

Inference in a class of optimization problems: confidence regions and finite sample bounds on errors in coverage probabilities

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INFERENCE IN A CLASS OF OPTIMIZATION PROBLEMS: CONFIDENCE REGIONS AND FINITE SAMPLE BOUNDS ON ERRORS IN COVERAGE PROBABILITIES

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ABSTRACT. This paper describes a method for carrying out inference on partially identified parameters that are solutions to a class of optimization problems. The optimization problems arise in applications in which grouped data are used for estimation of a model's structural parameters. The parameters are characterized by restrictions that involve the unknown population means of observed random variables in addition to the structural parameters of interest. Inference consists of finding confidence intervals for the structural parameters. Our theory provides a finite-sample bound on the difference between the true and nominal probabilities with which a confidence interval contains the true but unknown value of a parameter. We contrast our method with an alternative inference method based on the median-of-means estimator of Minsker (2015). The results of Monte Carlo experiments and empirical examples illustrate the usefulness of our method.

KEYWORDS: partial identification, normal approximation, finite-sample bounds

1. INTRODUCTION

We present a method for carrying out inference about a partially identified function of structural parameters of an econometric model. Our method applies to models that impose shape restrictions (e.g., Freyberger and Horowitz, 2015; Horowitz and Lee, 2017), a variety of partially identified models (e.g., Manski, 2007a; Tamer, 2010), and models in which a continuous function is inferred from the average values of variables in a finite number of discrete groups (e.g., Blundell, Duncan, and Meghir, 1998; Kline and Tartari, 2016). The specific inference problem consists of finding upper and lower

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bounds on the partially identified function $f(\psi)$ under the restrictions $g_1(\psi, \mu) \leq 0$ and $g_2(\psi, \mu) = 0$, where ψ is a vector of structural parameters; μ is a vector of unknown population means of observable random variables; f is a known, real-valued function; and g_1 and g_2 are known possibly vector-valued functions. The inequality $g_1(\psi, \mu) \leq 0$ holds component-wise.

Most existing methods for inference in our framework are based on asymptotic theory. They provide correct inference in the limit $n \rightarrow \infty$ but do not provide information about the accuracy of finite-sample inference. Our theory also relies on asymptotic approximations to form confidence intervals. However, in contrast to existing asymptotic methods, it provides a finite-sample bound on the difference between the true and nominal coverage probabilities of a confidence interval for $f(\psi)$. Thus, our results provide information about the accuracy of finite-sample inference.

There are several approaches to carrying out non-asymptotic inference in our framework. In some cases, a statistic with a known finite-sample distribution makes finite-sample inference possible. For example, the Clopper–Pearson (1934) confidence interval for a population probability is obtained by inverting the binomial probability distribution function. We use the Clopper–Pearson confidence interval in the empirical example presented in Section 5 of this paper. Manski (2007b) used the Clopper–Pearson interval to construct finite-sample confidence sets for counterfactual choice probabilities. A second method consists of using a finite-sample concentration inequality to obtain a confidence interval. This method is useful for applications only if the inequality provides a bound that does not depend on unknown population parameters. Hoeffding’s inequality for the mean of a scalar random variable with known bounded support provides such a bound. Syrgkanis, Tamer, and Ziani (2018) used Hoeffding’s inequality to construct a confidence interval for a partially identified population moment. Hoeffding’s inequality cannot be used if the (bounded) support of the underlying random variable is unknown. The generalization of Hoeffding’s inequality to confidence intervals for sub-Gaussian random variables requires information about a certain parameter of the distribution of the underlying random variable that is typically unavailable in applications. Minsker (2015) developed a confidence set for a vector of population means using a method called “median of means.” This method depends on certain tuning parameters. There are no data-based, efficient ways to choose these parameters in applications. Section 4 of this paper presents the results of Monte Carlo experiments comparing the widths of confidence intervals obtained by using Minsker’s (2015) method and our method.

