What's New in Econometrics **NBER, Summer 2007**

Lecture 9, Tuesday, July 31th, 3.15-4.15pm

Partial Identification

1. INTRODUCTION

Traditionally in constructing statistical or econometric models researchers look for models that are *(point-)identified*: given a large (infinite) data set, one can infer without uncertainty what the values are of the objects of interest, the estimands. Even though the fact that a model is identified does not necessarily imply that we do well in finite samples, it would appear that a model where we cannot learn the parameter values even in infinitely large samples would not be very useful. Traditionally therefore researchers have stayed away from models that are not (point-)identified, often adding assumptions beyond those that could be justified using substantive arguments. However, it turns out that even in cases where we cannot learn the value of the estimand *exactly* in large samples, in many cases we can still learn a fair amount, even in finite samples. A research agenda initiated by Manski (an early paper is Manski (1990), monographs include Manski (1995, 2003)), referred to as *partial identification*, or earlier as *bounds*, and more recently adopted by a large number of others, notably Tamer in a series papers (Haile and Tamer, 2003, Ciliberto and Tamer, 2007; Aradillas-Lopez and Tamer, 2007), has taken this perspective. In this lecture we focus primarily on a number of examples to show the richness of this approach. In addition we discuss some of the theoretical issues connected with this literature, and some practical issues in implementation of these methods.

The basic set up we adopt is one where we have a random sample of units from some population. For the typical unit, unit i, we observe the value of a vector of variables Z_i. Sometimes it is useful to think of there being in the background a latent variable variable W_i. We are interested in some functional θ of the joint distribution of Z_i and W_i, but, not observing W_i for any units, we may not be able to learn the value of θ even in infinite samples because the estimand cannot be written as a functional of the distribution of Z_i alone. The

three key questions are (i) what we can learn about θ in large samples (identification), (ii) how do we estimate this (estimation), and (iii) how do we quantify the uncertainty regarding θ (inference).

The solution to the first question will typically be a set, the *identified set*. Even if we can characterize estimators for these sets, computing them can present serious challenges. Finally, inference involves challenges concerning uniformity of the coverage rates, as well as the question whether we are interested in coverage of the entire identified set or only of the parameter of interest.

There are a number of cases of general interest. I will discuss two leading cases in more detail. In the first case the focus is on a scalar, with the identified set equal to an interval with lower and upper bound a smooth, $\sqrt{N}$-estimable functional of the data. A second case of interest is that where the information about the parameters can be characterized by moment restrictions, often arising from revealed preference comparisons between utilities at actions taken and actions not taken. I refer to this as the generalized inequality restrictions (GIR) setting. This set up is closely related to the generalized method of moments framework.

2. Partial Identification: Examples

Here we discuss a number of examples to demonstrate the richness of the partial identification approach.

2.1 Missing Data

This is a basic example, see e.g., Manski (1990), and Imbens and Manski (2004). It is substantively not very interesting, but it illustrates a lot of the basic issues. Suppose the observed variable is the pair $Z_i = (D_i, D_i \cdot Y_i)$, and the unobserved variable is $W_i = Y_i$. D_i is a binary variable. This corresponds to a missing data case. If $D_i = 1$, we observe Y_i, and if $D_i = 0$ we do not observe Y_i. We always observe the missing data indicator D_i. We assume the quantity of interest is the population mean $\theta = \mathbb{E}[Y_i]$.

In large samples we can learn $p = \mathbb{E}[D_i]$ and $\mu_1 = \mathbb{E}[Y_i|D_i = 1]$. The data contain no

information about $\mu_0 = \mathbb{E}[Y_i|D_i = 0]$. It can be useful, though not always possible, to write the estimand in terms of parameters that are point-identified and parameters that the data are not informative about. In this case we can do so:

$$\theta = p \cdot \mu_1 + (1 - p) \cdot \mu_0.$$

Since even in large samples we learn nothing about μ_0, it follows that without additional information there is no limit on the range of possible values for θ. Even if p is very close to 1, this small probability that $D_i = 0$ combined with the possibility that μ_0 is very large or very small allows for a wide range of values for θ.

Now suppose we know that the variable of interest is binary: $Y_i \in \{0, 1\}$. Then natural (not data-informed) lower and upper bounds for μ_0 are 0 and 1 respectively. This implies bounds on θ:

$$\theta \in [\theta_{\mathrm{LB}}, \theta_{\mathrm{UB}}] = [p \cdot \mu_1, p \cdot \mu_1 + (1 - p)].$$

These bounds are *sharp*, in the sense that without additional information we can not improve on them. Formally, for all values θ in $[\theta_{\mathrm{LB}}, \theta_{\mathrm{UB}}]$, we can find a joint distribution of (Y_i, W_i) that is consistent with the joint distribution of the observed data and with θ. Even if Y is not binary, but has some natural bounds, we can obtain potentially informative bounds on θ.

