\right)=\alpha, \tag{19.28}
\end{equation*}
$$

as desired, where the first equality holds by the law of iterated expectations and the second equality holds because $\bar{r}_{h}$ is a constant conditional on $\bar{D}_{h}$ and $c\left(1-\alpha, \bar{r}_{h}\right)$ is the $1-\alpha$ quantile of the conditional distribution of $\operatorname{clr}\left(\bar{r}_{h}\right)$ given $\bar{D}_{h}$, which equals that of $\overline{C L R} h$ given $\bar{D}_{h}$.
(ii) However, when $\bar{r}_{h}:=r_{h}\left(\bar{D}_{h}, \bar{M}_{h}\right)$ depends on $\bar{M}_{h}$, the distribution of $\bar{r}_{h}$ conditional on $\bar{D}_{h}$ is not a pointmass distribution. Rather, conditional on $\bar{D}_{h}, \bar{r}_{h}$ is a random variable that is not
independent of $\overline{L M}_{h}, \bar{J}_{h}$, and $\overline{C L R}_{h}$. In consequence, the second equality in 19.28) does not hold and the asymptotic null rejection probability of Kleibergen's CLR test may be larger or smaller than $\alpha$ depending upon the sequence $\left\{\lambda_{n, h} \in \Lambda_{K C L R}: n \geq 1\right\}$ (or $\left\{\lambda_{w_{n}, h} \in \Lambda_{K C L R}: n \geq 1\right\}$ ) when $q^{\dagger}<p$.

Next, we use Lemma 19.5 to provide an expression for the asymptotic size of Kleibergen's CLR test based on the Robin and Smith (2000) rank statistic with Jacobian-variance weighting.

Theorem 19.6 Let the parameter space for $F$ be $\mathcal{F}_{K C L R}$. Suppose the variance matrix estimator $\widetilde{V}_{D n}$ employed by the rank statistic $r k_{n}^{\dagger}$ (defined in 19.3) satisfies Assumption VD. Then, the asymptotic size of Kleibergen's CLR test based on rkn is

$$
\text { AsySz }=\max \left\{\alpha, \sup _{h \in H} P\left(\overline{C L R}_{h}>c\left(1-\alpha, \bar{r}_{h}\right)\right)\right\}
$$

provided $P\left(\overline{C L R}_{h}=c\left(1-\alpha, \bar{r}_{h}\right)\right)=0$ for all $h \in H$.
Comments: (i) Comment (i) to Theorem 19.1 also applies to Theorem 19.6 .
(ii) Theorem 19.6 and Lemma 19.2 combine to prove Theorem 19.1 .
(iii) A CS version of Theorem 19.6 holds with the parameter space $\mathcal{F}_{\Theta, K C L R}$ in place of $\mathcal{F}_{K C L R}$, see Comment (v) to Theorem 19.1 and the Comment to Proposition 10.1 .

### 19.5 Correct Asymptotic Size of Equally-Weighted CLR Tests Based on the Robin-Smith Rank Statistic

In this subsection, we consider equally-weighted CLR tests, a special case of which is considered in Section 6. By definition, an equally-weighted CLR test is a CLR test that is based on a $r k_{n}$ statistic that depends on $\widehat{D}_{n}$ only through $\widetilde{W}_{n} \widehat{D}_{n}$ for some general $k \times k$ weighting matrix $\widetilde{W}_{n}$. We show that such tests have correct asymptotic size when they are based on the rank statistic of Robin and Smith (2000) and employ a general weight matrix $\widetilde{W}_{n} \in R^{k \times k}$ that satisfies certain conditions. In contrast, the results in Section 6 consider the specific weight matrix $\widehat{\Omega}_{n}^{-1 / 2} \in R^{k \times k}$. The reason for considering these tests in this section is that the asymptotic results can be obtained as a relatively simple by-product of the results in Section 19.4. All that is required is a slight change in Assumption VD.

The rank statistic that we consider here is

$$
\begin{equation*}
r k_{n}^{\dagger}:=\lambda_{\min }\left(n \widehat{D}_{n}^{\prime} \widetilde{W}_{n}^{\prime} \widetilde{W}_{n} \widehat{D}_{n}\right) \tag{19.29}
\end{equation*}
$$

We replace Assumption VD in Section 19.4 by the following assumption.

Assumption W: For any sequence $\left\{\lambda_{n, h} \in \Lambda_{K C L R}: n \geq 1\right\}$, the random $k \times k$ weight matrix $\widetilde{W}_{n}$ is such that $n^{1 / 2}\left(\widetilde{W}_{n}-W_{F_{n}}^{\dagger}\right) \rightarrow_{d} \bar{W}_{h}$ for some non-random $k \times k$ matrices $\left\{W_{F_{n}}^{\dagger}: n \geq 1\right\}$ and some random $k \times k$ matrix $\bar{W}_{h} \in R^{k \times k}, W_{F_{n}}^{\dagger} \rightarrow W_{h}^{\dagger}$ for some nonrandom $\operatorname{pd} k \times k$ matrix $W_{h}^{\dagger}$, the convergence is joint with the convergence in (19.25), and ( $\bar{g}_{h}, \bar{D}_{h}, \bar{W}_{h}$ ) has a mean zero multivariate normal distribution with pd variance matrix. The same condition holds for any subsequence $\left\{w_{n}\right\}$ and any sequence $\left\{\lambda_{w_{n}, h} \in \Lambda_{K C L R}: n \geq 1\right\}$ with $w_{n}$ in place of $n$ throughout.

If one takes $\widetilde{M}_{n}\left(=\widetilde{V}_{D n}^{-1 / 2}\right)=I_{p} \otimes \widetilde{W}_{n}$ in Assumption VD, then $\widehat{D}_{n}^{\dagger}=\widetilde{W}_{n} \widehat{D}_{n}$ and the rank statistics in 19.3 and 19.29) are the same. Thus, Assumption W is analogous to Assumption VD with $\widetilde{M}_{n}=I_{p} \otimes \widetilde{W}_{n}$ and $M_{F_{n}}=I_{p} \otimes W_{F_{n}}^{\dagger}$. Note, however, that the latter matrix does not typically satisfy the condition in Assumption VD that $M_{F_{n}}$ is defined in 19.6, i.e., the condition that $M_{F_{n}}=\left(\Phi_{F_{n}}^{v e c\left(G_{i}\right)}\right)^{-1 / 2}$. Nevertheless, the results in Section 19.4 hold with Assumption VD replaced by Assumption W and with $M_{F}=I_{p} \otimes W_{F}^{\dagger}, D_{F}^{\dagger}=W_{F}^{\dagger} E_{F} G_{i}$, and $\bar{M}_{h}=I_{p} \otimes \bar{W}_{h}$. With these changes, $\bar{D}_{h}^{\dagger}=W_{h}^{\dagger} \bar{D}_{h}$ in 19.14 (because $\left(\Phi_{h}^{v e c\left(G_{i}\right)}\right)^{-1 / 2}$ is replaced by $\left.I_{p} \otimes W_{h}^{\dagger}\right), \bar{\Delta}_{h}^{\dagger}$ is defined as in 19.15 with $\bar{D}_{h}^{\dagger}$ as just given, and $\bar{M}_{h}^{\dagger}$ is defined as in 19.19 with $\bar{M}_{h, p-q^{\dagger}}^{\dagger}=\bar{W}_{h} h_{4} h_{2, p-q^{\dagger}}^{\dagger}$.

Below we show the key result that $\bar{M}_{h, p-q^{\dagger}}^{\dagger}=0^{k \times\left(p-q^{\dagger}\right)}$ for $r k_{n}^{\dagger}$ defined in 19.29. By 19.20, this implies that

$$
\begin{equation*}
r_{h}\left(\bar{D}_{h}, \bar{M}_{h}\right):=\lambda_{\min }\left(\left(\bar{\Delta}_{h, p-q^{\dagger}}^{\dagger}\right)_{3, k-q^{\dagger}}^{\dagger} h_{3, k-q^{\dagger}}^{\dagger \prime}\left(\bar{\Delta}_{h, p-q^{\dagger}}^{\dagger}\right)\right) \tag{19.30}
\end{equation*}
$$

when $q^{\dagger}<p$. Note that the rhs in 19.30 does not depend on $\bar{M}_{h}$ and, hence, is a function only of $\bar{D}_{h}$. That is, $r_{h}\left(\bar{D}_{h}, \bar{M}_{h}\right)=r_{h}\left(\bar{D}_{h}\right)$. Given that $r_{h}\left(\bar{D}_{h}, \bar{M}_{h}\right)$ does not depend on $\bar{M}_{h}$, Comment (i) to Lemma 19.5 implies that $P\left(\overline{C L R}_{h}>c\left(1-\alpha, \bar{r}_{h}\right)\right)=\alpha$ under all subsequences $\left\{w_{n}\right\}$ and all sequences $\left\{\lambda_{w_{n}, h} \in \Lambda_{K C L R}: n \geq 1\right\}$. This and Theorem 19.6 give the following result.

Corollary 19.7 Let the parameter space for $F$ be $\mathcal{F}_{K C L R}$. Suppose the rank statistic rk $k_{n}^{\dagger}$ (defined in 19.29) is based on a weight matrix $\widetilde{W}_{n}$ that satisfies Assumption W. Then, the asymptotic size of the corresponding equally-weighted version of Kleibergen's CLR test (defined in Section 5 with $\left.r k_{n}(\theta)=r k_{n}^{\dagger}\right)$ equals $\alpha$.

