Identification, weak instruments, and statistical inference in econometrics

Jean-Marie Dufour

CIRANO, CIREQ, and Département de sciences économiques, Université de Montréal

Abstract. We discuss statistical inference problems associated with identification and testability in econometrics. We consider inference in non-parametric models and weakly identified structural models (weak instruments). We point out that many ill-defined statistical problems, such as non-testable hypotheses, occur in these areas and are typically associated with asymptotic approximations. In non-parametric models, such problems include testing moments and inference under heteroscedasticity or serial dependence of unknown form. For weakly identified structural models, difficulties are typically associated with improper pivots, and we review recent developments aimed at proposing more reliable procedures, including alternative proposed statistics, bounds, projection, split-sampling, conditioning, Monte Carlo tests. JEL classification: C1, C12, C14, C15, C3, C5

Identification, instruments faibles, et inférence statistique en économie. Nous analysons les problèmes d’inférence associés à l’identification et à la testabilité en économie. Nous considérons l’inférence dans les modèles non-paramétriques et les modèles structurels faiblement identifiés (instruments faibles). Nous remarquons que beaucoup de problèmes mal posés, tels que des hypothèses non testables, apparaissent dans ces domaines et que ceux-ci sont typiquement associés à l’emploi d’approximations asymptotiques. Dans les modèles non-paramétriques, de tels problèmes incluent les tests sur les moments et l’inférence sous hétéroscédasticité ou dépendance sérielle de forme non spécifiée. Dans les modèles structurels faiblement identifiés, ces difficultés sont

This paper is based on the author’s Presidential Address to the Canadian Economics Association given on 31 May 2003, at Carleton University (Ottawa). The author thanks Bryan Campbell, Tarek Jouini, Lynda Khalaf, William McCausland, Nour Meddahi, Benoît Perron, and Mohamed Taamouti for several useful comments. This work was supported by the Canada Research Chair Program (Chair in Econometrics, Université de Montréal), the Alexander-von-Humboldt Foundation (Germany), the Canadian Network of Centres of Excellence (program on Mathematics of Information Technology and Complex Systems (MITACS)), the Canada Council for the Arts (Killam Fellowship), the Natural Sciences and Engineering Research Council of Canada, the Social Sciences and Humanities Research Council of Canada, the Fonds de recherche sur la société et la culture (Québec), and the Fonds de recherche sur la nature et les technologies (Québec). Email: jean.marie.dufour@umontreal.ca
habituellement associées à l’emploi de fonctions pivotales impropres et nous présentons un survol des méthodes récentes ayant pour objectif d’obtenir des procédures plus fiables, ce qui comprend les différentes statistiques proposées, l’emploi de bornes, la subdivision d’échantillon, les techniques de projection, le conditionnement et les tests de Monte Carlo.

1. Introduction

The main objective of econometrics is to supply methods for analysing economic data, building models, and assessing alternative theories. Over the last 25 years, econometric research has led to important developments in many areas, such as: (1) new fields of applications linked to the availability of new data, financial data, micro-data, panels, qualitative variables; (2) new models: multivariate time series models, GARCH-type processes; (3) a greater ability to estimate non-linear models that require an important computational capacity; (4) methods based on simulation: bootstrap, indirect inference, Markov chain Monte Carlo techniques; (5) methods based on weak distributional assumptions: non-parametric methods, asymptotic distributions based on ‘weak regularity conditions’; (6) discovery of various non-regular problems that require non-standard distributional theories, such as unit roots and unidentified (or weakly identified) models.

An important component of this work is the development of procedures for testing hypotheses (or models). Indeed, a view widely held by both scientists and philosophers is that testability or the formulation of testable hypotheses constitutes a central feature of scientific activity – a view we share. With the exception of mathematics, it is not clear that a discipline should be viewed as scientific if it does not lead to empirically testable hypotheses. But this requirement leaves open the question of formulating operational procedures for testing models and theories. To date, the only coherent – or, at least, the only well developed – set of methods are those supplied by statistical and econometric theory.

Last year, on the same occasion, MacKinnon (2002) discussed the use of simulation-based inference methods in econometrics, specifically bootstrapping, as a way of getting more reliable tests and confidence sets. In view of the importance of the issue, we also consider questions associated with the development of reliable inference procedures in econometrics. But our exposition will be, in a way, more specialized and in another way more general – and critical. Specifically, we shall focus on general statistical issues raised by identification in econometric models and, more specifically, on weak instruments in the context of structural models e.g., simultaneous equations models (SEM). We will find it useful to bring together two separate streams of literature: namely, results (from mathematical statistics and econometrics) on testability in non-parametric models and the recent econometric research on weak instruments.1 In particular, we shall

1 By a non-parametric model (or hypothesis), we mean a set of possible data distributions such that a distribution (e.g., the ‘true’ distribution) cannot be singled out by fixing a finite number of parameter values.
emphasize that identification problems arise in both literatures and have similar consequences for econometric methodology. Further, the literature on non-parametric testability sheds light on various econometric problems and their solutions.

Simultaneous equations models (SEM) are related in a natural way to the concept of equilibrium postulated by economic theory, both in microeconomics and macroeconomics. So it is not surprising that SEM were introduced and most often employed in the analysis of economic data. Methods for estimating and testing such models constitute a hallmark of econometric theory and represent one of its most remarkable achievements. The problems involved are difficult, raising, among various issues, the possibility of observational equivalence between alternative parameter values (non-identification) and the use of instrumental variables (IV). Further, the finite-sample distributional theory of estimators and test statistics is very complex, so inference is typically based on large-sample approximations. (For reviews, see Phillips 1983; Taylor 1983.)

IV methods have become a routine part of econometric analysis and, despite a lot of loose ends (often hidden by asymptotic distributional theory), the topic of SEM was dormant until a few years ago. Roughly speaking, an instrument should have two basic properties: first, it should be independent of (or, at least, uncorrelated with) the disturbance term in the equation of interest (exogeneity); second, it should be correlated with the included endogenous explanatory variables for which it is supposed to serve as an instrument (relevance). The exogeneity requirement has been well known from the very beginning of IV methods. The second one was also known from the theory of identification, but its practical importance was not well appreciated and was often hidden from attention by lists of instruments relegated to footnotes (if not simply absent) in research papers. It returned to centre stage with the discovery of so-called weak instruments, which can be interpreted as instruments with little relevance (i.e., weakly correlated with endogenous explanatory variables). Weak instruments lead to poor performance of standard econometric procedures and cases where they have pernicious effects may be difficult to detect.2 Interest in the problem also goes far beyond IV regressions and SEM, because it underscores the pitfalls in using large-sample approximations, as well as the importance of going back to basic statistical theory when developing econometric methods.

A parameter (or a parameter vector) in a model is not identified when it is not possible to distinguish between alternative values of the parameter. In parametric models, this is typically interpreted by stating that the postulated distribution of the data – as a function of the parameter vector (the likelihood

---

2 Early papers in which attention was called to the problem include Nelson and Startz (1990a,b), Buse (1992), Choi and Phillips (1992), Maddala and Jeong (1992), and Bound, Jaeger, and Baker (1993, 1995).
function) – can be the same for different values of the parameter vector. An important consequence of this sort of situation is a statistical impossibility: we cannot design a data-based procedure for distinguishing between equivalent parameter values (unless additional information is introduced). In particular, no reasonable test can be produced. In non-parametric models, identification is more difficult to characterize because a likelihood function (involving a finite number of parameters) is not available, and parameters are often introduced through more abstract techniques (e.g., functionals of distribution functions). But the central problem is the same: can we distinguish between alternative values of the parameter? So, quite generally, an identification problem can be viewed as a special form of non-testability. Specifically,

- identification involves the possibility of distinguishing different parameter values on the basis of the corresponding data distributions, while
- testability refers to the possibility of designing procedures that can discriminate between subsets of parameter values.

Alternatively, a problem of non-testability can be viewed as a form of non-identification (or underidentification). These problems are closely related. Furthermore, it is well known that one can create a non-identified model by introducing redundant parameters, and conversely identification problems can be eliminated by transforming the parameter space (e.g., by reducing the number of parameters). Problems of non-identification are associated with bad parameterizations, inappropriate choices of parameter representations. We will see below that the same remark applies quite generally to non-testability problems, especially in non-parametric set-ups.

In this paper, we pursue two main objectives: first, we analyse the statistical problems associated with non-identification within the broader context of testability; second, we review the inferential issues linked to the possible presence of weak instruments in structural models. More precisely, regarding the issue of testability, the following points will be emphasized:

1. many models and hypotheses are formulated in ways that make them fundamentally non-testable; in particular, this tends to be the case in non-parametric set-ups;
2. such difficulties arise in basic apparently well-defined problems, such as: (a) testing a hypothesis about a mean when the observations are independent and identically distributed (i.i.d.); (b) testing a hypothesis about a

---


4 By a reasonable test, we mean here a test that both satisfies a level constraint and may have power superior to the level when the tested hypothesis (the null hypothesis) does not hold. This will be discussed in greater detail below.
3. some parameters tend to be non-testable (badly identified) in non-parametric models while others are not; in particular, non-testability easily occurs for moments (e.g., means, variances) while it does not for quantiles (e.g., medians); from this viewpoint, moments are not a good way of representing the properties of distributions in non-parametric set-ups, while quantiles are so;

4. these phenomena underscore parametric non-separability problems: statements about a given parameter (often interpreted as the parameter of interest) are not empirically meaningful without information about other parameters (often called nuisance parameters); but hypotheses that set both the parameter of interest and some nuisance parameters may well be testable in such circumstances, so that the development of appropriate inference procedures should start from a joint approach;

5. to the extent that asymptotic distributional theory is viewed as a way of producing statistical methods that are valid under ‘weak’ distributional assumptions, it is fundamentally misleading because, under non-parametric assumptions, such approximations are arbitrarily bad in finite samples.

Concerning weak instruments, we will review the associated problems and proposed solutions, with an emphasis on finite-sample properties and the development of tests and confidence sets that are robust to the presence of weak instruments. In particular, the following points will be stressed:

1. in accordance with basic statistical theory, one should always look for pivots as the fundamental ingredient for building tests and confidence sets; this principle appears to be especially important when identification problems are present;

2. parametric non-separability arises in striking ways when some parameters may not be identified, so that proper pivots may easily involve many more parameters than the parameter of interest; this also indicates that the common distinction between parameters of interest and nuisance parameters can be quite arbitrary, if not misleading;

3. important additional criteria for evaluating procedures in such contexts include various forms of invariance (or robustness), such as: (a) robustness to weak instruments; (b) robustness to instrument exclusion; (c) robustness to the specification of the model for the endogenous explanatory variables in the equation(s) of interest;

4. weak instrument problems underscore in a striking way the limitations of large-sample arguments for deriving and evaluating inference procedures;

5. very few informative pivotal functions have been proposed in the context of simultaneous equations models;
6. the early statistic proposed by Anderson and Rubin (1949, AR) constitutes one of the (very rare) truly pivotal functions proposed for SEM; furthermore, it satisfies all the invariance properties listed above, so that it may reasonably viewed as a fundamental building block for developing reliable inference procedures in the presence of weak instruments;
7. a fairly complete set of inference procedures that allow one to produce tests and confidence sets for all model parameters can be obtained through projection techniques;
8. various extensions and improvements over the AR method are possible, especially in improving power; however, it is important to note that these often come at the expense of using large-sample approximations or giving up robustness.