We can also obtain informative bounds if we modify the object of interest a little bit. Suppose we are interested in quantiles of the distribution of Y_i. To make this specific, suppose we are interested in the median of Y_i, $\theta_{0.5} = \mathrm{med}(Y_i)$. The largest possible value for the median arises if all the missing value of Y_i are large. Define $q_\tau(Y_i|D_i = d)$ to be the τ quantile of the conditional distribution of Y_i given $D_i = d$. Then the median cannot be larger than $q_{1/(2p)}(Y_i|D_i = 1)$ because even if all the missing values were large, we know that at least $p \cdot (1/(2p)) = 1/2$ of the units have a value less than or equal to $q_{1/(2p)}(Y_i|D_i = 1)$. Similarly, the smallest possible value for the median corresponds to the case where all the

missing values are small, leading to a lower bound of $q_{(2p-1)/(2p)}(Y_i|D_i = 1)$. Then, if $p > 1/2$, we can infer that the median must satisfy

$$\theta_{0.5} \in [\theta_{\text{LB}}, \theta_{\text{UB}}] = \left[q_{(2p-1)/(2p)}(Y_i|D_i = 1), q_{1/(2p)}(Y_i|D_i = 1)\right],$$

and we end up with a well defined, and, depending on the data, more or less informative identified interval for the median. If fewer than 50% of the values are observed, or $p < 1/2$, then we cannot learn anything about the median of Y_i without additional information (for example, a bound on the values of Y_i), and the interval is $(-\infty, \infty)$. More generally, we can obtain bounds on the τ quantile of the distribution of Y_i, equal to

$$\theta_\tau \in [\theta_{\text{LB}}, \theta_{\text{UB}}] = \left[q_{(\tau-(1-p))/p}(Y_i|D_i = 1), q_{\tau/p}(Y_i|D_i = 1)\right].$$

which is bounded if the probability of Y_i being missing is less than $\min(\tau, 1 - \tau)$.

2.2 Returns to Schooling

Manski and Pepper (2000, MP) are interested in estimating returns to schooling. They start with an individual level response function $Y_i(w)$, where $w \in \{0, 1, \dots, 20\}$ is years of schooling. Let

$$\Delta(s, t) = \mathbb{E}[Y_i(t) - Y_i(s)],$$

be the difference in average outcomes (log earnings) given t rather than s years of schooling. Values of $\Delta(s, t)$ at different combinations of (s, t) are the object of interest. Let W_i be the actual years of school, and $Y_i = Y_i(W_i)$ be the actual log earnings. If one makes an unconfoundedness type assumption that

$$Y_i(w) \perp\!\!\!\perp W_i \,\Big|\, X_i,$$

for some set of covariates, one can estimate $\Delta(s, t)$ consistently given some support conditions. MP relax this assumption. Dropping this assumption entirely without additional

assumptions one can derive the bounds using the missing data results in the previous section. In this case most of the data would be missing, and the bounds would be wide. More interestingly MP focus on a number of alternative, weaker assumptions, that do not allow for point-identification of $\Delta(s, t)$, but that nevertheless may be able to narrow the range of values consistent with the data to an informative set. One of their assumptions requires that increasing education does not lower earnings:

Assumption 1 (MONOTONE TREATMENT RESPONSE)
If $w' \geq w$, then $Y_i(w') \geq Y_i(w)$.

Another assumption states that, on average, individuals who choose higher levels of education would have higher earnings at each level of education than individuals who choose lower levels of education.

Assumption 2 (MONOTONE TREATMENT SELECTION)
If $w'' \geq w'$, then for all w, $\mathbb{E}[Y_i(w)|W_i = w''] \geq \mathbb{E}[Y_i(w)|W_i = w']$.

Both assumptions are consistent with many models of human capital accumulation. They also address the main concern with the exogenous schooling assumption, namely that higher ability individuals who would have had higher earnings in the absence of more schooling, are more likely to acquire more schooling.

Under these two assumptions, the bound on the average outcome given w years of schooling is

$$\mathbb{E}[Y_i|W_i = w] \cdot \Pr(W_i \geq w) + \sum_{v<w} \mathbb{E}[Y_i|W_i = v] \cdot \Pr(W_i = v)$$

$$\leq \mathbb{E}[Y_i(w)] \leq$$

$$\mathbb{E}[Y_i|W_i = w] \cdot \Pr(W_i \leq w) + \sum_{v>w} \mathbb{E}[Y_i|W_i = v] \cdot \Pr(W_i = v).$$

Using data from the National Longitudinal Study of Youth MP a point estimator for the upper bound on the the returns to four years of college, $\Delta(12, 16)$ to be 0.397, with a 0.95 upper quantile of 0.450. Translated into an average yearl returns this gives us 0.099, which is in fact lower than some estimates that have been reported in the literature. This analysis suggests that the upper bound is in this case reasonably informative, given a remarkably weaker set of assumptions.

2.3 CHANGES IN INEQUALITY AND SELECTION

There is a large literature on the changes in the wage distribution and the role of changes in the returns to skills that drive these changes. One concern is that if one compares the wage distribution at two points in time, any differences may be partly or wholly due to differences in the composition of the workforce. Blundell, Gosling, Ichimura, and Meghir (2007, BGHM) investigate this using bounds. They study changes in the wage distribution in the United Kingdom for both men and women. Even for men at prime employment ages employment in the late nineties is less than 0.90, down from 0.95 in the late seventies. The concern is that the 10% who do not work are potentially different, both from those who work, as well as from those who did not work in the seventies, corrupting comparisons between the wage distributions in both years. Traditionally such concerns may have been ignored by implicitly assuming that the wages for those not working are similar to those who are working, possibly conditional on some observed covariates, or they may have been addressed by using selection models. The type of selection models used ranges from very parametric models of the type originally developed by Heckman (1978), to semi- and non-parametric versions of this (Heckman, 1990). The concern that BGHM raise is that those selection models rely on assumptions that are difficult to motivate by economic theory. They investigate what can be learned about the changes in the wage distributions without the final, most controversial assumptions of those selection models.