The approach that we use here consists of making a normal approximation to the unknown distribution of a sample average. A variety of results provide finite-sample upper bounds on the errors made by normal approximations. The Berry-Esséen inequality for the average of a scalar random variable is a well-known example of such a bound. Bentkus (2003) provides a bound for the error of a multivariate normal approximation to the distribution of the sample average of a random vector. Other normal approximations are given by Spokoiny and Zhilova (2015); Chernozhukov, Chetverikov, and Kato (2017); and Zhilova (2020); among others. The method described in this paper uses the normal approximation of Raič (2019), which is a refined version of Bentkus (2003) that does not require boundedness of the random variables involved and treats random vectors. When μ is high-dimensional, Chernozhukov, Chetverikov, and Kato (2017) may provide a tighter bound; however, the bound of Bentkus-Raič is narrower than that of Chernozhukov, Chetverikov, and Kato (2017) when the dimension of μ is fixed, which is the case we treat in this paper. In contrast to conventional asymptotic inference approaches, our theory provides a finite-sample bound on the difference between the true and nominal coverage probabilities of a confidence interval for the partially identified function $f(\psi)$. In sum, our approach is a hybrid of normal approximations and non-asymptotic theory in that critical values are determined by normal approximations and a lower bound on the finite-sample coverage probability of a confidence interval is available. In general, the lower bound depends on unknown population quantities and can be loose because it is a worst-case bound. Nonetheless, this non-asymptotic lower bound provides information about inference based on asymptotic approximations.

Our work is broadly related to the literature on inference in partial identified models. See Tamer (2010), Canay and Shaikh (2017), Ho and Rosen (2017), and Molinari (2020) for recent surveys. Chen, Christensen, and Tamer (2018) describe a Monte Carlo method for carrying out asymptotic inference for a class of models that includes our framework. Bugni, Canay, and Shi (2017) and Kaido, Molinari, and Stoye (2019) develop asymptotic inference methods for subvectors of partially identified parameters in moment inequality models. Chernozhukov, Chetverikov, and Kato (2019) and Belloni, Bugni, and Chernozhukov (2018) construct confidence regions by inverting pointwise tests on a hypothesis about the (sub)vector of parameters that are partially identified by a large number of moment inequalities. The inference problem we treat is different from those in the aforementioned papers in that we focus on the inference on parameters that are solutions to a class of optimization problems. A more

closely related working paper is Hsieh, Shi, and Shum (2017), who propose a method for asymptotic inference about estimators defined by mathematical programs. However, their inference method is different from ours because they recast their inference problem into one based on a set of inequalities with pre-estimated coefficients.

Our work is also related to the econometrics literature on finite-sample inference. Syrgkanis, Tamer, and Ziani (2018) consider finite-sample inference in auction models. Their framework and method are very different from those in this paper. In different context, Chernozhukov, Hansen, and Jansson (2009) and Rosen and Ura (2019) propose finite-sample inference for quantile regression models and for the maximum score estimand, respectively. Their methods are distinct from ours.

The remainder of this paper is organized as follows. Section 2 describes our method for obtaining confidence intervals and describes three empirical studies that illustrate how the inferential problem the method addresses arises in applications. Section 3 describes computational procedures for implementing our method. Section 4 reports the results of a Monte Carlo investigation of the numerical performance of our method, and Section 5 presents two empirical applications of the method. Section 6 gives concluding comments. Appendix A presents the proofs of theorems. Appendix B provides additional details on our computational procedures. Appendix C describes Minsker’s (2015) median of means method.

2. THE METHOD

Section 2.1 presents an informal description of inferential problem we address, and Section 2.2 generalizes the form of the objective function. Section 2.3 gives two examples of empirical applications in which the inferential problem arises. Section 2.4 provides a formal description of the method for constructing confidence intervals. Section 2.5 treats the possibility that g_1 and g_2 depend on a continuous covariate in addition to (ψ, μ) .