The literature on weak instruments is growing rapidly, and we cannot provide here a complete review. In particular, we will not discuss in any detail results on estimation, the detection of weak instruments, or asymptotic theory in this context. For that purpose, we refer the reader to the excellent survey recently published by Stock, Wright, and Yogo (2002).

The paper is organized as follows. In the next two sections, we review succinctly some basic notions concerning models (section 2) and statistical theory (section 3), which are important for our discussion. In section 4, we study testability problems in non-parametric models. In section 5, we review the statistical difficulties associated with weak instruments. In section 6, we examine a number of possible solutions in the context of linear SEM, while extensions to non-linear or non-Gaussian models are considered in section 7. We conclude in section 8.

2. Models

The purpose of econometric analysis is to develop mathematical representations of data, which we call models or hypotheses (models subject to restrictions). A hypothesis should have two basic features.

1. It must restrict the expected behaviour of observations, be informative. A non-restrictive hypothesis says nothing and, consequently, does not teach us anything: it is empirically empty, void of empirical content. The more restrictive a model is, the more informative it is, and the more interesting it is.
2. It must be compatible with available data; ideally, we would like it to be true.

However, these two criteria are not always compatible:

1. the information criterion suggests the use of parsimonious models that usually take the form of parametric models based on strong assumptions; note the information criterion is emphasized by an influential view in
philosophy of science that stresses falsifiability as a criterion for the scientific character of a theory (Popper 1968);

2. in contrast, compatibility with observed data is most easily satisfied by vague models that impose few restrictions; vague models may take the form of parametric models with a large number of free parameters or non-parametric models that involve an infinite set of free parameters and thus allow for weak assumptions.

Models can be classified as being either deterministic or stochastic. Deterministic models, which claim to make arbitrarily precise predictions, are highly falsifiable but always inconsistent with observed data. Accordingly, most models used in econometrics are stochastic. Such models are unverifiable: as with any theory that makes an indefinite number of predictions, we can never be sure that the model will not be put in jeopardy by new data. Moreover, they are logically unfalsifiable: in contrast with deterministic models, a probabilistic model is usually logically compatible with all possible observation vectors.

Given these facts, it is clear that any criterion for assessing whether a hypothesis is acceptable must involve a conventional aspect. The purpose of hypothesis testing theory is to supply a coherent framework for accepting or rejecting probabilistic hypotheses. It is a probabilistic adaptation of the falsification principle. (For further discussion on the issues in this section, see Dufour 2000.)

3. Statistical notions

In this section, we review succinctly basic statistical notions which are essential for understanding the rest of our discussion. The general outlook follows modern statistical testing theory, derived from the Neyman-Pearson approach and described in standard textbooks, such as Lehmann (1986).

3.1. Hypotheses

Consider an observational experiment whose result can be represented by a vector of observations

$$X^{(n)} = (X_1, \ldots, X_n)'$$

where $X_i$ takes real values, and let

$$F(x) = F(x_1, \ldots, x_n) = P[X_1 \leq x_1, \ldots, X_n \leq x_n]$$

be its distribution, where $x = (x_1, \ldots, x_n)$. We denote by $\mathcal{F}_n$ the set of possible distribution functions on $\mathbb{R}^n[\bar{F} \in \mathcal{F}_n]$.

For various reasons, we prefer to represent distributions in terms of parameters. There are two ways of introducing parameters in a model. The first is to
define a function from a space of probability distributions to a vector in some Euclidean space:

\[ \theta : \mathcal{F}_n \rightarrow \mathbb{R}^p. \] (3)

Examples of such parameters include the moments of a distribution (mean, variance, kurtosis, etc.) and its quantiles (median, quartiles, etc.). Such functions are also called functionals. The second approach is to define a family of distribution functions that are indexed by a parameter vector \( \theta \):

\[ F(x) = F_0(x | \theta), \] (4)

where \( F_0 \) is a distribution function with a specific form. For example, if \( F_0(x | \theta) \) represents a Gaussian distribution with mean \( \mu \) and variance \( \sigma^2 \) (e.g., corresponding to a Gaussian law), we have \( \theta = (\mu, \sigma^2) \).

A model is parametric if the distribution of the data is specified up to a finite number of (scalar) parameters. Otherwise, it is non-parametric. A hypothesis \( H_0 \) on \( X^{(n)} \) is an assertion of the type

\[ \mathcal{H}_0 \equiv \{ F(\cdot): F(x) = F_0(x | \theta_1, \theta_2) \text{ and } \theta_1 = \theta_1^0 \}. \] (6)

We usually abbreviate this as

\[ H_0: \theta_1 = \theta_1^0. \] (7)

In such a case, we call \( \theta_1 \) the parameter of interest and \( \theta_2 \) a nuisance parameter; the parameter of interest is set by \( H_0 \) but the nuisance parameter remains unknown. \( H_0 \) may be interpreted as follows: there is at least one distribution in \( H_0 \) that can be viewed as a representation compatible with the observed ‘behaviour’ of \( X^{(n)} \). Then we can say that

\[ H_0 \text{ is acceptable } \iff (\exists F \in \mathcal{H}_0) \text{ F is acceptable} \] (8)

or, equivalently,

\[ H_0 \text{ is unacceptable } \iff (\forall F \in \mathcal{H}_0) \text{ F is unacceptable}. \] (9)
Showing that \( H_0 \) is unacceptable requires one to show that all distributions in \( \mathcal{H}_0 \) are incompatible with the observed data.

### 3.2. Test level and size

A test for \( H_0 \) is a rule by which one decides to reject or accept the hypothesis (or to view it as incompatible with the data). It usually takes the form

\[
\begin{align*}
\text{reject } H_0 & \quad \text{if } S_n(X_1, \ldots, X_n) > c \\
\text{do not reject } H_0 & \quad \text{if } S_n(X_1, \ldots, X_n) \leq c.
\end{align*}
\]  

The test has level \( \alpha \) iff

\[
P_F[\text{Rejecting } H_0] \leq \alpha \text{ for all } F \in \mathcal{H}_0
\]

or, equivalently,

\[
\sup_{F \in \mathcal{H}_0} P_F[\text{Rejecting } H_0] \leq \alpha,
\]

where \( P_F[\cdot] \) is the function (probability measure) giving the probability of an event when the data distribution function is \( F \). The test has size \( \alpha \) if

\[
\sup_{F \in \mathcal{H}_0} P_F[\text{Rejecting } H_0] = \alpha.
\]

\( H_0 \) is testable if we can find a finite number \( c \) that satisfies the level restriction. Probabilities of rejecting \( H_0 \) for distributions outside \( \mathcal{H}_0 \) (i.e., for \( F \notin \mathcal{H}_0 \)) define the power function of the test.\(^5\) Power describes the ability of a test to detect a ‘false’ hypothesis. Alternative tests are typically assessed by comparing their powers: between two tests with the same level, the one with the highest power against a given alternative distribution \( F \notin \mathcal{H}_0 \) is deemed preferable (at least under this particular alternative). Among tests with the same level, we typically like to have a test with the highest possible power against ‘alternatives of interest.’

As the set \( \mathcal{H}_0 \) gets larger, the test procedure must satisfy a bigger set of constraints: the larger is the set of distributions compatible with a null hypothesis, the stronger are the restrictions on the test procedure. In other words, the less restrictive a hypothesis is, the more restricted will be the corresponding test procedure. It is easy to understand that imposing a large set of restrictions on a test procedure may reduce its power against specific alternatives. There may be

---

\(^5\) More formally, the power function can be defined as the function: \( P(F) = P_F[\text{Rejecting } H_0] \) for \( F \in \mathcal{H}_1 \setminus \mathcal{H}_0 \), where \( \mathcal{H}_1 \) is an appropriate subset of the set of all possible distributions \( \mathcal{F}_n \). Sometimes, it is also defined on the set \( \mathcal{H}_1 \cup \mathcal{H}_0 \), in which case it should satisfy the level constraint for \( F \in \mathcal{H}_0 \).
a point where the restrictions are no longer implementable, in the sense that no procedure that has some power can satisfy the level constraint: $H_0$ is non-testable. In such a case, we have an ill-defined test problem.

In a framework such as the one in (6), where we distinguish between a parameter of interest $\theta_1$ and a nuisance parameter $\theta_2$, this is typically due to heavy dependence of the distribution of $S_n$ on the nuisance parameter $\theta_2$. If the latter is specified, we may be able to find a (finite) critical value $c = c(\alpha, \theta_2)$ that satisfies the level constraint (11). But in ill-defined problems, $c(\alpha, \theta_2)$ depends heavily on $\theta_2$, so that it is not possible to find a useful (finite) critical value for testing $H_0$, that is, $\sup_{\theta_2} c(\alpha, \theta_2) = \infty$. Besides, even if this is the case, it does not imply that a hypothesis that would fix both $\theta_1$ and $\theta_2$ is not testable; that is, the hypothesis $H_0' : (\theta_1, \theta_2) = (\theta_1^0, \theta_2^0)$ may be perfectly testable. But only a complete specification of the vector $(\theta_1, \theta_2)$ does allow one to interpret the values taken by the test statistic $S_n$ (non-separability).

3.3. Confidence sets and pivots

If we consider a hypothesis of the form

$$H_0(\theta_1^0) : \theta_1 = \theta_1^0,$$

and if we can build a different test $S_n(\theta_1^0; X_1, \ldots, X_n)$ for each possible value of $\theta_1^0$, we can determine the set of values that can be viewed as compatible with the data according to the tests considered:

$$C = \{ \theta_1^0 : S_n(\theta_1^0; X_1, \ldots, X_n) \leq c(\theta_1^0) \}.$$  \hfill (15)

If

$$P_F[\text{Rejecting } H_0(\theta_1^0)] \leq \alpha \quad \text{for all } F \in \mathcal{H}(F_0, \theta_1^0),$$

we have

$$\inf_{\theta_1, \theta_2} P_F[\theta_1 \in C] \geq 1 - \alpha.$$  \hfill (17)

$C$ is a confidence set with level $1 - \alpha$ for $\theta_1$. The set $C$ covers the ‘true’ parameter value $\theta_1$ with probability at least $1 - \alpha$. The minimal probability of covering the true value of $\theta_1$, that is, $\inf_{\theta_1, \theta_2} P[\theta_1 \in C]$, is called the size of the confidence set.