BGHM focus on the interquartile range as their measure of dispersion in the wage distribution. As discussed in Section 2.1, this is convenient, because bounds on quantiles often exist in the presence of missing data. Let $F_{Y|X}(y|x)$ be the distribution of wages condi-

tional on some characteristics X. This is assumed to be well defined irrespective of whether an individual works or not. However, if an individual does not work, Y_i is not observed. Let D_i be an indicator for employment. Then we can estimate the conditional wage distribution given employment, $F_{Y|X,D}(y|x, d = 1)$, as well as the probability of employment, $p(x) = \text{pr}(D_i = 1|X_i = x)$. This gives us tight bounds on the (unconditional on employment) wage distribution

$$F_{Y|X,D}(y|x, d = 1) \cdot p(x) \le F_{Y|X,D}(y|x, d = 1) \le F_{Y|X,D}(y|x, d = 1) \cdot p(x) + (1 - p(x)).$$

We can convert this to bounds on the τ quantile of the conditional distribution of Y_i given $X_i = x$, denoted by $q_\tau(x)$:

$$q_{(\tau - (1 - p(x)))/p(x)}(Y_i|D_i = 1) \le q_\tau(x) \le q_{\tau/p(x)}(Y_i|D_i = 1),$$

Then this can be used to derive bounds on the interquartile range $q_{0.75}(x) - q_{0.25}(x)$:

$$q_{(0.75 - (1 - p(x)))/p(x)}(Y_i|D_i = 1) - q_{0.25/p(x)}(Y_i|D_i = 1)$$

$$\le q_{0.75}(x) - q_{0.25}(x) \le$$

$$q_{(0.25 - (1 - p(x)))/p(x)}(Y_i|D_i = 1) - q_{0.75/p(x)}(Y_i|D_i = 1).$$

So far this is just an application of the missing data bounds derived in the previous section. What makes this more interesting is the use of additional information short of imposing a full selection model that would point identify the interquartile range. The first assumption BGHM add is that of stochastic dominance of the wage distribution for employed individuals:

$$F_{Y|X,D}(y|x, d = 1) \le F_{Y|X,D}(y|x, d = 0).$$

One can argue with this stochastic dominance assumption, but within groups homogenous in background characteristics including education, it may be reasonable. This assumption tightens the bounds on the distribution function to:

$$F_{Y|X,D}(y|x, d = 1) \leq F_{Y|X,D}(y|x, d = 1) \leq$$

$$F_{Y|X,D}(y|x, d = 1) \cdot p(x) + (1 - p(x)).$$

Another assumption BGHM consider is a modification of an instrumental variables assumption that an observed covariate Z is excluded from the wage distribution:

$$F_{Y|X,Z}(y|X = x, Z = z) = F_{Y|X,Z}(y|X = x, Z = z'), \qquad \text{for all } x, z, z'.$$

This changes the bounds on the distribution function to:

$$\max_z F_{Y|X,Z,D}(y|x, z, d = 1) \cdot p(x, z)$$

$$\leq F_{Y|X,D}(y|x) \leq$$

$$\min_z F_{Y|X,Z,D}(y|x, z, d = 1) \cdot p(x) + (1 - p(x)).$$

(An alternative weakening of the standard instrumental variables assumption is in Hotz, Mullin and Sanders (1997), where a valid instrument exists, but is not observed directly.)

Such an instrument may be difficult to find, and BGHM argue that it may be easier to find a covariate that affects the wage distribution in one direction, using a monotone instrumental variables restriction suggested by Manski and Pepper (2000):

$$F_{Y|X,Z}(y|X = x, Z = z) \leq F_{Y|X,Z}(y|X = x, Z = z'), \qquad \text{for all } x, z < z'.$$

This discussion is somewhat typical of what is done in empirical work in this area. A number of assumptions are considered, with the implications for the bounds investigated. The results lay out part of the mapping between the assumptions and the bounds.

2.4 RANDOM EFFECTS PANEL DATA MODELS WITH INITIAL CONDITION PROBLEMS

Honoré and Tamer (2006) study dynamic random effects panel data models. We observe $(X_{i1}, Y_{i1}, \ldots, X_{iT}, Y_{iT})$, for $i = 1, \ldots, N$. The time dimension T is small relative to the cross-section dimension N. Large sample approximations are based on fixed T and large N. Inference would be standard if we specified a parametric model for the (components of the) conditional distribution of $(Y_{i1}, \ldots, Y_{iT})$ given $(X_{i1}, \ldots, X_{iT})$. In that case we could use maximum likelihood methods. However, it is difficult to specify this conditional distribution directly. Often we start with a model for the evolution of Y_{it} in terms of the present and past covariates and its lags. As an example, consider the model

$$Y_{it} = 1\{X'_{it}\beta + Y_{it-1}\gamma + \alpha_i + \epsilon_{it} \geq 0\},$$

with the ϵ_{it} independent over time and individuals, and normally distributed, $\epsilon_{it} \sim \mathcal{N}(0,1)$. The object of interest is the parameter governing the dynamics, γ. This model gives us the conditional distribution of $Y_{i2}, \ldots, Y_{iT}$ given Y_{i1}, α_i and given $X_{i1}, \ldots, X_{iT}$. Suppose we also postulate a parametric model for the random effects α_i:

$$\alpha|X_{i1}, \ldots, X_{iT} \sim G(\alpha|\theta),$$

(so in this case α_i is independent of the covariates). Then the model is (almost) complete, in the sense that we can almost write down the conditional distribution of $(Y_{i1}, \ldots, Y_{iT})$ given $(X_{i1}, \ldots, X_{iT})$. All that is missing is the conditional distribution of the initial condition:

$$p(Y_{i1}|\alpha_i, X_{i1}, \ldots, X_{iT}).$$

This is a difficult distribution to specify. One could directly specify this distribution, but one might want it to be internally consistent across different number of time periods, and that makes it awkward to choose a functional form. See for general discussions of this initial conditions problem Wooldridge (2002). Honoré and Tamer investigate what can be learned about γ without making parametric assumptions about this distribution. From the literature

it is known that in many cases γ is not point-identified (for example, the case with $T \leq 3$, no covariates, and a logistic distribution for ϵ_{it}). Nevertheless, it may be that the range of values of γ consistent with the data is very small, and it might reveal the sign of γ.