Comment: A CS version of Corollary 19.7 holds with the parameter space $\mathcal{F}_{\Theta, K C L R}$ in place of $\mathcal{F}_{K C L R}$, see Comment (v) to Theorem 19.1 and the Comment to Proposition 10.1.

Now, we establish that $\bar{M}_{h, p-q^{\dagger}}^{\dagger}\left(=\bar{W}_{h} h_{4} h_{2, p-q^{\dagger}}^{\dagger}\right)=0^{k \times\left(p-q^{\dagger}\right)}$. We have

$$
\begin{equation*}
W_{h}^{\dagger} h_{4}:=\lim W_{F_{n}}^{\dagger} E_{F_{n}} G_{i}=\lim C_{F_{n}}^{\dagger} \Upsilon_{F_{n}}^{\dagger} B_{F_{n}}^{\dagger \prime}=h_{3}^{\dagger} \lim \Upsilon_{F_{n}}^{\dagger} h_{2}^{\dagger \prime} \tag{19.31}
\end{equation*}
$$

where $C_{F_{n}}^{\dagger} \Upsilon_{F_{n}}^{\dagger}\left(B_{F_{n}}^{\dagger}\right)^{\prime}$ is the singular value decomposition of $W_{F_{n}}^{\dagger} E_{F_{n}} G_{i}, \Upsilon_{F_{n}}^{\dagger}$ is the $k \times p$ matrix with the singular values of $W_{F_{n}}^{\dagger} E_{F_{n}} G_{i}$, denoted by $\left\{\tau_{j F_{n}}^{\dagger}: n \geq 1\right\}$ for $j \leq p$, on the main diagonal
and zeroes elsewhere, and $C_{F_{n}}^{\dagger}$ and $B_{F_{n}}^{\dagger}$ are the corresponding $k \times k$ and $p \times p$ orthogonal matrices of singular vectors, as defined in 19.7. Hence, $\lim \Upsilon_{n}^{\dagger}$ exists, call it $\Upsilon_{h}^{\dagger}$, and equals $h_{3}^{\dagger \prime} h_{4} h_{2}^{\dagger}$. That is, the singular value decomposition of $W_{h}^{\dagger} h_{4}$ is

$$
\begin{equation*}
W_{h}^{\dagger} h_{4}=h_{3}^{\dagger} \Upsilon_{h}^{\dagger} h_{2}^{\dagger \prime} \tag{19.32}
\end{equation*}
$$

The $k \times p$ matrix $\Upsilon_{h}^{\dagger}$ has the limits of the singular values of $W_{F_{n}}^{\dagger} E_{F_{n}} G_{i}$ on its main diagonal and zeroes elsewhere. Let $\tau_{h, j}^{\dagger}$ for $j \leq p$ denote the limits of these singular values. By the definition of $q^{\dagger}, \tau_{h, j}^{\dagger}=0$ for $j=q^{\dagger}+1, \ldots, p$ (because $n^{1 / 2} \tau_{j F_{n}}^{\dagger} \rightarrow h_{1, j}^{\dagger}<\infty$ ). In consequence, $\Upsilon_{h}^{\dagger}$ can be written as

$$
\Upsilon_{h}^{\dagger}=\left[\begin{array}{cc}
\Upsilon_{h, q^{\dagger}}^{\dagger} & 0^{q^{\dagger} \times\left(p-q^{\dagger}\right)}  \tag{19.33}\\
0^{\left(k-q^{\dagger}\right) \times q^{\dagger}} & 0^{\left(k-q^{\dagger}\right) \times\left(p-q^{\dagger}\right)}
\end{array}\right], \text { where } \Upsilon_{h, q^{\dagger}}^{\dagger}:=\operatorname{Diag}\left\{\tau_{h, 1}^{\dagger}, \ldots, \tau_{h, q^{\dagger}}^{\dagger}\right\} .
$$

In addition,

$$
\begin{equation*}
h_{2}^{\dagger \prime} h_{2, p-q^{\dagger}}^{\dagger}=\binom{0^{q^{\dagger} \times\left(p-q^{\dagger}\right)}}{I_{p-q^{\dagger}}} . \tag{19.34}
\end{equation*}
$$

Thus, we have

$$
\begin{align*}
\bar{M}_{h, p-q^{\dagger}}^{\dagger} & :=\bar{W}_{h}\left(W_{h}^{\dagger}\right)^{-1} W_{h}^{\dagger} h_{4} h_{2, p-q^{\dagger}}^{\dagger}=\bar{W}_{h}\left(W_{h}^{\dagger}\right)^{-1} h_{3}^{\dagger} \Upsilon_{h}^{\dagger} h_{2}^{\dagger} h_{2, p-q^{\dagger}}^{\dagger} \\
& =\bar{W}_{h}\left(W_{h}^{\dagger}\right)^{-1} h_{3}^{\dagger}\left[\begin{array}{cc}
\Upsilon_{h, p-q^{\dagger}}^{\dagger} & 0^{q^{\dagger} \times\left(p-q^{\dagger}\right)} \\
0^{\left(k-q^{\dagger}\right) \times q^{\dagger}} & 0^{\left(k-q^{\dagger}\right) \times\left(p-q^{\dagger}\right)}
\end{array}\right]\binom{0^{q^{\dagger} \times\left(p-q^{\dagger}\right)}}{I_{p-q^{\dagger}}} \\
& =0^{k \times\left(p-q^{\dagger}\right)}, \tag{19.35}
\end{align*}
$$

where the first equality holds by the paragraph following Assumption W and uses the condition in Assumption W that $W_{h}^{\dagger}$ is pd and the second equality holds by 19.33 and 19.34. This completes the proof of Corollary 19.7.

### 19.6 Proofs of Results Stated in Sections 19.2 and 19.4

For notational simplicity, the proofs in this section are for the sequence $\{n\}$, rather than a subsequence $\left\{w_{n}: n \geq 1\right\}$. The same proofs hold for any subsequence $\left\{w_{n}: n \geq 1\right\}$.

Proof of Theorem 19.1. Theorem 19.1 follows from Theorem 19.6, which imposes Assumption VD, and Lemma 19.2, which verifies Assumption VD when $\widetilde{V}_{D n}$ is defined by 5.3).

Proof of Lemma 19.2. Consider any sequence $\left\{\lambda_{n, h} \in \Lambda_{K C L R}: n \geq 1\right\}$. By the CLT result in 19.11), the linear expansion of $n^{1 / 2}\left(\widehat{D}_{n}-E_{F_{n}} G_{i}\right)$ in 15.1, and the definitions of $\bar{g}_{h}$ and $\bar{D}_{h}$ in
(19.13), we have

$$
\begin{equation*}
n^{1 / 2}\left(\widehat{g}_{n}, \widehat{D}_{n}-E_{F_{n}} G_{i}\right) \rightarrow_{d}\left(\bar{g}_{h}, \bar{D}_{h}\right) . \tag{19.36}
\end{equation*}
$$

Next, we apply the delta method to the CLT result in 19.11) and the function $a(\cdot)$ defined in 19.16. The mean component in the lhs quantity in 19.11 is $\left(0^{(p+1) k \prime}, \operatorname{vech}\left(E_{F_{n}} f_{i}^{*} f_{i}^{* \prime}\right)^{\prime}\right)^{\prime}$. We have

$$
\begin{align*}
& a\left(\binom{0^{(p+1) k}}{\operatorname{vech}\left(E_{F_{n}} f_{i}^{*} f_{i}^{* \prime}\right)}\right) \\
& =\operatorname{vech}\left(\left(E_{F_{n}} \operatorname{vec}\left(G_{i}-E_{F_{n}} G_{i}\right) \operatorname{vec}\left(G_{i}-E_{F_{n}} G_{i}\right)^{\prime}-\Gamma_{F_{n}}^{v e c\left(G_{i}\right)} \Omega_{F_{n}}^{-1} \Gamma_{F_{n}}^{v e c}\left(G_{i}\right)^{\prime}\right)^{-1 / 2}\right) \\
& =\operatorname{vech}\left(\left(\Phi_{F_{n}}^{v e c\left(G_{i}\right)}\right)^{-1 / 2}\right)=\operatorname{vech}\left(M_{F_{n}}\right), \tag{19.37}
\end{align*}
$$

where $\Gamma_{F_{n}}^{v e c\left(G_{i}\right)}$ and $\Omega_{F_{n}}$ are defined in 3.6, the first equality uses the definitions of $a(\cdot)$ and $f_{i}^{*}$ (given in 19.16) and 5.6, respectively), the second equality holds by the definition of $\Phi_{F_{n}}^{v e c\left(G_{i}\right)}$ in 10.15, and the third equality holds by the definition of $M_{F_{n}}$ in 19.6. Also, $E_{F_{n}} f_{i}^{*} f_{i}^{* \prime} \rightarrow$ $h_{10, f^{*}}$ and $h_{10, f^{*}}$ is pd. Hence, $a(\cdot)$ is well defined and continuously partially differentiable at $\lim \left(0^{(p+1) k \prime}, \operatorname{vech}\left(E_{F_{n}} f_{i}^{*} f_{i}^{* \prime}\right)^{\prime}\right)^{\prime}=\left(0^{(p+1) k \prime}, \operatorname{vech}\left(h_{\left.10, f^{*}\right)^{\prime}}\right)^{\prime}\right.$, as required for the application of the delta method.

