VALIDITY OF SUBSAMPLING AND “PLUG-IN ASYMPTOTIC” INFEERENCE FOR PARAMETERS DEFINED BY MOMENT INEQUALITIES

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This paper considers inference for parameters defined by moment inequalities and equalities. The parameters need not be identified. For a specified class of test statistics, this paper establishes the uniform asymptotic validity of subsampling, \( m \) out of \( n \) bootstrap, and “plug-in asymptotic” tests and confidence intervals for such parameters. Establishing uniform asymptotic validity is crucial in moment inequality problems because the pointwise asymptotic distributions of the test statistics of interest have discontinuities as functions of the true distribution that generates the observations.

The size results are quite general because they hold without specifying the particular form of the moment conditions—only \( 2 + \delta \) moments finite are required. The results allow for independent and identically distributed (i.i.d.) and dependent observations and for preliminary consistent estimation of identified parameters.

1. INTRODUCTION

In this paper, we consider a confidence set (CS) for a true parameter \( \theta_0 \in \Theta \subset \mathbb{R}^d \) whose value is restricted by moment inequalities and equalities. The true parameter need not be identified. There are now numerous examples in the literature that fit into this framework. One way that moment inequalities arise in economic models is from the necessary conditions for Nash equilibria; see, e.g., Ciliberto.
and Tamer (2003), Andrews, Berry, and Jia (2004), Pakes, Porter, Ho, and Ishii (2004), and Bajari, Benkard, and Levin (2007). Moment inequalities also can arise from sufficient conditions for Nash equilibria; see, e.g., Ciliberto and Tamer (2003). Another way they arise is from data censoring, e.g., when a continuous variable is only observed to lie in an interval; see Manski and Tamer (2002). Moon and Schorfheide (2004) provide a macroeconomic example in which a moment inequality appears.

We consider a CS that is obtained by inverting a test that is based on a generalized method of moments–type (GMM-type) criterion function. The method is of Anderson–Rubin type and was first considered in the moment inequality context by Chernozhukov, Hong, and Tamer (2007). Chernozhukov et al. (2007) obtain critical values via subsampling. The present paper shows that for a broad class of test statistics subsampling CSs for the true parameter are uniformly asymptotically valid. The results hold for any specification of the moment functions (subject to $2 + \delta$ moments being finite). The paper also shows that subsampling CSs are not asymptotically conservative. Conditions on the form of the test statistic are given such that uniform validity holds. For example, the results hold for statistics given by the sum of squared negative parts of the normalized sample moment conditions, Gaussian quasi–likelihood ratio (QLR) statistics (also referred to as modified minimum distance statistics), generalized empirical likelihood (GEL) statistics, and a number of other statistics considered in the literature. The results of this paper are the first results available in the literature that establish uniform asymptotic validity of a method of inference for a general class of partially identified models. (A discussion of other results and methods in the literature appears later in this section.)

The results of this paper apply to CSs for the true parameter, as in Imbens and Manski (2004), rather than for the identified set (i.e., the set of points that are consistent with the population moment inequalities), as in Chernozhukov et al. (2007). The reason for this focus is that policy questions based on a structural model in which parameters are restricted by moment inequalities depend on the true parameter, rather than on the identified set. A CS for the identified set typically leads to conservative inference when interest is in the true parameter.

We say that a CS has asymptotically valid size if the limit as $n \to \infty$ of the finite-sample size of the CS is the nominal size. The definition of the finite-sample (or exact) size of a CS is the minimum coverage probability over the parameter space; see (3.5) later in this paper. Hence, the term size necessarily requires uniformity (and the alternative term uniformly asymptotically valid size is redundant). We stress that a CS has asymptotically valid size only if there is uniformity in the asymptotics. In regular models uniformity is often ignored (rightfully) because it holds under reasonable conditions and hence verification is just a technical exercise. This is not the case in the moment inequality model. The reason is that test statistics in this model have pointwise asymptotic distributions that are discontinuous in the true distribution that generates the data—a moment inequality enters the pointwise asymptotic distribution if and only if it holds as an equality.
But, this sharp discontinuity is not a feature of the finite-sample distribution. Discontinuities of this type are responsible for the asymptotic size problems that arise with weak instruments, near integrated processes, post-model selection inference, and parameters that are near a boundary. In such cases, (standard) bootstrap methods typically are not asymptotically valid in a pointwise or uniform sense. Furthermore, Andrews and Guggenberger (2005, 2009a, 2009b, 2010) Mikusheva (2007), and Guggenberger (2008, suppl.) show that even subsampling and \( m \) out of \( n \) bootstrap methods often fail to deliver correct asymptotic size. The results of this paper show that such problems do not arise in the moment inequality example provided subsampling is applied to an appropriate test statistic (and suitable moments exist). This is not true for all statistics. For example, subsampling the endpoints of the estimated set in the moment inequality model is not uniformly asymptotically valid.

The standard method in the literature for obtaining critical values for tests for multivariate one-sided null hypotheses is to use the least favorable asymptotic null distribution evaluated at a consistent estimator of the asymptotic variance matrix. We refer to such tests as plug-in asymptotic (PA) tests. We show that a CS based on a PA test has asymptotically valid size under similar conditions to those for subsampling. The PA critical values are at least as large as the subsampling critical values asymptotically, and in some cases strictly larger, which implies that subsampling CSs can be smaller than PA CSs. The PA CS is not asymptotically conservative provided there are no restrictions on the moment inequalities such that satisfaction of one inequality implies violation of another. But, such restrictions do arise in some examples; see, e.g., Rosen (2008). (In such cases, one could adjust the definition of the PA critical values to take account of the restrictions to obtain a CS that is not conservative.)

Model specification tests are easily constructed based on subsampling or PA CSs. One rejects correct model specification if the CS is empty. Asymptotic validity of the size of such a test follows immediately from the properties of the CS. But, these tests may be asymptotically conservative. See Guggenberger, Hahn, and Kim (2008) for a different test of model specification based on moment inequalities.

Under stated high-level conditions, our results also apply to the case where preliminary consistent estimators of identified parameters are plugged into the sample moment functions. This can be quite useful to reduce the dimension of the parameter under test. For brevity we do not verify the high-level conditions. See Soares (2005) for more primitive conditions.

The asymptotic results given in this paper for subsampling tests with subsample size \( b \) also apply to \( m \) out of \( n \) bootstrap tests with \( m = b \) and independent and identically distributed (i.i.d.) observations provided \( b^2/n \rightarrow 0 \). This is because subsampling based on subsamples of size \( b \) can be viewed as bootstrapping without replacement, which is not too different from bootstrapping with replacement when \( b^2/n \) is small. The subsampling results apply to both i.i.d. and time series observations, whereas the \( m \) out of \( n \) bootstrap results apply only to i.i.d. observations.
We now discuss the related literature. The closest papers are those of Romano and Shaikh (2005, 2008). The present paper and Romano and Shaikh (2008) provide uniformity results for subsampling CSs for the true parameter, whereas Romano and Shaikh (2005) does likewise for CSs for the identified set. Hence, we focus on a comparison of the present paper with Romano and Shaikh (2008).

First, in terms of priority, Romano and Shaikh (2008) established uniformity results for subsampling in two simple examples, a one-sided mean and a two-sided mean, before the moment inequality results of the present paper were established. Next chronologically, the subsampling uniformity results of this paper were obtained. These results apply to any moment condition model, allow for arbitrary moment functions, and apply to general test statistics. Subsequent to the results of the present paper, Romano and Shaikh (2008) established uniformity results for general models and moment functions for a particular nonscale equivariant test statistic.

Next, the present paper covers the following items that are not covered in Romano and Shaikh (2008); (i) the results apply to a general class of scale equivariant test statistics including the QLR statistic, generalized empirical likelihood statistics, and statistics that may depend on a preliminary estimator of an identified parameter; (ii) the results apply to models with both moment inequalities and equalities; (iii) the results apply to time series observations; (iv) explicit expressions are given for the probability limit of the subsampling critical value and the degree of asymptotic nonsimilarity of the subsampling test, and (v) the approach is conducive to the establishment of asymptotic local power results; as has been done in Andrews and Soares (2007). With regard to point (i), Romano and Shaikh (2008) consider a single test statistic that is not scale equivariant, which is highly undesirable in our view. (It implies that an arbitrary rescaling of one moment condition and not another yields a different test statistic from the original one.) The results of Romano and Shaikh (2008) do not cover the QLR statistic, which is the most widely used test statistic in the statistics literature on multidimensional one-sided tests.

Besides Romano and Shaikh (2005, 2008) the only other results in the literature that establish uniform validity of a method for inference with moment inequalities are those of Imbens and Manski (2004), Moon and Schorfheide (2004), Woutersen (2006), and Stoye (2007), and those of Soares (2005), Andrews and Soares (2007), and Fan and Park (2007) using moment selection methods. The results of Imbens and Manski (2004) and Woutersen (2006) are quite restrictive because their Assumption 1 requires (i) superefficiency of the implicit estimator of the length of the identified interval, which holds only in quite special cases (see Stoye, 2007), and (ii) joint asymptotic normality of lower and upper bound estimators \((\hat{\theta}_L, \hat{\theta}_U)\) of the identified interval. Joint estimation of the identified set typically does not yield estimators \((\hat{\theta}_L, \hat{\theta}_U)\) that satisfy asymptotic normality (even univariate asymptotic normality); see Andrews and Han (2009, Sects. 5.1 and 6.1) for simple examples. In consequence, their results do not apply to parameters defined by moment inequalities in general. The results of Stoye (2007)
that circumvent the superefficiency condition also are restrictive because they assume asymptotic normality of \( (\hat{\theta}_r, \hat{\theta}_u) \). Moon and Schorfheide (2004) consider a model in which moment equalities and a moment inequality appear and the parameter of interest is assumed to be point identified by the moment equalities. The Soares (2005), Andrews and Soares (2007), and Fan and Park (2007) uniformity results are obtained using the approach in this paper. Research on the power of tests in the moment inequality model is under way; see Andrews and Soares (2007).

Other papers in the literature that consider inference with moment inequalities include Chernozhukov et al. (2007), Andrews et al. (2004), Pakes et al. (2004), Rosen (2008), Bontemps, Magnac, and Maurin (2007), Bugni (2007), Canay (2007), Beresteau and Molinari (2008), and Galichon and Henry (2008). To date, none of these methods has been shown to be uniformly asymptotically valid. Some of these methods have the disadvantage of being asymptotically conservative (which leads to a larger CS than desired) either all of the time or some of the time. This is true of the methods in Andrews et al. (2004), Pakes et al. (2004), Rosen (2008), and Galichon and Henry (2008). The computational requirements for the different methods vary. For some methods computational simplicity is a comparative advantage.

The results in this paper use the general results given in Andrews and Guggenberger (2010) and generalize these results in two directions that are useful in the moment inequality model and in other models. First, we relax the (partial) product space assumption on the parameter space that is employed in Andrews and Guggenberger (2010) (see Assumption A in that paper). By doing so, the results applied to the moment inequality model allow for cases in which different moment conditions are related, e.g., one moment inequality cannot hold as an equality if some other one does. Restrictions of this type arise frequently in models with data censoring; see, e.g., Rosen (2008). Second, the results provide a larger lower bound on the asymptotic size (defined to be the limit of finite-sample size) of a CS than the results in Andrews and Guggenberger (2010). In many models, both bounds reduce to the same value and equal the upper bound. However, in the moment inequality model, the lower bound given in Andrews and Guggenberger (2010) is not sharp, whereas the lower bound given here is. Finally, the results of Andrews and Guggenberger (2010) are for tests, whereas the results given here are for CSs. This requires uniformity of the results with respect to the parameter of interest and also with respect to nuisance parameters. For tests the former is not required because the parameter of interest is fixed by the null hypothesis.

The general approach to uniformity given here and the way of setting up the moment inequality model to establish uniform results are useful for analyzing the asymptotic size of CSs that employ critical values that are not based on subsampling.

The remainder of the paper is organized as follows. Section 2 discusses the issue of uniformity. Section 3 describes the moment inequality/equality model.
Section 4 states the assumptions. Sections 5 and 6 introduce subsampling CSs and PA CSs, respectively, and show that these CSs are uniformly asymptotically valid for a specified class of test statistics. Section 7 introduces model specification tests. Section 8 discusses extensions to GEL ratio test statistics and test statistics based on preliminary consistent estimators of identified parameters. Section 9 provides general results for the asymptotic size of subsampling CSs. An Appendix contains proofs of the results.

For notational simplicity, throughout the paper we write partitioned column vectors as \( h = (h_1, h_2) \), rather than \( h = (h_1', h_2')' \). Let \( R_+ = \{ x \in R : x \geq 0 \} \), \( R_{+\infty} = R_+ \cup \{+\infty\} \), \( R_{[+\infty]} = R \cup \{+\infty\} \), \( K^p = K \times \cdots \times K \) (with \( p \) copies) for any set \( K \), \( \infty^p = (+\infty, \ldots, +\infty)' \) (with \( p \) copies). Let \( 0_k \) denote a \( k \)-vector of zeros. All limits are as \( n \to \infty \).

Let “pd” abbreviate “positive definite.”

2. UNIFORMITY

We are interested in a CS whose exact (i.e., finite-sample) size is close to its nominal level. By definition, the exact size of the CS is the supremum of its coverage probability over distributions that may generate the data. We use asymptotics to provide an approximation to the exact size. An asymptotic approximation is not necessarily accurate for the exact size if the asymptotic results are not uniform over the distributions that may generate the data. Thus, pointwise asymptotic results are insufficient to asymptotically validate the size of a CS unless they hold uniformly.

When a statistic has a discontinuity in its asymptotic distribution (as a function of a parameter or, more generally, as a function of the distribution generating the data), but not in its finite-sample distribution, pointwise asymptotics do not hold uniformly. The manifestation of this is that asymptotic distributions arise under drifting sequences of parameters that do not arise under pointwise asymptotics. Furthermore, data-dependent critical values may have probability limits under drifting sequences that are different from their probability limits under pointwise asymptotics. This is exactly what happens for tests and CSs in the moment inequality model. Given that pointwise asymptotics do not consider the full range of asymptotic behavior of the CS (which reflects the full range of its finite-sample behavior), asymptotic validity of the size of the CS cannot be established by its behavior under pointwise asymptotics. To determine the limit of the exact size and establish uniform validity, one needs to consider drifting sequences of parameters.

How serious are uniformity issues when a test statistic has a pointwise asymptotic distribution that is discontinuous in the distribution that generates the data? The answer is that they can be very serious. For example, in the weak instrument context, Dufour (1997) has shown that the exact size of the usual nominal 5% test based on the two-stage least squares (2SLS) estimator equals 100%. In a first-order autoregressive (AR(1)) model, the nominal 95% two-sided confidence interval for the autoregressive coefficient \( \rho \in (-1, 1) \) based on the usual normal
critical value has asymptotic size equal to 70% when an intercept is included in the model and 39% if an intercept and a time trend are included; see Andrews and Guggenberger (2009a). In post–model selection inference, Kabaila (1995) shows that a standard nominal 95% confidence interval based on a post–model selection estimator utilizing a consistent model selection procedure has asymptotic size 0%; see Leeb and Pötscher (2005) for related results. All of these problems with standard inference are due to a lack of uniformity.

In problems in which a lack of uniformity arises the (standard) bootstrap typically is even pointwise inconsistent. For example, for the parameter near a boundary case, see Andrews (2000). In the literature on the bootstrap, the usual prescription when the bootstrap is pointwise inconsistent is to use the \( m \) out of \( n \) bootstrap or subsampling; see Andrews and Guggenberger (2010) for references. Politis and Romano (1994) show that subsampling is pointwise consistent under very weak conditions, also see Politis, Romano, and Wolf (1999). Similarly, the \( m \) out of \( n \) bootstrap is pointwise consistent under weak conditions. These results, however, are pointwise asymptotic results. They are not uniform results.

Andrews and Guggenberger (2005, 2009a, 2009b, 2010) show that subsampling and the \( m \) out of \( n \) bootstrap are not necessarily asymptotically valid in a uniform sense. Also see Mikusheva (2007). Furthermore, the problem can be serious. For example, in the weak instrumental variables (IV) case, a nominal 5% equal-tailed two-sided subsampling test based on the 2SLS estimator has “adjusted” asymptotic size of 30% and exact size of 29% when \( n = 120 \), the subsample size \( b \) is 12, and 5 IVs are used; see Andrews and Guggenberger (2005). Furthermore, the exact size gets worse as \( n \to \infty \) and the (unadjusted) asymptotic size is 82%. Similarly, in the AR(1) model, a nominal 95% equal-tailed two-sided subsampling confidence interval has adjusted asymptotic size of 86% and exact size of 87% under normality of the errors when \( n = 130 \), the subsample size is \( b = 12 \), and an intercept is included in the model; see Andrews and Guggenberger (2009a). Again, the exact size gets worse as \( n \to \infty \) and the (unadjusted) asymptotic size is 60%. Subsampling in the post–consistent model selection example does not solve the uniformity problem. Andrews and Guggenberger (2009b) shows that the asymptotic size of a nominal 95% confidence interval in a simple location model is actually 0%.

In the moment inequality model, uniformity issues arise for some procedures when the identified set is sufficiently small that there is a nonnegligible probability of obtaining an estimated set that consists of a singleton. This scenario is of considerable empirical relevance. For example, this is the situation that arises in Andrews et al. (2004) and in both examples in Pakes et al. (2004). Note that the identified set does not have to be a singleton for the problem to arise, it just has to be sufficiently small relative to the sample size. Problems of this sort arise with the bootstrap applied to the interval endpoints (see Andrews, 2005), with subsampling applied to the interval endpoints, and with the procedure in Pakes et al. (2004) based on the pointwise asymptotic distribution of interval endpoints.
3. CONFIDENCE SETS BASED ON MOMENT INEQUALITIES

The moment inequality/equality model is defined as follows. We suppose there exists a true value $\theta_0 \in \Theta \subset \mathbb{R}^d$ that satisfies the moment conditions

$$
E_{F_0} m_j(W_i, \theta_0) \geq 0 \quad \text{for } j = 1, \ldots, p \quad \text{and}
$$

$$
E_{F_0} m_j(W_i, \theta_0) = 0 \quad \text{for } j = p + 1, \ldots, p + v,
$$

(3.1)

where $\{m_j(\cdot, \theta) : j = 1, \ldots, p + v\}$ are (known) real-valued moment functions and $\{W_i : i \geq 1\}$ are observed i.i.d. or stationary random vectors with joint distribution $F_0$. The true value $\theta_0$ is not necessarily identified. Thus, knowledge of $E_{F_0} m_j(W_i, \theta)$ for all $\theta \in \Theta_1$ does not necessarily imply knowledge of $\theta_0$. Furthermore, even knowledge of $F_0$ itself does not necessarily imply knowledge of the true value $\theta_0$. It may require more information than is available in the observed sample $\{W_i : i \leq n\}$ to identify the true parameter $\theta_0$. We are interested in CSs for the true value $\theta_0$.

Let

$$
m(W_i, \theta) = (m_1(W_i, \theta), \ldots, m_k(W_i, \theta))^T,
$$

(3.2)

where $k = p + v$. Let $(\theta, F)$ denote generic values of the parameters. For i.i.d. observations, the parameter space $\mathcal{F}$ for $(\theta, F)$ is the set of all $(\theta, F)$ that satisfy

(i) $\theta \in \Theta$,

(ii) $E_{F} m_j(W_i, \theta) \geq 0$ for $j = 1, \ldots, p$,

(iii) $E_{F} m_j(W_i, \theta) = 0$ for $j = p + 1, \ldots, k$,

(iv) $\{W_i : i \geq 1\}$ are i.i.d. under $F$,

(v) $\sigma_{F,j}^2(\theta) = \text{Var}_{F}(m_j(W_i, \theta)) \in (0, \infty)$ for $j = 1, \ldots, k$,

(vi) $\text{Corr}_{F}(m(W_i, \theta)) \in \Psi$, and

(vii) $E_{F} |m_j(W_i, \theta)/\sigma_{F,j}(\theta)|^{2+\delta} \leq M$ for $j = 1, \ldots, k$,

(3.3)

where $\Psi$ is a specified set of $k \times k$ correlation matrices (see the discussion that follows) and $M < \infty$ and $\delta > 0$ are fixed constants.\textsuperscript{10} For expository convenience, we specify $\mathcal{F}$ for dependent observations in the Appendix; see Section A1.1.