Honoré and Tamer study the case with a discrete distribution for α, with a finite and known set of support points. They fix the support to be $-3, -2.8, \ldots, 2.8, 3$, with unknown probabilities. Given that the ϵ_{it} are standard normal, this is very flexible. In a computational exercise they assume that the true probabilities make this discrete distribution mimic the standard normal distribution. In addition they set $\Pr(Y_{i1} = 1|\alpha_i) = 1/2$. In the case with $T = 3$ they find that the range of values for γ consistent with the data generating process (the identified set) is very narrow. If γ is in fact equal to zero, the width of the set is zero. If the true value is $\gamma = 1$, then the width of the interval is approximately 0.1. (It is largest for γ close to, but not equal to, -1.) See Figure 1, taken from Honoré and Tamer (2006).

The Honoré-Tamer analysis, in the context of the literature on initial conditions problems, shows very nicely the power of the partial identification approach. A problem that had been viewed as essentially intractable, with many non-identification results, was shown to admit potentially precise inferences despite these non-identification results.

2.5 AUCTION DATA

Haile and Tamer (2003, HT from hereon), in what is one of the most influential applications of the partial identification approach, study English or oral ascending bid auctions. In such auctions bidders offer increasingly higher prices until only one bidder remains. HT focus on a symmetric independent private values model. In auction t, for $t = 1, \ldots, T$, bidder i has a value ν_{it}, drawn independently from the value for bidder j. Large sample results refer to the number of auctions getting large. HT assume that the value distribution is the same in each auction (after adjusting for observable auction characteristics). A key object of interest, is the value distribution. Given that one can derive other interesting objects, such as the optimal reserve price.

One can imagine a set up where the researcher observes, as the price increases, for each
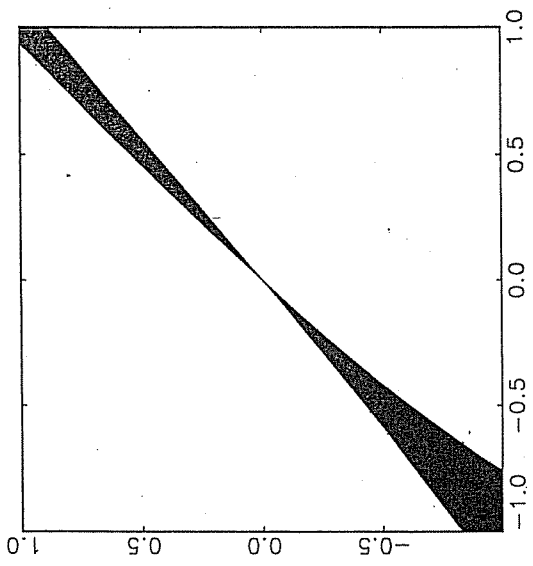
FIGURE 1.—Identified region for γ as a function of its true value.

bidder whether that bidder is still participating in the auction. (Milgrom and Weber (1982) assume that each bidder continuously confirms their participation by holding down a button while prices rise continuously.) In that case one would be able to infer for each bidder their valuation, and thus directly estimate the value distribution.

This is not what is typically observed. Instead of prices rising continuously, there are jumps in the bids, and for each bidder we do not know at any point in time whether they are still participating unless they subsequently make a higher bid. HT study identification in this, more realistic, setting. They assume that no bidder ever bids more than their valuation, and that no bidder will walk away and let another bidder win the auction if the winning bid is lower than their own valuation. Under those two assumptions, HT show that one can derive bounds on the value distribution.

One set of bounds they propose is as follows. Let the highest bid for participant i in auction t be b_{it}. The number of participants in auction t is n_t. Ignoring any covariates, let the distribution of the value for individual i, ν_{it}, be $F_\nu(v)$. This distribution function is the same for all auctions. Let $F_b(b) = \Pr(b_{it} \leq b)$ be the distribution function of the bids (ignoring variation in the number of bidders by auction). This distribution can be estimated because the bids are observed. The winning bid in auction t is $B_t = \max_{i=1,\ldots,n_t} b_{it}$. First HT derive an upper bound on the distribution function $F_\nu(v)$. Because no bidder ever bids more than their value, it follows that $b_{it} \leq \nu_{it}$. Hence, without additional assumptions,

$$F_\nu(v) \leq F_b(v), \qquad \text{for all } v.$$

For a lower bound on the distribution function one can use the fact that the second highest of the values among the n participants in auction t must be less than or equal to the winning bid. This follows from the assumption that no participant will let someone else win with a bid below their valuation. Let $F_{\nu,m:n}(v)$ denote the mth order statistic in a random sample of size n from the value distribution, and let $F_{B,n:n}(b)$ denote the distribution of the

winning bid in auctions with n participants. Then

$$F_{B,n:n}(v) \leq F_{\nu,n-1:n}(v).$$

The distribution of the any order statistic is monotonically related to the distribution of the parent distribution, and so a lower bound on $F_{\nu,n-1:n}(v)$ implies a lower bound on $F_\nu(v)$.

HT derive tighter bounds based on the information in other bids and the inequalities arising from the order statistics, but the above discussion illustrates the point that outside of the Milgrom-Weber button auction model one can still derive bounds on the value distribution in an English auction even if one cannot point-identify the value distribution. If in fact the highest bid for each individual was equal to their value (other than for the winner for whom the bid is equal to the second highest value), the bounds would collaps and point-identification would be obtained.