The delta method gives

$$
\begin{align*}
n^{1 / 2}\left(A_{n}-\operatorname{vech}\left(M_{F_{n}}\right)\right) & =n^{1 / 2}\left(a\left(n^{-1} \sum_{i=1}^{n}\binom{f_{i}^{*}}{\operatorname{vech}\left(f_{i}^{*} f_{i}^{* \prime}\right)}\right)-a\binom{0^{(p+1) k}}{\operatorname{vech}\left(E_{F_{n}} f_{i}^{*} f_{i}^{* \prime}\right)}\right) \\
& \rightarrow{ }_{d} \bar{A}_{h} \bar{L}_{h} \tag{19.38}
\end{align*}
$$

where the first equality holds by (19.37) and the definitions of $a(\cdot)$ and $A_{n}$ in 19.16), the convergence holds by the delta method using the CLT result in 19.11) and the definition of $\bar{A}_{h}$ following 19.16.

Applying the inverse $\operatorname{vech}(\cdot)$ operator, namely, $\operatorname{vech}_{k p, k p}^{-1}(\cdot)$, to both sides of 19.38 gives the reconfigured convergence result

$$
\begin{equation*}
\left.n^{1 / 2}\left(\operatorname{vech}_{k p, k p}^{-1}\left(A_{n}\right)\right)-M_{F_{n}}\right) \rightarrow_{d} \operatorname{vech}_{k p, k p}^{-1}\left(\bar{A}_{h} \bar{L}_{h}\right)=\bar{M}_{h} \tag{19.39}
\end{equation*}
$$

where the last equality holds by the definition of $\bar{M}_{h}$ in 19.18.
The convergence results in 19.36) and 19.39 hold jointly because both rely on the convergence result in 19.11.

We show below that

$$
\begin{equation*}
n^{1 / 2}\left(\widetilde{V}_{D n}-\left(\operatorname{vech}_{k p, k p}^{-1}\left(A_{n}\right)\right)^{-2}\right)=o_{p}(1) . \tag{19.40}
\end{equation*}
$$

This and the delta method applied again (using the function $\ell(A)=A^{-1 / 2}$ for a pd $k p \times k p$ matrix A) give

$$
\begin{equation*}
n^{1 / 2}\left(\widetilde{V}_{D n}^{-1 / 2}-\operatorname{vech}_{k p, k p}^{-1}\left(A_{n}\right)\right)=o_{p}(1) \tag{19.41}
\end{equation*}
$$

because $v e c h_{k p, k p}^{-1}\left(A_{n}\right)=\left(\Phi_{h}^{v e c\left(G_{i}\right)}\right)^{-1 / 2}+o_{p}(1)$ and $\Phi_{h}^{v e c\left(G_{i}\right)}$ is pd (because $h_{10, f^{*}}$ is pd and $\Phi_{h}^{v e c\left(G_{i}\right)}=$ $Q h_{10, f^{*} Q^{\prime}}$ for some full row rank matrix $Q$ ). Equations 19.36, 19.39, and 19.41) establish the result of the lemma.

Now we prove 19.40 . We have

$$
\begin{align*}
\widetilde{V}_{D n}:= & n^{-1} \sum_{i=1}^{n} \operatorname{vec}\left(G_{i}-\widehat{G}_{n}\right) \operatorname{vec}\left(G_{i}-\widehat{G}_{n}\right)^{\prime}-\widehat{\Gamma}_{n} \widehat{\Omega}_{n}^{-1} \widehat{\Gamma}_{n}^{\prime} \\
= & \left(n^{-1} \sum_{i=1}^{n} \operatorname{vec}\left(G_{i}-E_{F_{n}} G_{i}\right) \operatorname{vec}\left(G_{i}-E_{F_{n}} G_{i}\right)^{\prime}\right)-\left(\operatorname{vec}\left(\widehat{G}_{n}-E_{F_{n}} G_{i}\right) \operatorname{vec}\left(\widehat{G}_{n}-E_{F_{n}} G_{i}\right)^{\prime}\right) \\
& -\left(\widetilde{\Gamma}_{n}-\operatorname{vec}\left(\widehat{G}_{n}-E_{F_{n}} G_{i}\right) \widehat{g}_{n}^{\prime}\right)\left(\widetilde{\Omega}_{n}-\widehat{g}_{n} \widehat{g}_{n}^{\prime}\right)^{-1}\left(\widetilde{\Gamma}_{n}-\operatorname{vec}\left(\widehat{G}_{n}-E_{F_{n}} G_{i}\right) \widehat{g}_{n}^{\prime}\right)^{\prime} \\
= & n^{-1} \sum_{i=1}^{n} \operatorname{vec}\left(G_{i}-E_{F_{n}} G_{i}\right) \operatorname{vec}\left(G_{i}-E_{F_{n}} G_{i}\right)^{\prime}-\widetilde{\Gamma}_{n} \widetilde{\Omega}_{n}^{-1} \widetilde{\Gamma}_{n}^{\prime}+O_{p}\left(n^{-1}\right), \tag{19.42}
\end{align*}
$$

where the second equality holds by subtracting and adding $E_{F_{n}} G_{i}$ and some algebra, by the definitions of $\widehat{\Omega}_{n}$ and $\widehat{\Gamma}_{n}$ in 3.1, 3.2, and 5.3, and by the definitions of $\widetilde{\Omega}_{n}$ and $\widetilde{\Gamma}_{n}$ in 19.16 and the third equality holds because (i) the second summand on the lhs of the third equality is $O_{p}\left(n^{-1}\right)$ because $n^{1 / 2} \operatorname{vec}\left(\widehat{G}_{n}-E_{F_{n}} G_{i}\right)=O_{p}(1)$ (by the CLT using the moment conditions in $\mathcal{F}$, defined in (3.3) and (ii) $n^{1 / 2} \widehat{g}_{n}=O_{p}(1)$ (by Lemma 10.3), $n^{1 / 2} \operatorname{vec}\left(\widehat{G}_{n}-E_{F_{n}} G_{i}\right)=O_{p}(1)$, and $\widehat{\Gamma}_{n}=O_{p}(1)$, $\widehat{\Omega}_{n}^{-1}=O_{p}(1), \widetilde{\Gamma}_{n}=O_{p}(1)$, and $\widetilde{\Omega}_{n}^{-1}=O_{p}(1)$ (by the justification given for 15.1) .

Excluding the $O_{p}\left(n^{-1}\right)$ term, the rhs in 19.42 equals $\left(\text { vech }_{k p, k p}^{-1}\left(A_{n}\right)\right)^{-2}$. Hence, 19.40 holds and the proof is complete.

Proof of Theorem 19.3. The proof is similar to that of Lemma 10.3 in Section 10 with $\widehat{W}_{n}=W_{n}=I_{k}, \widehat{U}_{n}=U_{n}=I_{p}$, and the following quantities $q, \widehat{D}_{n}, D_{n}\left(=E_{F_{n}} G_{i}\right), B_{n, q}, \Upsilon_{n, q}, C_{n}$, and $\Upsilon_{n}$ replaced by $q^{\dagger}, \widehat{D}_{n}^{\dagger}, D_{n}^{\dagger}\left(=D_{F_{n}}^{\dagger}\right), B_{n, q^{\dagger}}^{\dagger}, \Upsilon_{n, q^{\dagger}}^{\dagger}, C_{n}^{\dagger}$, and $\Upsilon_{n}^{\dagger}$, respectively. The proof employs the notational simplifications in (9.1). We can write

$$
\begin{equation*}
\widehat{D}_{n}^{\dagger} B_{n, q^{\dagger}}^{\dagger}\left(\Upsilon_{n, q^{\dagger}}^{\dagger}\right)^{-1}=D_{n}^{\dagger} B_{n, q^{\dagger}}^{\dagger}\left(\Upsilon_{n, q^{\dagger}}^{\dagger}\right)^{-1}+n^{1 / 2}\left(\widehat{D}_{n}^{\dagger}-D_{n}^{\dagger}\right) B_{n, q^{\dagger}}^{\dagger}\left(n^{1 / 2} \Upsilon_{n, q^{\dagger}}^{\dagger}\right)^{-1} . \tag{19.43}
\end{equation*}
$$

By the singular value decomposition, $D_{n}^{\dagger}=C_{n}^{\dagger} \Upsilon_{n}^{\dagger} B_{n}^{\dagger^{\prime}}$. Thus, we obtain

$$
\begin{align*}
D_{n}^{\dagger} B_{n, q^{\dagger}}^{\dagger}\left(\Upsilon_{n, q^{\dagger}}^{\dagger}\right)^{-1} & =C_{n}^{\dagger} \Upsilon_{n}^{\dagger} B_{n}^{\dagger^{\prime}} B_{n, q^{\dagger}}^{\dagger}\left(\Upsilon_{n, q^{\dagger}}^{\dagger}\right)^{-1}=C_{n}^{\dagger} \Upsilon_{n}^{\dagger}\binom{I_{q^{\dagger}}}{0^{\left(p-q^{\dagger}\right) \times q^{\dagger}}}\left(\Upsilon_{n, q^{\dagger}}^{\dagger}\right)^{-1} \\
& =C_{n}^{\dagger}\binom{I_{q^{\dagger}}}{0^{\left(k-q^{\dagger}\right) \times q^{\dagger}}}=C_{n, q^{\dagger}}^{\dagger} . \tag{19.44}
\end{align*}
$$