In practice, confidence regions (or confidence intervals) were made possible by the discovery of pivotal functions (or pivots): a pivot for $\theta_1$ is a function $S_n(\theta_1; X_1, \ldots, X_n)$ whose distribution does not depend on unknown parameters
(nuisance parameters); in particular, the distribution does not depend on \( \theta_2 \). More generally, the function \( S_n(\theta_1; X_1, \ldots, X_n) \) is *boundedly pivotal* if its distribution function may depend on \( \theta \) but is bounded over the parameter space (see Dufour 1997). When we have a pivotal function (or a boundedly pivotal function), we can find a point \( c \) such that

\[
P_F[S_n(\theta_1; X_1, \ldots, X_n) \geq c] \leq \alpha, \quad \forall \theta_1.
\]

For example, if \( X_1, \ldots, X_n \overset{i.i.d.}{\sim} N[\mu, \sigma^2] \), the \( t \) statistic

\[
t_n(\mu) = \sqrt{n}(\bar{X}_n - \mu)/s_X,
\]

where \( \bar{X}_n = \Sigma_{i=1}^n X_i/n \) and \( s_X = \Sigma_{i=1}^n (X_i - \bar{X}_n)/(n - 1) \), follows a Student \( t(n-1) \) distribution, which does not depend on the unknown values of \( \mu \) and \( \sigma \); hence, it is a pivot. By contrast, \( \sqrt{n}(\bar{X}_n - \mu) \) is not a pivot because its distribution depends on \( \sigma \). More generally, in the classical linear model with several regressors, the \( t \) statistics for individual coefficients (say, \( t(\beta_i) = \sqrt{n}(\hat{\beta}_i - \beta_i)/\hat{\sigma}_\beta \)) constitute pivots because their distributions do not depend on unknown nuisance parameters; in particular, the values of the other regression coefficients disappear from the distribution.

### 3.4. Testability and identification

When formulating and trying to solve test problems, two types of basic difficulties can arise. First, there is no valid test that satisfies reasonable properties (such as depending upon the data): in such a case, we have a non-testable hypothesis, an empirically empty hypothesis. Second, the proposed statistic cannot be pivotal for the model considered: its distribution varies too much under the null hypothesis to determine a finite critical point satisfying the level restriction (18).

If a hypothesis is non-testable, we are not able to design a reasonable procedure for deciding whether it holds (without the introduction of additional data or information). This difficulty is closely related to the concept of identification in econometrics. A parameter \( \theta \) is identifiable iff

\[
\theta(F_1) \neq \theta(F_2) \Rightarrow F_1 \neq F_2.
\]

For \( \theta_1 \neq \theta_2 \), we can, in principle, design a procedure for deciding whether \( \theta = \theta_1 \) or \( \theta = \theta_2 \). The values of \( \theta \) are testable. More generally, a parametric transformation \( g(\theta) \) is identifiable iff

\[
g[\theta(F_1)] \neq g[\theta(F_2)] \Rightarrow F_1 \neq F_2.
\]
Intuitively, these definitions mean that different values of the parameter imply different distributions of the data, so that we may expect to be able to ‘tell’ the difference by looking at the data. This is certainly the case when a unique distribution is associated with each parameter value (e.g., we may use the Neyman-Pearson likelihood ratio test to make the decision), but this may not be the case when a parameter covers several distributions. In the next section, we examine several cases where this happens.

4. Testability, non-parametric models, and asymptotic methods

We will now discuss three examples of test problems that look perfectly well defined and sensible at first sight, but turn out to be ill defined when we look at them more carefully. These include (1) testing a hypothesis about a mean when the observations are independent and identically distributed (i.i.d.); (2) testing a hypothesis about a mean (or a median) with heteroscedasticity of unknown form; (3) testing the unit root hypothesis on an autoregressive model whose order can be arbitrarily large.6

4.1. Procedures robust to non-normality

One of the most basic problems in econometrics and statistics consists in testing a hypothesis about a mean, for example, its equality to zero. For instance, hypothesis tests on regression coefficients in linear regressions or, more generally, on parameters of models that are estimated by the generalized method of moments (GMM) can be viewed as extensions of this fundamental problem. If the simplest versions of the problem have no reasonable solution, the situation will not improve when we consider more complex versions (as is done routinely in econometrics).

The problem of testing a hypothesis about a mean has a very well-known and neat solution when the observations are independent and identically (i.i.d.) distributed according to a normal distribution: we can use a $t$ test. The normality assumption, however, is often considered to be too ‘strong.’ So it is tempting to consider a weaker (less restrictive) version of this null hypothesis, such as

$$H_0(\mu_0): X_1, \ldots, X_n \text{ are i.i.d. observations with } E(X_1) = \mu_0.$$ (22)

In other words, we would like to test the hypothesis that the observations have mean $\mu_0$, under the general assumption that $X_1, \ldots, X_n$ are i.i.d. Here $H_0(\mu_0)$ is a non-parametric hypothesis because the distribution of the data cannot be
completely specified by fixing a finite number of parameters. The set of possible data distributions (or data-generating processes) compatible with this hypothesis, that is,

\[ \mathcal{H}(\mu_0) = \{ \text{Distribution functions } f_n \in \mathcal{F}_n \text{ such that } H_0(\mu_0) \text{ is satisfied} \} \]  

is much larger here than in the Gaussian case and imposes very strong restrictions on the test. Indeed, the set \( \mathcal{H}(\mu_0) \) is so large that the following property must hold.

**THEOREM 1. MEAN NON-TESTABILITY IN NON-PARAMETRIC MODELS.** If a test has level \( \alpha \) for \( H_0(\mu_0) \), that is,

\[ P_{F_n}[\text{Rejecting } H_0(\mu_0)] \leq \alpha \text{ for all } F_n \in \mathcal{H}(\mu_0), \]  

then, for any \( \mu_1 \neq \mu_0 \),

\[ P_{F_n}[\text{Rejecting } H_0(\mu_0)] \leq \alpha \text{ for all } F_n \in \mathcal{H}(\mu_1). \]  

Further, if there is at least one value \( \mu_1 \neq \mu_0 \) such that

\[ P_{F_n}[\text{Rejecting } H_0(\mu_0)] \geq \alpha \text{ for at least one } F_n \in \mathcal{H}(\mu_1), \]  

then, for all \( \mu_1 \neq \mu_0 \),

\[ P_{F_n}[\text{Rejecting } H_0(\mu_0)] = \alpha \text{ for all } F_n \in \mathcal{H}(\mu). \]  

**Proof.** See Bahadur and Savage (1956).

In other words (by (25)), if a test has level \( \alpha \) for testing \( H_0(\mu_0) \) the probability of rejecting \( H_0(\mu_0) \) should not exceed the level irrespective how far the ‘true’ mean is from \( \mu_0 \). Further (by (27)), if ‘by luck’ the power of the test gets as high as the level, then the probability of rejecting should be uniformly equal to the level \( \alpha \). Here, the restrictions imposed by the level constraint are so strong that the test cannot have power exceeding its level: it should be insensitive to cases where the null hypothesis does not hold! An optimal test (say, at level 0.05) in such a problem can be run as follows: (1) ignore the data; (2) using a random number generator, produce a realization of a variable \( U \) according to a uniform distribution on the interval (0, 1); that is, \( U \sim U(0,1) \); (3) reject \( H_0 \) if \( U \leq 0.05 \). Clearly, this is not an interesting procedure. It is also easy to see that a similar result will hold if we add various non-parametric restrictions on the distribution, such as a finite variance assumption.
The above theorem also implies that tests based on the ‘asymptotic distribution’ of the usual $t$ statistic for $\mu = \mu_0$ ($t_n(\mu_0)$ defined in (19)) has size one under $H_0(\mu_0)$:

$$\sup_{F_n \in \mathcal{F}(\mu_0)} P_{F_n} [ |t_n(\mu_0)| > c ] = 1$$  \hspace{1cm} (28)

for any finite critical value $c$. In other words, procedures based on the asymptotic distribution of a test statistic have sizes that deviate arbitrarily from their nominal size.

A way to interpret what happens here is through the distinction between pointwise convergence and uniform convergence. Suppose, to simplify, that the probability of rejecting $H_0(\mu_0)$ when it is true depends on a single nuisance parameter $\gamma$ in the following way:

$$P_n(\gamma) \equiv P_{\gamma} [ |t_n(\mu_0)| > c ] = 0.05 + (0.95)e^{-|\gamma|n},$$  \hspace{1cm} (29)

where $\gamma \neq 0$. Then, for each value of $\gamma$, the test has level 0.05 asymptotically; that is,

$$\lim_{n \to \infty} P_n(\gamma) = 0.05,$$  \hspace{1cm} (30)

but the size of the test is one for all sample sizes:

$$\sup_{\gamma > 0} P_n(\gamma) = 1, \text{ for all } n.$$  \hspace{1cm} (31)

$P_n(\gamma)$ converges to a level of 0.05 pointwise (for each $\gamma$), but the convergence is not uniform, so that the probability of rejection is arbitrarily close to one for $\gamma$ sufficiently close to zero (for all sample sizes $n$).

Many other hypotheses lead to similar difficulties. Examples include

1. hypotheses about various moments of $X_t$:

   - $H_0(\sigma^2)$: $X_1, \ldots, X_n$ are i.i.d. observations such that $\text{Var} (X_t) = \sigma^2$
   - $H_0(\mu_p)$: $X_1, \ldots, X_n$ are i.i.d. observations such that $E(X_t^p) = \mu_p$;

2. most hypotheses on the coefficients of a regression (linear or non-linear), a structural equation (as in SEM), or a more general estimating function (Godambe 1960):

   - $H_0(\theta_0)$: $g_t(X_t, \theta_0) = u_t$, $t = 1, \ldots, T$, where $u_1, \ldots, u_T$ are i.i.d.
In econometrics, models of the form $H_0(\theta_0)$ are typically estimated and tested through a variant of the generalized method of moments (GMM), usually with weaker assumptions on the distribution of $u_1, \ldots, u_T$; see Hansen (1982), Newey and West (1987a), Newey and McFadden (1994) and Hall (1999). To the extent that GMM methods are viewed as a way to allow for ‘weak assumptions,’ it follows from the above discussion that they constitute pseudo-solutions of ill-defined problems.

It is important to observe that the above discussion does not imply that all non-parametric hypotheses are non-testable. In the present case, the problem of non-testability could be eliminated by choosing another measure of central tendency, such as a median:

$$H_0^{0.5}(m_0): X_1, \ldots, X_n \text{ are i.i.d. continuous r.v.s such that } \text{Med}(X_t) = m_0, \ t = 1, \ldots, T.$$  

$H_0^{0.5}(m_0)$ can be easily tested with a sign test (see Pratt and Gibbons 1981, chap. 2). More generally, hypotheses on the quantiles of the distribution of observations in random sample remain testable non-parametrically:

$$H_0^{p}(Q_{p0}): X_1, \ldots, X_n \text{ are i.i.d. observations such that } \Pr[X_t \leq Q_{p0}] = p, \ t = 1, \ldots, T.$$  

Moments are not empirically meaningful functionals in non-parametric models (unless strong distributional assumptions are added), though quantiles are so.