2.6 Entry Models and Inequality Conditions

Recently a number of papers has studied entry models in settings with multiple equilibria. In such settings traditionally researchers have added *ad hoc* equilbrium selection mechanisms. In the recent literature a key feature is the avoidance of such assumptions, as these are often difficult to justify on theoretical grounds. Instead the focus is on what can be learned in the absence of such assumptions. In this section I will discuss some examples from this literature. An important feature of these models is that they often lead to inequality restrictions, where the parameters of interest θ satisfy

$$\mathbb{E}[\psi(Z,\theta)] \geq 0,$$

for known $\psi(z,\theta)$. This relates closely to the standard (Hansen, 1983) generalized method of moments (GMM) set up where the functions $\psi(Z,\theta)$ would have expectation equal to zero at the true values of the parameters. We refer to this as the generalized inequality restrictions (GIR) form. These papers include Pakes, Porter, Ho, and Ishii (2006), Cilberto and Tamer (2004, CM from hereon), Andrews, Berry and Jia (2004). Here I will discuss a simplified

version of the CM model. Suppose two firms, A and B, contest a set of markets. In market m, $m = 1, \ldots, M$, the profits for firms A and B are

$$\pi_{Am} = \alpha_A + \delta_A \cdot d_{Bm} + \varepsilon_{Am}, \qquad \text{and} \quad \pi_{Bm} = \alpha_B + \delta_B \cdot d_{Am} + \varepsilon_{Bm}.$$

where $d_{Fm} = 1$ if firm F is present in market m, for $F \in \{A, B\}$, and zero otherwise. The more realistic model CM consider also includes observed market and firm characteristics. Firms enter market m if their profits in that market are positive. Firms observe all components of profits, including those that are unobserved to the econometrician, $(\varepsilon_{Am}, \varepsilon_{Bm})$, and so their decisions satisfy:

$$d_{Am} = 1\{\pi_{Am} \geq 0\}, \qquad d_{Bm} = 1\{\pi_{Bm} \geq 0\}. \tag{1}$$

(Pakes, Porter, Ho, and Ishii allow for incomplete information where expected profits are at least as high for the action taken as for actions not taken, given some information set.) The unobserved (to the econometrician) components of profits, ε_{Fm}, are independent accross markets and firms. For ease of exposition we assume here that they have a normal $\mathcal{N}(0, 1)$ distribution. (Note that we only observe indicators of the sign of profits, so the scale of the unobserved components is not relevant for predictions.) The econometrician observes in each market only the pair of indicators d_A and d_B. We focus on the case where the effect of entry of the other firm on a firm's profits, captured by the parameters δ_A and δ_B is negative, which is the case of most economic interest.

An important feature of this model is that given the parameters $\theta = (\alpha_A, \delta_A, \alpha_B, \delta_B)$, for a given set of $(\varepsilon_{Am}, \varepsilon_{Bm})$ there is not necessarily a unique solution (d_{Am}, d_{Bm}). For pairs of values $(\varepsilon_{Am}, \varepsilon_{Bm})$ such that

$$-\alpha_A < \varepsilon_A \leq -\alpha_A - \delta_A, \qquad -\alpha_B < \varepsilon_B \leq -\alpha_B - \delta_B,$$

both $(d_A, d_B) = (0, 1)$ and $(d_A, d_B) = (1, 0)$ satisfy the profit maximization condition (1). In the terminology of this literature, the model is not *complete*. It does not specify the

outcomes given the inputs. Figure 1, adapted from CM, shows the different regions in the $(\varepsilon_{Am}, \varepsilon_{Bm})$ space.

The implication of this is that the probability of the outcome $(d_{Am}, d_{Bm}) = (0, 1)$ cannot be written as a function of the parameters of the model, $\theta = (\alpha_A, \delta_A, \alpha_B, \delta_B)$, even given distributional assumptions on $(\varepsilon_{Am}, \varepsilon_{Bm})$. Instead the model implies a lower and upper bound on this probability:

$$H_{L,01}(\theta) \leq \Pr\left((d_{Am}, d_{Bm}) = (0, 1)\right) \leq H_{U,01}(\theta).$$

Inspecting Figure 1 shows that

$$H_{L,01}(\theta) = \Pr(\varepsilon_{Am} < -\alpha_A, -\alpha_B < \varepsilon_{Bm})$$

$$+\Pr(-\alpha_A \leq \varepsilon_{Am} < -\alpha_A - \delta_A, -\alpha_B - \delta_B < \varepsilon_{Bm}),$$

and

$$H_{U,01}(\theta) = \Pr(\varepsilon_{Am} < -\alpha_A, \alpha_B < \varepsilon_{Bm})$$

$$+\Pr(-\alpha_A \leq \varepsilon_{Am} < -\alpha_A - \delta_A, -\alpha_B - \delta_B < \varepsilon_{Bm}),$$

$$+\Pr(-\alpha_A \leq \varepsilon_{Am} < -\alpha_A - \delta_A, -\alpha_B < \varepsilon_{Bm} < -\alpha_B - \delta_B),$$

Similar expressions can be derived for the probability $\Pr\left((d_{Am}, d_{Bm}) = (1, 0)\right)$. Thus in general we can write the information about the parameters in large samples as

$$\begin{pmatrix} H_{L,00}(\theta) \\ H_{L,01}(\theta) \\ H_{L,10}(\theta) \\ H_{L,11}(\theta) \end{pmatrix} \leq \begin{pmatrix} \Pr\left((d_{Am}, d_{Bm}) = (0, 0)\right) \\ \Pr\left((d_{Am}, d_{Bm}) = (0, 1)\right) \\ \Pr\left((d_{Am}, d_{Bm}) = (1, 0)\right) \\ \Pr\left((d_{Am}, d_{Bm}) = (1, 1)\right) \end{pmatrix} \leq \begin{pmatrix} H_{U,00}(\theta) \\ H_{U,01}(\theta) \\ H_{U,11}(\theta) \\ H_{U,11}(\theta) \end{pmatrix}.$$