Let $\widehat{D}_{n}=\left(\widehat{D}_{1 n}, \ldots, \widehat{D}_{p n}\right) \in R^{k \times p}$ and $\bar{D}_{h}=\left(\bar{D}_{1 h}, \ldots, \bar{D}_{p h}\right) \in R^{k \times p}$. We have

$$
\begin{align*}
n^{1 / 2}\left(\widehat{D}_{n}^{\dagger}-D_{n}^{\dagger}\right)= & n^{1 / 2} \sum_{j=1}^{p}\left(\widetilde{M}_{1 j n} \widehat{D}_{j n}-M_{1 j F_{n}} E_{F_{n}} G_{i j}, \ldots, \widetilde{M}_{p j n} \widehat{D}_{j n}-M_{p j F_{n}} E_{F_{n}} G_{i j}\right) \\
= & \sum_{j=1}^{p}\left[\widetilde{M}_{1 j n} n^{1 / 2}\left(\widehat{D}_{j n}-E_{F_{n}} G_{i j}\right)+n^{1 / 2}\left(\widetilde{M}_{1 j n}-M_{1 j F_{n}}\right) E_{F_{n}} G_{i j}, \ldots\right. \\
& \left.\widetilde{M}_{p j n} n^{1 / 2}\left(\widehat{D}_{j n}-E_{F_{n}} G_{i j}\right)+n^{1 / 2}\left(\widetilde{M}_{p j n}-M_{p j F_{n}}\right) E_{F_{n}} G_{i j}\right] \\
\rightarrow & d \sum_{j=1}^{p}\left(M_{1 j h} \bar{D}_{j h}+\bar{M}_{1 j h} h_{4, j}, \ldots, M_{p j h} \bar{D}_{j h}+\bar{M}_{p j h} h_{4, j}\right), \tag{19.45}
\end{align*}
$$

where the convergence holds by Lemma 10.2 in Section 10. Assumption VD, and $E_{F_{n}} G_{i j} \rightarrow h_{4, j}$ (by the definition of $h_{4, j}$ ).

Combining (19.43)-19.45) gives

$$
\begin{equation*}
\widehat{D}_{n}^{\dagger} B_{n, q^{\dagger}}^{\dagger}\left(\Upsilon_{n, q^{\dagger}}^{\dagger}\right)^{-1}=C_{n, q^{\dagger}}^{\dagger}+o_{p}(1) \rightarrow_{p} h_{3, q^{\dagger}}^{\dagger}=\bar{\Delta}_{h, q^{\dagger}}^{\dagger}, \tag{19.46}
\end{equation*}
$$

where the equality uses $n^{1 / 2} \tau_{j F_{n}}^{\dagger} \rightarrow \infty$ for all $j \leq q^{\dagger}$ by the definition of $q^{\dagger}$ and $B_{n, q^{\dagger}}^{\prime} B_{n, q^{\dagger}}=I_{q^{\dagger}}$, the convergence holds by the definition of $h_{3, q^{\dagger}}^{\dagger}$, and the last equality holds by the definition of $\bar{\Delta}_{h, q^{\dagger}}^{\dagger}$ in 19.15.

Using the singular value decomposition $D_{n}^{\dagger}=C_{n}^{\dagger} \Upsilon_{n}^{\dagger} B_{n}^{\dagger \prime}$ again, we obtain

$$
\begin{align*}
& n^{1 / 2} D_{n}^{\dagger} B_{n, p-q^{\dagger}}^{\dagger}=n^{1 / 2} C_{n}^{\dagger} \Upsilon_{n}^{\dagger} B_{n}^{\dagger \prime} B_{n, p-q^{\dagger}}^{\dagger}=n^{1 / 2} C_{n}^{\dagger} \Upsilon_{n}^{\dagger}\binom{0^{q^{\dagger} \times\left(p-q^{\dagger}\right)}}{I_{p-q^{\dagger}}} \\
= & C_{n}^{\dagger}\left(\begin{array}{c}
0^{q^{\dagger} \times\left(p-q^{\dagger}\right)} \\
n^{1 / 2} \Upsilon_{n, p-q^{\dagger}}^{\dagger} \\
0^{(k-p) \times\left(p-q^{\dagger}\right)}
\end{array}\right) \rightarrow h_{3}^{\dagger}\left(\begin{array}{c}
0^{q^{\dagger} \times\left(p-q^{\dagger}\right)} \\
\operatorname{Diag}\left\{h_{1, q^{\dagger}+1}^{\dagger}, \ldots, h_{1, p}^{\dagger}\right\} \\
0^{(k-p) \times\left(p-q^{\dagger}\right)}
\end{array}\right)=h_{3}^{\dagger} h_{1, p-q^{\dagger}}^{\dagger}, \tag{19.47}
\end{align*}
$$

where the second equality uses $B_{n}^{\dagger \prime} B_{n}^{\dagger}=I_{p}$, the convergence holds by the definitions of $h_{3}^{\dagger}$ and $h_{1, j}^{\dagger}$ for $j=1, \ldots, p$, and the last equality holds by the definition of $h_{1, p-q^{\dagger}}^{\dagger \diamond}$ in the paragraph following
(19.10), which uses 10.17).

By 19.45 and $B_{n, p-q^{\dagger}}^{\dagger} \rightarrow h_{2, p-q^{\dagger}}^{\dagger}$, we have

$$
\begin{equation*}
n^{1 / 2}\left(\widehat{D}_{n}^{\dagger}-D_{n}^{\dagger}\right) B_{n, p-q^{\dagger}}^{\dagger} \rightarrow_{d} \bar{D}_{h}^{\dagger} h_{2, p-q^{\dagger}}^{\dagger}+\bar{M}_{h, p-q^{\dagger}}^{\dagger} \tag{19.48}
\end{equation*}
$$

using the definitions of $\bar{D}_{h}^{\dagger}$ and $\bar{M}_{h, p-q^{\dagger}}^{\dagger}$ in 19.14 and 19.19, respectively.
Using (19.47) and (19.48), we get

$$
\begin{align*}
n^{1 / 2} \widehat{D}_{n}^{\dagger} B_{n, p-q^{\dagger}}^{\dagger} & =n^{1 / 2} D_{n}^{\dagger} B_{n, p-q^{\dagger}}^{\dagger}+n^{1 / 2}\left(\widehat{D}_{n}^{\dagger}-D_{n}^{\dagger}\right) B_{n, p-q^{\dagger}}^{\dagger} \\
& \rightarrow{ }_{d} h_{3}^{\dagger} h_{1, p-q^{\dagger}}^{\dagger \diamond}+\bar{D}_{h}^{\dagger} h_{2, p-q^{\dagger}}^{\dagger}+\bar{M}_{h, p-q^{\dagger}}^{\dagger}=\bar{\Delta}_{h, p-q^{\dagger}}^{\dagger}+\bar{M}_{h, p-q^{\dagger}}^{\dagger} \tag{19.49}
\end{align*}
$$

where the last equality holds by the definition of $\bar{\Delta}_{h, p-q^{\dagger}}^{\dagger}$ in 19.15 .
Equations 19.46) and 19.49) combine to give

$$
\begin{align*}
n^{1 / 2} \widehat{D}_{n}^{\dagger} T_{n}^{\dagger} & =n^{1 / 2} \widehat{D}_{n}^{\dagger} B_{n}^{\dagger} S_{n}^{\dagger}=\left(\widehat{D}_{n}^{\dagger} B_{n, q^{\dagger}}^{\dagger}\left(\Upsilon_{n, q^{\dagger}}^{\dagger}\right)^{-1}, n^{1 / 2} \widehat{D}_{n}^{\dagger} B_{n, p-q^{\dagger}}^{\dagger}\right) \\
& \rightarrow{ }_{d}\left(\bar{\Delta}_{h, q^{\dagger}}^{\dagger}, \bar{\Delta}_{h, p-q^{\dagger}}^{\dagger}+\bar{M}_{h, p-q^{\dagger}}^{\dagger}\right)=\bar{\Delta}_{h}^{\dagger}+\bar{M}_{h}^{\dagger} \tag{19.50}
\end{align*}
$$

using the definitions of $S_{n}^{\dagger}$ and $T_{n}^{\dagger}$ in 19.26, $\bar{\Delta}_{h}^{\dagger}$ in 19.15, and $\bar{M}_{h}^{\dagger}$ in 19.19.
By Lemma $10.2, n^{1 / 2}\left(\widehat{g}_{n}, \widehat{D}_{n}-E_{F_{n}} G_{i}\right) \rightarrow_{d}\left(\bar{g}_{h}, \bar{D}_{h}\right)$. This convergence is joint with that in 19.50) because the latter just relies on the convergence of $n^{1 / 2}\left(\widehat{D}_{n}-E_{F_{n}} G_{i}\right)$, which is part of the former, and of $n^{1 / 2}\left(\widetilde{M}_{n}-M_{F_{n}}\right) \rightarrow_{d} \bar{M}_{h}$, which holds jointly with the former by Assumption VD. This establishes the convergence result of Theorem 19.3 .