4.2. Procedures robust to heteroscedasticity of unknown form

Another common problem in econometrics consists in developing methods which remain valid in making inference on regression coefficients when the variances of the observations are not identical (heteroscedasticity). In particular, this may go as far as looking for tests that are ‘robust to heteroskedasticity of unknown form.’ But it is not widely appreciated that this involves very strong restrictions on the procedures that can satisfy this requirement. To see this, consider the problem that consists in testing whether $n$ observations are independent with common zero median, namely:

$$H_0: X_1, \ldots, X_n \text{ are independent random variables each with a distribution symmetric about zero.}$$  

Equivalently, $H_0$ states that the joint distribution $F_n$ of the observations belongs to the (huge) set $\mathcal{H}_0 = \{F_n \in \mathcal{F}_n: F_n \text{satisfies } H_0\}$. $H_0$ allows heteroscedasticity of unknown form.
THEOREM 2. CHARACTERIZATION OF HETEROSCEDASTICITY ROBUST TESTS. If a test has level $\alpha$ for $H_0$, where $0 < \alpha < 1$, then it must satisfy the condition

$$P[\text{Rejecting } H_0 \mid |X_1|, \ldots, |X_n|] \leq \alpha \text{ under } H_0.$$  \hfill (33)

Proof: See Pratt and Gibbons (1981, sec. 5.10) and Lehmann and Stein (1949).

In other words, a valid test with level $\alpha$ must be a sign test – or, more precisely, its level must be equal to $\alpha$ conditional on the absolute values of the observations (which amounts to considering a test based on the signs of the observations). From this, the following remarkable property follows.

COROLLARY 3. If, for all $0 < \alpha < 1$, the condition (33) is not satisfied, then the size of the test is equal to one; that is,

$$\sup_{F_n \in \mathcal{H}_0} P_{F_n}[\text{Rejecting } H_0] = 1.$$  \hfill (34)

In other words, if a test procedure does not satisfy (33) for all levels $0 < \alpha < 1$, then its true size is one irrespective of its nominal size. Most so-called heteroscedasticity robust procedures based on “corrected” standard errors (see White 1980; Newey and West 1987b; Davidson and MacKinnon 1993, chap. 16; Cushing and McGarvey 1999) do not satisfy condition (33) and consequently have size one. (For examples of size distortion, see Dufour 1981; Campbell and Dufour 1995, 1997.)

4.3. Procedures robust to autocorrelation of arbitrary form

As a third illustration, let us now examine the problem of testing the unit root hypothesis in the context of an autoregressive model whose order is infinite or is not bounded by a prespecified maximal order:

$$X_t = \beta_0 + \sum_{k=1}^{p} \lambda_k X_{t-k} + u_t, \quad u_t \overset{i.i.d.}{\sim} N[0, \sigma^2], \quad t = 1, \ldots, n,$$

where $p$ is not bounded a priori. This type of problem has attracted a lot of attention in recent years. (For Reviews of this huge literature, see Banerjee et al. 1993; Stock 1994; Tanaka 1996; Maddala and Kim 1998.) We wish to test

$$\hat{H}_0: \sum_{k=1}^{p} \lambda_k = 1$$  \hfill (36)
or, more precisely,

\[ \tilde{H}_0: X_t = \beta_0 + \sum_{k=1}^{p} \lambda_k X_{t-k} + u_t, \ t = 1, \ldots, n, \text{for some } p \geq 0, \]

\[ \sum_{k=1}^{p} \lambda_k = 1 \text{ and } u_t \overset{\text{i.i.d.}}{\sim} N[0, \sigma^2]. \quad (37) \]

About this problem, we can show the following theorem and corollary.

**THEOREM 4. UNIT ROOT NON-TESTABILITY IN NON-PARAMETRIC MODELS.** If a test has level \( \alpha \) for \( \tilde{H}_0 \), that is,

\[ P_{F_n}[\text{Rejecting } \tilde{H}_0] \leq \alpha \text{ for all } F_n \text{ satisfying } \tilde{H}_0, \]

then

\[ P_{F_n}[\text{Rejecting } \tilde{H}_0] \leq \alpha \text{ for all } F_n. \quad (39) \]


**COROLLARY 5.** If, for all \( 0 < \alpha < 1 \), the condition (39) is not satisfied, then the size of the test is equal to one; that is,

\[ \sup_{F_n \in \mathcal{H}_0} P_{F_n}[\text{Rejecting } \tilde{H}_0] = 1, \]

where \( \mathcal{H}_0 \) is the set of all data distributions \( F_n \) that satisfy \( \tilde{H}_0 \).

As in the mean problem, the null hypothesis is simply too ‘large’ (unrestricted) to allow testing from a finite data set. Consequently, all procedures that claim to offer corrections for very general forms of serial dependence (e.g., Phillips 1987; Phillips and Perron 1988) are affected by these problems; irrespective of the nominal level of the test, the true size under the hypothesis \( \tilde{H}_0 \) is equal to one.

To get a testable hypothesis, it is essential to fix jointly the order of the AR process (i.e., a numerical upper bound on the order) and the sum of the coefficients: for example, we could consider the following null hypothesis where the order of the autoregressive process is equal to 12:

\[ H_0(12): X_t = \beta_0 + \sum_{k=1}^{12} \lambda_k X_{t-k} + u_t, \ t = 1, \ldots, n, \]

\[ \sum_{k=1}^{12} \lambda_k = 1 \text{ and } u_t \overset{\text{i.i.d.}}{\sim} N[0, \sigma^2]. \quad (40) \]
The order of the autoregressive process is an essential part of the hypothesis: it is not possible to separate inference on the unit root hypothesis from inference on the order of the process. Similar difficulties will also occur for most other hypotheses on the coefficients of (37). For further discussion of this topic, the reader may consult Sims (1971a,b), Blough (1992), Faust (1996, 1999), and Pötscher (2002).

5. Structural models and weak instruments

Several authors in the past have noted that usual asymptotic approximations are not valid or lead to very inaccurate results when parameters of interest are close to regions where these parameters are no longer identifiable. The literature on this topic is now considerable. In this section, we shall examine these issues in the context of SEM.

5.1. Standard simultaneous equations model

Let us consider the standard simultaneous equations model:

\[ y = Y\beta + X_1\gamma + u \]  
\[ Y = X_1\Pi_1 + X_2\Pi_2 + V, \]

where \( y \) and \( Y \) are \( T \times 1 \) and \( T \times G \) matrices of endogenous variables, \( X_1 \) and \( X_2 \) are \( T \times k_1 \) and \( T \times k_2 \) matrices of exogenous variables, \( \beta \) and \( \gamma \) are \( G \times 1 \) and \( k_1 \times 1 \) vectors of unknown coefficients, \( \Pi_1 \) and \( \Pi_2 \) are \( k_1 \times G \) and \( k_2 \times G \) matrices of unknown coefficients, \( u = (u_1, \ldots, u_T)' \) is a \( T \times 1 \) vector of structural disturbances, and \( V = [V_1, \ldots, V_T]' \) is a \( T \times G \) matrix of reduced-form disturbances. Further, \( X = [X_1, X_2] \) is a full-column rank \( T \times k \) matrix,

---

where \( k = k_1 + k_2 \). Finally, to get a finite-sample distributional theory for the test statistics, we shall use the following assumptions on the distribution of \( u \):

\[
    u \text{ and } X \text{ are independent;}
\]

\[
    u \sim N[0, \sigma_u^2 I_T].
\]

(44) may be interpreted as the strict exogeneity of \( X \) with respect to \( u \).

Note that the distribution of \( V \) is not otherwise restricted; in particular, the vectors \( V_1, \ldots, V_T \) need not follow a Gaussian distribution and may be heteroscedastic. Below, we shall also consider the situation where the reduced-form equation for \( Y \) includes a third set of instruments \( X_3 \), which are not used in the estimation:

\[
    Y = X_1 \Pi_1 + X_2 \Pi_2 + X_3 \Pi_3 + V,
\]

where \( X_3 \) is a \( T \times k_3 \) matrix of explanatory variables (not necessarily strictly exogenous); in particular, \( X_3 \) may be unobservable. We view this situation as important because, in practice, it is quite rare that one can consider all the relevant instruments that could be used. Even more generally, we could also assume that \( Y \) obeys a general non-linear model of the form:

\[
    Y = g(X_1, X_2, X_3, V, \Pi),
\]

where \( g(\cdot) \) is a possibly unspecified non-linear function and \( \Pi \) is an unknown parameter matrix.

The model presented in (41)–(42) can be rewritten in reduced form as

\[
    y = X_1 \pi_1 + X_2 \pi_2 + v
\]

\[
    Y = X_1 \Pi_1 + X_2 \Pi_2 + V,
\]

where \( \pi_1 = \Pi_1 \beta + \gamma, \ v = u + V\beta \) and

\[
    \pi_2 = \Pi_2 \beta.
\]

Suppose, now, that we are interested in making inference about \( \beta \).

Equation (50) is the crucial equation governing identification in this system: we need to be able to recover \( \beta \) from the values of the regression coefficients \( \pi_2 \) and \( \Pi_2 \). The necessary and sufficient condition for identification is the well-known rank condition for the identification of \( \beta \):
\( \beta \) is identifiable iff \( \text{rank}(\Pi_2) = G \).  \hspace{1cm} (51)

We have a weak instrument problem when either \( \text{rank}(\Pi_2) < k_2 \) (non-identification), or \( \Pi_2 \) is close to having deficient rank (i.e., \( \text{rank}(\Pi_2) = k_2 \) with strong linear dependence between the rows (or columns) of \( \Pi_2 \)). There is no compelling definition of the notion near-identification, but reasonable characterizations include the condition that \( \det(\Pi_2'\Pi_2) \) is ‘close to zero,’ or that \( \Pi_2'\Pi_2 \) has one or several eigenvalues ‘close to zero’.

Weak instruments are notorious for causing serious statistical difficulties on several fronts: (1) parameter estimation; (2) confidence interval construction; (3) hypothesis testing. We now consider these problems in greater detail.

5.2. Statistical problems associated with weak instruments

The problems associated with weak instruments were originally discovered through their consequences for estimation. Work in this area includes

1. theoretical work on the exact distribution of two-stage least squares (2SLS) and other ‘consistent’ structural estimators and test statistics (Phillips 1983, 1984, 1985, 1989; Rothenberg 1984, Hillier 1990; Nelson and Startz 1990a, b; Buse 1992; Maddala and Jeong 1992; Choi and Phillips 1992; Dufour 1997);
2. weak-instrument (local to non-identification) asymptotics (Staiger and Stock 1997; Stock and Wright 2000);
3. empirical examples (Bound, Jaeger, and Baker 1995).

The main conclusions of this research can be summarized as follows.