As is standard, we consider a CS obtained by inverting a test. The test is based on a test statistic $T_n(\theta_0)$ for testing $H_0 : \theta = \theta_0$. The nominal level $1 - \alpha$ CS for $\theta$ is

$$
CS_n = \{\theta \in \Theta : T_n(\theta) \leq c_{1-\alpha}(\theta)\},
$$

(3.4)

where $c_{1-\alpha}(\theta)$ is a critical value. We consider subsampling and PA critical values subsequently.

The exact and asymptotic confidence sizes of $CS_n$ are

$$
\text{ExCS}_n = \inf_{(\theta, F) \in \mathcal{F}} P_F(T_n(\theta) \leq c_{1-\alpha}(\theta)) \quad \text{and} \quad \text{AsyCS} = \liminf_{n \to \infty} \text{ExCS}_n,
$$

(3.5)
respectively. The exact maximum coverage probability and the asymptotic maximum coverage probability are

\[
ExMaxCP_n = \sup_{(\theta, F) \in \mathcal{F}} P_F(T_n(\theta) \leq c_{1-\alpha}(\theta)) \quad \text{and} \quad AsyMaxCP = \lim sup_{n \to \infty} ExMaxCP_n. \tag{3.6}
\]

The difference \( AsyMaxCP - AsyCS \) measures the magnitude of asymptotic non-similarity of the CS. For a CS with \( AsyCS \geq 1 - \alpha \), a large value of \( AsyMaxCP \) indicates that the CS may be larger than desirable.

The definition of \( AsyCS \) in (3.5) takes the “inf” before the “lim inf.” In consequence, uniformity over \((\theta, F)\) is built into the definition of \( AsyCS \). Uniformity is necessary for the asymptotic size to provide a good approximation to the finite-sample size of CSs. Andrews and Guggenberger (2005, 2009a, 2009b, 2010) show that when a test statistic has a discontinuity in its limit distribution, as occurs in the moment inequality model, pointwise asymptotics (in which one takes the “lim inf” before the “inf”) can be very misleading in some models.

We consider a general class of test statistics \( T_n(\theta) \) that are defined as follows. The sample moment functions are

\[
m_{n, j}(\theta) = n^{-1} \sum_{i=1}^{n} m_j(W_i, \theta) \quad \text{for} \quad j = 1, \ldots, k. \tag{3.7}
\]

Let \( \hat{\Sigma}_n(\theta) \) be an estimator of the asymptotic variance matrix, \( \Sigma(\theta) \), of \( n^{1/2} \bar{m}_n(\theta) \). When the observations are i.i.d., we take

\[
\hat{\Sigma}_n(\theta) = n^{-1} \sum_{i=1}^{n} (m(W_i, \theta) - \bar{m}_n(\theta))(m(W_i, \theta) - \bar{m}_n(\theta))'. \tag{3.8}
\]

When the observations are dependent, \( \hat{\Sigma}_n(\theta) \) must take this into account. A heteroskedasticity and autocorrelation consistent (HAC) estimator may be required.

The statistic \( T_n(\theta) \) is defined to be of the form

\[
T_n(\theta) = S(n^{1/2} \bar{m}_n(\theta)), \hat{\Sigma}_n(\theta), \tag{3.9}
\]

where \( S \) is a real function on \( \mathbb{R}_{[+\infty]}^p \times \mathbb{R}^v \times \mathcal{V}_{k\times k} \), where \( \mathcal{V}_{k\times k} \) is the space of \( k \times k \) variance matrices. (The set \( \mathbb{R}_{[+\infty]}^p \times \mathbb{R}^v \) contains \( k \)-vectors whose first \( p \) elements are either real or \(+\infty\) and whose last \( v \) elements are real.) The function \( S \) is required to satisfy Assumptions 1–4 stated in Section 4. Examples of functions that do so are now defined.

The first test function \( S \) that we consider is

\[
S_1(m, \Sigma) = \sum_{j=1}^{p} [m_j/\sigma_j]^2 + \sum_{j=p+1}^{p+v} (m_j/\sigma_j)^2, \quad \text{where} \quad [x]_- = \begin{cases} x & \text{if} \ x < 0 \\ 0 & \text{if} \ x \geq 0 \end{cases}, \quad m = (m_1, \ldots, m_k)', \tag{3.10}
\]
and \( \sigma_j^2 \) is the \( j \)th diagonal element of \( \Sigma \). With this function, the parameter space \( \Psi \) for the correlation matrices in condition (vi) of (3.3) is not restricted. That is, (3.3) holds with \( \Psi = \Psi_1 \), where \( \Psi_1 \) contains all \( k \times k \) correlation matrices. The function \( S_1 \) leads to the test statistic

\[
T_n(\theta) = n \sum_{j=1}^{p} \left[ \frac{m_{n,j}(\theta)}{\hat{\sigma}_{n,j}(\theta)} \right]^2 + n \sum_{j=p+1}^{p+v} \left[ \frac{m_{n,j}(\theta)}{\hat{\sigma}_{n,j}(\theta)} \right]^2,
\]

where \( \hat{\sigma}_{n,j}(\theta) = \left[ \hat{\Sigma}_n(\theta) \right]_{jj} \). This is an Anderson–Rubin-type GMM statistic that gives positive weight to moment inequalities only when they are violated. This type of statistic has been considered in Chernozhukov et al. (2007).

The second test function is a Gaussian QLR (or minimum distance) function defined by

\[
S_2(m, \Sigma) = \inf_{t=(t_1,0)v: t_1 \in \mathbb{R}^p_{+}, c} (m - t)' \Sigma^{-1} (m - t). \tag{3.12}
\]

With this function, we restrict the parameter space \( \Psi \) in (3.3). In particular, we take \( \Psi = \Psi_2 \), where \( \Psi_2 \) contains all \( k \times k \) correlation matrices whose determinant is greater than or equal to \( \varepsilon \) for some \( \varepsilon > 0 \). This type of statistic has been considered in numerous papers on tests of inequality constraints (see, e.g., Kudo, 1963; Silvapulle and Sen, 2005, Sect. 3.8) and also in papers in the moment inequality literature (see Manski and Tamer, 2002; and Rosen, 2008).

The following function yields a test with particularly good power against alternatives with \( p_1 (\leq p) \) moment inequalities violated:

\[
S_3(m, \Sigma) = \sum_{j=1}^{p_1} \left[ \frac{m_{j}}{\sigma_j} \right]^2 + \sum_{j=p+1}^{p+v} \left( \frac{m_{j}}{\sigma_j} \right)^2, \tag{3.13}
\]

where \( \left[ \frac{m_{j}}{\sigma_j} \right]^2 \) denotes the \( j \)th largest value among \( \{ \left[ \frac{m_{\ell}}{\sigma_{\ell}} \right]^2 : \ell = 1, \ldots, p \} \) and \( p_1 \) is some specified integer. The function \( S_3 \) satisfies (3.3) with \( \Psi = \Psi_1 \). The function \( S_3 \) is considered in Andrews and Jia (2008). Note that the function \( S_1 \) is a special case of \( S_3 \).

Other test functions \( S \) can be considered that satisfy Assumptions 1–4. For example, one could alter \( S_1 \) or \( S_3 \) by replacing the step function \( [x]_+ \) by a smooth function, by replacing the square by the absolute value to a different positive power (such as one), or by adding weights.

Generally it is not possible to compare the performance of one test function/statistic with that of another without specifying the critical values to be used. The reason is that most critical values, such as the subsampling and PA critical values considered here, are data-dependent and have limits as \( n \to \infty \) that depend on the distribution of the observations. Hence, a given test statistic generates different tests depending on the critical values employed, and the differences do not vanish asymptotically.
The test statistics based on the functions $S_1$ and $S_3$ are easier to compute than those based on $S_2$ because the former are simple functions of the data, whereas the latter involve minimization over $t_1 \in R_{p,\infty}^p$. Computation of $S_2$ requires solving quadratic programming problems. This can be done quickly. But, many computations of the test statistic are required to construct a CS, especially if one is using resampling methods, because (i) one needs to compute tests for an arbitrarily large number of null parameter values $\theta_0$ to construct a CS; (ii) in most cases a different critical value is needed for each null value; and (iii) each critical value requires numerous computations of the test statistic if resampling methods are employed. On the other hand, the function $S_2$ employs information about the correlation matrix $\Omega_1 = D^{-1/2} \Sigma_1 D^{-1/2}$ where $D = \text{Diag}(\Sigma)$, which has power advantages in some cases, whereas $S_1$ and $S_3$ do not.\textsuperscript{13}

One also could consider a test statistic that is the same as $S_1$ but without the division by $\sigma_j$ in each summand. Pakes et al. (2004) and Romano and Shaikh (2008) consider a test statistic of this form. In this case, the uniform asymptotic validity results given subsequently for subsampling and for PA methods can be shown to hold provided $\sigma^2_{F,j}(\theta)$ is bounded away from zero in condition (v) of (3.3). This test statistic is not recommended, however, because it is not equivariant to rescaling of the moment conditions and, hence, is not likely to have good properties in terms of the volume of CSs in general. (In fact, in their empirical applications, Pakes et al., 2004, find that it is desirable to consider an alternative test statistic to the one they first propose that roughly standardizes the variances of the moment conditions.)

4. ASSUMPTIONS

In this section we state Assumptions 1–4 concerning the function $S$ and show that the functions $S_1–S_3$ satisfy them. We also state some assumptions that are not needed for the main results given subsequently but are used for some peripheral results.

Let $B \subset R^w$. We say that a real function $G$ on $R_{[+\infty]}^p \times B$ is continuous at $x \in R_{[+\infty]}^p \times B$ if $y \to x$ for $y \in R^p \times B$ implies that $G(y) \to G(x)$. In the assumptions that follow, the set $\Psi$ is as in condition (vi) of (3.3).\textsuperscript{14} For $p$-vectors $m_1, m_1^*, m_1 < m_1^*$ means that $m_1 \leq m_1^*$ and at least one inequality in the $p$-vector of inequalities holds strictly.

**Assumption 1.**

(a) $S(m_1, m_2, \Sigma)$ is nonincreasing in $m_1$, for all $m_1 \in R^p$, $m_2 \in R^w$, and variance matrices $\Sigma \in R^{k \times k}$.

(b) $S(m, \Sigma) = S(Dm, D\Sigma D)$ for all $m \in R^k$, $\Sigma \in R^{k \times k}$, and pd diagonal $D \in R^{k \times k}$.

(c) $S(m, \Omega) \geq 0$ for all $m \in R^k$ and $\Omega \in \Psi$.

(d) $S(m, \Omega)$ is continuous at all $m \in R_{[+\infty]}^p \times R^w$ and $\Omega \in \Psi$. 


Assumption 2. For all $h_1 \in R^+_{+\infty}$, all $\Omega \in \Psi$, and $Z \sim N(0_k, \Omega)$, the distribution function (df) of $S(Z + (h_1, 0_v), \Omega)$ at $x \in R$ is

(a) continuous for $x > 0$,
(b) strictly increasing for $x > 0$ unless $v = 0$ and $h_1 = \infty^p$, and
(c) less than or equal to $\frac{1}{2}$ at $x = 0$ whenever $v \geq 1$ or $h_1 = 0^p$.

Assumption 3. For some finite $\zeta \leq 0$, $S(m, \Omega) > 0$ if and only if $m_j < \zeta$ for some $j = 1, \ldots, p$ or $m_j \neq 0$ for some $j = p + 1, \ldots, k$, where $m = (m_1, \ldots, m_k)'$ and $\Omega \in \Psi$.

Assumption 4.

(a) The df of $S(Z, \Omega)$ is continuous at its $1 - \alpha$ quantile, $c(\Omega, 1 - \alpha)$, for all $\Omega \in \Psi$, where $Z \sim N(0_k, \Omega)$ and $\alpha \in (0, 1/2)$.
(b) $c(\Omega, 1 - \alpha)$ is continuous in $\Omega$ uniformly for $\Omega \in \Psi$.

In Assumption 2, if an element of $h_1$ equals $+\infty$, then by definition the corresponding element of $Z + (h_1, 0_v)$ equals $+\infty$.

Assumptions 1–4 are shown in Lemma 1, which follows, not to be restrictive.

Assumption 1(a) is the key assumption that is needed to ensure that subsampling CSs have correct asymptotic size. Assumption 1(b) is a natural assumption that specifies that the test statistic is invariant to the scale of each sample moment. Assumptions 1(b) and 1(d) are conditions that enable one to determine the asymptotic properties of $T_n(\theta)$. Assumption 1(c) normalizes the test statistic to be non-negative. Assumptions 2 and 3 are used to show that certain asymptotic dfs satisfy suitable continuity/strictly increasing properties. These properties ensure that the subsampling critical value converges in probability to a constant and the CS has asymptotic size that is not affected by a jump in a df. Assumption 3 implies that $S(\infty^p, \Sigma) = 0$ when $v = 0$. Assumption 4 is a mild continuity assumption.

Lemma 1. The functions $S_1(m, \Sigma) – S_3(m, \Sigma)$ satisfy Assumptions 1–4 with $\Psi = \Psi_1$ for $S_1(m, \Sigma)$ and $S_3(m, \Sigma)$ and with $\Psi = \Psi_2$ for $S_2(m, \Sigma)$.

Remark. In Lemma 1, the function $S_2$ requires the correlation matrices to be bounded away from singularity, whereas none of the other functions requires this.

Next we introduce three conditions that are not needed to show that subsampling and PA CSs have correct asymptotic size (i.e., $\text{AsyCS} \geq 1 - \alpha$). Rather, the first and third conditions are used to show that subsampling and PA CSs, respectively, are not asymptotically conservative (i.e., $\text{AsyCS} \neq 1 - \alpha$). The second condition is used when showing that subsampling CSs have $\text{AsyMaxCP} = 1$ when $v = 0$. 
For \((\theta, F) \in \mathcal{F}\), define \(h_{1,j}(\theta, F) = \infty\) if \(E_F m_j(W_i, \theta) > 0\) and \(h_{1,j}(\theta, F) = 0\) if \(E_F m_j(W_i, \theta) = 0\) for \(j = 1, \ldots, p\). Let \(h_1(\theta, F) = (h_{1,1}(\theta, F), \ldots, h_{1,p}(\theta, F))^\prime\) and \(\Omega(\theta, F) = \lim_{n \to \infty} \text{Corr}_F(n^{1/2} \bar{m}_n(\theta)).\)

**Assumption C1.** For some \((\theta, F) \in \mathcal{F}\), the df of \(S(Z + (h_1(\theta, F), 0)), \Omega(\theta, F))\) is continuous at its \(1 - \alpha\) quantile, where \(Z \sim N(0_\kappa, \Omega(\theta, F)).\)

**Assumption C2.** \(v = 0\) and for some \((\theta, F) \in \mathcal{F}\), \(E_F m_j(W_i, \theta) > 0\) for all \(j = 1, \ldots, p\).

**Assumption C3.** For some \((\theta, F) \in \mathcal{F}\) with \(h_1(\theta, F) = 0_p\), the df of \(S(Z, \Omega(\theta, F))\) is continuous at its \(1 - \alpha\) quantile, where \(Z \sim N(0_\kappa, \Omega(\theta, F)).\)

Assumption C1 is a very weak continuity condition. (Hence, subsampling CSs typically are not asymptotically conservative.) Assumption C2 typically holds if the identified set is not a singleton. Assumption C3 holds quite generally if there are no restrictions relating the expectation of one moment function to that of another. But, if such restrictions exist, then Assumption C3 fails and the PA CS is asymptotically conservative. (Assumption C3 fails when there are restrictions because there is no \((\theta, F) \in \mathcal{F}\) with \(h_1(\theta, F) = 0_p\).) For example, Assumption C3 fails in a regression model in which one only observes the integer part of a latent dependent variable.

### 5. Subsampling Confidence Sets

We now define subsampling critical values and CSs. Let \(b\) denote the subsample size when the full-sample size is \(n\). We assume \(b \to \infty\) and \(b/n \to 0\) as \(n \to \infty\) (throughout the paper). The choice of \(b\) is discussed in the subsampling literature; e.g., see Politis et al. (1999). We do not discuss it further here. (It is beyond the scope of this paper.) The number of different subsamples of size \(b\) is \(q_n\). With i.i.d. observations, there are \(q_n = n!/(n-b)!b!\) different subsamples of size \(b\). With time series observations, there are \(q_n = n-b+1\) subsamples each consisting of \(b\) consecutive observations.

The subsample statistics used to construct the subsampling critical value are \(\{T_{n,b,j}(\theta) : j = 1, \ldots, q_n\}\), where \(T_{n,b,j}(\theta)\) is a subsample statistic defined exactly as \(T_n(\theta)\) is defined but based on the \(j\)th subsample of size \(b\) rather than the full sample. The empirical df and \(1 - \alpha\) sample quantile of \(\{T_{n,b,j}(\theta) : j = 1, \ldots, q_n\}\) are

\[
U_{n,b}(\theta, x) = q_n^{-1} \sum_{j=1}^{q_n} 1(T_{n,b,j}(\theta) \leq x) \quad \text{for} \quad x \in R \quad \text{and}
\]

\[
c_{n,b}(\theta, 1 - \alpha) = \inf\{x \in R : U_{n,b}(\theta, x) \geq 1 - \alpha\}.
\]

The subsampling test rejects \(H_0 : \theta = \theta_0\) if \(T_n(\theta_0) > c_{n,b}(\theta_0, 1 - \alpha)\). The nominal level \(1 - \alpha\) subsampling CS is given by (3.4) with \(c_{1-\alpha}(\theta) = c_{n,b}(\theta, 1 - \alpha)\).
The following theorem applies to i.i.d. observations, in which case \( F \) is defined in (3.3), and to dependent observations, in which case for brevity \( F \) is defined in (A.2)–(A.3) in the Appendix.

**THEOREM 1.** Suppose Assumptions 1–3 hold and \( 0 < \alpha < 1/2 \). Then, the nominal level \( 1 - \alpha \) subsampling CS based on \( T_n(\theta) \) satisfies

(i) \( \text{AsyCS} \geq 1 - \alpha \),

(ii) \( \text{AsyCS} = 1 - \alpha \) if Assumption C1 also holds, and

(iii) \( \text{AsyMaxCP} = 1 \) if \( v = 0 \) (i.e., no moment equalities appear) and Assumption C2 also holds.

**Remarks.**

1. An important feature of Theorem 1 is that no assumptions are placed on the moment functions \( m(W_i, \theta) \) beyond the existence of mild moment conditions (e.g., \( 2 + \delta \) moments finite in the i.i.d. case) that appear in the definition of \( F \) and Assumption C2 that is used in Theorem 1(iii). Thus, the results apply to moment conditions based on instruments that are weak. (The reason is that the test statistics considered are of the Anderson–Rubin type.)

2. The asymptotic distribution of \( T_n(\theta) \) differs depending on the sequence of true values \( \{(\theta_n, F_n) \in F : n \geq 1\} \) considered. The Appendix provides explicit expressions for \( \text{AsyCS} \) and \( \text{AsyMaxCP} \) in terms of these asymptotic distributions; see (A.13) and below (A.13) and the definitions in (A.10) and (9.3). Hence, \( \text{AsyMaxCP} \) can be evaluated in cases in which \( v \geq 1 \).