2.1. The Inferential Problem. Let $\{X_i : i = 1, \dots, n\}$ be a random sample from the distribution of the random vector $X \in \mathbb{R}^p$ for some finite $p \geq 1$. Define $\mu = \mathbb{E}(X)$ and $\Sigma = \text{cov}(X)$. Let ψ be a finite-dimensional parameter and $f(\psi)$ be a real-valued, known function. We assume throughout this section that $f(\psi)$ is only partially identified by the sampling process, though our results also hold if $f(\psi)$ is point identified. We seek a confidence interval for $f(\psi)$, which we define as an interval that contains $f(\psi)$ with probability exceeding a known value. Let $g_1(\psi, \mu)$ and $g_2(\psi, \mu)$ be possibly vector valued known functions satisfying $g_1(\psi, \mu) \leq 0$ and

$g_2(\psi, \mu) = 0$. Define

$$(2.1) \quad J_+ := \max_{\psi} f(\psi) \quad \text{and} \quad J_- := \min_{\psi} f(\psi)$$

subject to the component-wise constraints:

$$(2.2a) \quad g_1(\psi, \mu) \leq 0,$$

$$(2.2b) \quad g_2(\psi, \mu) = 0,$$

$$(2.2c) \quad \psi \in \Psi,$$

where Ψ is a compact parameter set.

In this setting, we are interested in the identification interval $J_- \leq f(\psi) \leq J_+$. However, this interval cannot be calculated in applications because μ is unknown. Therefore, we estimate μ by the sample average $\bar{X} = n^{-1} \sum_{i=1}^n X_i$, and we estimate J_+ and J_- by

$$(2.3) \quad \hat{J}_+(\bar{X}) := \max_{\psi, m} f(\psi) \quad \text{and} \quad \hat{J}_-(\bar{X}) := \min_{\psi, m} f(\psi)$$

subject to

$$(2.4a) \quad g_1(\psi, m) \leq 0,$$

$$(2.4b) \quad g_2(\psi, m) = 0,$$

$$(2.4c) \quad \psi \in \Psi,$$

and

$$(2.4d) \quad n^{1/2}(\bar{X} - m) \in \mathcal{S},$$

where \mathcal{S} is a set, specified in Section 2.3, for which $n^{1/2}(\bar{X} - \mu) \in \mathcal{S}$ with high probability. Since μ is unknown, we replace it with a variable of optimization in (2.3)–(2.4) but restrict that variable to \mathcal{S} . The resulting confidence interval for $f(\psi)$ is

$$(2.5) \quad \hat{J}_-(\bar{X}) \leq f(\psi) \leq \hat{J}_+(\bar{X}).$$

Section 2.4 provides a finite-sample lower bound on the probability that this interval contains $f(\psi)$. That is, Section 2.4 provides a finite-sample lower bound on

$$(2.6) \quad \mathbb{P} \left[\hat{J}_-(\bar{X}) \leq J_- \leq f(\psi) \leq J_+ \leq \hat{J}_+(\bar{X}) \right].$$

2.2. Generalizing the Form of the Objective Function. It is straightforward to allow the objective function $f(\psi)$ to depend on μ . For the lower bound $\hat{J}_-(\bar{X})$, we introduce an auxiliary variable t that performs as an upper bound on $f(\psi, \mu)$ and solve: $\min_{\psi, m, t} t$ subject to $f(\psi, m) \leq t$ and (2.4). Analogously, for the upper bound $\hat{J}_+(\bar{X})$, we introduce a lower bound on the objective and solve: $\max_{\psi, m, s} s$ subject to $f(\psi, m) \geq s$ and (2.4). We focus on the original formulation in the previous section since the new formulation can be rewritten as the original form by redefining f , g_1 and g_2 .