1. Theoretical results show that the distributions of various estimators depend in a complicated way on unknown nuisance parameters. Thus, they are difficult to interpret.
2. When identification conditions are not satisfied, standard asymptotic theory for estimators and test statistics typically collapses.
3. With weak instruments,
   a) the 2SLS estimator becomes heavily biased (in the same direction as ordinary least squares (OLS);
   b) the distribution of the 2SLS estimator is quite far from the normal distribution (e.g., bimodal).
4. A striking illustration of these problems appears in the reconsideration by Bound, Jaeger, and Baker (1995) of a study on returns to education by Angrist and Krueger (1991). Using 329,000 observations, these authors found that replacing the instruments used by Angrist and Krueger (1991)
with randomly generated (totally irrelevant) instruments produced very similar point estimates and standard errors. This result indicates that the original instruments were weak.

For a more complete discussion of estimation with weak instruments, the reader may consult Stock, Wright, and Yogo (2002).

5.3. Characterization of valid tests and confidence sets

Weak instruments also lead to very serious problems when one tries to perform tests or build confidence intervals on the parameters of the model. Consider a situation where we have two parameters \( \theta_1 \) and \( \theta_2 \) (i.e., \( \theta = (\theta_1, \theta_2) \)) such that \( \theta_2 \) is no longer identified when \( \theta_1 \) takes a certain value, say, \( \theta_1 = \theta_1^0 \):

\[
L(y | \theta_1, \theta_2) \equiv \mathcal{L}(y | \theta_1^0).
\]

**Theorem 6.** If \( \theta_2 \) is a parameter whose value is not bounded, then the confidence region \( C \) with level \( 1 - \alpha \) for \( \theta_2 \) must have the following property:

\[
P_\theta[\text{C is unbounded}] > 0,
\]

and, if \( \theta_1 = \theta_1^0 \),

\[
P_\theta[\text{C is unbounded}] \geq 1 - \alpha.
\]

**Proof.** See Dufour (1997).

**Corollary 7.** If \( C \) does not satisfy the property given in the previous theorem, its size must be zero.

This will be the case, in particular, for any Wald-type confidence interval, obtained by assuming that

\[
t_{\theta_2} = \frac{\hat{\theta}_2 - \theta_2}{\hat{\sigma}_{\theta_2}} \approx N(0, 1),
\]

which yields confidence intervals of the form \( \hat{\theta}_2 - c\hat{\sigma}_{\theta_2} \leq \theta_2 \leq \hat{\theta}_2 + c\hat{\sigma}_{\theta_2} \), where \( P[|N(0, 1)| > c] \leq \alpha \). By the above corollary, this type of interval has level zero, irrespective of the critical value \( c \) used:

\[
\inf_{\theta} P_\theta \left[ \hat{\theta}_2 - c\hat{\sigma}_{\theta_2} \leq \theta_2 \leq \hat{\theta}_2 + c\hat{\sigma}_{\theta_2} \right] = 0.
\]

In such situations, the notion of standard error loses its usual meaning and does not constitute a valid basis for building confidence intervals. In SEM, for
example, this applies to standard confidence intervals based on 2SLS estimators and their asymptotic ‘standard errors.’

Correspondingly, if we wish to test a hypothesis of form $H_0: \theta_2 = \theta_2^0$, the size of any test of the form

$$\left| t_{\hat{\theta}_2}(\theta_2^0) \right| = \left| \frac{\hat{\theta}_2 - \theta_2^0}{\sigma_{\theta_2}} \right| > c(\alpha)$$

will deviate arbitrarily from its nominal size. No unique large-sample distribution for $t_{\hat{\theta}_2}$ can provide valid tests and confidence intervals based on the asymptotic distribution of $t_{\hat{\theta}_2}$. From a statistical viewpoint, this means that $t_{\hat{\theta}_2}$ is not a pivotal function for the model considered. More generally, this type of problem affects the validity of all Wald-type methods, which are based on comparing parameter estimates with their estimated covariance matrix.

By contrast, in models of the form (41)–(45), the distribution of the LR statistics for most hypotheses on model parameters can be bounded and cannot move arbitrarily: likelihood ratios are boundedly pivotal functions and provide a valid basis for testing and confidence set construction (see Dufour 1997).

The central conclusion here is tests and confidence sets on the parameters of a structural model should be based on proper pivots.

### 6. Approaches to weak instrument problems

What should the features of a satisfactory solution to the problem of making inference in structural models? We shall emphasize here four properties: (1) the method should be based on proper pivotal functions (ideally, a finite-sample pivot); (2) robustness to the presence of weak instruments; (3) robustness to excluded instruments; (4) robustness to the formulation of the model for the explanatory endogenous variables $Y$ (which is desirable in many practical situations).

In the light of these criteria, we shall discuss first, the Anderson-Rubin procedure, which in our view is the reference method for dealing with weak instruments in the context of standard SEM; second, the projection technique that provides a general way of making a wide spectrum of tests and confidence sets; and third, several recent proposals aimed at getting improvements over the AR procedure.

#### 6.1. Anderson-Rubin statistic

A solution to testing in the presence of weak instruments has been available for more than 50 years (Anderson and Rubin 1949) and is now centre stage again (Dufour 1997; Staiger and Stock 1997). It is interesting that the AR method can be viewed as an alternative way of exploiting ‘instruments’ for inference on a structural model, although it predates the introduction of 2SLS methods in
SEM (Theil 1953; Basmann 1957), which later became the most widely used method for estimating linear structural equations models. The basic problem considered consists in testing the hypothesis

$$H_0(\beta_0): \beta = \beta_0$$

in model (41)–(44). In order to do that, we consider an auxiliary regression obtained by subtracting $Y_0$ from both sides of (41) and expanding the right-hand side in terms of the instruments. This yields the regression

$$y - Y_0 = X_1\theta_1 + X_2\theta_2 + \varepsilon,$$

where $\theta_1 = \gamma + \Pi_1(\beta - \beta_0)$, $\theta_2 = \Pi_2(\beta - \beta_0)$ and $\varepsilon = u + V(\beta - \beta_0)$. Under the null hypothesis $H_0(\beta_0)$, this equation reduces to

$$y - Y_0 = X_1\theta_1 + \varepsilon,$$

so we can test $H_0(\beta_0)$ by testing $H_0(\beta_0): \theta_2 = 0$, in the auxiliary regression (59). This yields the following F-statistic – the Anderson-Rubin statistic – which follows a Fisher distribution under the null hypothesis:

$$AR(\beta_0) = \frac{[SS_0(\beta_0) - SS_1(\beta_0)]/k_2}{SS_1(\beta_0)/(T - k)} \sim F(k_2, T - k)$$

where $SS_0(\beta_0) = (y - Y_0)'M(X_1)(y - Y_0)$ and $SS_1(\beta_0) = (y - Y_0)'M(X)(y - Y_0)$; for any full-rank matrix $A$, we denote $P(A) = A(A'A)^{-1}A'$ and $M(A) = I - P(A)$. What plays the crucial role here is the fact that we have instruments ($X_2$) that can be related to $Y$ but are excluded from the structural equation. To draw inference on the structural parameter $\beta$, we ‘hang’ on the variables in $X_2$: if we add $X_2$ to the constrained structural equation (60), its coefficient should be zero. For these reasons, we shall call the variables in $X_2$ auxiliary instruments.

Since the latter statistic is a proper pivot, it can be used to build confidence sets for $\beta$:

$$C_\beta(\alpha) = \{\beta_0: AR(\beta_0) \leq F_\alpha(k_2, T - k)\},$$

where $F_\alpha(k_2, T - k)$ is the critical value for a test with a level $\alpha$ based on the $F(k_2, T - k)$ distribution. When there is only one endogenous explanatory

8 The basic ideas for using instrumental variables for inference on structural relationships appear to go back to Working (1927) and Wright (1928). For an interesting discussion of the origin of IV methods in econometrics, see Stock and Trebbi (2003).
variable \((G = 1)\), this set has an explicit solution involving a quadratic inequa-
tion; that is,

\[
C_\beta(\alpha) = \{\beta_0 : a\beta_0^2 + b\beta_0 + c \leq 0\},
\]  

(63)

where \(a = Y'HY\), \(H \equiv M(X_1) - M(X) [1 + k_2 F_\alpha(k_2, T - k)/(T - k)]\), \(b = -2Y'H_y\), and \(c = y'^HY_y\). The set \(C_\beta(\alpha)\) may easily be determined by finding the roots of the quadratic polynomial in equation (63); see Dufour and Jasiak (2001) and Zivot, Startz, and Nelson (1998) for details.

When \(G > 1\), the set \(C_\beta(\alpha)\) is not in general an ellipsoid, but it remains fairly manageable by using the theory of quadrics (Dufour and Taamouti 2000). When the model is correct and its parameters are well identified by the instruments used, \(C_\beta(\alpha)\) is a closed bounded set close to an ellipsoid. In other cases, it can be unbounded or empty. Unbounded sets are highly likely when the model is not identified, so they point to lack of identification. Empty confidence sets can occur (with a non-zero probability) when we have more instruments than parameters in the structural equation (41), that is, the model is over-identified. An empty confidence set means that no value of the parameter vector \(\beta\) is judged to be compatible with the available data, which indicates that the model is misspecified. So the procedure provides as an interesting byproduct a specification test. (For further discussion of this point, see Kleibergen 2002b.)

It is also easy to see that the above procedure remains valid even if the extended reduced form (46) is the correct model for \(Y\). In other words, we can leave out a subset of the instruments \((X_3)\) and use only \(X_2\); the level of the procedure will not be affected. Indeed, this will also hold if \(Y\) is determined by the general – possibly non-linear – model (47). The procedure is robust to excluded instruments as well as to the specification of the model for \(Y\). The power of the test may be affected by the choice of \(X_2\), but its level is not. Since it is quite rare that an investigator can be sure relevant instruments have not been left out, this is an important practical consideration.

The AR procedure can be extended easily to deal with linear hypotheses that involve \(\gamma\) as well. For example, to test a hypothesis of the form

\[
H_0(\beta_0, \gamma_0) : \beta = \beta_0 \text{ and } \gamma = \gamma_0,
\]  

(64)

we can consider the transformed model

\[
y - Y\beta_0 - X_1\gamma_0 = X_1\theta_1 + X_2\theta_2 + \varepsilon.
\]  

(65)

Since, under \(H_0(\beta_0, \gamma_0)\),

\[
y - Y\beta_0 - X_1\gamma_0 = \varepsilon,
\]  

(66)
we can test \( H_0(\beta_0, \gamma_0) \) by testing \( H_0'(\beta_0, \gamma_0) : \theta_1 = 0 \) and \( \theta_2 = 0 \) in the auxiliary regression (65); see Maddala (1974). Tests for more general restrictions of the form

\[
H_0(\beta_0, \nu_0) : \beta = \beta_0 \text{ and } R\gamma = \nu_0, \tag{67}
\]

where \( R \) is an \( r \times K \) fixed full-rank matrix, are discussed in Dufour and Jasiak (2001).

The AR procedure thus enjoys several remarkable features. Namely, it is (1) pivotal in finite samples; (2) robust to weak instruments; (3) robust to excluded instruments; (4) robust to the specification of the model for \( Y \) (which can be non-linear with an unknown form); further, (5) the AR method provides asymptotically ‘valid’ tests and confidence sets under quite weak distributional assumptions (basically, the assumptions that cover the usual asymptotic properties of linear regression); and (6) it can be extended easily to test restrictions and build confidence sets that also involve the coefficients of the exogenous variables, such as \( H_0(\beta_0, \nu_0) \) in (67).