3. The results of Theorem 1 hold even when there are restrictions on the moment inequalities such that when one moment inequality holds as an equality then another moment inequality cannot. Restrictions of this sort arise in a variety of models. For example, they arise in a location model with interval outcomes. In this model, \( y_i \) is observed, \( y_i^* \) and \( u_i \) are not observed, \( y_i^* = \theta_0 + u_i \) for \( i = 1, \ldots, n \), \( y_i = [y_i^*] \) (i.e., \( y_i \) equals the integer part of \( y_i^* \)), and \( u_i \) has mean zero. The interval outcome \( [y_i, y_i + 1] \) necessarily includes the unobserved outcome variable \( y_i^* \). Two moment inequalities that place bounds on \( \theta_0 \) are (i) \(-E_{\theta_0, y_i} + \theta_0 \geq 0 \) and (ii) \( E_{\theta_0, y_i} + 1 - \theta_0 \geq 0 \). Obviously, both inequalities cannot simultaneously hold as equalities. Subsampling automatically takes this into account and generates a (data-dependent) critical value that is smaller than what one would obtain if no functional relationship existed between the two moment functions. This yields a CS that is smaller than otherwise, as is desirable.

The subsample statistic can be defined using a recentering, and the uniform asymptotic validity results go through with some additional effort. In fact, based on finite-sample simulations of size and power (not reported here), this is the version of subsampling that we recommend for moment inequality models. The recentered subsample statistic \( \hat{T}_{n,b,j}(\theta) \) is defined to be
\[ \hat{T}_{n,b,j}(\theta) = S(b^{1/2}(\hat{m}_{n,b,j}(\theta) - \overline{m}_n(\theta)), \hat{\Sigma}_{n,b,j}(\theta)), \]  
\tag{5.2}

where \( \overline{m}_{n,b,j}(\theta) \) is the sample average based on the observations in the \( j \)th sub-sample and \( \hat{\Sigma}_{n,b,j}(\theta) \) is the variance matrix estimator based on the observations in the \( j \)th subsample.

Chernozhukov and Fernandez-Val (2005) consider a recentered subsampling method in the context of inference for quantile processes. Simulation results in Linton, Maasoumi, and Whang (2005) for a different testing problem (viz., tests of stochastic dominance) find that the recentered subsampling test performs substantially worse in terms of both size and power than the subsampling test without recentering. Based on the simulations we have done, however, this is not the case in moment inequality models.

6. PLUG-IN ASYMPTOTIC CONFIDENCE SETS

Next, we discuss CSs based on an asymptotic critical value. The least favorable asymptotic null distributions of the statistic \( T_n(\theta) \) are shown to be those for which the moment inequalities hold as equalities. These distributions depend on the (asymptotic) correlation matrix, \( \Omega \), of the moment functions. We consider PA critical values that are obtained from the least favorable asymptotic null distribution evaluated at a consistent estimator of \( \Omega \). Critical values of this type have long been considered in the literature on multivariate one-sided tests; see Silvapulle and Sen (2005) for references. They have been considered in the moment inequality literature by Rosen (2008). We exploit results in Andrews and Guggenberger (2009a) for “plug-in size-corrected fixed critical values” to obtain the asymptotic results given in this section.

As before, let \( c(\Omega, 1 - \alpha) \) denote the \( 1 - \alpha \) quantile of \( S(Z, \Omega) \), where \( Z \sim N(0_k, \Omega) \). This is the \( 1 - \alpha \) quantile of the asymptotic null distribution of \( T_n(\theta) \) when the moment inequalities hold as equalities. Define

\[ \hat{\Omega}_n(\theta) = \hat{D}_n^{-1/2}(\theta) \hat{\Sigma}_n(\theta) \hat{D}_n^{-1/2}(\theta), \]  
\tag{6.1}

where \( \hat{D}_n(\theta) = \text{Diag}(\hat{\Sigma}_n(\theta)) \) and \( \hat{\Sigma}_n(\theta) \) is defined in (3.8) for i.i.d. observations and is a consistent estimator of \( \lim_{n \to \infty} \text{Var}(n^{1/2}\hat{m}_n(\theta)) \) for dependent observations.

The nominal \( 1 - \alpha \) PA CS is given by (3.4) with critical value \( c_{1-\alpha}(\theta) \) equal to \( c(\hat{\Omega}_n(\theta), 1 - \alpha) \).

\[ \text{THEOREM 2.} \quad \text{Suppose Assumptions 1 and 4 hold and} \ 0 < \alpha < 1/2. \ \text{Then, the nominal level} \ 1 - \alpha \ \text{PA CS based on} \ T_n(\theta) \ \text{satisfies} \]

(i) \( \text{AsyCS} \geq 1 - \alpha \) and
(ii) \( \text{AsyCS} = 1 - \alpha \) provided Assumption C3 also holds.
Remark. Theorem 2(i) holds even when there are restrictions on the moment inequalities such that when one moment inequality holds as an equality then another moment inequality cannot. However, Theorem 2(ii) does not hold in this case because Assumption C3 fails. The PA critical value does not automatically take functional relationships between the moment functions into account as the subsampling critical value does. The PA critical value is larger than necessary and the PA CS is asymptotically conservative in this scenario. Thus, subsampling CSs have advantages over PA CSs in this scenario.

7. MODEL SPECIFICATION TESTS

Tests of model specification can be constructed using the subsampling and PA CSs introduced in Sections 5 and 6. The null hypothesis of interest is that there exists a parameter \( \theta_0 \in \Theta \) such that (3.1) holds (with additional conditions specified by the parameter space for \((\theta, F)\), such as those in (3.3) or those given in the Appendix for temporally dependent observations). The idea of such specification tests is the same as for the \( J \) test of overidentifying restrictions in GMM; see Hansen (1982). With the \( J \) test, one rejects the null hypothesis of correct model specification if the GMM criterion function evaluated at the GMM estimator is sufficiently large. In the moment inequality/equality model, the analogue of the GMM criterion function is the test statistic \( T_n(\theta) \). By definition, the subsampling test rejects the model specification if \( T_n(\theta) \) exceeds the subsampling critical value \( c_{n,b}(\theta, 1 - \alpha) \) for all \( \theta \in \Theta \). Equivalently, it rejects if the subsampling CS is empty. The PA model specification test is analogous with the PA critical value in place of the subsampling critical value.

If the model specified in (3.1) is correctly specified, then the subsampling CS and the PA CS contain the true value with asymptotic probability \( 1 - \alpha \) (or greater) uniformly over the parameter space. Hence, under the null hypothesis of correct model specification, the limit as \( n \to \infty \) of the finite-sample size of the subsampling and PA model specification tests are \( \leq \alpha \) under the assumptions of Theorems 1(i) and 2(i), respectively. Note that the model specification tests may be asymptotically conservative (i.e., have asymptotic size \( < \alpha \)) even when the assumptions of part (ii) of those theorems hold. As discussed earlier, it is crucial that the asymptotic sizes of these tests are shown to be valid uniformly over the parameter space because the present testing scenario is one in which the test statistic \( T_n(\theta) \) has a limit distribution that is discontinuous in the parameters.

8. EXTENSIONS


Here we consider CSs for parameters in the moment inequality/equality model based on a GEL test statistic, \( T_{nGEL}(\theta) \), rather than a test statistic of the form \( T_n(\theta) \) in (3.9). In the context of moment equalities, Smith (1997) considers GEL statistics. In the context of moment inequalities, Soares (2006) considers GEL
statistics, and Moon and Schorfheide (2004), Otsu (2006), and Canay (2007) consider empirical likelihood (EL) statistics. In the Appendix, we show that the asymptotic distribution of $T_n^{GEL}(\theta)$ is the same as that of the statistic $T_n(\theta)$ in (3.9) based on $S_2(m, \Sigma)$ in (3.12) when the observations come from a row-wise i.i.d. triangular array. In consequence, by the same argument as for the latter statistic, GEL-based subsampling and PA CSs based on $T_n^{GEL}(\theta)$ have correct asymptotic size.

For $t \in \mathbb{R}^p$, define

$$m_i(t, \theta) = m(W_i, \theta) - (t, 0)_n.$$  \hfill (8.1)

The vector $t$ can be viewed as an additional nuisance parameter that captures the slackness in the first $p$ moment inequalities. The minimum distance (MD) formulation of the EL statistic for inference under moment equalities and inequalities is given by

$$L_{EL}(\theta, t) = \sup_{\pi} \left\{ \prod_{i=1}^n \pi_i : \pi_i \geq 0, \sum_{i=1}^n \pi_i = 1, \sum_{i=1}^n \pi_i m_i(t, \theta) = 0 \right\},$$  \hfill (8.2)

where $\pi = (\pi_1, \ldots, \pi_n)'$. Under weak additional assumptions, the MD formulation of the EL estimator $\hat{\theta}_{EL} = \arg\max_{\theta \in \Theta} \sup_{t \in \mathbb{R}^p} L_{EL}(\theta, t)$ can be reexpressed (equivalently) as the solution to a saddlepoint problem $\hat{\theta}_{EL} = \arg\min_{\theta \in \Theta} \inf_{t \in \mathbb{R}^p} \sup_{\lambda \in \hat{\Lambda}_n(t, \theta)} n \sum_{i=1}^n \ln(1 - \lambda'_i m_i(t, \theta))$, where $\hat{\Lambda}_n(t, \theta)$ is defined in (8.3). We consider the GEL generalization of this saddlepoint problem and work with the statistic

$$T_n^{GEL}(\theta) = \inf_{t \in \mathbb{R}^p} \sup_{\lambda \in \hat{\Lambda}_n(t, \theta)} n \hat{P}_\rho(t, \theta, \lambda),$$

where

$$\hat{P}_\rho(t, \theta, \lambda) = 2n^{-1} \sum_{i=1}^n (\rho(\lambda'_i m_i(t, \theta)) - \rho(0)),$$

and

$$\hat{\Lambda}_n(t, \theta) = \{ \lambda \in \mathbb{R}^k : \lambda'_i m_i(t, \theta) \in Q \text{ for } i = 1, \ldots, n \}.$$  \hfill (8.3)

$Q$ is an open interval of the real line that contains 0, and $\rho : Q \rightarrow \mathbb{R}$ is a concave function that is twice continuously differentiable on a neighborhood of 0 with first and second derivatives at 0 normalized to equal $-1$. For $\rho(x) = \ln(1 - x)$ we obtain the EL estimator; for $\rho(x) = -(1 + x)^2/2$, we obtain the continuous updating estimator; and for $\rho(x) = -\exp(-x)$, we obtain the exponential tilting estimator. In the Appendix we show that under the i.i.d. setup of (3.3) and Assumption GEL the equivalents of Theorems 1 and 2 hold for CSs based on $T_n^{GEL}(\theta)$ rather than $T_n(\theta)$ in (3.9).

8.2. Preliminary Estimation of Identified Parameters

Suppose the population moment functions are of the form $E_F m_j(W_i, \theta_0, \tau_0) \geq 0$ for $j = 1, \ldots, p$ and $E_F m_j(W_i, \theta_0, \tau_0) = 0$ for $j = p + 1, \ldots, k$, where $\tau_0$ is a
parameter for which a preliminary asymptotically normal estimator \( \hat{\tau}_n(\theta_0) \) exists. Of course, this typically requires that \( \tau_0 \) is identified. Soares (2005) considers this scenario in some detail. The sample moment functions in this case are of the form 
\[
m_{n,j}(\theta) = n^{-1} \sum_{i=1}^n m_{j}(W_i, \theta, \hat{\tau}_n(\theta)).
\]
The asymptotic variance of \( n^{1/2} m_{n,j}(\theta) \) is different when \( \tau_0 \) is replaced by \( \hat{\tau}_n(\theta) \), and hence \( \hat{\Sigma}_n(\theta) \) needs to take this into account, but otherwise the theoretical treatment of this model is the same. In fact, Theorems 1 and 2 hold in this case using the conditions given in (A.3) of the Appendix. These are high-level conditions that essentially just require that \( m_{n,j}(\theta, \hat{\tau}_n(\theta)) \) is asymptotically normal (after suitable normalization) under certain drifting sequences of parameters.

9. GENERAL RESULTS FOR SUBSAMPLING CONFIDENCE SETS

This section provides general results for CSs. These results are used in the Appendix to prove Theorem 1 for subsample CSs in the moment inequality model. Let \( R_\infty = R \cup \{\pm \infty\} \).

9.1. Definition of Confidence Sets

We consider CSs for a parameter \( \theta \in R^d \) when nuisance parameters \( \eta \in R^s \) and \( \gamma_3 \in T_3 \) may appear, where \( T_3 \) is an arbitrary, possibly infinite-dimensional, space. We obtain CSs for \( \theta \) by inverting tests based on a test statistic \( T_n(\theta_0) \) for testing the null hypothesis \( H_0 : \theta = \theta_0 \). Fixed and subsampling critical values are considered. Let \( \Theta \subseteq R^p \) denote the parameter space for \( \theta \). The CS for \( \theta \) is defined as in (3.4). The focus of this section is on the behavior of CSs when the asymptotic distribution of \( T_n(\theta) \) depends on the parameters \( (\theta, \eta) \) and is discontinuous at some value(s) of \( (\theta, \eta) \).

We partition \( \theta \) and \( \eta \) into \( (\theta_1, \theta_2) \) and \( (\eta_1, \eta_2) \), where \( \theta_j \in R^{d_j} \) and \( \eta_j \in R^{s_j} \) for \( j = 1, 2 \). By definition (made precise subsequently), \( \gamma_1 = (\theta_1, \eta_1) \) are parameters that determine how close the asymptotic distribution of \( T_n(\theta) \) is to a point of discontinuity and \( \gamma_2 = (\theta_2, \eta_2) \) are parameters that do not do so but still may affect the asymptotic distribution of \( T_n(\theta) \). The parameter \( \gamma_3 \) does not affect the asymptotic distribution of \( T_n(\theta) \). Define \( \gamma = (\gamma_1, \gamma_2, \gamma_3) \), where \( \gamma_1 \in R^p, \gamma_2 \in R^q, \gamma_3 \in T_3 \), \( p = d_1 + s_1 \), and \( q = d_2 + s_2 \). Let \( \Gamma \) denote the parameter space for \( \gamma \). In most models, either no parameter \( \theta_1 \) or \( \theta_2 \) appears (i.e., \( d_1 = 0 \) or \( d_2 = 0 \)). For example, in the moment inequality model, \( d_1 = 0 \).

The terms \( ExCS_n, AsyCS, ExMaxCP_n, \) and \( AsyMaxCP \) are defined as in (3.5) and (3.6) with \( \gamma \in \Gamma \) in place of \( (\theta, F) \in \mathcal{F} \).

9.2. Critical Values

A test rejects the null hypothesis when \( T_n(\theta_0) \) exceeds some critical value. We consider two types of critical values for use with the test statistic \( T_n(\theta_0) \). The first is a fixed critical value (FCV) and is denoted \( c_{Fix}(\theta_0, 1 - \alpha) \), where \( \alpha \in (0, 1) \) is the nominal size of the FCV test. The FCV test rejects \( H_0 \) when \( T_n(\theta_0) > c_{Fix}(\theta_0, 1 - \alpha) \).
A common choice is $c_{F_{1S}}(\theta_0, 1-\alpha)$, where $c_{F_{1S}}(\theta_0, 1-\alpha)$ denotes the $1-\alpha$ quantile of $J_{1S}$ and $J_{1S}$ is the asymptotic null distribution of $T_n(\theta_0)$ when $\gamma$ is fixed and is not a point of discontinuity. (Of course, this choice only applies when $T_n(\theta_0)$ has the same asymptotic distribution for all fixed $\gamma$ that are not points of discontinuity.) Another choice is $c_{F_{1S}}(\theta_0, 1-\alpha)$ equals the $1-\alpha$ quantile of a least favorable asymptotic null distribution, which may occur at a point of discontinuity.

Let $b$ and $q_n$ be as in Section 5. The subsample statistics that are used to construct the subsampling critical value are denoted by $\{T_{n,b,j}(\theta_0) : j = 1, \ldots, q_n\}$ when testing $H_0 : \theta = \theta_0$.

Let $\{T_{n,b,j}(\theta) : j = 1, \ldots, q_n\}$ be subsample statistics that are defined exactly as $T_n(\theta)$. In most cases, the subsample statistics $\{\hat{T}_{n,b,j}(\theta_0) : j = 1, \ldots, q_n\}$ are defined to satisfy one or the other of the following assumptions.

**Assumption Sub1.** $\hat{T}_{n,b,j}(\theta_0) = T_{n,b,j}(\hat{\theta}_n)$ for all $j \leq q_n$, where $\hat{\theta}_n$ is an estimator of $\theta$.

**Assumption Sub2.** $\hat{T}_{n,b,j}(\theta_0) = T_{n,b,j}(\theta_0)$ for all $j \leq q_n$.

The estimator $\hat{\theta}_n$ in Assumption Sub1 usually is chosen to be an estimator that is consistent under both the null and alternative hypotheses. In the moment inequality example, the subsample statistics are defined such that Assumption Sub2 holds—because we do not assume that $\theta$ is identified and hence no consistent estimator $\hat{\theta}_n$ is available.

Let $L_{n,b}(\theta, x)$ and $c_{n,b}(\theta, 1-\alpha)$ denote the empirical df and $1-\alpha$ sample quantile, respectively, of the subsample statistics $\{\hat{T}_{n,b,j}(\theta) : j = 1, \ldots, q_n\}$. By definition,

$\begin{align*}
L_{n,b}(\theta, x) &= q_n^{-1} \sum_{j=1}^{q_n} 1(\hat{T}_{n,b,j}(\theta) \leq x) \quad \text{for } x \in \mathbb{R} \quad \text{and} \\
c_{n,b}(\theta, 1-\alpha) &= \inf\{x \in \mathbb{R} : L_{n,b}(\theta, x) \geq 1-\alpha\}.
\end{align*}$

The subsampling test rejects $H_0 : \theta = \theta_0$ if $T_n(\theta_0) > c_{n,b}(\theta_0, 1-\alpha)$.

### 9.3. Parameter Space

The parameter space for $\gamma$ is $\Gamma$.

**Assumption A0.** $\Gamma \subset \{(\gamma_1, \gamma_2, \gamma_3) : \gamma_1 \in \mathbb{R}^p, \gamma_2 \in \mathbb{R}^q, \gamma_3 \in \mathbb{T}_3\}$.

In contrast to Assumption A of Andrews and Guggenberger (2010), Assumption A0 does not require $\gamma_1$ to lie in a product space in $\mathbb{R}^p$ and does not require $\gamma_1, \gamma_2$ to lie in a product space of the form $\Gamma_1 \times \Gamma_2$ for some $\Gamma_1 \subset \mathbb{R}^p$ and $\Gamma_2 \subset \mathbb{R}^q$. The latter product space condition is typically violated in the moment inequality model. The former product space condition is sometimes violated.
in the moment inequality model. The relaxation of Assumption A of Andrews and Guggenberger (2010) to Assumption A0 is a substantial contribution of this paper. It is useful in a variety of models beyond the moment inequality model.

### 9.4. Convergence Assumption

For an arbitrary distribution $G$, let $G(\cdot)$ denote the df of $G$, let $G(x \leftarrow \cdot)$ denote the limit from the left of $G(\cdot)$ at $x$, and let $C(G)$ denote the set of continuity points of $G(\cdot)$. Define the $1 - \alpha$ quantile, $q(1 - \alpha)$, of a distribution $G$ by $q(1 - \alpha) = \inf\{x \in R : G(x) \geq 1 - \alpha\}$. The distribution $J_h$ considered subsequently is the distribution of a proper random variable that is finite with probability one.