(For $(d_A, d_B) = (0, 0)$ or $(d_A, d_B) = (1, 1)$ the lower and upper bound coincide, but for ease of exposition we treat all four configurations symmetrically.) The $H_{L,ij}(\theta)$ and $H_{U,ij}(\theta)$ are

Figure 1

(d_A, d_B)
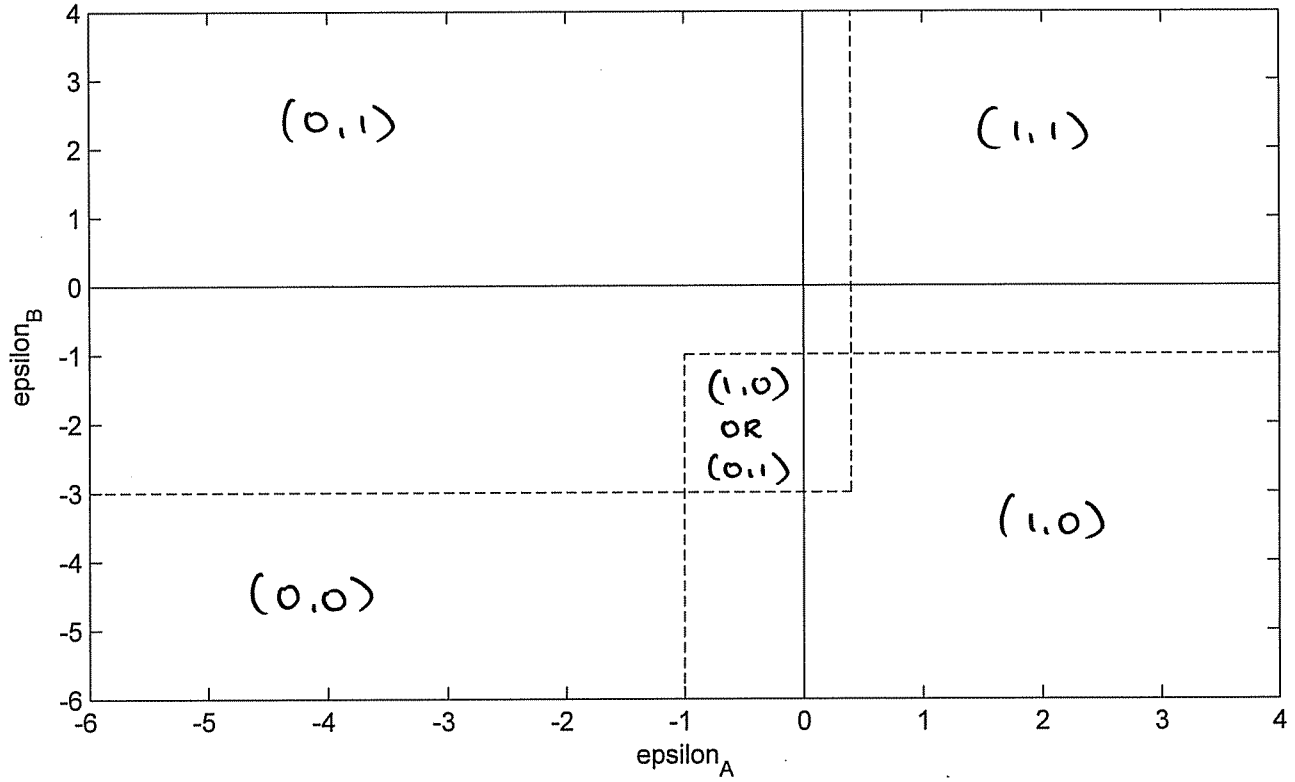
$\alpha_A = 1 \qquad \delta_A = -1.4$

$\alpha_B = 3 \qquad \delta_B = -2$

known functions of θ. The data allow us to estimate the foru probabilities, which contain only three separate pieces of information because the probabilities add up to one. Given these probabilities, the identified set is the set of all θ that satisfy all eight inequalities. In the simple model above, there are four parameters. Even in the case with the lower and upper bounds for the probabilities coinciding, these would in general not be identified.

We can write this in the GIR form by defining

$$\psi(d_A, d_B | \alpha_A, \alpha_B, \delta_A, \delta_B) = \begin{pmatrix} H_{U,00}(\theta) - (1 - d_A) \cdot (1 - d_B) \\ (1 - d_A) \cdot (1 - d_B) - H_{L,00}(\theta) \\ H_{U,01}(\theta) - (1 - d_A) \cdot d_B \\ (1 - d_A) \cdot d_B - H_{L,01}(\theta) \\ H_{U,10}(\theta) - d_A \cdot (1 - d_B) \\ d_A \cdot (1 - d_B) - H_{L,10}(\theta) \\ H_{U,11}(\theta) - d_A \cdot d_B \\ d_A \cdot d_B - H_{L,11}(\theta) \end{pmatrix},$$

so that the model implies that at the true values of the parameters

$$\mathbb{E}\left[\psi(d_A, d_B | \alpha_A, \alpha_B, \delta_A, \delta_B)\right] \geq 0.$$

3. ESTIMATION

Chernozhukov, Hong, and Tamer (2007, CHT) consider, among other things, the case with moment inequality conditions,

$$\mathbb{E}[\psi(Z, \theta)] \geq 0,$$

where $\psi(z, \theta)$ is a known vector of functions, of dimension M, and the unknown parameter θ is of dimension K. Let Θ be the parameter space, a subset of $\mathbb{R}^K$.