But the method also has its drawbacks. The main ones are (1) the tests and confidence sets obtained in this way apply only to the full vector \( \beta \) [or \( (\beta', \gamma')' \)]; what can we do, if \( \beta \) has more than one element? (2) power may be low if too many instruments are added (\( X_2 \) has too many variables) to perform the test, especially if the instruments are irrelevant; (3) the error normality assumption is restrictive and we may wish to consider other distributional assumptions; (4) the structural equations are assumed to be linear. We will now discuss a number of methods that have been proposed to circumvent these drawbacks.

6.2. Projections and inference on parameter subsets

Suppose, now, that \( \beta \) [or \( (\beta', \gamma')' \)] has more than one component. The fact that a procedure with a finite-sample theory has been obtained for ‘joint hypotheses’ of the form \( H_0(\beta_0) \) (or \( H_0(\beta_0, \gamma_0) \)) is not due to chance: since the distribution of the data is determined by the full parameter vector, there is no reason in general why one should be able to decide on the value of a component of \( \beta \) independently of the others. Such a separation is feasible only in special situations, for example, in the classical linear model (without exact multicollinearity). Lack of identification is precisely a situation where the value of a parameter may be determined only after various restrictions (e.g., the values of other parameters) have been imposed. So parametric non-separability arises here, and inference should start from a simultaneous approach. If the data-generating process corresponds to a model where parameters are well identified, precise inferences on individual coefficients may be feasible. This raises the question of how one can move from a joint inference on \( \beta \) to its components.
A general approach to this problem consists in using a projection technique. If
\[ P[\beta \in C_\beta(\alpha)] \geq 1 - \alpha, \]  
then, for any function \( g(\beta) \),
\[ P[g(\beta) \in g[C_\beta(\alpha)]] \geq 1 - \alpha. \]  
If \( g(\beta) \) is a component of \( \beta \) or (more generally) a linear transformation \( g(\beta) = w/\beta \), the confidence set for a linear combination of the parameters, say \( w/\beta \) takes the usual form \( [w/\beta - \hat{\sigma}_\alpha, w/\beta + \hat{\sigma}_\alpha] \), with \( \hat{\beta} \) a k-class type estimator of \( \beta \); see Dufour and Taamouti (2000).9

Another interesting feature comes from the fact that the confidence sets obtained in this way are simultaneous in the sense of Scheffé. More precisely, if \( \{g_a(\beta) : a \in A\} \) is a set of functions of \( \beta \), then
\[ P[g_a(\beta) \in g[C_\beta(\alpha)] \text{ for all } a \in A] \geq 1 - \alpha. \]  
If these confidence intervals are used to test different hypotheses, an unlimited number of hypotheses can be tested without losing control of the overall level.

6.3. Alternatives to the AR procedure
With a view to improving the power of AR procedures, a number of alternative methods have been recently suggested. We will now discuss several of them.

6.3.1. Generalized auxiliary regression
A general approach to the problem of testing \( H_0(\beta_0) \) consists in replacing \( X_2 \) in the auxiliary regression
\[ y - Y\beta_0 = X_1\theta_1 + X_2\theta_2 + \varepsilon \]  
with an alternative set of auxiliary instruments, say, \( Z \) of dimension \( T \times k_2 \). In other words, we consider the generalized auxiliary regression,
\[ y - Y\beta_0 = X_1\theta_1 + Z\tilde{\theta}_2 + \varepsilon, \]  
where \( \tilde{\theta}_2 = 0 \) under \( H_0(\beta_0) \). So we can test \( H_0(\beta_0) \) by testing \( \tilde{\theta}_2 = 0 \) in (72). Then the problem consists in selecting \( Z \) so that the level can be controlled and power may be improved with respect to the AR auxiliary regression (71). For

---

9 \( g[C_\beta(\alpha)] \) takes the form of a bounded confidence interval as soon as the confidence set \( g[C_\beta(\alpha)] \) is unbounded. For further discussion of projection methods, the reader may consult Dufour (1990, 1997), Campbell and Dufour (1997), Abdelkhalek and Dufour (1998), Dufour, Hallin, and Mizera (1998), Dufour and Kiviet (1998), and Dufour and Jasiak (2001).
example, it is easy to see that the power of the AR test could become low if a large set of auxiliary instruments is used, especially if the latter are weak. So several alternative procedures can be generated by reducing the number of auxiliary instruments (the number of columns in $Z$).

At the outset, we should note that, if (42) were the correct model and $\Pi = [\Pi_1, \Pi_2]$ were known, then an optimal choice from the viewpoint of power consists in choosing $Z = X_2 \Pi_2$; see Dufour and Taamouti (2001b). The practical problem, of course, is that $\Pi_2$ is unknown. This suggests that we replace $X_2 \Pi_2$ with an estimate, such as

$$Z = X_2 \hat{\Pi}_2,$$

where $\hat{\Pi}_2$ is an estimate of the reduced-form coefficient $\Pi_2$ in (42). The problem then consists in choosing $\hat{\Pi}$. For that purpose, it is tempting to use the least squares estimator $\hat{\Pi} = (X'X)^{-1}X'Y$. However, $\hat{\Pi}$ and $\varepsilon$ are not independent, and we continue to face a simultaneity problem with messy distributional consequences. Ideally, we would like to select an estimate $\hat{\Pi}_2$ which is independent of $\varepsilon$.

6.3.2. Split-sample optimal auxiliary instruments
If we can assume that the error vectors $(u_t, V_t)'$, $t = 1, \ldots, T$, are independent, this approach to estimating $\Pi$ may be feasible if a split-sample technique is used: a fraction of the sample is used to obtain $\hat{\Pi}$ and the rest to estimate the auxiliary regression (72) with $Z = X_2 \hat{\Pi}_2$. Under such circumstances, by conditioning on $\hat{\Pi}$, we can easily see that the standard $F$ test for $\theta_2 = 0$ is valid. Further, this procedure is robust to weak instruments, excluded instruments as well as the specification of the model for $Y$ (that is, under the general assumptions (46) and (47)], as long as the independence between $\hat{\Pi}$ and $\varepsilon$ can be maintained. Of course, using a split-sample may involve a loss of the effective number of observations, and there will be a trade-off between the efficiency gain from using a smaller number of auxiliary instruments and the observations that are ‘sacrificed’ to get $\hat{\Pi}$. Better results tend to be obtained by using a relatively small fraction of the sample to obtain $\hat{\Pi}$ – 10% for example – and the rest for the main equation. For further details on this procedure, the reader may consult Dufour and Jasiak (2001) and Kleibergen (2002a).\(^{10}\)

A number of alternative procedures can be cast in the framework of equation (72).

\(^{10}\) Split-sample techniques often lead to important distributional simplifications; for further discussion of this type of method, see Angrist and Krueger (1995) and Dufour and Torrès (1998, 2000).
6.3.3. LM-type GMM-based statistic

If we take \( Z = Z_{WZ} \) with

\[
Z_{WZ} = P[M(X_1)X_2] Y = P[M(X_1)X_2] M(X_1) Y = [M(X_1)X_2] \hat{\Pi}_2
\]

(74)

\[
\hat{\Pi}_2 = [X'_2M(X_1)X_2]^{-1} X'_2M(X_1) Y,
\]

(75)

the F-statistic (say, \( F_{GMM}(\beta_0) \)) for \( \bar{\theta}_2 = 0 \) is a monotonic transformation of the LM-type statistic \( LMGMM(\beta_0) \) proposed by Wang and Zivot (1998). Namely,

\[
F_{GMM}(\beta_0) = \left( \frac{T - k_1 - G}{GT} \right) \frac{LMGMM(\beta_0)}{1 - (1/T)LMGMM(\beta_0)},
\]

(76)

where

\[
LMGMM(\beta_0) = \frac{(y - Y\beta_0)' P[Z_{WZ}] (y - Y\beta_0)}{(y - Y\beta_0)' M(X_1)(y - Y\beta_0)/T}.
\]

(77)

Note that \( \hat{\Pi}_2 \) above is the ordinary least squares (OLS) estimator of \( \Pi_2 \) from the multivariate regression (42), so that \( F_{GMM}(\beta_0) \) can be obtained by computing the F-statistic for \( \bar{\theta}_2^\ast = 0 \) in the regression

\[
y - Y\beta_0 = X_1\theta_1^\ast + (X_2\hat{\Pi}_2)\theta_2^\ast + u.
\]

(78)

When \( k_2 \geq G \), the statistic \( F_{GMM}(\beta_0) \) can also be obtained by testing \( \bar{\theta}_2^{**} = 0 \) in the auxiliary regression

\[
y - Y\beta_0 = X_1\theta_1^{**} + \hat{Y}\theta_2^{**} + u,
\]

(79)

where \( \hat{Y} = X\hat{\Pi} \). It is also interesting to note that the OLS estimates of \( \theta_1^{**} \) and \( \theta_2^{**} \), obtained by fitting the latter equation, are identical to the 2SLS estimates of \( \theta_1^{**} \) and \( \theta_2^{**} \) in the equation

\[
y - Y\beta_0 = X_1\theta_1^{**} + Y\theta_2^{**} + u.
\]

(80)

The \( LMGMM(\beta_0) \) test may thus be interpreted as an approximation to the optimal test based on replacing the optimal auxiliary instrument \( X_2\Pi_2 \) by \( X_2\hat{\Pi}_2 \). The statistic \( LMGMM(\beta_0) \) is also numerically identical to the corresponding LR-type and Wald-type tests, based on the same GMM estimator (in this case, the 2SLS estimator of \( \beta \)).
As mentioned above, the distribution of this statistic will be affected by the fact that $X_2\hat{\Pi}_2$ and $u$ are not independent. In particular, it is influenced by the presence of weak instruments. But Wang and Zivot (1998) showed that the distribution of $LM_{GMM}(\beta_0)$ is bounded by the $\chi^2(k_2)$ asymptotically. When $k_2 = G$ (usually deemed the ‘just-identified’ case, although the model may be under-identified in that case), we see easily (from (78)) that $FGMM(\beta_0)$ is (almost surely) identical with the AR statistic, that is,

$$F_{GMM}(\beta_0) = AR(\beta_0) \text{ if } k_2 = G,$$

so that $FGMM(\beta_0)$ follows an exact $F(G, T - k)$ distribution, while for $k_2 > G$,

$$G \cdot FGMM(\beta_0) \leq \left( \frac{T - k_1 - G}{T - k_1 - k_2} \right) k_2 \cdot AR(\beta_0),$$

so that the distribution of $LM_{GMM}(\beta_0)$ can be bounded in finite samples by the distribution of a monotonic transformation of a $F(k_2, T - k)$ variable (which, for $T$ large, is very close to the $\chi^2(k_2)$ distribution). But, for $T$ reasonably large, $AR(\beta_0)$ will always reject when $FGMM(\beta_0)$ rejects (at a given level), so the power of the AR test is uniformly superior to that of the $LM_{GMM}$ bound test.$^{11}$