Let $r > 0$ denote a rate of convergence index such that when the true parameter $\gamma_1$ satisfies $n^{r}\gamma_1 \to h_1$, then the test statistic $T_n(\theta_0)$ has an asymptotic distribution that depends on the localization parameter $h_1$ (see Assumption B0, which follows). In most examples, including the moment inequality example, $r = 1/2$. For a given model, we assume there is a single fixed $r > 0$.

Let $\{w_n : n \geq 1\}$ denote some subsequence of $\{n\}$. Given $\{w_n\}$, we consider sequences of parameters with the following properties.

**Definition of $\{\gamma_{w_n,h} : n \geq 1\}$.** Given $r > 0$ and $h = (h_1, h_2) \in R^2_\infty \times R^2_\infty$, let $\{\gamma_{w_n,h} = (\gamma_{w_n,h,1}, \gamma_{w_n,h,2}, \gamma_{w_n,h,3}) : n \geq 1\}$ denote a sequence of parameters in $\Gamma$ for which $w_n^{r}\gamma_{w_n,h,1} \to h_1$, $\gamma_{w_n,h,2} \to h_2$, $\gamma_{w_n,h} = (\theta_{w_n,h,1}, \eta_{w_n,h,1}), (\theta_{w_n,h,2}, \eta_{w_n,h,2}), (\theta_{w_n,h,3})$, and $\theta_{w_n,h} = (\theta_{w_n,h,1}, \theta_{w_n,h,2})$ if such a sequence exists.

Define

$$H = \{h \in R^2_\infty \times R^2_\infty : \exists \text{ a subsequence } \{w_n\} \text{ and a sequence } \{\gamma_{w_n,h} : n \geq 1\}\}.$$  

(9.2)

The sequence $\{\gamma_{w_n,h} : n \geq 1\}$ is defined such that under $\{\gamma_{w_n,h} : n \geq 1\}$, the asymptotic distribution of $T_{w_n}(\theta_{w_n,h})$ depends on $h$ and only $h$.

**Assumption B0.** For some $r > 0$, all $h \in H$, all subsequences $\{w_n\}$ of $\{n\}$, all sequences $\{\gamma_{w_n,h} : n \geq 1\}$, and some distributions $J_h$, $T_{w_n}(\theta_{w_n,h}) \to_d J_h$ under $\{\gamma_{w_n,h} : n \geq 1\}$, where $\gamma_{w_n,h} = (\theta_{w_n,h,1}, \eta_{w_n,h,1}), (\theta_{w_n,h,2}, \eta_{w_n,h,2}), (\theta_{w_n,h,3})$ and $\theta_{w_n,h} = (\theta_{w_n,h,1}, \theta_{w_n,h,2})$.

Assumption B0 is a strengthening of Assumption B of Andrews and Guggenberger (2010) to cover subsequences $\{w_n\}$ rather than just sequences $\{n\}$. Also it differs slightly from Assumption B of Andrews and Guggenberger (2010) because it applies to CSs rather than tests. Although more complicated, Assumption B0 is usually not more difficult to verify than Assumption B. When Assumption A of Andrews and Guggenberger (2010) holds, Assumptions B and B0 are equivalent; see the proof of (8.6) of Andrews and Guggenberger (2010). Assumption B0 holds in a wide variety of examples of interest; see the Appendix for the moment inequality model and Andrews and Guggenberger (2005, 2009a, 2009b, 2010) for other models.
9.5. Subsampling Assumptions

Theorem 3 in section 9.6 shows that the asymptotic size of a subsampling CS is determined by the asymptotic distributions of the full-sample statistic \( T_{w_h}(\theta_{w_h}) \) and the subsample statistic \( T_{w_h,b_{w_h}}(\theta_{w_h}) \) under certain parameter sequences \( \{\gamma_{w_n}, g, h : n \geq 1\} \). By Assumption B0, the asymptotic distribution of \( T_{w_h}(\theta_{w_h}) \) is \( J_h \). The asymptotic distribution of \( T_{w_h,b_{w_h}}(\theta_{w_h}) \) under \( \{\gamma_{w_n}, g, h : n \geq 1\} \) is shown to be \( J_g \) for \( g \in H \).

**Definition of \( \{\gamma_{w_n}, g, h : n \geq 1\} \)**. Given \( r > 0 \), \( g = (g_1, g_2) \in R^p_\infty \times R^p_\infty \), and \( h = (h_1, h_2) \in R^p_\infty \times R^p_\infty \) with \( g_2 = h_2 \), let \( \{\gamma_{w_n}, g, h = (\gamma_{w_n}, g, h, 1, \gamma_{w_n}, g, h, 2, \gamma_{w_n}, g, h, 3) : n \geq 1\} \) denote a sequence of parameters in \( \Gamma \) for which \( \gamma_{w_n}, g, h, 1 \to h_1 \), \( \gamma_{w_n}, g, h, 2 \to h_2 \), and \( \gamma_{w_n}, g, h = (\theta_{w_n}, g, h, 1, \eta_{w_n}, g, h, 1), (\theta_{w_n}, g, h, 2, \eta_{w_n}, g, h, 2), \gamma_{w_n}, g, h, 3) \), and \( \theta_{w_n}, g, h = (\theta_{w_n}, g, h, 1, \theta_{w_n}, g, h, 2) \) if such a sequence exists.

By definition, a sequence \( \{\gamma_{w_n}, g, h : n \geq 1\} \) also is of the form \( \{\gamma_{w_h}, g, h : n \geq 1\} \). The index set of the asymptotic distributions of \( T_{w_h}(\theta_{w_h}) \) and \( T_{w_h,b_{w_h}}(\theta_{w_h}) \) under sequences \( \{\gamma_{w_n}, g, h : n \geq 1\} \) is denoted by \( GH \). By definition,

\[
GH = \{(g, h) \in (R^p_\infty \times R^p_\infty)^2 : \\
\exists \text{ a subsequence } \{w_n\} \text{ and a sequence } \{\gamma_{w_n}, g, h : n \geq 1\} \}. \tag{9.3}
\]

By definition of \( \{\gamma_{w_n}, g, h : n \geq 1\} \) and Assumption C, which follows (i.e., \( b/n \to 0 \)), for all \( (g, h) = ((g_1, g_2), (h_1, h_2)) \in GH \), we have \( g_2 = h_2 \) and \( |g_1, j| \leq |h_1, j| \) for \( j = 1, \ldots, p \), where \( g_1 = (g_{1,1}, \ldots, g_{1,p})' \) and \( h_1 = (h_{1,1}, \ldots, h_{1,p})' \).

For subsampling CSs, we require the following additional assumptions.

**Assumption C.** (a) \( b \to \infty \) and (b) \( b/n \to 0 \).

**Assumption D.**

(a) \( \{T_{n,b,j}(\theta) : j = 1, \ldots, q_n\} \) are identically distributed under any \( \gamma \in \Gamma \) for all \( n \geq 1 \) and

(b) \( T_{n,b,j}(\theta) \) and \( T_{b}(\theta) \) have the same distribution under any \( \gamma \in \Gamma \) for all \( n \geq 1 \), where \( \theta = (\theta_1, \theta_2) \) and \( \gamma = ((\theta_1, \eta_1), (\theta_2, \eta_2), \gamma_3) \).

**Assumption E0.** For all subsequences \( \{w_n\} \) of \( \{n\} \) and all sequences \( \{\gamma_{w_n}, g, h \in \Gamma : n \geq 1\} \), \( U_{w_n,b_{w_n}}(\theta_{w_n}, g, h, x) \to -E_{\gamma_{w_n}, g, h} U_{w_n,b_{w_n}}(\theta_{w_n}, g, h, x) \to 0 \) under \( \{\gamma_{w_n}, g, h : n \geq 1\} \) for all \( x \in R \), where \( \theta_{w_n, g, h} = ((\theta_{w_n, g, h, 1}, \theta_{w_n, g, h, 2}) \text{ and } \gamma_{w_n, g, h} = ((\theta_{w_n, g, h, 1}, \eta_{w_n, g, h, 1}), (\theta_{w_n, g, h, 2}, \eta_{w_n, g, h, 2}), \gamma_{w_n, g, h, 3}) \).

**Assumption F.** For all \( \epsilon > 0 \) and \( h \in H \), \( J_h(c_h(1-\alpha) + \epsilon) > 1 - \alpha \), where \( c_h(1-\alpha) \) is the \( 1 - \alpha \) quantile of \( J_h \).
Assumption G0. For all \( h = (h_1, h_2) \in H \), all subsequences \( \{w_n\} \) of \( [n] \), and all sequences \( \{\gamma_{w_n,g,h} : n \geq 1\} \), if \( U_{w_n,b_{w_n}}(\theta_{w_n,g,h}, x) \to_p J_g(x) \) under \( \{\gamma_{w_n,g,h} : n \geq 1\} \) for all \( x \in C(J_g) \), then \( L_{w_n,b_{w_n}}(\theta_{w_n,g,h}, x) - U_{w_n,b_{w_n}}(\theta_{w_n,g,h}, x) \to_p 0 \) under \( \{\gamma_{w_n,g,h} : n \geq 1\} \) for all \( x \in C(J_g) \), where \( \theta_{w_n,g,h} = (\theta_{w_n,g,h,1}, \theta_{w_n,g,h,2}) \) and \( \gamma_{w_n,g,h} = ((\theta_{w_n,g,h,1}, \eta_{w_n,g,h,1}), (\theta_{w_n,g,h,2}, \eta_{w_n,g,h,2}), \gamma_{w_n,g,h,3}) \).

Assumptions C, D, and F are the same as in Andrews and Guggenberger (2010). Assumptions E0 and G0 are extensions of Assumptions E and G of Andrews and Guggenberger (2010) to cover subsequences \( \{w_n\} \) rather than just full sequences \( [n] \). Assumptions C and D are standard assumptions in the subsampling literature and are not restrictive. Assumption D necessarily holds when the observations are i.i.d. or stationary and the subsamples are constructed in the usual way. Assumption E0 holds automatically when the observations are i.i.d. for each fixed \( \gamma \in \Gamma \) or are stationary strong mixing for each fixed \( \gamma \in \Gamma \) and \( \sup_{\gamma \in \Gamma} a_{\gamma}(j) \to 0 \) as \( j \to \infty \), where \( \{a_{\gamma}(j) : j \geq 1\} \) are the strong mixing numbers of the observations when the true parameter is \( \gamma \); see Andrews and Guggenberger (2010). Assumption F is not restrictive. It holds in all of the examples considered in Andrews and Guggenberger (2005, 2009a, 2009b, 2010). Assumption G0 holds automatically when Assumption Sub2 holds (because \( L_n,h(\cdot, \cdot) = U_n,h(\cdot, \cdot) \) for all \( n, b \)), as occurs with the moment inequality subsample statistics considered in this paper. In Andrews and Guggenberger (2010), sufficient conditions for Assumption G are given when Assumption Sub1 holds. These can be extended to provide sufficient conditions for Assumption G0.

9.6. Asymptotic Results

Theorem 3 in this section is a CS analogue of the testing results of Theorem 1 of Andrews and Guggenberger (2010) but with two improvements that are needed in the moment inequality example. The first improvement is that the product space form of \( \Gamma_1 \) and \( \Gamma_1 \times \Gamma_2 \) is eliminated (Assumption A of Andrews and Guggenberger, 2010, is replaced by Assumption A0). This extension is useful in many models. The price to pay for this extension is the more complicated form of \( GH \) here than in Andrews and Guggenberger (2010) and the more complicated forms of Assumptions B0, E0, and G0, which involve subsequences \( \{w_n\} \), than Assumptions B, E, and G of Andrews and Guggenberger (2010).

The second improvement is that Theorem 3 provides a larger lower bound on \( A_{yCS} \) than does the straight analogue of Theorem 1 of Andrews and Guggenberger (2010). In most examples, continuity of \( J_h(x) \) at suitable values of \( (h, x) \) yields the lower and upper bounds given in Theorem 1 of Andrews and Guggenberger (2010) to be equal, and, hence, the latter delivers the precise value of asymptotic size. This continuity does not hold in the moment inequality example when \( v = 0 \). We introduce an improvement that is applicable in models in which \( J_h(x) \) has a discontinuity at \( x = c_g(1 - \alpha) \) for some \( (g, h) \in GH \) and the test statistic and the subsample statistics have a common lower bound on their
support for all \( n \geq 1 \). The improvement is possible because the test statistic and the subsampling critical values cannot be smaller than the lower bound.

Let \( GH^* \) be the set of points \((g, h) \in GH\) such that for all sequences \( \gamma_{w_n,g,h} : n \geq 1 \) we have

\[
\liminf_{n \to \infty} P_{\gamma_{wn,g,h}}(T_{\gamma_{wn,g,h}}(\theta_{w_n,g,h}, 1 - \alpha)) \geq J_h(c_g(1 - \alpha)). \tag{9.4}
\]

The improved lower bound on AsyCS for subsampling CSs is

\[
\text{Min}^-_{CS,Sub}(\alpha) = \min \left\{ \inf_{(g,h) \in GH} J_h(c_g(1 - \alpha)) - \right\}
\]

(9.5)

(where infimum over a null set is defined to be \( \infty \)). Clearly, \( \text{Min}^-_{CS,Sub}(\alpha) \geq \inf_{(g,h) \in GH} J_h(c_g(1 - \alpha)) \), where the latter is the lower bound on AsyCS without the improvement. Define \( \text{Max}^-_{CS,Sub}(\alpha) \) analogously to \( \text{Min}^-_{CS,Sub}(\alpha) \) with min and inf replaced by max and sup.

Sufficient conditions for \((g, h)\) to be in \( GH^* \) are that for all sequences \( \gamma_{w_n,g,h} : n \geq 1 \), (a) there exists a finite nonstochastic lower bound \( LB_h \) such that the subsample statistics are \( \geq LB_h \) almost surely (a.s.) under \( \gamma_{w_n,g,h} : n \geq 1 \), (b) \( J_h(LB_h) \geq J_h(c_g(1 - \alpha)) \), and (c) \( \liminf_{n \to \infty} P_{\gamma_{wn,g,h}}(T_{\gamma_{wn,g,h}}(\theta_{w_n,g,h}) \leq LB_h) \geq J_h(LB_h) \). (Conditions (a)–(c) imply (9.4) because \( \liminf_{n \to \infty} P_{\gamma_{wn,g,h}}(T_{\gamma_{wn,g,h}}(\theta_{w_n,g,h}) \leq LB_h) \geq J_h(c_g(1 - \alpha)) \).

The main results of this section are as follows.

**THEOREM 3.**

(i) Suppose Assumptions A0 and B0 hold. Then, an FCV CS satisfies

\[
\text{AsyCS} \in \left[ \inf_{h \in H} J_h(c_{F_{Fix}}(1 - \alpha)), \inf_{h \in H} J_h(c_{F_{Fix}}(1 - \alpha)) \right]
\]

and

\[
\text{AsyMaxCP} \in \left[ \sup_{h \in H} J_h(c_{F_{Fix}}(1 - \alpha)), \sup_{h \in H} J_h(c_{F_{Fix}}(1 - \alpha)) \right].
\]

(ii) Suppose Assumptions A0, B0, C, D, E0, F, and G0 hold. Then, a subsampling CS satisfies

\[
\text{AsyCS} \in \left[ \text{Min}^-_{CS,Sub}(\alpha), \inf_{(g,h) \in GH} J_h(c_g(1 - \alpha)) \right]
\]

and

\[
\text{AsyMaxCP} \in \left[ \text{Max}^-_{CS,Sub}(\alpha), \sup_{(g,h) \in GH} J_h(c_g(1 - \alpha)) \right].
\]
Remarks.

1. Theorem 3(ii) is used in the Appendix to prove Theorem 1.
2. When the parameter space $\Gamma$ takes on a partial product-space form as in Assumption A of Andrews and Guggenberger (2010), then the forms of the localization parameter spaces $H$ and $GH$ can be made more explicit, and the results of Theorem 3 hold under a somewhat simpler assumption than Assumption B0. See the Appendix for details.

NOTES

1. The lack of a uniformly most powerful test even in a Gaussian location testing problem with a multivariate one-sided null hypothesis (which is a special case of the nonlinear moment inequality model considered here) indicates that it is not possible to unambiguously rank the different statistics that one might use. However, some choices have better all around properties than others.
2. The uniformity issue arises whether one is interested in a confidence set for the true parameter or for the identified set. The issue is uniformity over the true distribution generating the data, not uniformity of coverage of all the points in the identified set. See Romano and Shaikh (2005) for some results concerning the uniformity of subsampling confidence sets for the identified set.
We note that correct asymptotic size depends on the specification of the parameter space. A procedure that does not have correct asymptotic size for one parameter space may have correct asymptotic size for a subset of this parameter space. Hence, it is important to exclude from the parameter space points that are not empirically relevant, especially if these are points that cause problems of uniformity.
3. Such critical values can be calculated by computing the appropriate bound from a weighted chi-square distribution or by simulating from the least favorable asymptotic distribution given the estimated variance matrix.
4. The $m$ out of $n$ bootstrap uses a bootstrap sample of size $m$ when the full-sample size is $n$, where $m \to \infty$ and $m/n \to 0$ as $n \to \infty$.
5. In an i.i.d. scenario, the distribution of a subsample of size $b$ is the same as the conditional distribution of a nonparametric bootstrap sample of size $b$ conditional on there being no duplicates of observations in the bootstrap sample. If $b^2/n \to 0$, then the probability of no duplicates goes to one as $n \to \infty$; see Politis et al. (1999, p. 48). In consequence, $b$ out of $n$ bootstrap tests and subsampling tests have the same first-order asymptotic properties.
6. Moment equalities cannot be handled in the setup of Romano and Shaikh (2008) by writing a moment equality as two moment inequalities because this approach yields a singular variance matrix for the resulting moment inequalities and the latter is not covered by their Lemma 3.1, which is used in their proof of uniformity.
7. These results and the ones given elsewhere in the paper are based on simulation of the formula for asymptotic size and hence are accurate up to simulation error.
8. The “adjusted” asymptotic size is defined in Andrews and Guggenberger (2009a). It is based on a formula for the asymptotic size that is adjusted to take into account the ratio of the subsample size, $b$, to the full-sample size, $n$, that is actually used in a given problem. In many cases, the adjusted asymptotic size is found to be more accurate than the usual “unadjusted” asymptotic size.
9. Also note that the probability of obtaining a singleton set does not have to be large to have adverse effects on some procedures because errors in tests or confidence intervals with probability .05 are what is typically relevant.
10. The moment condition (vii) could be relaxed slightly to uniform integrability of second moments.
11. With dependent observations, $\Psi$ is the parameter space for the limiting correlation matrix, $\lim_{n \to \infty} \text{Corr}_F(n^{1/2} \pi_n(\theta))$. 

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12. If $\Sigma$ is singular, we define $S_2$ using the Moore–Penrose inverse $\Sigma^+$ in place of $\Sigma^{-1}$ in (3.12). With some work, it may be possible to extend the results given in the text for the function $\delta = S_2$ to the case where $\Psi = \Psi_1$.

The definition of $S_2(m, \Sigma)$ takes the infimum over $t_1 \in R^p_{+\infty}$, rather than over $t_1 \in R^p_{+\infty}$. For calculation of the test statistic based on $S_2$, using the latter gives an equivalent value. To obtain the correct asymptotic distribution, however, the former definition is required because it leads to continuity at infinity of $S_2$ when some elements of $m$ may be infinity. For example, suppose $k = p = 1$. In this case, when $m \in R_+$, $\inf_{t_1 \in R_{+\infty}} (m-t_1)^2 = \inf_{t_1 \in R_+} (m-t_1)^2 = 0$. However, when $m = \infty$, $\inf_{t_1 \in R_{+\infty}} (m-t_1)^2 = 0$, but $\inf_{t_1 \in R_+} (m-t_1)^2 = \infty$.