Define for a vector x the vector $(x)_+$ to be the component-wise non-negative part, and $(x)_-$ to be the component-wise non-positive part, so that for all x, $x = (x)_- + (x)_+$. For a given $M \times M$ non-negative definite weight matrix W, CHT consider the population objective function

$$Q(\theta) = \mathbb{E}[\psi(Z, \theta)]'_- W \mathbb{E}[\psi(Z, \theta)]_-.$$

For all θ in the identified set, denoted by $\Theta_I \subset \Theta$, we have $Q(\theta) = 0$.

The sample equivalent to this population objective function is

$$Q_N(\theta) = \left(\frac{1}{N} \sum_{i=1}^{N} \psi(Z_i, \theta) \right)'_{-} W \left(\frac{1}{N} \sum_{i=1}^{N} \psi(Z_i, \theta) \right)_{-}.$$

We cannot simply estimate the identified set as

$$\tilde{\Theta}_I = \{ \theta \in \Theta \, | \, Q_N(\theta) = 0 \},$$

The reason is that even for θ in the identified set $Q_N(\theta)$ may be positive with high probability. A simple way to see that is to consider the standard GMM case with equalities and over-identification. If $\mathbb{E}[\psi(Z, \theta)] = 0$, the objective function will not be zero in finite samples in the case with over-identification. As a result, $\tilde{\Theta}_I$ can be empty when Θ_I is not, even in large samples.

This is the reason CHT estimate the set Θ_I as

$$\hat{\Theta}_I = \{ \theta \in \Theta \, | \, Q_N(\theta) \leq a_N \},$$

where $a_N \to 0$ at the appropriate rate. In most regular problems $a_N = c/N$, leading to an estimator $\hat{\Theta}_I$ that is consistent for Θ_I, by which we mean that the two sets get close to each other, in the Hausdorf sense that

$$\sup_{\theta \in \Theta_I} \inf_{\theta' \in \hat{\Theta}_I} d(\theta, \theta') \longrightarrow 0, \qquad \text{and} \quad \sup_{\theta' \in \hat{\Theta}_I} \inf_{\theta \in \Theta_I} d(\theta, \theta') \longrightarrow 0,$$

where $d(\theta, \theta') = ((\theta - \theta)'(\theta - \theta'))^{1/2}$.

3. INFERENCE: GENERAL ISSUES

There is a rapidly growing literature concerned with developing methods for inference in partially identified models, including Beresteanu and Molinari (2006), Chernozhukov, Hong, and Tamer (2007), Imbens and Manski (2004), Rosen (2006), and Romano and Shaikh

(2007ab). In many cases the partially identified set itself is difficult to characterize. In the scalar case this can be much simpler. There it often is an interval, $[\theta_{\text{LB}}, \theta_{\text{UB}}]$. There are by now a number of proposals for constructing confidence sets. They differ in implementation as well as in their goals. One issue is whether one wants a confidence set that includes each element of the identified set with fixed probability, or the entire identified set with that probability. Formally, the first question looks for a confidence set CI_α^θ that satisfies

$$\inf_{\theta \in [\theta_{\text{LB}}, \theta_{\text{UB}}]} \Pr\left(\theta \in \text{CI}_\alpha^\theta\right) \geq \alpha.$$

In the second case we look for a set $\text{CI}_\alpha^{[\theta_{\text{LB}}, \theta_{\text{UB}}]}$ such that

$$\Pr\left([\theta_{\text{LB}}, \theta_{\text{UB}}] \subset \text{CI}_\alpha^\theta\right) \geq \alpha.$$

The second requirement is stronger than the first, and so generally $\text{CI}_\alpha^\theta \subset \text{CI}_\alpha^{[\theta_{\text{LB}}, \theta_{\text{UB}}]}$. Here we follow Imbens and Manski (2004) and Romano and Shaikh (2007a) who focus on the first case. This seems more in line with teh traditional view of confidence interval in that they should cover the true value of the parameter with fixed probability. It is not clear why the fact that the object of interest is not point-identified should change the definition of a confidence interval. CHT and Romano and Shaikh (2007b) focus on the second case.

Next we discuss two specific examples to illustrate some of the issues that can arise, in particular the uniformity of confidence intervals.

3.1 INFERENCE: A MISSING DATA PROBLEM

Here we continue the missing data example from Section 2.1. We have a random sample of $(W_i, W_i \cdot Y_i)$, for $i = 1, \ldots, N$. Y_i is known to lie in the interval $[0, 1]$, interest is in $\theta = \mathbb{E}[Y]$, and the parameter space is $\Theta = [0, 1]$. Define $\mu_1 = \mathbb{E}[Y|W = 1]$, $\lambda = \mathbb{E}[Y|W = 0]$, $\sigma^2 = \mathbb{V}(Y|W = 1)$, and $p = \mathbb{E}[W]$. For ease of exposition we assume p is known. The identified set is

$$\Theta_I = [p \cdot \mu_1, p \cdot \mu_1 + (1 - p)].$$

Imbens and Manski (2004) discuss confidence intervals for this case. The key feature of this problem, and similar ones, is that the lower and upper bounds are well-behaved functionals of the joint distribution of the data that can be estimated at the standard parametric $\sqrt{N}$ rate with an asymptotic normal distribution. In this specific example the lower and upper bound are both functions of a single unknown parameter, the conditional mean μ_1. The first step is a 95% confidence interval for μ_1. Let $N_1 = \sum_i W_i$ and $\overline{Y}_1 = \sum_i W_i \cdot Y_i / N_1$. The standard confidence interval is

$$\mathrm{CI}_\alpha^{\mu_1} = \left[\overline{Y} - 1.96 \cdot \sigma / \sqrt{N_1}, \overline{Y} + 1.96 \cdot \sigma / \sqrt{N_1} \right].$$