The independence of $\bar{g}_{h}$ and $\left(\bar{D}_{h}, \bar{\Delta}_{h}^{\dagger}\right)$ follows from the independence of $\bar{g}_{h}$ and $\bar{D}_{h}$, which holds by Lemma 10.2, and the fact that $\bar{\Delta}_{h}^{\dagger}$ is a nonrandom function of $\bar{D}_{h}$.

Proof of Lemma 19.4. The proof of Lemma 19.4 is analogous to the proof of Theorem 10.4 with $\widehat{W}_{n}=W_{n}=I_{k}, \widehat{U}_{n}=U_{n}=I_{p}$, and the following quantities $q, \widehat{D}_{n}, D_{n}\left(=E_{F_{n}} G_{i}\right), \widehat{\kappa}_{j n}, B_{n}, B_{n, q}$, $S_{n}, S_{n, q}, \tau_{j F_{n}}$, and $h_{3, q}$ replaced by $q^{\dagger}, \widehat{D}_{n}^{\dagger}, D_{n}^{\dagger}\left(=D_{F_{n}}^{\dagger}\right), \widehat{\kappa}_{j n}^{\dagger}, B_{n}^{\dagger}, B_{n, q^{\dagger}}^{\dagger}, S_{n}^{\dagger}, S_{n, q^{\dagger}}^{\dagger}, \tau_{j F_{n}}^{\dagger}$, and $h_{3, q^{\dagger}}^{\dagger}$, respectively. Theorem 19.3, rather than Lemma 10.3 , is employed to obtain the results in 17.37). In consequence, $\bar{\Delta}_{h, q}$ and $\bar{\Delta}_{h, p-q}$ are replaced by $\bar{\Delta}_{h, q^{\dagger}}^{\dagger}+\bar{M}_{h, q^{\dagger}}^{\dagger}$ and $\bar{\Delta}_{h, p-q^{\dagger}}^{\dagger}+\bar{M}_{h, p-q^{\dagger}}^{\dagger}$, respectively, where $\bar{\Delta}_{h, q^{\dagger}}^{\dagger}+\bar{M}_{h, q^{\dagger}}^{\dagger}=\bar{\Delta}_{h, q^{\dagger}}^{\dagger}$ (because $\bar{M}_{h, q^{\dagger}}^{\dagger}:=0^{k \times q^{\dagger}}$ by 19.19). The quantities $\bar{\Delta}_{h, q}$ and $\bar{\Delta}_{h, p-q}$ are replaced by $\bar{\Delta}_{h, q^{\dagger}}^{\dagger}$ and $\bar{\Delta}_{h, p-q^{\dagger}}^{\dagger}+\bar{M}_{h, p-q^{\dagger}}^{\dagger}$ in 17.37) and in the rest of the proof of Theorem 10.4. Note that 17.39 holds with $h_{3, q}$ replaced by $h_{3, q^{\dagger}}^{\dagger}$ because $\bar{\Delta}_{h, q^{\dagger}}^{\dagger}=h_{3, q^{\dagger}}^{\dagger}$ by 19.15 (just as $\overline{\bar{\Delta}}_{h, q}=h_{3, q}$ ). Because $\widehat{U}_{n}=U_{n}$, the matrices $\widehat{A}_{n}$ and $A_{j n}$ for $j=1,2,3$ (defined in 17.39 ) are all zero matrices, which simplifies the expressions in (17.41)-17.44) considerably.

The proof of Theorem 10.4 uses Lemma 17.1 to obtain 17.42 . Hence, an analogue of Lemma 17.1 is needed, where the changes listed in the first paragraph of this proof are made and $h_{6, j}$ and $C_{n}$ are replaced by $h_{6, j}^{\dagger}$ and $C_{n}^{\dagger}$, respectively. In addition, $\mathcal{F}_{W U}$ is replaced by $\mathcal{F}_{K C L R}$ (because $\mathcal{F}_{K C L R} \subset \mathcal{F}_{W U}$ for $\delta_{W U}$ sufficiently small and $M_{W U}$ sufficiently large using the facts that $\mathcal{F}_{0} \cap \mathcal{F}_{W U}$ equals $\mathcal{F}_{0}$ for $\delta_{W U}$ sufficiently small and $M_{W U}$ sufficiently large by the argument following (10.5) and $\mathcal{F}_{K C L R} \subset \mathcal{F}_{0}$ by the argument following 19.5). Because $\widehat{U}_{n}=U_{n}$, the matrices $\widehat{A}_{j n}$ for $j=1,2,3$ (defined in (17.2) are all zero matrices, which simplifies the expressions in (17.9)-(17.12) considerably. For $(17.3)$ to go through with the changes listed above (in particular, with $\widehat{W}_{n}, \widehat{D}_{n}$, $D_{n}$, and $U_{n}$ replaced by $I_{k}, \widehat{D}_{n}^{\dagger}, D_{n}^{\dagger}$, and $I_{p}$, respectively), we need to show that

$$
\begin{equation*}
n^{1 / 2}\left(\widehat{D}_{n}^{\dagger}-D_{n}^{\dagger}\right)=O_{p}(1) \tag{19.51}
\end{equation*}
$$

By (5.4) with $\theta=\theta_{0}$ (and with the dependence of various quantities on $\theta_{0}$ suppressed for notational simplicity), we have

$$
\widehat{D}_{n}^{\dagger}=\sum_{j=1}^{p}\left(\widetilde{M}_{1 j n} \widehat{D}_{j n}, \ldots, \widetilde{M}_{p j n} \widehat{D}_{j n}\right), \text { where } \widetilde{M}_{n}=\left[\begin{array}{ccc}
\widetilde{M}_{11 n} & \cdots & \widetilde{M}_{1 p n}  \tag{19.52}\\
\vdots & \ddots & \vdots \\
\widetilde{M}_{p 1 n} & \cdots & \widetilde{M}_{p p n}
\end{array}\right]:=\widetilde{V}_{D n}^{-1 / 2} \in R^{k p \times k p} .
$$

By 19.6), we have

$$
\begin{equation*}
D_{n}^{\dagger}=\sum_{j=1}^{p}\left(M_{1 j F_{n}} D_{j n}, \ldots, M_{p j F_{n}} D_{j n}\right) \tag{19.53}
\end{equation*}
$$

$\operatorname{using} D_{n}=\left(D_{1 n}, \ldots, D_{p n}\right)$, and $D_{j n}:=E_{F_{n}} G_{i j}$ for $j=1, \ldots, p$.
For $s=1, \ldots, p$, we have

$$
\begin{equation*}
n^{1 / 2}\left(\widetilde{M}_{s j n} \widehat{D}_{j n}-M_{s j F_{n}} D_{j n}\right)=\widetilde{M}_{s j n} n^{1 / 2}\left(\widehat{D}_{j n}-D_{j n}\right)+n^{1 / 2}\left(\widetilde{M}_{s j n}-M_{s j F_{n}}\right) D_{j n}=O_{p}(1), \tag{19.54}
\end{equation*}
$$

where $n^{1 / 2}\left(\widehat{D}_{j n}-D_{j n}\right)=O_{p}(1)($ by Lemma 10.2 $), n^{1 / 2}\left(\widetilde{M}_{s j n}-M_{s j F_{n}}\right)=O_{p}(1)$ (because $n^{1 / 2}\left(\widetilde{M}_{n}-\right.$ $\left.M_{F_{n}}\right) \rightarrow{ }_{d} \bar{M}_{h}$ by Assumption VD), $M_{s j F_{n}}=O(1)$ (because $M_{F}=\left(\Phi_{F}^{v e c\left(G_{i}\right)}\right)^{-1 / 2}, \Phi_{F}^{v e c\left(G_{i}\right)}$ defined in 10.15 satisfies $\Phi_{F}^{\operatorname{vec}\left(G_{i}\right)}:=\operatorname{Var}_{F}\left(\operatorname{vec}\left(G_{i}\right)-\Gamma_{F}^{\operatorname{vec}\left(G_{i}\right)} \Omega_{F}^{-1} g_{i}\right)=\left[-E_{F} \operatorname{vec}\left(G_{i}\right) g_{i}^{\prime} \Omega_{F}^{-1}: I_{p k}\right] \operatorname{Var} F\left(f_{i}^{*}\right)$, and $\lambda_{\min }\left(\operatorname{Var}_{F}\left(f_{i}^{*}\right)\right) \geq \delta_{2}$ by the definition of $\mathcal{F}_{K C L R}$ in 19.5) , and $D_{j n}=O(1)$ (by the moment conditions in $\mathcal{F}$, defined in (3.3).

Hence,

$$
\begin{equation*}
n^{1 / 2}\left(\widehat{D}_{n}^{\dagger}-D_{n}^{\dagger}\right)=\sum_{j=1}^{p} n^{1 / 2}\left[\left(\widetilde{M}_{1 j n} \widehat{D}_{j n}, \ldots, \widetilde{M}_{p j n} \widehat{D}_{j n}\right)-\left(M_{1 j F_{n}} D_{j n}, \ldots, M_{p j F_{n}} D_{j n}\right)\right]=O_{p}(1) . \tag{19.55}
\end{equation*}
$$

This completes the proof of the analogue of Lemma 17.1, which completes the proof of parts (a)-(d) of Lemma 19.4 ,

For part (e) of Lemma 19.4 , the results of parts (a)-(d) hold jointly with those in Theorem 19.3 , rather than those in Lemma 10.3, because Theorem 19.3 is used to obtain the results in 17.37, rather than Lemma 10.3 . This completes the proof.