13. Note that the tests based on $S_1$ and $S_2$ depend on $\Omega$ through the subsampling critical value. However, the subsampling critical value converges in probability to a constant (as is shown explicitly in Theorem 3). In consequence, the form of the test (by which we mean the shape of the rejection region in the $k$-dimensional space of outcomes of the sample moment conditions) does not depend on $\Omega$ for large $n$, and the test fails to exploit some information provided by $\Omega$.

14. For temporally dependent observations, $\Psi$ is as in condition (v) of (A.2) in the Appendix.

15. In Assumptions 1(d) and 4(b), $S(m, \Omega)$ and $c(\Omega, 1-\alpha)$ are viewed as functions defined on the space of all correlation matrices $\Psi_1$. By definition, $c(\Omega, 1-\alpha)$ is continuous in $\Omega$ uniformly for $\Omega \in \Psi$ if for all $\eta > 0$ there exists $\delta > 0$ such that whenever $||\Omega^* - \Omega|| < \delta$ for $\Omega^* \in \Psi_1$ and $\Omega \in \Psi$ we have $|c(\Omega^* (1-\alpha) - c(\Omega (1-\alpha))| < \eta$.

16. This result holds under Assumption GEL (stated in the Appendix) under sequences of parameters $\{\gamma_{n,h} : n \geq 1\}$ as in Assumption B0 in Section 9.4.

REFERENCES


References


APPENDIX

In the Appendix, we first show how the moment inequality model fits into the general framework for CSs introduced in Section 9. Next, we prove the main results stated in the paper for the moment inequality model, and, in particular, we use the general result Theorem 3 to prove Theorem 1. Then, we provide results for GEL test statistics. Finally, we prove Theorem 3.

A1. Moment Inequality Model.

A1.1. Specification of Parameters. In this section we specify a one-to-one mapping between the parameters \((\theta, F)\) in the moment inequality model and the parameter \(\gamma = (\gamma_1, \gamma_2, \gamma_3)\) that appears in the general results of Section 9. We define \(\gamma_1 = (\gamma_{1,1}, \ldots, \gamma_{1,p})^\prime \in R_p^p\) by writing the moment inequalities in (3.1) as moment equalities:

\[
\sigma_{F,j}^{-1}(\theta)E m_j(W_i, \theta) = \gamma_{1,j} = 0 \quad \text{and} \quad \gamma_{1,j} \geq 0 \quad \text{for } j = 1, \ldots, p, \tag{A.1}
\]

where \(\sigma_{F,j}(\theta) = \text{AsyVar}_F(n^{1/2} \bar{m}_{n,j}(\theta))\) denotes the variance of the asymptotic distribution of \(n^{1/2} \bar{m}_{n,j}(\theta)\) when the true parameter is \(\theta\) and the true distribution of the data is \(F\). Let \(\Omega = \Omega(\theta, F) = \text{AsyCorr}_F(n^{1/2} \bar{m}_n(\theta))\), where \(\text{AsyCorr}_F(n^{1/2} \bar{m}_n(\theta))\) denotes the correlation matrix of the asymptotic distribution of \(n^{1/2} \bar{m}_n(\theta)\) when the true parameter is \(\theta\) and the true distribution of the data is \(F\). (We only consider \((\theta, F)\) for which these asymptotic variances and correlation matrix exist; see conditions (iv) and (v) of (A.2) later in this section.) When no preliminary estimator \(\hat{\gamma}(\theta)\) appears, \(\sigma_{F,j}(\theta) = \lim_{n \to \infty} \text{Var}_F(n^{1/2} \bar{m}_{n,j}(\theta))\) and \(\Omega(\theta, F) = \lim_{n \to \infty} \text{Corr}_F(n^{1/2} \bar{m}_n(\theta))\), where \(\text{Var}_F(\cdot)\) and \(\text{Corr}_F(\cdot)\) denote finite-sample variance and correlation under \((\theta, F)\), respectively. Let \(\gamma_2 = (\gamma_{2,1}, \gamma_{2,2}) = (\theta, \text{vech}_d(\Omega(\theta, F))) \in R^{d}, \) where \(\text{vech}_d(\Omega)\) denotes the vector of elements of \(\Omega\) that lie below the main diagonal, \(q = d + k(k - 1)/2\), and \(\gamma_3 = F\).

For the case described in Section 8.2 (where the sample moment functions depend on a preliminary estimator \(\hat{\gamma}(\theta)\) of an identified parameter vector \(\gamma_0\)), we define \(m_j(W_i, \theta) = m_j(W_i, \theta, \tau_0)\), \(m(W_i, \theta) = (m_1(W_i, \theta, \tau_0), \ldots, m_k(W_i, \theta, \tau_0))^\prime\), \(\bar{m}_{n,j}(\theta) = n^{-1} \sum_{i=1}^n m_j(W_i, \theta, \tau_0)\), and \(\bar{m}_n(\theta) = (\bar{m}_{n,1}(\theta), \ldots, \bar{m}_{n,k}(\theta))^\prime\). (Hence, in this case, \(\bar{m}_n(\theta) \neq n^{-1} \sum_{i=1}^n m(W_i, \theta)\).)

For i.i.d. observations (and no preliminary estimator \(\hat{\gamma}(\theta)\)), the parameter space \(\Gamma\) for \(\gamma\) in the moment inequality example is defined by \(\Gamma = \{\gamma = (\gamma_1, \gamma_2, \gamma_3) : \text{for some } (\theta, F) \in F\}\), where \(F\) is defined in (3.3), \(\gamma_1\) satisfies (A.1), \(\gamma_2 = (\theta, \text{vech}_d(\Omega(\theta, F)))\), and \(\gamma_3 = F\).
For dependent observations and for sample moment functions that depend on preliminary estimators of identified parameters, we specify the parameter space $\Gamma$ for the moment inequality model using a set of high-level conditions. To verify the high-level conditions using primitive conditions one has to specify an estimator $\hat{\Sigma}_n(\theta)$ of the asymptotic variance matrix $\Sigma(\theta)$ of $n^{1/2}\bar{m}_n(\theta)$. For brevity, we do not do so here. Because there is a one-to-one mapping from $\gamma$ to $(\theta, F)$, $\Gamma$ also defines the parameter space $F$ of $(\theta, F)$. Let $\Psi$ be a specifically set of $k \times k$ correlation matrices. Let $\{\pi(j) : j \geq 1\}$ be a sequence of nonnegative numbers that satisfies $\pi(j) \to 0$ as $j \to \infty$. The parameter space $\Gamma$ is defined to include parameters $\gamma = (\gamma_1, \gamma_2, \gamma_3) = (\gamma_1, (\theta, \gamma_{2,2}), F)$ that satisfy

(i) $\theta \in \Theta$,
(ii) $\sigma_{F,j}^2(\theta)E_F m_j(W_i, \theta) - \gamma_{1,j} = 0$ and $\gamma_{1,j} \geq 0$ for $j = 1, \ldots, p$,
(iii) $E_F m_j(W_i, \theta) = 0$ for $j = p + 1, \ldots, k$,
(iv) $\sigma_{F,j}^2(\theta) = \text{AsyVar}_F\left(n^{1/2}\bar{m}_{n,j}(\theta)\right)$ exists and lies in $(0, \infty)$ for $j = 1, \ldots, k$,
(v) $\text{AsyCorr}_F\left(n^{1/2}\bar{m}_n(\theta)\right)$ exists and equals $\Omega_{\gamma_{2,2}} \in \Psi$,
(vi) $\{W_i : i \geq 1\}$ are stationary and strong mixing under $F$

with strong mixing numbers $\alpha_F(j) \leq \pi(j)$ for all $j \geq 1$, (A.2)

where $\gamma_1 = (\gamma_1, 1, \ldots, \gamma_1, \mu)$ and $\Omega_{\gamma_{2,2}}$ is the $k \times k$ correlation matrix determined by $\gamma_{2,2}$. Furthermore, $\Gamma$ must be restricted by enough additional conditions such that under any sequence $\{\gamma_{n,h} = (\gamma_{n,h}, 0_{n,h}, \text{vec} \hat{\gamma}_{n,h}, (\Omega_{n,h}), F_{n,h}) : n \geq 1\}$ of parameters in $\Gamma$ that satisfies $n^{1/2}\gamma_{n,h,1} \to h_1$ and $(\theta_{n,h}, \text{vec} \hat{\gamma}_{n,h}(\Omega_{n,h})) \to h_2 = (h_{2,1}, h_{2,2})$ for some $h = (h_1, h_2) \in R_p^\infty \times R_k^\infty$, we have

(vii) $A_n = (A_{n,1}, \ldots, A_{n,k}) \to_d Z_{h_{2,2}} \sim N(0_k, \Omega_{h_{2,2}})$ as $n \to \infty$, where

\[ A_{n,j} = n^{1/2}\left(\bar{m}_{n,j}(\theta_{n,h}) - n^{-1}\sum_{i=1}^n E_{n,h} m_j(W_i, \theta_{n,h})\right)/\sigma_{F,n,j}(\theta_{n,h}), \]

(viii) $\bar{\sigma}_{n,j}(\theta_{n,h})/\sigma_{F,n,j}(\theta_{n,h}) \to p_1$ as $n \to \infty$ for $j = 1, \ldots, k$,
(ix) $\hat{\sigma}_{n,j}(\theta_{n,h}) \bar{\bar{\sigma}}_n(\theta_{n,h}) \hat{\sigma}_{n,j}(\theta_{n,h}) \to p \Omega_{h_{2,2}}$ as $n \to \infty$, (A.3)

(x) conditions (vii)–(ix) hold for all subsequences $\{w_n\}$ in place of $\{n\}$,

where $\Omega_{h_{2,2}}$ is the $k \times k$ correlation matrix for which $\text{vec} \hat{\gamma}_{n,h}(\Omega_{h_{2,2}}) = h_{2,2}$, $\bar{\sigma}_{n,j}(\theta) = \{\Sigma_n(\theta)\}_{ij}$ for $j = 1, \ldots, k$, and $\bar{\sigma}_n(\theta) = \text{Diag}(\hat{\sigma}_{n,1}^2(\theta), \ldots, \hat{\sigma}_{n,k}^2(\theta)) = \text{Diag}(\Sigma_n(\theta))$.

(When a preliminary estimator $\hat{\gamma}_n(\theta)$ appears, $A_{n,j}$ can be written equivalently as $n^{1/2}(n^{-1}\sum_{i=1}^n m_j(W_i, \theta_{n,h}, \hat{\gamma}_n(\theta_{n,h}))) - n^{-1}\sum_{i=1}^n E_{n,h} m_j(W_i, \theta_{n,h}, \tau_0))/\sigma_{F,n,j}(\theta_{n,h})$, which typically is asymptotically normal with an asymptotic variance matrix $\Omega_{h_{2,2}}$ that reflects the fact that $\tau_0$ has been estimated. When a preliminary estimator $\hat{\gamma}_n(\theta)$ appears, $\Sigma_n(\theta)$ needs to be defined to take account of the fact that $\tau_0$ has been estimated. When no preliminary estimator $\hat{\gamma}_n(\theta)$ appears, $A_{n,j}$ can be written equivalently as $n^{1/2}(\bar{m}_{n,j}(\theta_{n,h}) - E_{n,h} \bar{m}_{n,j}(\theta_{n,h}))/\sigma_{F,n,j}(\theta_{n,h})$. Condition (x) of (A.3) requires that conditions (vii)–(ix) must hold under any sequence of parameters $\{\gamma_{w_n,h} : n \geq 1\}$ that satisfies the conditions preceding (A.3) with $n$ replaced by $w_n$.)
For example, for i.i.d. observations, conditions (i)–(vi) of (3.3) imply conditions (i)–(vi) of (A.2). Furthermore, conditions (i)–(vi) of (3.3) plus the definition of \( \hat{\Sigma}_n(\theta) \) in (3.8) and the additional condition (vii) of (3.3) imply conditions (vii)–(x) of (A.3).

**LEMMA 2.** The parameter space \( \Gamma \) for i.i.d. observations (that is, for \( \Gamma \) defined by the restrictions summarized in (3.3)) is such that conditions (i)–(x) of (A.2)–(A.3) hold when \( \hat{\Sigma}_n(\theta) \) is defined by (3.8).

For dependent observations, one needs to specify a particular variance estimator \( \hat{\Sigma}_n(\theta) \) before one can specify primitive “additional conditions” beyond conditions (i)–(vi) in (A.2) that ensure that \( \Gamma \) is such that any sequence \( \{\gamma_{n,h} : n \geq 1\} \) in \( \Gamma \) satisfies (A.3). For brevity, we do not do so here. Note that the strong mixing assumption in condition (vi) of (A.2) is used to verify Assumption E0.

**A1.2. Proofs for the Moment Inequality Model.**

**Proof of Theorem 1.** We prove Theorem 1 for the moment inequality/equality model by showing (a) Assumptions A0, B0, C, D, E0, F, and G0 hold and hence Theorem 3 applies, (b) \( \min_{CS,Sub}(\alpha) = \min_{(g,h) \in \mathcal{GH}} \min_{j} J_{h,c_{g}(1-\alpha)}(A_{n,j} + n^{1/2} \gamma_{n,h,1,j}) \) (using the fact that \( \gamma_{n,h} \in (\gamma_{n,h,1,1}, \ldots, \gamma_{n,h,1,p})^{\top} \) and by definition \( \gamma_{n,h,1,j} = 0 \) for \( j = p+1, \ldots, k \). Condition (vii) of (3.3) and the definition of \( \{\gamma_{n,h} : n \geq 1\} \) imply that if \( h_{1,j} = \infty \) and \( j \leq p \), then \( A_{n,j} + n^{1/2} \gamma_{n,h,1,j} \to \infty \) under \( \{\gamma_{n,h} : n \geq 1\} \). In consequence, if any element of \( h_{1} \) equals \( \infty \), \( \hat{\Sigma}_n(\theta) \) does not converge in distribution (to a proper finite random vector), and the continuous mapping theorem cannot be applied to obtain the asymptotic distribution of the right-hand side of (A.4).

To circumvent these problems, we consider a \( k \)-vector-valued function of \( \hat{\Sigma}_n(\theta) \) that converges in distribution whether or not some elements of \( h_{1} \) equal \( \infty \). Then, we write the right-hand side of (A.4) as a continuous function of this \( k \)-vector and apply the continuous mapping theorem. Let \( G(\cdot) \) be a strictly increasing continuous df on \( R \), such as the standard normal df. For \( j \leq k \), we have

\[
G_{n,j} = G\left(\hat{\sigma}_{n,j}^{-1}(\theta_{n,h})n^{1/2}m_{n,j}(\theta_{n,h})\right) = G\left(\hat{\sigma}_{n,j}^{-1}(\theta_{n,h})\sigma_{F_{n,j}}(\theta_{n,h})\left[A_{n,j} + n^{1/2} \gamma_{n,h,1,j}\right]\right).
\]  

(A.5)

Let \( Z_{h_{2,j}} = (Z_{h_{2,1}}, \ldots, Z_{h_{2,k}})^{\top} \sim N(0_{k}, \Omega_{h_{2}}) \). Define \( h_{1,j} = 0 \) for \( j = p+1, \ldots, k \). If \( j \leq p \) and \( h_{1,j} < \infty \) or if \( j = p+1, \ldots, k \), then

\[
G_{n,j} \to d G\left(Z_{h_{2,j}} + h_{1,j}\right)
\]  

(A.6)
by (A.5), conditions (vii) and (viii) of (A.3), and the continuous mapping theorem. If \( j \leq p \) and \( h_{1,j} = \infty \), then

\[
G_{n,j} = G \left( \frac{1}{\sqrt{n}} \left( \bar{\theta}_{n,h} \right) n^{1/2} \tilde{m}_{n,j}(\theta_{n,h}) \right) \rightarrow p 1 \tag{A.7}
\]

by (A.5), \( A_{n,j} = O_p(1) \), and \( G(x) \rightarrow 1 \) as \( x \rightarrow \infty \). The results in (A.6)–(A.7) hold jointly and combine to give

\[
G_n = (G_{n,1}, \ldots, G_{n,h})' \rightarrow_d (G(Z_{h2,2}, h_{1,1}), \ldots, G(Z_{h2,2,k} + h_{1,k}))' = G_\infty, \tag{A.8}
\]

where \( G(Z_{h2,2,j} + h_{1,j}) \) denotes \( G(\infty) = 1 \) when \( h_{1,j} = \infty \).

Let \( G^{-1} \) denote the inverse of \( G \). For \( x = (x_1, \ldots, x_k)' \in \mathbb{R}^p_{[+\infty]} \times \mathbb{R}^p \), let \( G(k)(x) = (G(x_1), \ldots, G(x_k))' \in (0,1]^p \times (0,1)^p \). For \( y = (y_1, \ldots, y_k)' \in (0,1]^p \times (0,1)^p \), let \( G^{-1}(y) = (G^{-1}(y_1), \ldots, G^{-1}(y_k))' \in \mathbb{R}^p_{[+\infty]} \times \mathbb{R}^p \). Define \( S^* \) as

\[
S^*(y, \Omega) = S(G^{-1}_k(y), \Omega) \tag{A.9}
\]

for \( y \in (0,1]^p \times (0,1)^p \) and \( \Omega \in \Psi \). Assumption 1(d) implies that \( S^*(y, \Omega) \) is continuous at all \((y, \Omega)\) for \( y \in (0,1]^p \times (0,1)^p \) and \( \Omega \in \Psi \). We now have

\[
T_n(\theta_{n,h}) = S \left( G^{-1}(G_n), \tilde{D}_n^{-1/2}(\theta_{n,h}) \tilde{\Sigma}(\theta_{n,h}) \tilde{D}_n^{-1/2}(\theta_{n,h}) \right)
= S^* \left( G_n, \tilde{D}_n^{-1/2}(\theta_{n,h}) \tilde{\Sigma}(\theta_{n,h}) \tilde{D}_n^{-1/2}(\theta_{n,h}) \right)
\rightarrow_d S^*(G_\infty, \Omega_{h2,2})
= S(G^{-1}_k(G_\infty), \Omega_{h2,2})
= S(Z_{h2,2} + (h_{1}, 0_0), \Omega_{h2,2})
\sim J_h, \tag{A.10}
\]

where the first equality holds by (A.4) and the definition of \( G^{-1}_k(G_n) \), the second and third equalities hold by the definition of \( S^* \), the convergence holds by (A.8), condition (ix) of (A.3), and the continuous mapping theorem, the last equality holds by the definitions of \( G^{-1}_k \) and \( G_\infty \) and the definition that if \( h_{1,j} = \infty \), then the corresponding element of \( Z_{h2,2} + (h_{1}, 0_0) \) equals \( \infty \), and the last line gives the definition of \( J_h \) (where \( h = (h_1, h_2) \), \( h_2 = (h_{2,1}, h_{2,2}) \), \( h_{2,1} \in \mathbb{R}^d \) is arbitrary because it does not appear in \( S(Z_{h2,2} + (h_1, 0_0), \Omega_{h2,2}) \), and \( h_{2,2} = \text{vech}_4(\Omega_{h2,2}) \)). By the same argument but using condition (x) of (A.3) in place of conditions (vii)–(ix), the result of (A.10) holds with \( \{w_n\} \) in place of \( \{n\} \) for any subsequence \( \{w_n\} \). Hence, Assumption B0 holds with \( J_h \) defined as in (A.10).

Assumption C is assumed in Section 5. Assumption D holds by stationarity and the standard definition of subsample statistics in the i.i.d. and dependent cases. Assumption E0 holds for i.i.d. and stationary strong mixing observations by the remarks at the end of Section 9.5 using condition (vi) of (A.2).