### 6.3.4. Kleibergen’s $K$ Test

If we take $Z = Z_K$ with

$$Z_K = P(X) \left[ Y - (y - Y\beta_0) \frac{s_{e\epsilon}V(\beta_0)}{s_{\epsilon\epsilon}(\beta_0)} \right] = X\hat{\Pi}(\beta_0) \equiv \hat{Y}(\beta_0)$$

$$\hat{\Pi}(\beta_0) = \hat{\Pi} - \hat{\pi}(\beta_0) \frac{s_{e\epsilon}V(\beta_0)}{s_{\epsilon\epsilon}(\beta_0)}, \quad \hat{\Pi} = (X'X)^{-1}X'Y$$

$$\hat{\pi}(\beta_0) = (X'X)^{-1}X'(y - Y\beta_0), \quad s_{e\epsilon}V(\beta_0) = \frac{1}{T - k} (y - Y\beta_0)'M(X)Y$$

$$s_{\epsilon\epsilon}(\beta_0) = \frac{(y - Y\beta_0)'M(X)(y - Y\beta_0)}{T - k},$$

$^{11}$ The $\chi^2(k_2)$ bound also follows in a straightforward way from (82). Note that Wang and Zivot (1998) do not provide the auxiliary regression interpretation (78)–(79) of their statistics. For details, see Dufour and Taamouti (2001b).
we obtain a statistic, which reduces to the one proposed by Kleibergen (2002a) for $k_1 = 0$. More precisely, with $k_1 = 0$, the F-statistic $F_K(\beta_0)$ for $\theta_2 = 0$ is equal to Kleibergen’s statistic $K(\beta_0)$ divided by $G$:

$$F_K(\beta_0) = K(\beta_0)/G.$$  \hfill (87)

This procedure tries to correct the simultaneity problem associated with the use of $\hat{Y}$ in the $L_{GMM}$ statistic by “purging” it from its correlation with $u$ (by subtracting the term $\hat{\pi}(\beta_0)_{s_{eY}}(\beta_0)/s_{eY}(\beta_0)$ in $Z_K$). In other words, $F_K(\beta_0)$ and $K(\beta_0) = G F_K(\beta_0)$ can be obtained by testing $\theta_2 = 0$ in the regression

$$y - Y \beta_0 = X_1 \theta_1 + \hat{Y}(\beta_0) \theta_2 + u,$$  \hfill (88)

where the fitted values $\hat{Y}$, which appear in the auxiliary regression (79) for the $L_{GMM}$ test, have been replaced by $\hat{Y}(\beta_0) = \hat{Y} - X \hat{\pi}(\beta_0)_{s_{eY}}(\beta_0)/s_{eY}(\beta_0)$, which are closer to being orthogonal with $u$.

If $k_2 = G$, we have $F_K(\beta_0) = AR(\beta_0) \sim F(G, T - k)$, while in the other cases ($k_2 \geq G$), we can see easily that the bound for $F_{G_{GMM}}(\beta_0)$ in (82) also applies to $F_K(\beta_0)$:

$$G F_K(\beta_0) \leq \left( \frac{T - k_1 - G}{T - k_1 - k_2} \right) k_2 AR(\beta_0).$$  \hfill (89)

Kleibergen (2002a) did not supply a finite-sample distributional theory but showed (assuming $k_1 = 0$) that $K(\beta_0)$ follows a $\chi^2(G)$ distribution asymptotically under $H_0(\beta_0)$, irrespective of the presence of weak instruments. This entails that the $K(\beta_0)$ test will have power higher than the one of $L_{GMM}$ test (based on the $\chi^2(k_2)$ bound), at least in the neighbourhood of the null hypothesis, although not necessarily far away from the null hypothesis.

It is also interesting to note that the inequality (89) indicates that the distribution of $K(\beta_0) = G F_K(\beta_0)$ can be bounded in finite samples by a $[k_2(T - k_1 - G)/(T - k)]F(k_2, T - k)$ distribution. However, because of the stochastic dominance of $AR(\beta_0)$, there would be no advantage in using the bound to get critical values for $K(\beta_0)$, since the AR test would then have better power.

In view of the fact that the above procedure is based on estimating the mean of $\chi^2$ (using $\chi^2$) and the covariances between the errors in the reduced form for $Y$ and $u$ (using $s_{eY}(\beta_0)$), it can become quite unreliable in the presence of excluded instruments.
6.3.5. Likelihood ratio test

The likelihood ratio (LR) statistic for \( \beta = \beta_0 \) was also studied by Wang and Zivot (1998). The LR test statistic in this case takes the form:

\[
LRLIML = T \left[ \ln \left( \kappa(\beta_0) \right) - \ln \left( \kappa(\hat{\beta}_{LIML}) \right) \right],
\]  

(90)

where \( \hat{\beta}_{LIML} \) is the limited information maximum likelihood estimator (LIML) of \( \beta \) and

\[
\kappa(\beta) = \frac{(y - Y\beta)' M(X_1) (y - Y\beta)}{(y - Y\beta)' M(X) (y - Y\beta)}. 
\]  

(91)

Like \( LMGMM \), the distribution of \( LRLIML \) depends on unknown nuisance parameters under \( H_0(\beta_0) \), but its asymptotic distribution is \( \chi^2(k_2) \) when \( k_2 = G \) and is bounded by the \( \chi^2(k_2) \) distribution in the other cases (a result in accordance with the general LR distributional bound given in Dufour 1997).

This bound can also be easily derived from the following inequality:

\[
LRLIML \leq \left( \frac{T}{T - k} \right) k_2 AR(\beta_0),
\]  

(92)

so that the distribution of \( LRLIML \) is bounded in finite samples by the distribution of a \( [Tk_2/(T - k)]F(k_2, T - k) \) variable; for details, see Dufour and Khalaf (2000). For \( T \) reasonably large, this entails that the \( AR(\beta_0) \) test will have power higher than the one of \( LRLIML \) test [based on the \( \chi^2(k_2) \) bound], at least in the neighbourhood of the null hypothesis. The power the AR test is uniformly superior to the one of the \( LRLIML \) bound test. Because the LR test depends heavily on the specification of the model for \( Y \), it is not robust to excluded instruments.

6.3.6. Conditional tests

A promising approach was recently proposed by Moreira (2003a). His suggestion consists in conditioning upon an appropriately selected portion of the sufficient statistics for a gaussian SEM. Assuming that the covariance matrix of the errors is known, the corresponding conditional distribution of various test statistics for \( H_0(\beta_0) \) does not involve nuisance parameters. The conditional distribution is typically not standard but may be established by simulation. Such an approach may lead to power gains. On the other hand, the assumption that error covariances are known is rather implausible, and the extension of the method to the case where the error covariance matrix is unknown is obtained at the expense of using a large-sample approximation. Like Kleibergen’s procedure, this method yields an asymptotically similar test. For further discussion, see Moreira and Poi (2001) and Moreira (2003b).
6.3.7. Instrument selection procedures

Systematic search methods for identifying relevant instruments and excluding unimportant instruments have been discussed by several authors (Hall, Rudebusch, and Wilcox 1996; Hall and Peixe 2000; Dufour and Taamouti 2001a; Donald and Newey 2001). In this set-up, the power of AR-type tests depends on a function of model parameters called the concentration coefficient. One way to approach instrument selection is to maximize the concentration coefficient towards maximizing test power. Robustness to instrument exclusion is very handy in this context. For further discussion, the reader may consult Dufour and Taamouti (2001a).

To summarize, in special situations, alternatives to the AR procedure may allow some power gains with respect to the AR test with an unreduced set of instruments. They themselves may have some important drawbacks. In particular, (1) only an asymptotic distributional theory is supplied; (2) the statistics used are not pivotal in finite samples, although Kleibergen’s and Moreira’s statistics are asymptotically pivotal; (3) they are not robust to instrument exclusion or to the formulation of the model for the explanatory endogenous variables. It is also of interest to note that finite-sample versions of several of these asymptotic tests may be obtained by using split-sample methods.

All the problems and techniques discussed above relate to sampling-based statistical methods. SEM can also be analysed through a Bayesian approach, which may alleviate the indeterminacies associated with identification via the introduction of a prior distribution on the parameter space. Bayesian inferences always depend on the choice of prior distribution (a property viewed as undesirable in the sampling approach), but this dependence becomes especially strong when identification problems are present (see Gleser and Hwang 1987). In this paper we aim only at discussing problems and solutions that arise within the sampling framework, and it is beyond its scope to debate the advantages and disadvantages of Bayesian methods under weak identification. For additional discussion on this issue, see Kleibergen and Zivot (2003) and Sims (2001).

7. Extensions

We will discuss succinctly some extensions of the above results to multivariate set-ups (where several structural equations may be involved), models with non-Gaussian errors, and non-linear models.

7.1. Multivariate regression, simulation-based inference, and non-normal errors

Another approach to inference on a linear structural equation model is based on observing that the structural model (41)–(44) can be put in the form of a multivariate linear regression (MLR):

\[ \bar{Y} = XB + U, \]  

(93)
where $\hat{Y} = [y, Y], B = \begin{bmatrix} \pi, \Pi \end{bmatrix}, U = \begin{bmatrix} u, V \end{bmatrix} = [\hat{U}_1, \ldots, \hat{U}_T]', \pi = [\pi_1', \pi_2'], \Pi = [\Pi_1', \Pi_2'], \pi_1 = \Pi_1\beta + \gamma$ and $\pi_2 = \Pi_2\beta$. (Most of this section is based on Dufour and Khalaf 2001 and Dufour 2002) This model is linear except for the nonlinear restriction $\pi_2 = \Pi_2\beta$. Let us now make the assumption that the errors in the different equations for each observation, $\hat{U}_t$, satisfy the property:

$$\hat{U}_t = JW_t, \ t = 1, \ldots, T,$$

(94)

where the vector $w = \text{vec}(W_1, \ldots, W_n)$ has a known distribution and $J$ is an unknown nonsingular matrix (which enters into the covariance matrix $\Sigma$ of the error vectors $\hat{U}_t$). This distributional assumption is, in a way, more restrictive than the one made in section 5.1 – because of the assumption on $V$ – and in another way, less restrictive, because the distribution of $u$ is not taken to be necessarily $N[0, \sigma_u^2 I_T]$.