Proof of Lemma 19.5. The proof of parts (a) and (b) is the same as the proof of Theorem 12.1 for the case where Assumption $\mathrm{R}(\mathrm{a})$ holds (which states that $r k_{n} \rightarrow_{p} \infty$ ) using Lemma 19.4(a), which shows that $r k_{n}^{\dagger} \rightarrow_{d} \infty$ if $q^{\dagger}=p$.

The proofs of parts (c) and (d) are the same as in 12.5 - 12.9 ) in the proof of Theorem 12.1 for the case where Assumption $\mathrm{R}(\mathrm{b})$ holds, using Theorem 19.3 and Lemma 19.4 (b) in place of Lemma 10.3. with $r_{h}\left(\bar{D}_{h}, \bar{M}_{h}\right)$ (defined in 19.20 ) in place of $r_{h}\left(\bar{D}_{h}\right)$, and for part (d), with the proviso that $P\left(\overline{C L R}_{h}=c\left(1-\alpha, \bar{r}_{h}\right)\right)=0$. (The proof in Theorem 12.1 that $P\left(\overline{C L R}_{h}=c\left(1-\alpha, \bar{r}_{h}\right)\right)=0$ does not go through in the present case because $\bar{r}_{h}=r_{h}\left(\bar{D}_{h}, \bar{M}_{h}\right)$ is not necessarily a constant conditional on $\bar{D}_{h}$ and alternatively, conditional on $\left(\bar{D}_{h}, \bar{M}_{h}\right), \overline{L M}_{h}$ and $\bar{J}_{h}$ are not necessarily independent and distributed as $\chi_{p}^{2}$ and $\chi_{k-p}^{2}$.) Note that 12.10 does not necessarily hold in the present case, because $\bar{r}_{h}=r_{h}\left(\bar{D}_{h}, \bar{M}_{h}\right)$ is not necessarily a constant conditional on $\bar{D}_{h}$.

The proof of Theorem 19.6 given below uses Corollary 2.1(a) of ACG, which is stated below as Proposition 19.8. It is a generic asymptotic size result. Unlike Proposition 10.1 above, Proposition 19.8 applies when the asymptotic size is not necessarily equal to the nominal size $\alpha$. Let $\left\{\phi_{n}: n \geq 1\right\}$ be a sequence of tests of some null hypothesis whose null distributions are indexed by a parameter $\lambda$ with parameter space $\Lambda$. Let $R P_{n}(\lambda)$ denote the null rejection probability of $\phi_{n}$ under $\lambda$. For a finite nonnegative integer $J$, let $\left\{h_{n}(\lambda)=\left(h_{1 n}(\lambda), \ldots, h_{J n}(\lambda)\right)^{\prime} \in R^{J}: n \geq 1\right\}$ be a sequence of functions on $\Lambda$. Define $H$ as in 10.1.

For a sequence of scalar constants $\left\{C_{n}: n \geq 1\right\}$, let $C_{n} \rightarrow\left[C_{1, \infty}, C_{2, \infty}\right]$ denote that $C_{1, \infty} \leq$ $\liminf _{n \rightarrow \infty} C_{n} \leq \limsup \operatorname{sum}_{n \rightarrow \infty} C_{n} \leq C_{2, \infty}$.

Assumption B: For any subsequence $\left\{w_{n}\right\}$ of $\{n\}$ and any sequence $\left\{\lambda_{w_{n}} \in \Lambda: n \geq 1\right\}$ for which $h_{w_{n}}\left(\lambda_{w_{n}}\right) \rightarrow h \in H, R P_{w_{n}}\left(\lambda_{w_{n}}\right) \rightarrow\left[R P^{-}(h), R P^{+}(h)\right]$ for some $R P^{-}(h), R P^{+}(h) \in[0,1]$.

Proposition 19.8 (ACG, Corollary 2.1(a)) Under Assumption B, the tests $\left\{\phi_{n}: n \geq 1\right\}$ have $A s y S z:=\limsup _{n \rightarrow \infty} \sup _{\lambda \in \Lambda} R P_{n}(\lambda) \in\left[\sup _{h \in H} R P^{-}(h), \sup _{h \in H} R P^{+}(h)\right]$.

Comments: (i) Corollary 2.1(a) of ACG is stated for CS's, rather than tests. But, following Comment 4 to Theorem 2.1 of ACG, with suitable adjustments (as in Proposition 19.8 above) it applies to tests as well.
(ii) Under Assumption B, if $R P^{-}(h)=R P^{+}(h)$ for all $h \in H$, then $A s y S z=\sup _{h \in H} R P^{+}(h)$. We use this to prove Theorem 19.6. The result of Proposition 19.8 for the case where $R P^{-}(h) \neq$ $R P^{+}(h)$ for some $h \in H$ is used when proving Comment (i) to Theorem 19.1 and the Comment to Theorem 19.6

Proof of Theorem 19.6. Theorem 19.6 follows from Lemma 19.5 and Proposition 19.8 because Lemma 19.5 verifies Assumption B with $R P^{-}(h)=R P^{+}(h)=\alpha$ when $q^{\dagger}=p$ and with $R P^{-}(h)=$ $R P^{+}(h)=P\left(\overline{C L R}_{h}>c\left(1-\alpha, \bar{r}_{h}\right)\right)$ when $q^{\dagger}<p$.

### 19.7 Proof of Lemma 5.2

Proof of Lemma 5.2. Define $J_{n}(\theta)$ by the decomposition $A R_{n}(\theta)=L M_{n}(\theta)+J_{n}(\theta)$. Under all subsequences $\left\{w_{n}\right\}$ and all sequences $\left\{\lambda_{w_{n}, h}: n \geq 1\right\}$ with $\lambda_{w_{n}, h} \in \Lambda_{0}$ in 10.10, 11.2 and 12.6) imply that

$$
\begin{equation*}
\binom{J_{w_{n}}\left(\theta_{0}\right)}{L M_{w_{n}}\left(\theta_{0}\right)} \rightarrow_{d}\binom{\bar{J}_{h}}{L M_{h}} \sim\binom{\bar{g}_{h}^{\prime} h_{5, g}^{-1 / 2} M_{\bar{\Delta}_{h}} h_{5, g}^{-1 / 2} \bar{g}_{h}}{\bar{g}_{h}^{\prime} h_{5, g}^{-1 / 2} P_{\bar{\Delta}_{h}} h_{5, g}^{-1 / 2} \bar{g}_{h},}, \tag{19.56}
\end{equation*}
$$

where $\bar{\Delta}_{h}$ is defined in 10.17). Note that the parameter space $\Lambda_{0}$ for $\lambda$ defined in 10.9 is equivalent to the parameter space $\mathcal{F}_{0}$, see Comment (i) to Theorem 11.1 .

Equation 19.56) and the CMT imply that the test statistic in (5.7) converges in distribution to

$$
\begin{equation*}
\sup _{r \in[0, \infty]}\left[\frac{1}{2}\left(\overline{L M}_{h}+\bar{J}_{h}-r+\sqrt{\left(\overline{L M}_{h}+\bar{J}_{h}-r\right)^{2}+4{\overline{L M_{h}}}_{h} r}\right)-c(1-\alpha, r)\right] . \tag{19.57}
\end{equation*}
$$

Conditional on $\bar{\Delta}_{h}, \overline{L M}_{h}$ and $\bar{J}_{h}$ are independent and distributed as $\chi_{p}^{2}$ and $\chi_{k-p}^{2}$, respectively. Therefore, the conditional distribution of the random variable in 19.57 given $\bar{\Delta}_{h}$ is the same as the distribution of the quantity in 5.8). Since the latter does not depend on $\bar{\Delta}_{h}$, the same statement holds for the unconditional distribution of the random variable in (19.57).

The results of the previous paragraph verify Assumption B* (stated just above Proposition 10.1) with the limit of the rejection probabilities in Assumption $\mathrm{B}^{*}$, i.e., $\lim _{n \rightarrow \infty} R P_{w_{n}}\left(\lambda_{w_{n}}\right)$, equal to the probability that the random variable in (5.8) is positive. The asymptotic size result of the Lemma now follows by Proposition 10.1.

## 20 Proof of Theorem 7.1

Theorem 7.1 of AG1. Suppose the LM test, the CLR test with moment-variance weighting, and when $p=1$ the CLR test with Jacobian-variance weighting are defined as in this section, the parameter space for $F$ is $\mathcal{F}_{T S, 0}$ for the first two tests and $\mathcal{F}_{T S, J V W, p=1}$ for the third test, and

Assumption V holds. Then, these tests have asymptotic sizes equal to their nominal size $\alpha \in(0,1)$ and are asymptotically similar (in a uniform sense). Analogous results hold for the corresponding CS's for the parameter spaces $\mathcal{F}_{\Theta, T S, 0}$ and $\mathcal{F}_{\Theta, T S, J V W, p=1}$.