Next, we verify Assumption F. When \( \nu = 0 \) and \( h_1 = \infty, \) the limit random variable in (A.10) is \( S(Z_{h2,2} + \infty^p, \Omega_{h2,2}) = S(\infty^p, \Omega_{h2,2}) = 0 \) using Assumption 3. In consequence,
Next, we show that $v \geq \{\text{ch} that AsyCS given in (A.10) and $J_h$ where the first and third equalities hold by (A.4) and Assumption 3, the second equality and Assumption F holds, (b) if $J_h$ holds, i.e., $J_h(x) = 1 - \alpha$ for all $x > 0$ and Assumption F holds (otherwise, $J_h(x) = 1 - \alpha$ for some $x > 0$ and $J_h(x/2) = 1 - \alpha$ because $J_h$ is nondecreasing, which contradicts the fact that $J_h(x)$ is strictly increasing for $x > 0$). Hence, Assumption F holds. Assumption G0 holds automatically because the subsampling procedure satisfies Assumption Sub2.

Given that Assumptions A0, B0, C, D, E0, F, and G0 hold, the result of Theorem 3(ii) holds, i.e., $\text{AsyCS} \in [\text{Min}_{\text{CS,Sub}}(\alpha), \inf_{(g,h) \in GH} J_h(c_g(1-\alpha))]$.

We now prove Theorem 1(i) by showing that $\text{Min}_{\text{CS,Sub}}(\alpha) \geq 1 - \alpha$. First, by Assumption 1(a), for $0 \leq g_1 \leq h_1 \in R^p_+ \cap \infty$ and $(g, h) \in GH$, we have

$$S(Z_{h,2,2} + (g_1, 0_h), \Omega_{h,2,2}) \geq S(Z_{h,2,2} + (h_1, 0_h), \Omega_{h,2,2}),$$

$$c_g(1 - \alpha) = c_h(1 - \alpha), \quad \text{and}$$

$$J_h(c_g(1 - \alpha)) = J_h(c_h(1 - \alpha)).$$

(A.11)

Next, we show that $GH^* = GH$. Given $(g, h) \in GH$, suppose $c_g(1 - \alpha) > 0$. Then, $J_h(c_g(1 - \alpha)) = J_h(c_h(1 - \alpha))$ because $J_h(x)$ is continuous for all $x > 0$ by Assumption 2(a). This and Lemma 6(vi) of Andrews and Guggenberger (2010) (which holds under Assumptions A0, B0, C, D, E0, F, and G0 by the proof of Theorem 3 in Section A2) establish (9.4). Hence, $(g, h) \in GH^*$.

Now, suppose $c_g(1 - \alpha) = 0$. (Assumption 1(c) rules out $c_g(1 - \alpha) < 0$.) This implies that $c_h(1 - \alpha) = 0$ by (A.11). The conditions $c_h(1 - \alpha) = 0$ and $0 < \alpha < 1/2$ are consistent with Assumption 2(c) only if $v = 0$. Given $v = 0$, we show $(g, h) \in GH^*$ by verifying conditions (a)–(c) in the paragraph preceding Theorem 3. Condition (a) holds with $LB_h = 0$ by Assumption 1(c). Condition (b) holds because $LB_h = c_g(1 - \alpha) = 0$. Next we show condition (c). Under $\gamma_{n, h} : n \geq 1$ and with $v = 0$, we have

$$P_{\gamma_{n, h}}(T_n(\theta_{n, h}) \leq 0)$$

$$= P_{\gamma_{n, h}}(n^{1/2} \bar{m}_{n, j}(\theta_{n, h})/\sigma_{F_{n, h}, j}(\theta_{n, h}) \geq \zeta \quad \text{for all} \quad j = 1, \ldots, p)$$

$$= P_{\gamma_{n, h}}(A_{n, j} + n^{1/2} \gamma_{n, h, 1, j} \geq \zeta \quad \text{for all} \quad j = 1, \ldots, p)$$

$$\Rightarrow P(Z_{h,2,2, j} + h_{1, j} \geq \zeta \quad \text{for all} \quad j = 1, \ldots, p)$$

$$= P(S(Z_{h,2,2} + h_{1}, \Omega_{h,2,2}) \leq 0)$$

$$= J_h(0),$$

(A.12)

where the first and third equalities hold by (A.4) and Assumption 3, the second equality and the convergence hold by (A.3), and the last equality holds by the definition of $J_h$ given in (A.10) and $v = 0$. The same argument holds with $\gamma_{n, h} : n \geq 1$ in place of $\gamma_{n, h} : n \geq 1$. Hence, (A.12) completes the verification of condition (c) and concludes the proof that $GH^* = GH$.

For subsampling CSs, we now have

$$\text{AsyCS} = \inf_{(g, h) \in GH} J_h(c_g(1 - \alpha)) \geq \inf_{h \in H} J_h(c_h(1 - \alpha)) \geq 1 - \alpha,$$

(A.13)
where the equality holds by Theorem 3 and \( GH^* = GH \), the first inequality holds by (A.11), and the second inequality holds by the definition of \( c_h(1 - \alpha) \). This establishes Theorem 1(i). (Note that \( \text{AsyMaxCP} \) is given by the second expression in (A.13) with “sup” in place of “inf.”)

Next, let \( (\theta^*, F^*) \) be an element of \( \mathcal{F} \) for which Assumption C1 applies and let \( \gamma^* \) be the value in \( \Gamma \) that corresponds to \( (\theta^*, F^*) \in \mathcal{F} \). Define \( h^* = (h^*_1, h^*_2) \), where \( h^*_1 = h_1(\theta^*, F^*) \), \( h^*_2,1 = \theta^* \in \mathbb{R}^d \), and \( h^*_2,2 = \vech_\sigma(\Omega(\theta^*, F^*)) \). We have \( (h^*, h^*) \in GH \) because the sequence \( \{\gamma_{w_q,g,h} : n \geq 1\} \) defined by \( \gamma_{w_q,g,h} = \gamma^* \) for all \( n \geq 1 \) leads to the point \( (h^*, h^*) \in GH \) by the definition of \( GH \) given in (9.3). By Assumption C1, \( J_{h^*}(c_{h^*}(1 - \alpha)) = 1 - \alpha \). In consequence, we have

\[
\text{AsyCS} = \inf_{(g, h) \in GH} J_h(c_h(1 - \alpha)) \leq J_{h^*}(c_{h^*}(1 - \alpha)) = 1 - \alpha. \tag{A.14}
\]

Combining (A.13) and (A.14) completes the proof of Theorem 1(ii).

Now, we prove Theorem 1(iii). By assumption, \( \nu = 0 \). Assumption C2 guarantees the existence of \( (\theta^*, F^*) \in \mathcal{F} \) for which \( E_f m_j(W_i, \theta^*)/\sigma F_{-j}(\theta^*) > 0 \) for all \( j = 1, \ldots, p \). The sequence of constant true values \( \{(\theta^*, F^*) \in \mathcal{F} : n \geq 1\} \) satisfies \( n^{1/2} E_f m_j(W_i, \theta^*)/\sigma F_{-j}(\theta^*) \to \infty \) and \( b^{1/2} E_f m_j(W_i, \theta^*)/\sigma F_{-j}(\theta_h) \to \infty \) for all \( j = 1, \ldots, p \). Let \( \gamma^* = (\gamma^*_1, (\theta^*, \gamma^*_2,2), F^*) \in \Gamma \) correspond to \( (\theta^*, F^*) \in \mathcal{F} \). Define \( g^* = h^* = (\infty^p, (\theta^*, \gamma^*_2,2)) \). Then, \( (g^*, h^*) \in GH \) if and only if \( \phi_h^*(x) = 1 \) for \( x \geq 0 \) because \( S(\Omega_{h^*_2,2}^g + \infty^p, \Omega_{h^*_2,2}^g) = S(\infty^p, \Omega_{h^*_2,2}^g) = 0 \) using Assumption 3 and \( c_{g^*}(1 - \alpha) \geq 0 \) by Assumption 1(c). Hence, \( J_{h^*}(c_{g^*}(1 - \alpha)) = 1 \). Given the previous result that \( GH^* = GH \), we have \( \text{MaxCS,Sub}(\alpha) = \sup_{(g, h) \in GH} J_h(c_h(1 - \alpha)) \geq J_{h^*}(c_{g^*}(1 - \alpha)) = 1. \)

**Proof of Theorem 2.** We prove Theorem 2 using the confidence set analogue of Theorem 3 of Andrews and Guggenberger (2009a) discussed in Section 3.3 of that paper. (Note that the PA CS considered in this paper is an example of the plug-in size-corrected (PSC) FCV CS considered in Andrews and Guggenberger, 2009a.) Theorem 3 of Andrews and Guggenberger (2009a) can be extended to hold with Assumptions A0, B0, E0, and G0 in place of Assumptions A, B, E, and G (defined in Andrews and Guggenberger, 2009a, and used in their Thm. 3) using the same arguments as in the proof of Theorem 3 in Section A2. The definitions of \( H \) and \( GH \) are then given by (9.2) and (9.3) of the current paper. For the case of PSC-FCV tests, which are relevant here, only Assumptions A, B, L(i), N, and O(a), defined in Andrews and Guggenberger (2009a), are needed for Theorem 3 of Andrews and Guggenberger (2009a). Hence, we only need to verify Assumptions A0, B0, L(i), N, and O(a) here and the set \( GH \) is not relevant.

In the present case, the quantity \( cv_{h_{12}}(1 - \alpha) \) in (3.5) of Andrews and Guggenberger (2009a) satisfies

\[
sv_{h_{12}}(1 - \alpha) = \sup_{h_1 \in H_1} c(h_1, h_2)(1 - \alpha) \leq c(0_2, h_2)(1 - \alpha), \tag{A.15}
\]

where \( H_1 = \{h_1 \in R_\infty^p : h = (h_1, h_2) \in H \text{ for some } h_2 \in R_\infty^q \} \), the equality is by definition, and the inequality holds by (A.11) because (A.11) holds for all \( g = (0_2, h_2) \), \( h = (h_1, h_2) \), and \( h_1 \in R_\infty^p \) whether or not \( (g, h) \in GH \) (which is not necessarily the case here). (Note that the inequality in (A.15) is not necessarily an equality because \( 0_2 \) is not necessarily in \( H_1 \).) Hence, the critical value \( c(\Omega_h(\theta), 1 - \alpha) \) in (6.2) is greater than or equal to the
critical value \(c_{\alpha^2}(1-\alpha)\) of Andrews and Guggenberger (2009a), and Theorem 3 of the same paper yields \(\text{AsyCS} \geq 1-\alpha\). Under Assumption C3, the inequality in (A.15) holds as an equality because \((0_p, h_2) \in H\) for some \(h_2\) (by a similar argument to that preceding (A.14)) and so \(0_p \in H_1\). In this case, \(c(\Omega_0(\theta), 1-\alpha)\) equals the critical value \(c_{\alpha^2}(1-\alpha)\) of Andrews and Guggenberger (2009a), and Theorem 3 of that paper yields \(\text{AsyCS} = 1-\alpha\).

Now, it suffices to show that Assumptions L(i), N, and O(a) of Andrews and Guggenberger (2009a) hold because Assumptions A0 and B0 hold by the proof of Theorem 1. Assumption L(i) holds, i.e., \(\sup_{\Omega \in \Psi} c(\Omega, 1-\alpha) < \infty\) because \(c(\Omega, 1-\alpha)\) is a uniformly continuous function (by Assumption 4(b)) on the subset \(\Psi\) of the compact set \(\Psi_1\) of all \(k \times k\) correlation matrices and hence is bounded on \(\Psi\). For dependent observations, Assumption N holds by condition (ix) of (A.3) (which holds for i.i.d. observations by Lemma 2) and the definition of \(\{\gamma_{n,h} : n \geq 1\}\)—which implies that \(\gamma_{n,h,2} \rightarrow h_2\). Assumptions O(a)(i) and O(a)(ii) hold by Assumptions 2(b) and 2(a), respectively. Assumption O(a)(iii) holds with \(h_1^* = 0_p\).

**Proof of Lemma 1.** For \(S_1\), Assumptions 1 and 3 hold immediately with \(\xi = 0\) in Assumption 3. Assumption 2(a) holds because (a) if \(v \geq 1\), the summand \(\sum_{j=p+1}^{v} Z_j^2\) is absolutely continuous, where \(Z = (Z_1, \ldots, Z_k)^\prime\), (b) if \(v = 0\) and \(h_1 \neq \infty^p\), the summands \([Z_j + h_1, j]_{1}^2\) are absolutely continuous for \(x > 0\) for all \(j = 1, \ldots, p\) such that \(h_1,j < \infty\), and (c) if \(v = 0\) and \(h_1 = \infty^p\), \(S_1(Z + h_1, \Omega) = 0\) and the sum does also. Assumptions 2(c) holds because (a) if \(v \geq 1\), the summand \(\sum_{j=p+1}^{v} Z_j^2\) has positive density on \(R_+\), each summand \([Z_j + h_1, j]_{1}^2\) for which \(h_1,j < \infty\) (of which there may be none) has positive density on \(R_+\) and so does the sum and (b) if \(v = 0\) and \(h_1 \neq \infty^p\), each summand \([Z_j + h_1, j]_{1}^2\) for which \(h_1,j < \infty\) (of which there is at least one) has positive density on \(R_+\) and the sum does also. Assumption 2(c) holds because if \(v \geq 1\), \(P(S_1(Z + (h_1, 0_n), \Omega) \leq 0) \leq P(\sum_{j=p+1}^{\infty} Z_j^2 \leq 0) = 0\), and if \(h_1 = 0\) and \(v = 0\), \(P(S_1(Z, \Omega) \leq 0) \leq P([Z]^2 \leq 0) = P(Z \geq 0) = 1/2\) where the inequality holds for any \(j \leq p\). Assumption 4(a) holds by the same argument as for Assumption 2(a). Assumption 4(b) holds because (a) if \(v \geq 1\), \(c(\Omega, 1-\alpha)\) is continuous at each \(\Omega \in \Psi\) and \(\Psi = \Psi_1\) is compact. To see the former, let \((\Omega_N : N \geq 1)\) be a sequence of correlation matrices such that \(\Omega_N \rightarrow \Omega\) as \(N \rightarrow \infty\). We need to show that \(c(\Omega_N, 1-\alpha) \rightarrow c(\Omega, 1-\alpha)\). Denote by \(f_N\) and \(f\) the dfs of \(S_1(Z_N, \Omega_N)\) and \(S_1(Z, \Omega)\), respectively, where \(Z_N \sim N(0_n, \Omega_N)\) and \(Z \sim N(0_n, \Omega)\). By Assumption 2(b), \(f\) is increasing for \(x > 0\) (because \(h_1 = 0_p\), not \(\infty^p\), in this case). By Assumption 2(c) we have \(c(\Omega, 1-\alpha) > 0\), and it follows that \(f\) is increasing at \(c(\Omega, 1-\alpha)\). This implies that \(c(\Omega_N, 1-\alpha) \rightarrow c(\Omega, 1-\alpha)\) because \(S_1(Z, \Omega) = \sum_{j=1}^{p} [Z_j]^2 + \sum_{j=p+1}^{\infty} Z_j^2\) have \(\sup_{x \in R} [f_N(x) - f(x)] \rightarrow 0\).

For \(S_2\), Assumptions 1(b) and (c) and 3 hold immediately with \(\xi = 0\) in Assumption 3. Assumption 1(d) holds straightforwardly using the specification of \(\Psi_2\), which bounds the determinant of the correlation matrix \(\Omega\) away from zero. Assumption 1(a) holds because for \(x \in R^p\) with \(x \geq 0\), we have

\[
S_2((m_1 + x, m_2), \Sigma) = \inf_{t_1 \in R_{+}^n, m_2} \left( m_1 + x - t_1 \right)^\prime \Sigma^{-1} \left( m_1 + x - t_1 \right) \\
= \inf_{t_1 \in R_{+}^n, m_2} \left( m_1 - t_1 \right)^\prime \Sigma^{-1} \left( m_1 - t_1 \right)
\]
\[ \frac{1}{m_1 - t_1} \left( \begin{array}{c} m_1 - t_1 \\ m_2 - t_1 \end{array} \right)^{\prime} \Sigma^{-1} \left( \begin{array}{c} m_1 - t_1 \\ m_2 - t_1 \end{array} \right) \]

\[ = S_2((m_1, m_2, \Sigma)). \]  

(A.16)

To show Assumption 2(c), first suppose \( h_1 = 0 \); then

\[ S_2(Z, \Omega) = \inf_{t=(t_1, 0_0); t_1 \in R^2_{+, \infty}} (Z - t)_\prime \Omega^{-1} (Z - t) \]  

(A.17)

\[ = \inf_{t=(t_1, 0_0); t_1 \in R^2_{+, \infty}} (Z - Bt)_\prime (Z - Bt) = \inf_{t_1 \in R^2_{+, \infty}} (Z^* - Bt_1)_\prime (Z^* - Bt_1), \]

where \( Z \sim N(0_k, \Omega), B' B = \Omega^{-1}, Z^* = BZ \sim N(0_k, I_k), B = [B_1 : B_2], \) and \( B_1 \) is \( k \times p \) and full rank \( p \leq k \). The right-hand side of (A.17) is zero only if \( Z^* = Bt_1 \) for some \( t_1 \in R^2_{+, \infty} \). The latter holds with probability zero if \( k > p \) and with probability \( \leq 1/2 \) if \( k = p \), which verifies Assumption 2(c) for \( h_1 = 0 \). Next, suppose \( v \geq 1 \); without loss of generality (wlog) we can assume \( \| B_1 \| < \infty \) (because if some element of \( h_1 \) equals infinity then the infimum in \( S_2(Z, \Omega) \) is obtained by taking the corresponding element of \( t_1 \) equal to infinity). Then, using (A.17), we have

\[ S_2(Z + (h_1, 0_0), \Omega) = \inf_{t_1 \in R^2_{+, \infty}} (Z^* - Bt_1)_\prime (Z^* - Bt_1), \]  

(A.18)

where \( Z^* = B(Z + (h_1, 0_0)) \sim N(B_1 h_1, I_k) \) and \( B \) and \( B_1 \) are as before. As previously, \( S_2(Z, \Omega) = 0 \) is zero only if \( Z^* = Bt_1 \) for some \( t_1 \in R^2_{+, \infty} \). The support of \( Z^* \) is \( R^k \), whereas \( [B_1 t_1 : t_1 \geq 0] \) is a subset of a \( p \)-dimensional linear subspace of \( R^k \). Because \( v = k - p > 0 \), \( S_2(Z, \Omega) = 0 \) with probability zero.

Next, we show that Assumptions 2(a) and (b) hold for \( S_2 \). If \( v = 0 \) and \( h_1 = \infty^P \), then \( S_2(Z + h_1, \Omega) = 0 \), \( J_v(x) = 1 \) for all \( x > 0 \), Assumption 2(a) holds, and Assumption 2(b) does not impose any restriction. Otherwise, \( v \geq 1 \) or \( h_1 \neq \infty^P \). As before, wlog we can assume \( \| B_1 \| < \infty \) (because \( v \geq 1 \) or \( h_1 \neq \infty^P \)) implies that at least one element of \( Z \) remains after setting to zero all those elements \( Z_j + h_{1,j} - t_{1,j} \) for which \( h_{1,j} = \infty \). Equation (A.18) holds in the present case (whether or not \( v \geq 1 \)). The random variable \( S_2(Z + (h_1, 0_0), \Omega) \) in (A.18) has support \( R^+ \) and is absolutely continuous. Hence, Assumptions 2(a) and 2(b) hold. Assumption 4(a) holds by the same argument as for Assumption 2(a).