Consider, now, a hypothesis of the form

$$H_0: \text{RBC} = D,$$

(95)

where $R$, $C$, and $D$ are fixed matrices. This is called a uniform linear (UL) hypothesis; for example, the hypothesis $\beta = \beta_0$ tested by the AR test can be written in this form (see Dufour and Khalaf 2000). The corresponding gaussian LR statistic is

$$LR(H_0) = T \ln(|\hat{\Sigma}_0|/|\hat{\Sigma}|),$$

(96)

where $\hat{\Sigma} = \hat{U}'\hat{U}/T$ and $\hat{\Sigma}_0 = \hat{U}_0'\hat{U}_0/T$ are, respectively, the unrestricted and restricted estimates of the error covariance matrix. The AR test can also be obtained as a monotonic transformation of a statistic of the form $LR(H_0)$. An important feature of $LR(H_0)$ in this case is that its distribution under $H_0$ does not involve nuisance parameters and may be easily simulated (it is a pivot); see Dufour and Khalaf (2002). In particular, its distribution is completely invariant to the unknown $J$ matrix (or the error covariance matrix). In such a case, even though this distribution may be complicated, we can use Monte Carlo test techniques – a form of parametric bootstrap – to obtain exact test procedures. (For further discussion see Dufour and Khalaf 2001; Dufour 2002.) Multivariate extensions of AR tests, which impose restrictions on several structural equations, can be obtained in this way. Further, this approach allows one to consider any (possibly non-gaussian) distribution on $w$.

More generally, it is of interest to note that the LR statistic for about any hypothesis on $B$ can be bounded by a LR statistic for an appropriately selected UL hypothesis: setting $b = \text{vec}(B)$ and

$$H_0: Rb \in \Delta_0,$$

(97)
where $R$ an arbitrary $q \times k(G + 1)$ matrix and $\Delta_0$ is an arbitrary subset of $\mathbb{R}^q$, the distribution of the corresponding LR statistic can be bounded by the LR statistic for a UL hypothesis (which is pivotal). This covers as special cases all restrictions on the coefficients of SEM (as long as they are written in the MLR form). To avoid the use of such bounds (which may be quite conservative), it is also possible to use maximized Monte Carlo tests (Dufour 2002).

All the above procedures are valid for parametric models that specify the error distribution up to an unknown linear transformation (the $J$ matrix), which allows an unknown covariance matrix. It is easy to see that these (including the exact procedures discussed in section 6) yield ‘asymptotically valid’ procedures under much weaker assumptions than those used to obtain finite-sample results. However, in view of the discussion in section 4, the pitfalls and limitations of such arguments should be remembered: there is no substitute for a provably exact procedure.

If we aim at getting tests and confidence sets for non-parametric versions of the SEM (where the error distribution involves an infinite set of nuisance parameters), this may be achievable by looking at distribution-free procedures based on permutations, ranks or signs. There is very little work on this topic in the SEM. For an interesting first look, however, the reader should look at a recent paper by Bekker (2002).

7.2. Non-linear models

It is relatively difficult to characterize identification and study its consequences in non-linear structural models. But problems similar to those noted for linear SEM do arise. Non-linear structural models typically take the form:

$$f_t(y_t, x_t, \theta) = u_t, \quad E_{\theta}[u_t | Z_t] = 0, \quad t, \ldots, T, \tag{98}$$

where $f_t(\cdot)$ is a vector of possibly non-linear relationships, $y_t$ is a vector endogenous variables, $x_t$ is a vector of exogenous variables, $\theta$ is vector of unknown parameters, $Z_t$ is a vector of conditioning variables (or instruments) – usually with a number of additional distributional assumptions – and $E_{\theta}[\cdot]$ refers to the expected value under a distribution with parameter value $\theta$. In such models, $\theta$ can be viewed as identifiable if there is only one value of $\theta$ (say, $\theta = \bar{\theta}$) that satisfies (98), and we have weak identification (or weak instruments) when the expected values $E_{\theta}[f_t(y_t, x_t, \theta) | Z_t] = 0, t, \ldots, T$, are ‘weakly sensitive’ to the value of $\theta$.

Work on weak identification in non-linear models remains scarce. Non-linearity makes it difficult to construct finite-sample procedures even in models where identification difficulties do not occur. So it is not surprising that results in this area have been based mostly on large-sample approximations. Stock and Wright (2000) studied the asymptotic distributions of GMM-based estimators and test statistics under conditions of weak identification (and weak ‘high level’ asymptotic distributional assumptions). While GMM estimators of
θ have non-standard asymptotic distributions, the objective function minimized by the GMM procedure follows an asymptotic distribution that is unaffected by the presence of weak instruments: it is asymptotically pivotal. So tests and confidence sets based on the objective function can be asymptotically valid irrespective of the presence of weak instruments. These results are achieved for the full parameter vector θ, that is, for hypotheses of the form θ = θ₀ and the corresponding joint confidence sets. This is not surprising: parametric non-separability arises here for two reasons: model non-linearity and the possibility of non-identification. Of course, once a joint confidence set for θ has been built, inference on individual parameters can be drawn via projection methods. Other contributions in this area include papers by Kleibergen (2001a, 2003), who proposed an extension of the $K(\beta_0)$ test, and Wright (2003), who proposed tests of underidentification and identification.

In view of the discussion in section 4, the fact that all these methods are based on large-sample approximations without a finite-sample theory remains a concern. However, a first attempt at deriving finite-sample procedures is available in Dufour and Taamouti (2001b). Under parametric assumptions on the errors, the hypothesis $\theta = \theta_0$ is tested by testing $\gamma = 0$ in an auxiliary regression of the form:

$$f_t(y_t, x_t, \theta_0) = z_t(\theta_0, \theta_1)\gamma + \varepsilon_t, \quad t, \ldots, T,$$

where the $z_t(\theta_0, \theta_1)$ are instruments in a way that maximizes power against a reference alternative (point-optimal instruments). In this way one gets point-optimal tests (see King 1988; Dufour and King 1991). Inferences on non-linear regressions are also covered by this set-up. As in the case of linear SEM, sample-split techniques may be exploited to approximate optimal instruments, and projection methods can be used to draw inferences on subvectors of θ.

8. Conclusion

By way of conclusion, we will summarize the main points made in this paper.

1. There are basic pitfalls and limitations faced in developing inference procedures in econometrics. If we are not careful, we can easily be led into ill-defined problems and find ourselves

   a) trying to test to test a non-testable hypothesis, such as a hypothesis on a moment in the context of an insufficiently restrictive non-parametric model, or a hypothesis (e.g., a unit root hypothesis) on a dynamic model while allowing a dynamic structure with an unlimited (not necessarily infinite) number of parameters;
b) trying to solve an inference problem using a technique that cannot deliver a solution because of the very structure of the technique, as in (i) testing a hypothesis on a mean (or median) under heteroscedasticity of unknown form, via standard least-square-based ‘heteroscedasticity-robust’ standard errors, or (ii) building a confidence interval for a parameter that is not identifiable in a structural model, via the usual technique based on standard errors. In particular, this type of difficulty arises for Wald-type statistics in the presence of weak instruments (or weakly identified models)

2. In many econometric problems (e.g., inference on structural models), several of the intuitions derived from the linear regression model and standard asymptotic theory can easily be misleading.

   a) Standard errors do not constitute a valid way of assessing parameter uncertainty and building confidence intervals.

   b) In many models, such as structural models where parameters may be underidentified, individual parameters in statistical models generally are not meaningful, but parameter vectors can be so (at least in parametric models). We called this phenomenon parametric nonseparability. As a result, restrictions on individual coefficients may not be testable, while restrictions on the whole parameter vector are so. This feature should play a central role in designing for dealing with weakly identified models.

3. The above difficulties underscore the pitfalls of large-sample approximations, which typically are based on pointwise (rather than uniform) convergence results and may be arbitrarily inaccurate in finite samples.

4. Concerning solutions to such problems, and more specifically in the context of weakly identified models, we have emphasized the following points.

   a) In accordance with basic statistical theory, one should always look for pivots as the fundamental ingredient for building tests and confidence sets.

   b) Pivots generally are not available for individual parameters, but they can be obtained in a much a wider set of cases for appropriately selected vectors of parameters.

   c) Given a pivot for a parameter vector, we can construct valid tests and confidence sets for the parameter vector.

   d) Inference on individual coefficients may then be derived through projection methods.

5. In the specific example of SEM, the following general remarks are, in our view, important.

   a) Besides being pivotal, the AR statistic enjoys several remarkable robustness properties, such as robustness to the presence of weak instruments, to excluded instruments, or to the specification of a model for the endogenous explanatory variables.
b) It is possible to improve the power of AR-type procedures (especially by reducing the number of instruments), but power improvements may come at the expense of using a possibly unreliable large-sample approximation or losing robustness (such as robustness to excluded instruments). As usual, there is a trade-off between power (which typically is increased by considering more restrictive models) and robustness (which involves considering a wider hypothesis).

c) Trying to adapt and improve AR-type procedures (without ever forgetting basic statistical principles) constitutes the most promising avenue for dealing with weak instruments.

References


Chao, John, and Norman R. Swanson (2001) ‘Bias and MSE analysis of the IV estimator under weak identification with application to bias correction,’ Technical Report, Department of Economics, Rutgers University, New Brunswick, NJ


— (1990) ‘Exact tests and confidence sets in linear regressions with autocorrelated errors,’ *Econometrica* 58, 475–94
— (1997) ‘Some impossibility theorems in econometrics, with applications to structural and dynamic models,’ *Econometrica* 65, 1365–89


Dufour, Jean-Marie, and Maxwell L. King (1991) ‘Optimal invariant tests for the autocorrelation coefficient in linear regressions with stationary or nonstationary AR(1) errors,’ *Journal of Econometrics* 47, 115–43


— (2001b) ‘Point-optimal instruments and generalized Anderson-Rubin procedures for nonlinear models,’ Technical Report, CRDE, Université de Montréal


— (2002c) ‘Weak instruments: diagnosis and cures in empirical econometrics,’ Technical Report, Department of Economics, Massachusetts Institute of Technology


Kleibergen, Frank (2001a) ‘Testing parameters in GMM without assuming that they are identified,’ Technical Report TI 01-067/4, Tinbergen Institute, Amsterdam

— (2001b) ‘Testing subsets of structural coefficients in the IV regression model,’ Technical Report, Department of Quantitative Economics, University of Amsterdam


— (2002b) ‘Two independent statistics that test location and misspecification and add up to the Anderson-Rubin statistic,’ Technical Report, Department of Quantitative Economics, University of Amsterdam

— (2003) ‘Expansions of GMM statistics that indicate their properties under weak and/or many instruments,’ Technical Report, Department of Quantitative Economics, University of Amsterdam


Moreira, Marcelo J. (2001) ‘Tests with correct size when instruments can be arbitrarily weak,’ Technical Report, Department of Economics, Harvard University, Cambridge, MA

Identification, weak instruments, statistical inference 807


Nelson, C.R., and R. Startz (1990a) ‘The distribution of the instrumental variable estimator and its t-ratio when the instrument is a poor one,’ Journal of Business 63, 125–40

— (1990b) ‘Some further results on the exact small properties of the instrumental variable estimator,’ Econometrica 58, 967–76


Sims, Christopher (1971a) ‘Distributed lag estimation when the parameter space is explicitly infinite-dimensional,’ Annals of Mathematical Statistics 42, 1622–36

— (1971b) ‘Discrete approximations to continuous time distributed lags in econometrics,’ Econometrica 39, 545–63