The proof of Theorem 7.1 is analogous to that of Theorems 4.1, 5.3, and 6.1. In the time series case, for tests, we define $\lambda=\left(\lambda_{1, F}, \ldots, \lambda_{9, F}\right)$ and $\left\{\lambda_{n, h}: n \geq 1\right\}$ as in 10.9) and 10.11), respectively, but with $\lambda_{5, F}$ defined differently than in the i.i.d. case. (For CS's in the time series case, we make the adjustments outlined in the Comment to Proposition 10.1.) We define

$$
\begin{equation*}
\lambda_{5, F}:=V_{F}=\sum_{m=-\infty}^{\infty} E_{F}\binom{g_{i}}{\operatorname{vec}\left(G_{i}-E_{F} G_{i}\right)}\binom{g_{i-m}}{\operatorname{vec}\left(G_{i-m}-E_{F} G_{i-m}\right)}^{\prime} . \tag{20.1}
\end{equation*}
$$

In consequence, $\lambda_{5, F_{n}} \rightarrow h_{5}$ implies that $V_{F_{n}} \rightarrow h_{5}$ and the condition in Assumption V holds with $V=h_{5}$. The difference in the definitions of $\lambda_{5, F}$ in the i.i.d. and time series cases reflects the difference in the definitions of $\Sigma_{F}^{v e c\left(G_{i}\right)}$ in these two cases. See the discussion following 7.1 of AG1 above regarding the latter.

The proof of Theorem 7.1 uses the CLT given in the following lemma.
Lemma 20.1 Let $f_{i}:=\left(g_{i}^{\prime} \text {, vec }\left(G_{i}\right)^{\prime}\right)^{\prime}$. We have: $w_{n}^{-1 / 2} \sum_{i=1}^{w_{n}}\left(f_{i}-E_{F_{n}} f_{i}\right) \rightarrow{ }_{d} N\left(0^{(p+1) k}, h_{5}\right)$ under all subsequences $\left\{w_{n}\right\}$ and all sequences $\left\{\lambda_{w_{n}, h}: n \geq 1\right\}$.

Proof of Theorem 7.1. The proof is the same as the proofs of Theorems 4.1, 5.3, and 6.1 (given in Sections 11, 12, and 13, respectively, above) and the proofs of Lemmas 10.2 and 10.3 and Theorem 10.4 (given in Sections 15,16 , and 17 above), upon which the former proofs rely, for the i.i.d. case with some modifications. The modifications affect the proofs of Lemmas 10.2 and 10.3 and the proof of Theorem 5.3. No modifications are needed elsewhere.

The first modification is the change in the definition of $\lambda_{5, F}$ described in (20.1).
The second modification is that $\widehat{\Omega}_{n}=\widehat{\Omega}_{n}\left(\theta_{0}\right) \rightarrow_{p} h_{5, g}$ not by the WLLN but by Assumption V and the definition of $\widehat{\Omega}_{n}(\theta)$ in 7.4). In the time series case, by definition, $\lambda_{5, F}:=V_{F}$, so $h_{5}:=\lim \lambda_{5, F_{n}}=\lim V_{F_{n}}$. By definition, $h_{5, g}$ is the upper left $k \times k$ submatrix of $h_{5}$ and $\Omega_{F}$ is the upper left $k \times k$ submatrix of $V_{F}$ by (7.1) and (20.1). Hence, $h_{5, g}=\lim \Omega_{F_{n}}$. By the definition of $\mathcal{F}_{T S}, \lambda_{\min }\left(\Omega_{F}\right) \geq \delta \forall F \in \mathcal{F}_{T S}$. Hence, $h_{5, g}$ is pd.

Let $h_{5, G_{j} g}$ be the $k \times k$ submatrix of $h_{5}$ that corresponds to the submatrix $\widehat{\Gamma}_{j n}(\theta)$ of $\widehat{V}_{n}(\theta)$ in (7.4) for $j=1, \ldots, p$. The third modification is that $\widehat{\Gamma}_{j n}=\widehat{\Gamma}_{j n}\left(\theta_{0}\right)=h_{5, G_{j} g}+o_{p}(1)$ in 15.1) in the proof of Lemma 10.2 (rather than $\left.\widehat{\Gamma}_{j n}=E_{F_{n}} G_{i j} g_{i}^{\prime}+o_{p}(1)\right)$ for $j=1, \ldots, p$ and this holds by Assumption V and the definition of $\widehat{\Gamma}_{j n}(\theta)$ in 7.4 (rather than by the WLLN).

We write

$$
h_{5}=\left(\begin{array}{cc}
h_{5, g} & h_{5, G g}^{\prime}  \tag{20.2}\\
h_{5, G g} & h_{5, G}
\end{array}\right) \text { for } h_{5, g} \in R^{k \times k}, h_{5, G g}=\left(\begin{array}{c}
h_{5, G_{1} g} \\
\vdots \\
h_{5, G_{p} g}
\end{array}\right) \in R^{p k \times k}, \text { and } h_{5, G} \in R^{p k \times p k} .
$$

The fourth modification is that $\widetilde{V}_{D n}$ in 13.1 in the proof of Theorem 5.3 is defined as described in Section 7, rather than as in 5.3. In addition, $\widetilde{V}_{D n} \rightarrow_{p} h_{7}$ in 13.1 holds with $h_{7}=h_{5, G}-$ $h_{5, G g}\left(h_{5, g}\right)^{-1} h_{5, G g}^{\prime}$ by Assumption V, rather than by the WLLN.

The fifth modification is the use of a WLLN and CLT for triangular arrays of strong mixing random vectors, rather than i.i.d. random vectors, for the quantities in the proof of Lemma 10.2 and elsewhere. For the WLLN, we use Example 4 of Andrews (1988), which shows that for a strong mixing row-wise-stationary triangular array $\left\{W_{i}: i \leq n\right\}$ we have $n^{-1} \sum_{i=1}^{n}\left(\xi\left(W_{i}\right)-E_{F_{n}} \xi\left(W_{i}\right)\right) \rightarrow_{p}$ 0 for any real-valued function $\xi(\cdot)$ (that may depend on $n$ ) for which $\sup _{n \geq 1} E_{F_{n}}\left\|\xi\left(W_{i}\right)\right\|^{1+\delta}<\infty$ for some $\delta>0$. For the CLT, we use Lemma 20.1 as follows. The joint convergence of $n^{1 / 2} \widehat{g}_{n}$ and $n^{1 / 2}\left(\widehat{D}_{n}-E_{F_{n}} G_{i}\right)$ in the proof of Lemma 10.2 is obtained from 15.1 , modified by the second and third modifications above, and the following result:

$$
\begin{align*}
& n^{-1 / 2} \sum_{i=1}^{n}\left(\zeta\left(W_{i}\right)-E_{F_{n}} \zeta\left(W_{i}\right)\right)=\left(\begin{array}{cc}
I_{k} & 0^{k \times p k} \\
-h_{5, G g} h_{5, g}^{-1} & I_{p k}
\end{array}\right) n^{-1 / 2} \sum_{i=1}^{n}\left(f_{i}-E_{F_{n}} f_{i}\right) \\
& \rightarrow{ }_{d} N\left(0^{(p+1) k}, L_{h_{5}}\right), \text { where } \\
& \zeta\left(W_{i}\right):=\binom{g_{i}}{\operatorname{vec}\left(G_{i}\right)-h_{5, G g} h_{5, g}^{-1} g_{i}}=\left(\begin{array}{cc}
I_{k} & 0^{k \times p k} \\
-h_{5, G g} h_{5, g}^{-1} & I_{p k}
\end{array}\right)\binom{g_{i}}{\operatorname{vec}\left(G_{i}\right)}, \tag{20.3}
\end{align*}
$$

$f_{i}=\left(g_{i}^{\prime}, \operatorname{vec}\left(G_{i}\right)^{\prime}\right)^{\prime}$, and the convergence holds by Lemma 20.1. Using 20.2), the variance matrix $L_{h_{5}}$ in (20.3) takes the form:

$$
\begin{align*}
L_{h_{5}} & =\left(\begin{array}{cc}
I_{k} & 0^{k \times p k} \\
-h_{5, G g} h_{5, g}^{-1} & I_{p k}
\end{array}\right)\left(\begin{array}{cc}
h_{5, g} & h_{5, G g^{\prime}} \\
h_{5, G g} & h_{5, G}
\end{array}\right)\left(\begin{array}{cc}
I_{k} & -h_{5, g}^{-1} h_{5, G g}^{\prime} \\
0^{p k \times k} & I_{p k}
\end{array}\right) \\
& =\left(\begin{array}{cc}
I_{k} & 0^{k \times p k} \\
-h_{5, G g} h_{5, g}^{-1} & I_{p k}
\end{array}\right)\left(\begin{array}{cc}
h_{5, g} & 0^{k \times p k} \\
h_{5, G g} & \Phi_{h}^{v e c\left(G_{i}\right)}
\end{array}\right)=\left(\begin{array}{cc}
h_{5, g} & 0^{k \times p k} \\
0^{p k \times k} & \Phi_{h}^{v e c\left(G_{i}\right)}
\end{array}\right), \text { where } \\
\Phi_{h}^{v e c\left(G_{i}\right)} & =h_{5, G}-h_{5, G g} h_{5, g}^{-1} h_{5, G g}^{\prime} . \tag{20.4}
\end{align*}
$$

Equations (15.1) (modified as described above), (20.3), and 20.4 combine to give the result of Lemma 10.2 for the time series case.