To show Assumption 4(b) for \( S_2 \), first we show continuity of \( c(\Omega, 1 - \alpha) \) at a fixed \( \Omega \in \Psi_2 \). Let \( \{ \Omega_N : N \geq 1 \} \) be a sequence of correlation matrices not necessarily in \( \Psi_2 \) such that \( \Omega_N \to \Omega \) as \( N \to \infty \). We need to show that \( c(\Omega_N, 1 - \alpha) \to c(\Omega, 1 - \alpha) \). Denote by \( f_N \) and \( f \) the dfs of \( S_2(Z_N, \Omega_N) \) and \( S_2(Z, \Omega) \), respectively, where \( Z_N \sim N(0_k, \Omega_N) \) and \( Z \sim N(0_k, \Omega) \). By Assumption 2(b), \( f \) is increasing for \( x > 0 \) (because \( h_1 = 0 \), not \( \infty^P \), in this case). By Assumption 2(c) we have \( c(\Omega, 1 - \alpha) > 0 \), and it follows that \( f \) is increasing at \( c(\Omega, 1 - \alpha) \). This implies that \( c(\Omega_N, 1 - \alpha) \to c(\Omega, 1 - \alpha) \) because by (A.18) with \( h_1 = 0 \), we have \( \sup_{x \in R} |f_N(x) - f(x)| \to 0 \).

Next, choose \( \delta > 0 \) small enough that for the compact set \( \Psi_2^\delta = \{ \Omega \in \Psi : \det(\Omega) \geq \delta \} \) (where \( \delta > 0 \) as in the definition of \( \Psi_2 \)); it holds that for every \( \Omega_2 \in \Psi_2 \) we have \( \Omega \in \Psi : \| \Omega_2 - \Omega \| < \delta \subset \Psi_2^\delta \). By the argument in the preceding paragraph, \( c(\Omega, 1 - \alpha) \) is continuous on \( \Psi_2^\delta \) as a function on \( \Psi_2^\delta \) and thus is uniformly continuous on \( \Psi_2^\delta \). This implies that \( c(\Omega, 1 - \alpha) \) is uniformly continuous on \( \Psi_2 \) as a function on \( c(\Omega, 1 - \alpha) \).

The proof for \( S_3 \) is essentially the same as that for \( S_1 \).
Proof of Lemma 2. Condition (ii) of (A.2) holds by the definition of $\gamma_{1,j}$ in (A.1) for $j = 1, \ldots, p$ using condition (ii) of (3.3). Conditions (i), (iii), (iv), (v), and (vi) of (A.2) hold by conditions (i), (iii), (iv) and (v), (iv) and (vi), and (iv) of (3.3), respectively.

Condition (vii) of (A.3) holds by the combination of the Cramér–Wold device and the Liapounov triangular array central limit theorem (CLT) for row-wise i.i.d. random variables with mean zero and variance one by conditions (iv) and (vii) of (3.3). Conditions (viii) and (ix) of (A.3) hold by standard arguments using a weak law of large numbers (LLN) for row-wise i.i.d. random variables with variance one by conditions (iv) and (vii) of (3.3). Condition (x) of (A.3) holds by the same argument as for conditions (vii)–(ix) of (A.3).

A1.3. Results for GEL Statistics. Here we prove that Theorems 1 and 2 hold for CSs based on the GEL statistic $T^\text{GEL}_n(\theta)$ rather than $T_n(\theta)$, provided Assumption GEL, which follows, holds, which requires that the observations are i.i.d. for each fixed $\theta, F \in \mathcal{F}$. It suffices to show that for any sequence $\{\gamma_{w_h,n} : n \geq 1\}$ for $h = (h_1, h_2) \in \mathbb{R}^p_+ \times \mathbb{R}_+^{d_\mathcal{A}}$ and corresponding $\{\theta_{w_h,n}, F_{w_h,n}\} \in \mathcal{F}$ for $n \geq 1$, we have

$$T^\text{GEL}_{w_h}(\theta_{w_h,n}, h) - T_{w_h}(\theta_{w_h,n}) = o_p(1). \quad (A.19)$$

The result in (A.19) implies that $T^\text{GEL}_{w_h}(\theta_{w_h,n}, h)$ satisfies Assumption B0. The remainder of the proofs are the same as the proofs of Theorems 1 and 2.

We use the following notation. Let

$$t_{n,h} = \mathbb{E}_{F_{w_h,n}} m(W_i, \theta_{n,h}),$$

$$\hat{t}_n = \arg\min_{t \in \mathbb{R}_+^p} \sup_{\lambda \in \Lambda_\mathcal{A}(t, \theta_{n,h})} n\hat{F}_p(t, \theta_{n,h}, \lambda), \quad \text{if it exists,}$$

$$\hat{m}_n(t) = n^{-1} \sum_{i=1}^n m_i(t, \theta_{n,h}),$$

$$\hat{\Lambda}(t) = n^{-1} \sum_{i=1}^n m_i(t, \theta_{n,h}) m_i(t, \theta_{n,h})' \quad \text{and}$$

$$\Lambda_n = \{ \lambda \in \mathbb{R}^k : ||\lambda|| \leq n^{-1/2+\delta/2} \} \quad (A.20)$$

for $\delta > 0$ as in condition (vii) of (3.3). Let $p_j(x) = (\partial^j / \partial x^j)(x)$ for $j = 1, 2$. Let w.p.a.1 denote “with probability that approaches one as $n \to \infty$.” We make the following assumption.

Assumption GEL.

(a) The observations are i.i.d. for each fixed $(\theta, F) \in \mathcal{F}$.

(b) Condition (v) of (3.3) is strengthened to $\text{Var}_F(m_j(W_i, \theta)) \in [\varepsilon_*, M_*]$ for some $\varepsilon_* > 0$ and $M_* < \infty$ for $j = 1, \ldots, k$.

(c) The parameter space $\Psi$ in (3.3) equals $\Psi_2$.

(d) For any sequence $\{\gamma_{w_h,n} : n \geq 1\}$ and corresponding $\{\theta_{w_h,n}, F_{w_h,n}\} \in \mathcal{F}$ for $n \geq 1$, $\tilde{t}_{w_h,n} = \arg\min_{t \in \mathbb{R}_+^p} \sup_{\lambda \in \Lambda_{w_h,n}(t, \theta_{w_h,n})} w_n \hat{F}_p(t, \theta_{w_h,n}, \lambda)$ and $\tilde{t}^*_{w_h,n} = \arg\min_{t \in \mathbb{R}_+^p} w_n (\Omega_{w_h,n}(\theta_{w_h,n}) - (t, 0_n))' \tilde{\Sigma}_{w_h,n}(\theta_{w_h,n} - (t, 0_n))$ exist and satisfy $\sup_{n \geq 1} ||\tilde{t}_{w_h,n}|| \leq K$ and $\sup_{n \geq 1} ||\tilde{t}^*_{w_h,n}|| \leq K$ w.p.a.1 for some constant $K < \infty$, where $\tilde{\Sigma}_{w_h,n}(\theta)$ is defined in (3.8).
The proof of (A.19) uses a similar approach to that in Newey and Smith (2004). The proof uses the following four lemmas.

**LEMMA 3.** Suppose Assumption GEL holds. For any sequence \(\{\gamma_{n,h} : n \geq 1\} \) and corresponding \(\{\theta_{w_n,h}, F_{w_n,h}\} \in \mathcal{F} : n \geq 1\), there exist constants \(K < \infty\) and \(\varepsilon > 0\) such that

(i) \(\lambda_{\min}(\hat{\Delta}(t_{w_n,h})) \geq \varepsilon\) w.p.a.1 and \(w_n^{-1} \sum_{i=1}^{w_n} ||m_i(t_{w_n,h}, \theta_{w_n,h})m_i(t_{w_n,h}, \theta_{w_n,h})'|| = O_p(1)\),

(ii) for every random sequence \(\{\tilde{t}_{w_n} \in R^p_n : n \geq 1\}\) with \(\sup_{n \geq 1} ||\tilde{t}_{w_n}|| \leq K\) w.p.a.1,

\[
\max_{1 \leq i \leq w_n} ||m_i(\tilde{t}_{w_n}, \theta_{w_n,h})|| = O_p(w_n^{1/(2+\delta)}) \quad \text{and} \quad \lambda_{\max}(\hat{\Delta}(\tilde{t}_{w_n})) \leq K \text{ w.p.a.1},
\]

(iii) \(\tilde{m}_{w_n}(t_{w_n,h}) = O_p(n^{-1/2})\), and

(iv) \(w_n^{-1} \sum_{i=1}^{w_n} ||m(W_i, \theta_{w_n,h})|| = O_p(1)\) and \(||t_{w_n,h}|| \leq K\).

The next three lemmas are analogues of Lemmas 7–9 of Guggenberger and Smith (2005). They all hold under a given sequence \(\{\gamma_{n,h} : n \geq 1\}\) and corresponding \(\{\theta_{n,h}, F_{n,h}\} \in \mathcal{F} : n \geq 1\) and Assumption GEL.

**LEMMA 4.** Assume that for a (possibly random) sequence \(\{t_n \in R^p_n : n \geq 1\}\) we have \(\max_{1 \leq i \leq n} ||m_i(t_n, \theta_{n,h})|| = O_p(n^{1/(2+\delta)})\). Then, \(\sup_{\lambda \in \Lambda_n} |\hat{\lambda}m_i(t_n, \theta_{n,h})| \rightarrow p 0\) and \(\Lambda_n \subset \Lambda \left(\hat{t}_n, \theta_{n,h}\right)\) w.p.a.1.

**LEMMA 5.** Assume that for a (possibly random) sequence \(\{t_n \in R^p_n : n \geq 1\}\) we have \(\max_{1 \leq i \leq n} ||m_i(t_n, \theta_{n,h})|| = O_p(n^{1/(2+\delta)})\), \(\hat{\lambda}_{\min}(\hat{\Delta}(t_n)) \geq \varepsilon\) w.p.a.1 for an \(\varepsilon > 0\), and \(\hat{m}(t_n) = O_p(n^{-1/2})\). Then, \(\hat{\lambda}(t_n, \theta_{n,h}) \in \Lambda_n(t_n, \theta_{n,h})\) satisfying \(\hat{P}_\rho(t_n, \theta_{n,h}, \lambda(t_n, \theta_{n,h})) = \sup_{\lambda \in \Lambda_n(t_n, \theta_{n,h})} \hat{P}_\rho(t_n, \theta_{n,h}, \lambda)\) exists w.p.a.1. \(\lambda(t_n, \theta_{n,h}) = O_p(n^{-1/2})\), and \(\hat{m}(t_n) = O_p(n^{-1})\).

**LEMMA 6.** Suppose \(\hat{t}_n\) (defined in (A.20)) exists w.p.a.1, \(\max_{1 \leq i \leq n} ||m_i(\hat{t}_n, \theta_{n,h})|| = O_p(n^{1/(2+\delta)})\), \(\lambda_{\max}(\hat{\Delta}(\hat{t}_n)) \leq K\) w.p.a.1 for some \(K < \infty\), and \(\sup_{\lambda \in \Lambda_n(t_n, \theta_{n,h})} \hat{P}(t_n, \theta_{n,h}, \lambda) = O_p(n^{-1})\). Then, \(\hat{m}(\hat{t}_n) = O_p(n^{-1/2})\).

**Remark.** Lemmas 4–6 hold for any subsequence \(\{w_n : n \geq 1\}\) in place of \(\{n\}\).

**Proof of (A.19).** For notational simplicity we use \(n\) instead of \(w_n\) in the proof. We use the abbreviations \(\hat{P}_\rho(t, \lambda) = \hat{P}_\rho(t, \theta_{n,h}, \lambda)\) and \(m_i(t) = m_i(t, \theta_{n,h})\). Using Lemma 3, Lemma 5 with \(t_n = t_{n,h}\) yields \(\sup_{\lambda \in \Lambda_n(t_n, \theta_{n,h})} \hat{P}_\rho(t_n, \theta_{n,h}, \lambda) = O_p(n^{-1})\). Using this result, Assumption GEL(d), and Lemma 3(ii), Lemma 6 gives \(\hat{m}(\hat{t}_n) = O_p(n^{-1/2})\). Therefore, \(O_p(n^{-1/2}) = \hat{m}(\hat{t}_n) = n^{-1} \sum_{i=1}^{n} m(W_i, \theta_{n,h}) - (t_{n,h}, 0) + (t_{n,h}, 0) - (\hat{t}_n, 0)\) (A.21) and hence \(\hat{t}_n - t_{n,h} \rightarrow p 0\). By Lemma 3(i) we have \(\hat{\lambda}_{\min}(\hat{\Delta}(t_{w_n,h})) \geq \varepsilon\) w.p.a.1, and it was just shown that \(t_n - t_{n,h} \rightarrow p 0\). These two statements imply \(\hat{\lambda}_{\min}(\hat{\Delta}(\hat{t}_n)) \geq \varepsilon\) w.p.a.1.
using the technicalities in Lemma 3(iv) and simply multiplying out. By Lemma 3(ii), the second part of Lemma 3(iv), and $\hat{t}_n - t_{n,h} \to p 0_{p}$, we have $\max_{1 \leq i \leq n} |m_i(\hat{G}_n, \theta_{n,h})| = O_p(n^{1/(2+\delta)})$. Using this, we apply Lemma 5 with $t_n = \hat{t}_n$ and conclude that $\lambda_n = \lambda(\hat{G}_n, \theta_{n,h}) \in \Lambda_n(\hat{G}_n, \theta_{n,h})$ satisfying $\bar{P}_\rho(\hat{G}_n, \lambda_n) = \sup_{\lambda \in \Lambda_n(\hat{G}_n, \theta_{n,h})} \bar{P}_\rho(\hat{G}_n, \lambda) \exists w.p.a.1$ and $\lambda_n = O_p(n^{-1/2})$. Therefore, the first-order conditions

\[
\begin{equation}
\tag{A.22}
n^{-1} \sum_{i=1}^n \rho_1(j^*_n m_i(\hat{G}_n))m_i(\hat{G}_n) \to 0_k
\end{equation}
\]

hold w.p.a.1. Expanding the first-order conditions in $\lambda$ around $0_k$, there exists a mean value $\lambda_n$ between $0_k$ and $\lambda_n$ (that may be different for each row) such that w.p.a.1

\[
0_k = -\hat{m}_n(\hat{G}_n) + \left[ n^{-1} \sum_{i=1}^n \rho_2(j^*_n m_i(\hat{G}_n))m_i(\hat{G}_n) \right] \lambda_n
\]

\[
= -\hat{m}_n(\hat{G}_n) - \Delta_n \lambda_n,
\]

(A.23)

where the matrix $\Delta_n$ has been defined implicitly. Because $\lambda_n = O_p(n^{-1/2})$, $\max_{1 \leq i \leq n} |m_i(\hat{G}_n, \theta_{n,h})| = O_p(n^{1/(2+\delta)})$, and $\rho_2(0) = -1$, we have $\max_{1 \leq i \leq n} |\rho_2(j^*_n m_i(\hat{G}_n))| + 1 \to p 0$. Thus, by Lemma 3(i), $\lambda_n - \Delta(\hat{G}_n) \to p 0_k \times k$. In addition, by the preceding argument, $\lambda_n - \Delta(\hat{G}_n) \geq \epsilon$ w.p.a.1. In consequence, $\Delta_n$ is invertible w.p.a.1, and

\[
\lambda_n = -\Delta_n^{-1} \hat{m}_n(\hat{G}_n)
\]

(A.24)

w.p.a.1. Inserting this into a second-order Taylor expansion for $\bar{P}_\rho(\hat{G}_n, \lambda)$ with mean value $\lambda_n$, it follows that w.p.a.1

\[
\bar{P}_\rho(\hat{G}_n, \lambda_n) = -2 \lambda_n^t \hat{m}_n(\hat{G}_n) + \lambda_n^t \left[ n^{-1} \sum_{i=1}^n \rho_2(j^*_n m_i(\hat{G}_n))m_i(\hat{G}_n) \right] \lambda_n
\]

\[
= 2 \hat{m}_n(\hat{G}_n)^t \Delta_n^{-1} \hat{m}_n(\hat{G}_n) - \hat{m}_n(\hat{G}_n)^t \Delta_n^{-1} \Delta_n \Delta_n^{-1} \hat{m}_n(\hat{G}_n),
\]

(A.25)

where $\Delta_n = n^{-1} \sum_{i=1}^n \rho_2(j^*_n m_i(\hat{G}_n))m_i(\hat{G}_n) m_i(\hat{G}_n)$ satisfies $\Delta_n \to p 0$ by the same argument as used earlier to show $\Delta_n - \Delta(\hat{G}_n) \to p 0_k \times k$. Therefore, up to $O_p(1)$ terms, we have

\[
\tau_n^{GEL}(\theta_{n,h}) = \frac{n \hat{m}_n(\hat{G}_n)^t \Delta_n^{-1} \hat{m}_n(\hat{G}_n)}{n \hat{m}_n(\hat{G}_n)^t \Delta_n^{-1} \hat{m}_n(\hat{G}_n)}
\]

\[
= n \left( \frac{\Delta_n}{\Delta_n(\theta_{n,h})} - \left( \hat{G}_n, 0_0 \right) \right)^t \Delta_n^{-1} \left( \frac{\Delta_n}{\Delta_n(\theta_{n,h})} - \left( \hat{G}_n, 0_0 \right) \right)
\]

\[
= n \left( \frac{\Delta_n}{\Delta_n(\theta_{n,h})} - \left( \hat{G}_n, 0_0 \right) \right)^t \hat{S}_n(\theta_{n,h})^{-1} \left( \frac{\Delta_n}{\Delta_n(\theta_{n,h})} - \left( \hat{G}_n, 0_0 \right) \right)
\]

\[
= \min_{t \in \mathbb{R}_+^\infty} n \left( \frac{\Delta_n}{\Delta_n(\theta_{n,h})} - \left( t, 0_0 \right) \right)^t \hat{S}_n(\theta_{n,h})^{-1} \left( \frac{\Delta_n}{\Delta_n(\theta_{n,h})} - \left( t, 0_0 \right) \right)
\]

\[
= \tau_n^{GEL}(\theta_{n,h}) = S_2 \left( n^{1/2} m_n(\theta_{n,h}), \hat{S}_n(\theta_{n,h}) \right),
\]

(A.26)

where the third equality holds because $\hat{S}_n(\theta_{n,h}) = n^{-1} \sum_{i=1}^n (m(W_i, \theta_{n,h}) - E m(W_i, \theta_{n,h}))(m(W_i, \theta_{n,h}) - E m(W_i, \theta_{n,h}))^t + O_p(1)$ and $t_{n,h} - t_n \to p 0_p$. The second to last equality holds by the following argument. Denote by $t_n^*$ the minimizing $t \in \mathbb{R}_+^\infty$ in the second to last equality of (A.26). We have to show that

\[
\begin{equation}
\tag{A.27}
n \left( \frac{\Delta_n}{\Delta_n(\theta_{n,h})} - \left( t_n^*, 0_0 \right) \right)^t \hat{S}_n(\theta_{n,h})^{-1} \left( \frac{\Delta_n}{\Delta_n(\theta_{n,h})} - \left( t_n^*, 0_0 \right) \right)
\end{equation}
\]
is $o_p(1)$. If this does not hold, then it could not be the case that $\hat{m}_n(t_n^*) = O_p(n^{-1/2})$ because if the latter were true, then the argument in (A.21)–(A.26) could be applied to $t_n^*$ instead of $t_n$, yielding $T_n^GEL(\theta_n, h) = n(\hat{m}_n(\theta_n, h) - (t_n^*, 0))\Sigma_n(\theta_n, h)^{-1}(m_n(\theta_n, h) - (t_n^*, 0))$ w.p.a.1, which is a contradiction. Therefore, $\hat{m}_n(t_n^*)$ is not $O_p(n^{-1/2})$. But then $T_n(\theta_n, h) = \min_{t \in R_n} n(\hat{m}_n(\theta_n, h) - (t, 0))\Sigma_n(\theta_n, h)^{-1}(m_n(\theta_n, h) - (t, 0))$ cannot be $O_p(1)$ because $\lambda_{\min}(\Sigma_n(\theta_n, h)^{-1}) \geq \varepsilon$ w.p.a.1 by the second part of Lemma 3(ii). Therefore, we get a contradiction to (A.10) (where the latter shows that $T_n(\theta_n, h) \rightarrow_d J_n$ and thus $T_n(\theta_n, h) = O_p(1)$). Therefore, the expression in (A.27) is indeed $o_p(1)$.