Consider the confidence interval for the lower and upper bound:

$$\mathrm{CI}_\alpha^{p \cdot \mu_1} = \left[p \cdot \left(\overline{Y} - 1.96 \cdot \sigma / \sqrt{N_1} \right), p \cdot \left(\overline{Y} + 1.96 \cdot \sigma / \sqrt{N_1} \right) \right],$$

and

$$\mathrm{CI}_\alpha^{p \cdot \mu_1 + (1-p)} = \left[p \cdot \left(\overline{Y} - 1.96 \cdot \sigma / \sqrt{N_1} \right) + (1-p), p \cdot \left(\overline{Y} + 1.96 \cdot \sigma / \sqrt{N_1} \right) + 1 - p \right].$$

A simple and valid confidence interval can be based on the lower confidence bound on the lower bound and the upper confidence bound on the upper bound:

$$\mathrm{CI}_\alpha^{\theta} = \left[p \cdot \left(\overline{Y} - 1.96 \cdot \sigma / \sqrt{N_1} \right), p \cdot \left(\overline{Y} + 1.96 \cdot \sigma / \sqrt{N_1} \right) + 1 - p \right].$$

This is generally conservative. For each θ in the interior of Θ_I, the asymptotic coverage rate is 1. For $\theta \in \{ \theta_{\mathrm{LB}}, \theta_{\mathrm{UB}} \}$, the coverage rate is $\alpha + (1 - \alpha)/2$.

The interval can be modified to give asymptotic coverage equal to α by changing the quantiles used in the confidence interval construction, essentially using one-sided critical values,

$$\mathrm{CI}_\alpha^{\theta} = \left[p \cdot \left(\overline{Y} - 1.645 \cdot \sigma / \sqrt{N_1} \right), p \cdot \left(\overline{Y} + 1.645 \cdot \sigma / \sqrt{N_1} \right) + 1 - p \right].$$

This has the problem that if $p = 0$ (when θ is point-identified), the coverage is only $\alpha - (1-\alpha)$. In fact, for values of p close to zero, the confidence interval would be shorter than the confidence interval in the point-identified case. Imbens and Manski (2004) suggest modifying the confidence interval to

$$\mathrm{CI}_\alpha^\theta = \left[p \cdot \left(\overline{Y} - C_N \cdot \sigma/\sqrt{N_1} \right), p \cdot \left(\overline{Y} + C_N \cdot \sigma/\sqrt{N_1} \right) + 1 - p \right],$$

where the critical value C_N satisfies

$$\Phi \left(C_N + \sqrt{N} \cdot \frac{1-p}{\sigma/\sqrt{p}} \right) - \Phi \left(-C_N \right) = \alpha.$$

and $C_N = 1.96$ if $p = 0$. This confidence interval has asymptotic coverage 0.95, uniformly over p.

3.2. INFERENCE: MULTIPLE INEQUALITIES

Here we look at inference in the Genereralized Inequality (GIR) setting. The example is a simplified version of the moment inequality type of problems discussed in CHT, Romano and Shaikh (2007ab), Pakes, Porter, Ho, and Ishii (2006), and Andrews and Guggenberger (2007). Suppose we have two moment inequalities,

$$\mathbb{E}[X] \geq \theta, \qquad \text{and} \ \ \mathbb{E}[Y] \geq \theta.$$

The parameter space is $\Theta = [0, \infty)$. Let $\mu_X = \mathbb{E}[X]$, and $\mu_Y = \mathbb{E}[Y]$. We have a random sample of size N of the pairs (X, Y). The identified set is

$$\Theta_I = [0, \min(\mu_X, \mu_Y)].$$

The key difference with the previous example is that the upper bound is no longer a smooth, well-behaved functional of the joint distribution. In the simple two-inequality example, if μ_X is close to μ_Y, the distribution of the estimator for the upper bound is not well approximated by a normal distribution. Suppose we estimate the means of X and Y by

$\overline{X}$, and $\overline{Y}$, and that the variances of X and Y are known to be equal to σ^2. A naive 95% confidence interval would be

$$C_\alpha^\theta = [0, \min(\overline{X}, \overline{Y}) + 1.645 \cdot \sigma/N].$$

This confidence interval essentially ignores the moment inequality that is not binding in the sample. It has asymptotic 95% coverage for all values of μ_X, μ_Y, as long as $\min(\mu_X, \mu_Y) > 0$, and $\mu_X \neq \mu_Y$. The first condition $(\min(\mu_X, \mu_Y) > 0)$ is the same as the condition in the Imbens-Manski example. It can be dealt with in the same way by adjusting the critical value slightly based on an initial estimate of the width of the identified set.

The second condition raises a different uniformity concern. The naive confidence interval essentially assumes that the researcher knows which moment conditions are binding. This is true in large samples, unless there is a tie. However, in finite samples ignoring uncertainty regarding the set of binding moment inequalities may lead to a poor approximation, especially if there are many inequalities. One possibility is to construct conservative confidence intervals (e.g., Pakes, Porter, Ho, and Ishii, 2007). However, such intervals can be unnecessarily conservative if there are moment inequalities that are far from binding.

One would like construct confidence intervals that asymptotically ignore irrelevant inequalities, and at the same time are valid uniformly over the parameter space. Bootstrapping is unlikely to work in this setting. One way of obtaining confidence intervals that are uniformly valid is based on subsampling. See Romano and Shaikh (2007a), and Andrews and Guggenberger (2007). Little is known about finite sample properties in realistic settings.

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