The sixth modification occurs in the proof of Lemma 10.3 (d) in Section 16 in this SM. In the time series case, the proof goes through as is, except that the calculations in (16.13) are not needed because $\Sigma_{F}^{a_{i}}$ (and, hence, $\Psi_{F}^{a_{i}}$ as well) is defined with its underlying components recentered at their means (which is needed to ensure that $\Sigma_{F}^{a_{i}}$ is a convergent sum). The latter implies that $\lim \Psi_{F_{n}}^{v e c\left(G_{i}\right)}=\Phi_{h}^{v e c\left(G_{i}\right)}$ automatically holds and $\lim \Psi_{F_{n}}^{v e c\left(C_{F_{n}, k-q}^{\prime} \Omega_{F_{n}}^{-1 / 2} G_{i} B_{F_{n}, p-q} \xi_{2}\right)}=$ $\Phi_{h}^{v e c\left(h_{3, k-q}^{\prime} h_{5, g}^{-1 / 2} G_{i} h_{2, p-q} \xi_{2}\right)}$ (which, in the i.i.d. case, is proved in 16.133.

This completes the proof of Theorem 7.1.
Proof of Lemma 20.1. For notational simplicity, we prove the result for the sequence $\{n\}$ rather than a subsequence $\left\{w_{n}: n \geq 1\right\}$. The same proof applies for any subsequence. By the CramérWold device, it suffices to prove the result with $f_{i}-E_{F_{n}} f_{i}$ and $h_{5}$ replaced by $s\left(W_{i}\right)=b^{\prime}\left(f_{i}-E_{F_{n}} f_{i}\right)$ and $b^{\prime} h_{5} b$, respectively, for arbitrary $b \in R^{(p+1) k}$. First, we show

$$
\begin{equation*}
\lim \operatorname{Var}_{F_{n}}\left(n^{-1 / 2} \sum_{i=1}^{n} s\left(W_{i}\right)\right)=b^{\prime} h_{5} b, \tag{20.5}
\end{equation*}
$$

where by assumption $\lambda_{5, F_{n}}=\sum_{m=-\infty}^{\infty} E_{F_{n}} s\left(W_{i}\right) s\left(W_{i-m}\right) \rightarrow h_{5}$. By change of variables, we have

$$
\begin{equation*}
\operatorname{Var}_{F_{n}}\left(n^{-1 / 2} \sum_{i=1}^{n} s\left(W_{i}\right)\right)=\sum_{m=-n+1}^{n-1} \operatorname{Cov}_{F_{n}}\left(s\left(W_{i}\right), s\left(W_{i-m}\right)\right)-\sum_{m=-n+1}^{n-1} \frac{|m|}{n} \operatorname{Cov}_{F_{n}}\left(s\left(W_{i}\right), s\left(W_{i-m}\right)\right) . \tag{20.6}
\end{equation*}
$$

This gives

$$
\begin{align*}
& \left\|\operatorname{Var}_{F_{n}}\left(n^{-1 / 2} \sum_{i=1}^{n} s\left(W_{i}\right)\right)-b^{\prime} \lambda_{5, F_{n}} b\right\| \\
\leq & 2 \sum_{m=n}^{\infty}\left\|\operatorname{Cov}_{F_{n}}\left(s\left(W_{i}\right), s\left(W_{i-m}\right)\right)\right\|+\sum_{m=-n+1}^{n-1} \frac{|m|}{n}\left\|\operatorname{Cov}_{F_{n}}\left(s\left(W_{i}\right), s\left(W_{i-m}\right)\right)\right\| . \tag{20.7}
\end{align*}
$$

By a standard strong mixing covariance inequality, e.g., see Davidson (1994, p. 212),

$$
\begin{equation*}
\sup _{F \in \mathcal{F}_{T S}}\left\|\operatorname{Cov}_{F}\left(s\left(W_{i}\right), s\left(W_{i-m}\right)\right)\right\| \leq C_{1} \alpha_{F}^{\gamma /(2+\gamma)}(m) \leq C_{1} C^{\gamma /(2+\gamma)} m^{-d \gamma /(2+\gamma)}, \text { where } d \gamma /(2+\gamma)>1, \tag{20.8}
\end{equation*}
$$

for some $C_{1}<\infty$, where the second inequality uses the definition of $\mathcal{F}_{T S}$ in (7.2). In consequence, both terms on the rhs of (20.7) converge to zero. This and $b^{\prime} \lambda_{5, F_{n}} b \rightarrow b^{\prime} h_{5} b$ establish 20.5).

When $b^{\prime} h_{5} b=0$, we have $\lim _{n \rightarrow \infty} \operatorname{Var}_{F_{n}}\left(n^{-1 / 2} \sum_{i=1}^{n} s\left(W_{i}\right)\right)=0$, which implies that $n^{-1 / 2} \sum_{i=1}^{n}$ $s\left(W_{i}\right) \rightarrow_{d} N\left(0, b^{\prime} h_{5} b\right)=0$. When $b^{\prime} h_{5} b>0$, we can assume $\sigma_{n}^{2}=\operatorname{Var}_{F_{n}}\left(n^{-1 / 2} \sum_{i=1}^{n} s\left(W_{i}\right)\right) \geq c$ for some $c>0 \forall n \geq 1$ without loss of generality. We apply the triangular array CLT in Corollary

1 of de Jong (1997) with (using de Jong's notation) $\beta=\gamma=0, c_{n i}:=n^{-1 / 2} \sigma_{n}^{-1}$, and $X_{n i}:=$ $n^{-1 / 2} s\left(W_{i}\right) \sigma_{n}^{-1}$. Now we verify conditions (a)-(c) of Assumption 2 of de Jong (1997). Condition (a) holds automatically. Condition (b) holds because $c_{n i}>0$ and $E_{F_{n}}\left|X_{n i} / c_{n i}\right|^{2+\gamma}=E_{F_{n}}\left|s\left(W_{i}\right)\right|^{2+\gamma} \leq$ $2\|b\|^{2+\gamma} M<\infty \forall F_{n} \in \mathcal{F}_{T S}$. Condition (c) holds by taking $V_{n i}=X_{n i}$ (where $V_{n i}$ is the random variable that appears in the definition of near epoch dependence in Definition 2 of de Jong (1997)), $d_{n i}=0$, and using $\alpha_{F_{n}}(m) \leq C m^{-d} \forall F_{n} \in \mathcal{F}_{T S}$ for $d>(2+\gamma) / \gamma$ and $C<\infty$. By Corollary 1 of de Jong (1997), we have $X_{n i} \rightarrow_{d} N(0,1)$. This and 20.5) give

$$
\begin{equation*}
n^{-1 / 2} \sum_{i=1}^{n} s\left(W_{i}\right) \rightarrow_{d} N\left(0, b^{\prime} h_{5} b\right), \tag{20.9}
\end{equation*}
$$

as desired.

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[^0]:    ${ }^{11}$ Note that $\sup _{j \geq 1, F \in \mathcal{F}_{W U}} \tau_{j F}<\infty$ by the conditions $\left\|W_{F}\right\| \leq M_{1}$ and $\left\|U_{F}\right\| \leq M_{1}$ in $\mathcal{F}_{W U}$ and the moment conditions in $\mathcal{F}$. Thus, $\left\{\tau_{j F_{n}}: n \geq 1\right\}$ does not diverge to infinity, and the "order of magnitude" of $\left\{\tau_{j F_{n}}: n \geq 1\right\}$ refers to whether this sequence converges to zero, and how slowly or quickly it does, when it does converge to zero.

[^1]:    ${ }^{12}$ When the matrix $M_{21 F} \neq 0^{k \times k}$, the argument in 19.2 does not go through because $n^{1 / 2} \widetilde{M}_{21 n}$ does not converge in distribution (since $n^{1 / 2}\left(\widetilde{M}_{21 n}-M_{21 F}\right) \rightarrow_{d} \bar{M}_{21 h}$ by assumption). In this case, one has to alter the definition of $T_{n}^{\dagger}$ so that it rotates the columns of $\widehat{D}_{n}$ before rescaling them. The rotation required depends on both $M_{F}$ and $E_{F} G_{i}$.

[^2]:    ${ }^{13}$ Note that $q^{\dagger}=0$ implies $q=0$ when $p=1$ because $n^{1 / 2} D_{F_{n}}^{\dagger}=n^{1 / 2} M_{F_{n}} E_{F_{n}} G_{i}=O(1)$ when $q^{\dagger}=0$ (by the definition of $\left.q^{\dagger}\right)$ and this implies that $n^{1 / 2} E_{F_{n}} G_{i}=O(1)$ using the first condition in $\mathcal{F}_{K C L R}$. In turn, the latter implies that $n^{1 / 2} \Omega_{F_{n}}^{-1 / 2} E_{F_{n}} G_{i}=O(1)$ using the last condition in $\mathcal{F}$. That is, $q=0$ (since $W_{F}=\Omega_{F}^{-1 / 2}$ and $U_{F}=I_{p}$ because $\widehat{W}_{n}=\widehat{\Omega}_{n}^{-1 / 2}$ and $\widehat{U}_{n}=I_{p}$ in the present case, see the Comment to Lemma 19.2.