Proof of Lemma 3. The first part of Lemma 3(i) holds because (a) $\hat{\Delta}(t_{wn}) - EF_{wn} \hat{\Delta}(t_{wn}) = o_p(1)$ by a weak LLN for row-wise i.i.d. random variables given Assumptions GEL(a) and (b) and condition (vii) of (3.3) and (b) $\lambda_{\min}(EF_{wn} \hat{\Delta}(t_{wn})) \geq \varepsilon$ for some $\varepsilon > 0$ for all $n$ by Assumptions GEL(b) and (c). The second part of Lemma 3(i) and the first part of Lemma 3(ii) hold by a weak LLNs, Assumptions GEL(a) and (b), and condition (vii) of (3.3). The first part of Lemma 3(ii) holds using Assumptions GEL(a) and (b) and condition (vii) of (3.3); e.g., see Guggenberger and Smith (2005, eqn. (2.4)). The second part of Lemma 3(ii) holds by a weak LLN, Assumptions GEL(a) and (b), and condition (vii) of (3.3). Lemma 3(iii) holds by a Liapounov CLT for a row-wise i.i.d. triangular array of random variables applied to $m_i(t_{wn}, \theta_{wn}, h): i = 1, \ldots, n; n \geq 1$ using Assumptions GEL(a) and (b), condition (vii) of (3.3), and the fact that $m_i(t_{wn}, \theta_{wn}, h)$ has mean zero given the definition of $t_{wn}$. The second part of Lemma 3(iv) holds by Assumption GEL(b) and condition (vii) of (3.3).

Proof of Lemma 4. The result of the lemma follows from $\sup_{\lambda \in \Lambda_n} \sup_{1 \leq i \leq n} |\lambda'|^2 m_i(t_n, \theta_n, h)] < O_p(n^{-1/(2+\delta/2)}n^{1/(2+\delta)}) = o_p(1)$.

Proof of Lemma 5. We use the abbreviations $\tilde{P}_\rho(\lambda) = \tilde{P}_\rho(t_n, \theta_n, h, \lambda)$ and $m_i = m_i(t_n, \theta_n, h)$ in this proof. Let $\lambda_{\rho} \in \Lambda_n$ be such that $\tilde{P}_\rho(\lambda_{\rho}) = \max_{\lambda \in \Lambda_n} \tilde{P}_\rho(\lambda)$. Such a $\lambda_{\rho} \in \Lambda_n$ exists w.p.a.1 because a continuous function takes on its maximum on a compact set and by Lemma 4 $\tilde{P}_\rho(\lambda)$ is twice continuously differentiable, i.e., $C^2$, on some open neighborhood of $\Lambda_n$ w.p.a.1. We now show that actually $\tilde{P}_\rho(\lambda_{\rho}) = \sup_{\lambda \in \Lambda_n, t_n, h} \hat{P}_\rho(\lambda)$ w.p.a.1, which then proves the first part of the lemma. By a second-order Taylor expansion around $\lambda = 0$, there is a $\lambda_{\rho}'$ on the line segment joining $0_k$ and $\lambda_{\rho}$ such that for some positive constants $C_1$ and $C_2$, we have

$\lambda_{\rho}' = \tilde{P}_\rho(0) \leq \tilde{P}_\rho(\lambda_{\rho})$

$= -2\lambda_{\rho}' m_n(t_n) + \lambda_{\rho}' n^{-1} \sum_{i=1}^n \rho_2(\lambda_{\rho}'m_i)m_i|\lambda_{\rho}'|\leq -2\lambda_{\rho}' m_n(t_n) - C_1\lambda_{\rho}' \hat{\Delta}(t_n)\lambda_{\rho} \leq 2||\lambda_{\rho}'||\parallel \hat{m}_n(t_n)\parallel - C_2||\lambda_{\rho}'\parallel^2$ (A.28)

w.p.a.1, where the second inequality follows because $\max_{1 \leq i \leq n} \rho_2(\lambda_{\rho}'m_i) < -1/2$ w.p.a.1 by Lemma 4, continuity of $\rho_2(\lambda)$ at zero, and $\rho_2(0) = -1$. The last inequality follows from $\lambda_{\rho}'(\hat{\Delta}(t_n)) \geq \varepsilon > 0$ w.p.a.1. Now, (A.28) implies that $(C_2/2)||\lambda_{\rho}'|| \leq ||\hat{m}_n(t_n)||$ w.p.a.1, the latter being $O_p(n^{-1/2})$. Therefore, $\lambda_{\rho} \in \text{int}(\Lambda_n)$ w.p.a.1. Hence, the first-order conditions for an interior maximum $\partial \tilde{P}_\rho(\lambda)/\partial \lambda = 0$ hold at $\lambda = \lambda_{\rho}$ w.p.a.1. By Lemma 4, $\lambda_{\rho} \in \Lambda_n(t_n, \theta_n, h)$ w.p.a.1, and thus by concavity of $\tilde{P}_\rho(\lambda)$ and convexity of
\( \widehat{\Lambda}_n(t_n, \theta_n, h) \) it follows that \( \widehat{P}_p(\lambda_n) = \sum_{\hat{\lambda} \in \widehat{\Lambda}_n(t_n, \theta_n, h)} \widehat{P}_p(\hat{\lambda}) \) w.p.a.1, which implies the first part of the lemma. From before, \( \lambda(t_n, \theta_n, h) = \hat{\lambda}_n = O_p(n^{-1/2}) \). Thus, the second part of the lemma holds. This, \( ||\hat{m}_n(t_n)|| = O_p(n^{-1/2}) \), and (A.28) give the third part of the lemma. 

Proof of Lemma 6. We use the abbreviations \( \widehat{P}_p(t, \lambda) = \widehat{P}_p(t, \theta_n, h, \lambda) \) and \( m_i(t) = m_i(t, \theta_n, h) \) in this proof. Wlog. \( \hat{m}_n(\hat{\gamma}_n) \neq 0 \) can be assumed. Define \( \hat{\lambda}_n = -n^{-1/2} \hat{m}_n(\hat{\gamma}_n)/||\hat{m}_n(\hat{\gamma}_n)|| \). Note that \( \hat{\lambda}_n \in \Lambda_n \) and thus \( \hat{\lambda}_n \in \widehat{\Lambda}_n(\hat{\gamma}_n) \) w.p.a.1 by Lemma 4 (applied with \( t_n = \hat{\gamma}_n \)). By a second-order Taylor expansion around \( \hat{\lambda} = 0 \), there is a \( \tilde{\lambda}_n \) on the line segment joining \( 0 \) and \( \hat{\lambda}_n \) such that for some positive constants \( C_1 \) and \( C_2 \), we have

\[
\hat{P}_p(\hat{\gamma}_n, \hat{\lambda}_n) = -2n^{-1/2} \hat{m}_n(\hat{\gamma}_n) + \hat{m}_n[n^{-1} \sum_{i=1}^{n} \rho_2(\hat{\gamma}_n, m_i(\hat{\gamma}_n))m_i(\hat{\gamma}_n)m_i(\hat{\gamma}_n)^t] \hat{\lambda}_n
\]

\[
\geq 2n^{-1/2} ||\hat{m}_n(\hat{\gamma}_n)|| - C_1 \hat{\lambda}_n \hat{\lambda}_n
\]

\[
\geq 2n^{-1/2} ||\hat{m}_n(\hat{\gamma}_n)|| - C_2 n^{-1}
\]

w.p.a.1, where the first inequality follows from \( \min_{1 \leq i \leq n} \rho_2(\hat{\gamma}_n, m_i(\hat{\gamma}_n)) \geq -1.5 \) w.p.a.1, which is implied by Lemma 4. The second inequality follows by \( \hat{\lambda}_n(\hat{\gamma}_n) \leq K < \infty \) w.p.a.1. The definition of \( t_n \) implies

\[
\hat{P}_p(\hat{\gamma}_n, \hat{\lambda}_n) \leq \sup_{\hat{\lambda} \in \widehat{\Lambda}_n(\hat{\gamma}_n, \theta_n, h)} \hat{P}_p(\hat{\gamma}_n, \hat{\lambda}) \leq \sup_{\hat{\lambda} \in \widehat{\Lambda}_n(t_n, h, \theta_n, h)} \hat{P}_p(t_n, h, \lambda) = O_p(n^{-1}).
\]

Combining equations (A.29) and (A.30) implies \( n^{-1/2} ||\hat{m}_n(\hat{\gamma}_n)|| = O_p(n^{-1}) \).

A2. General Results. This section is concerned with the general results of Section 9. First, we state a corollary to Theorem 3 that applies when the parameter space \( \Gamma \) takes on a partial product-space form, as in Assumption A of Andrews and Guggenberger (2010). In this case, the form of \( H \) and \( GH \) can be made more explicit, and the results of Theorem 3 hold under an assumption that eliminates the subsequences that appear in Assumption B0. Second, we prove Theorem 3.

Let \( \lfloor \cdot \rfloor \) denote the left endpoint of an interval that may be open or closed at the left end. Define \( \lfloor \cdot \rfloor \) analogously for the right endpoint. The following assumption implies Assumption A0. Let \( R_- = \{x \in R : x < 0 \} \) and \( R_{-\infty} = R_- \cup \{-\infty\} \).

Assumption A.

(a) For some \( \Gamma_1 \subset R^p, \Gamma_2 \subset R^q, \) and \( \Gamma_3(\gamma_1, \gamma_2) \subset T_3 \), which may depend on \( \gamma_1 \) and \( \gamma_2 \), \( \Gamma \) satisfies

\[
\Gamma = \{(\gamma_1, \gamma_2, \gamma_3) : \gamma_1 \in \Gamma_1, \gamma_2 \in \Gamma_2, \gamma_3 \in \Gamma_3(\gamma_1, \gamma_2)\}
\]

and

(b) \( \Gamma_1 = \prod_{m=1}^{p} \Gamma_{1,m} \) where \( \Gamma_{1,m} = [\gamma_{1,m}^0, \gamma_{1,m}^0] \) for some \( -\infty \leq \gamma_{1,m}^0 \leq \gamma_{1,m}^0 \leq \infty \) that satisfy \( \gamma_{1,m}^0 \leq 0 \leq \gamma_{1,m}^0 \) for \( m = 1, \ldots, p \).
Under Assumption A, it follows that

\[ H = H_1 \times H_2, \quad H_1 = \prod_{m=1}^{p} \begin{cases} R_{+\infty} & \text{if } \gamma_{1,m}^f = 0 \\ R_{-\infty} & \text{if } \gamma_{1,m}^f = 0 \\ R_{\infty} & \text{if } \gamma_{1,m}^f < 0 \quad \text{and} \quad \gamma_{1,m}^r > 0, \end{cases} \]

\[ H_2 = \operatorname{cl}(\Gamma_2), \]

where \( \operatorname{cl}(\Gamma_2) \) is the closure of \( \Gamma_2 \) with respect to \( R_q^{\infty} \). For example, if \( p = 1, \gamma_{1,1}^f = 0 \), and \( \Gamma_2 = R^d \), then \( H_1 = R_{+\infty}, H_2 = R_q^{\infty}, \) and \( H = R_{+\infty} \times R_q^{\infty} \).

Under Assumption A, the set \( GH \) reduces to

\[ GH = \{ (g, h) \in H \times H : g = (g_1, g_2), h = (h_1, h_2), \quad g_2 = h_2, \quad \text{and for} \}

\[ m = 1, \ldots, p, \] (i) \( g_{1,m} = 0 \) if \( |h_{1,m}| < \infty \), (ii) \( g_{1,m} \in R_{+\infty} \) if \( h_{1,m} = +\infty \), and (iii) \( g_{1,m} \in R_{-\infty} \) if \( h_{1,m} = -\infty \}, \]

\[ (A.32) \]

where \( g_1 = (g_1, \ldots, g_{1,p})' \in H_1 \) and \( h_1 = (h_1, \ldots, h_{1,p})' \in H_1 \). Note that for \( (g, h) \in GH \), we have \( |g_{1,m}| \leq |h_{1,m}| \) for all \( m = 1, \ldots, p \).

Given Assumption A, the following weakened version of Assumption B0 is sufficient.

**Assumption B'**. For some \( r > 0 \), all \( h \in H \), all sequences \( \{\gamma_n, n \geq 1\} \), and some distributions \( J_h, T_n(\theta_n,h) \to_d J_h \) under \( \{\gamma_n, n \geq 1\} \), where \( \gamma_n = ((\theta_{n,1}, \eta_{n,1}), (\theta_{n,2}, \eta_{n,2}), \gamma_{n,3}) \) and \( \theta_{n,h} = (\theta_{n,1}, \theta_{n,2}) \).

Assumption B' is the same as Assumption B of Andrews and Guggenberger (2010) except that it applies to CSs rather than tests, so that \( T_n(\theta) \) is evaluated at \( \theta_{n,h} \) rather than at the null value \( \theta_0 \) and \( \gamma_{n,h} \) is defined differently.

We have the following corollary to Theorem 3.

**Corollary 1.** Theorem 3 holds with \( H \) and \( GH \) defined in (A.32) and (A.33), respectively, with Assumptions A0 and B0 replaced by Assumptions A and B'.

**Remarks.**

1. Assumption B' is simpler and weaker than Assumption B0. But typically the work needed to verify these assumptions and the strength of the assumptions are almost the same. Hence, the main advantage of Corollary 1 is that when Assumption A holds one has the explicit forms for \( H \) and \( GH \) given in (A.32) and (A.33).

2. Corollary 1 is proved using the proof of Theorem 3 coupled with the argument given in (8.6)–(8.7) of the proof of Lemma 6 of Andrews and Guggenberger (2010).

**Proof of Theorem 3.** The proof of the results of Theorem 3 for AsyCS is analogous to that of Theorem 1 of Andrews and Guggenberger (2010) with the following changes: AsySz(\( \theta_0 \)) is replaced by \( 1 - \text{AsyCS} \), probabilities \( P_{\theta,\gamma}(\cdot) \) and expectations \( E_{\theta,\gamma}(\cdot) \) are
replaced by $P_\gamma(\cdot)$ and $E_\gamma(\cdot)$, respectively, because $\theta$ is a subvector of $\gamma$, the test statistic $T_n(\theta_0)$ is replaced by $T_n(\theta_0)$ throughout, where $\theta_0$ denotes the true value of $\theta$ which may depend on $n$, and one makes use of the fact that $\inf_{h \in H} J_h(c_{FG}(1 - \alpha)) = 1 - \Max_{\Fix}^\infty(\alpha)$, $\inf_{(g,h) \in GH} J_h(c_\gamma(1 - \alpha)) = 1 - \Max_{\Sub}^\infty(\alpha)$, etc., where $\Max_{\Fix}^\infty(\alpha)$ and $\Max_{\Sub}^\infty(\alpha)$ are defined in Andrews and Guggenberger (2010). The proof of the results for $\AsyMaxCP$ is quite similar to those for $\AsyCS$ and hence is not discussed.

The replacement of Assumptions A, B, E, and G of Andrews and Guggenberger (2010) by Assumptions A0, B0, E0, and G0 requires the following changes in the proof of Theorem 1 of that paper. First, we show that the results of Lemma 6(i)–(vi) of Andrews and Guggenberger (2010) hold with $\{\gamma_{\text{Asy}} : n \geq 1\}$ equal to $\{\gamma_{\text{Asy} \cdot \text{CS}, \text{Sub}} : n \geq 1\}$ (defined in this paper) under the assumptions of Theorem 3 of this paper. Lemma 6(i) of Andrews and Guggenberger (2010) (i.e., $(g, h) \in GH$) holds by the definition of $GH$ in this paper. The proof of Lemma 6(ii) is the same as in Andrews and Guggenberger (2010) (noting that $\{\gamma_{\text{Asy}} : n \geq 1\}$ in Lemma 6 is of the form $\{\gamma_{\text{Asy} \cdot \text{CS}, \text{Sub}} : n \geq 1\}$ considered in this paper) but with (8.6) holding by Assumption B0 rather than by the proof given in Andrews and Guggenberger (2010). The proof of Lemma 6(c) is much simpler than that in Andrews and Guggenberger (2010). By Assumption E0, $U_{\gamma_{\text{Asy}}}(x) - E_{\gamma_{\text{Asy}} \cdot \text{CS}, \text{Sub}}(x) \rightarrow p 0$ under $\{\gamma_{\text{Asy} \cdot \text{CS}, \text{Sub}} : n \geq 1\}$, and so (8.7) of Andrews and Guggenberger (2010) is not needed. The result and the result of Lemma 6(ii) yield Lemma 6(iii). Similarly, the proof of Lemma 6(d) is much simpler than that in Andrews and Guggenberger. The result of Lemma 6(iv) holds immediately by Assumption G0 and the result of Lemma 6(iii). The proof of Lemmas 6(v)–(vi) of Andrews and Guggenberger (2010) is the same as in that paper but with the result of (II) stated in the proof holding by Assumption B0 rather than by the proof given in Andrews and Guggenberger (2010).

Second, in the proof of Theorem 1 of Andrews and Guggenberger (2010), (8.10) holds by the definition of $GH$ (which guarantees that for each $(g, h) \in GH$ there is a sequence $\{\gamma_{\text{Asy} \cdot \text{CS}, \text{Sub}} : n \geq 1\}$ and the result of Lemma 6(v) of Andrews and Guggenberger (2010). The remainder of the proof holds using the modified version of Lemma 6 (which holds under Assumptions A0, B0, C, D, E0, F, and G0) with the only change being that $h_2 \in \overline{\{\theta \in \Gamma_1\}}$ in (8.14) is replaced by $h_2 \in \overline{\{\theta \in \Gamma_2\}}$. This concludes the adjustment of the proof of Theorem 1 of Andrews and Guggenberger (2010) to take account of the change in assumptions.

The improvement to the lower bound on $\AsyCS$ for subsampling CSs is obtained as follows. If the assumption is added to Lemma 5 of Andrews and Guggenberger (2010) that $\lim_{n \to \infty} P(T_n \leq c_\theta) = G_T(c_\infty)$, then the lemma yields the stronger conclusion that $P(T_n \leq c_\theta) \to G_T(c_\infty)$. This follows directly from the proof of Lemma 5(ii) of Andrews and Guggenberger (2010). Therefore, for any $(g, h) \in GH^*$ and any sequence $\{\gamma_{\text{Asy} \cdot \text{CS}, \text{Sub}} : n \geq 1\}$, the proof of Lemma 6(vi) of Andrews and Guggenberger (2010) yields the stronger conclusion that $P_{\gamma_{\text{Asy} \cdot \text{CS}, \text{Sub}}}(T_n(\theta_{\gamma_{\text{Asy} \cdot \text{CS}, \text{Sub}}, h) \leq c_{\gamma_{\text{Asy} \cdot \text{CS}, \text{Sub}}, he_{\gamma_{\text{Asy} \cdot \text{CS}, \text{Sub}}, h, 1 - \alpha}) \to J_h(c_\gamma(1 - \alpha))$. Combining this with the proof of Theorem 1(ii) of Andrews and Guggenberger (2010) establishes the lower bound $\Min_{\Asy}^\infty(\alpha)$ to $\AsyCS$ given in the theorem.