2.3. Examples of Empirical Applications.

Example 1. Blundell, Duncan, and Meghir (1998) use grouped data to estimate labor supply effects of tax reforms in the United Kingdom. To motivate our setup, we consider a simple model with which Blundell, Duncan, and Meghir (1998) describe how to use grouped data to estimate β in the labor supply model with no income effect:

$$(2.7) \quad h_{it} = \alpha + \beta \ln w_{it} + U_{it},$$

where h_{it} and w_{it} , respectively, are hours of work and the post-tax hourly wage rate of individual i in year t , and U_{it} is an unobserved random variable that satisfies certain conditions. The parameter β is identified by a relation of the form

$$\beta = \beta(h_{gt}, lw_{gt}),$$

where h_{gt} and lw_{gt} are the mean hours and log wages in year t of individuals in group g . There are 8 groups defined by four year-of-birth cohorts and level of education. The data span the period 1978-1992.

A nonparametric version of (2.7) is

$$(2.8) \quad h_{it} = f(w_{it}) + U_{it},$$

where $f \in \mathcal{F}$ is an unknown continuous function and \mathcal{F} is a function space. A nonparametric analog of β is the weighted average derivative

$$\tilde{\beta} = \int \frac{\partial f(u)}{\partial u} w(u) du,$$

where w is a non-negative weight function. The average derivative $\tilde{\beta}$ is not identified non-parametrically by the mean values of hours and wages for finitely many groups and time periods. It can be partially identified, however, by imposing a shape restriction such as weak monotonicity on the labor supply function f . Assume, for

example, that $\mathbb{E}[h_{it} - f(w_{it})|g, t] = 0$. Blundell, Duncan, and Meghir (1998) set $\mathbb{E}[h_{it} - f(w_{it})|g, t] = a_g + m_t$, where a_g and m_t , respectively, are group and time fixed effects. These are accommodated by our framework but we do not do this in the present discussion.

The identification interval for $\tilde{\beta}$ is $\tilde{\beta}_- \leq \tilde{\beta} \leq \tilde{\beta}_+$, where

$$(2.9) \quad \tilde{\beta}_+ = \max_{f \in \mathcal{F}} \int \frac{\partial f(u)}{\partial u} w(u) du \quad \text{and} \quad \tilde{\beta}_- = \min_{f \in \mathcal{F}} \int \frac{\partial f(u)}{\partial u} w(u) du$$

subject to

$$(2.10a) \quad f(w_{gt}) - f(w_{g't'}) \leq 0 \quad \text{if} \quad w_{gt} < w_{g't'},$$

$$(2.10b) \quad h_{gt} - f(w_{gt}) = 0.$$

The continuous mathematical programming problem (2.9)-(2.10) can be put into the finite-dimensional framework of (2.3)-(2.4) by observing that under mild conditions on \mathcal{F} , f can be approximated very accurately by the truncated infinite series

$$(2.11) \quad f(u) \approx \sum_{j=1}^J \psi_j \phi_j(u),$$

where the ψ_j 's are constant parameters, the ϕ_j 's are basis functions for \mathcal{F} , and J is a truncation point. In an estimation setting, J can be an increasing function of the sample size, though we do not undertake this extension here. The approximation error of (2.11) can be bounded. Here, however, we assume that J is sufficiently large to make the error negligibly small. The finite-dimensional analog of (2.9)-(2.10) is

$$(2.12) \quad J_+ = \max_{\psi_j: j=1, \dots, J} \sum_{j=1}^J \psi_j \int \frac{\partial \phi_j(u)}{\partial u} w(u) du \quad \text{and} \quad J_- = \min_{\psi_j: j=1, \dots, J} \sum_{j=1}^J \psi_j \int \frac{\partial \phi_j(u)}{\partial u} w(u) du$$

subject to

$$(2.13a) \quad \sum_{j=1}^J \psi_j [\phi_j(w_{gt}) - \phi_j(w_{g't'})] \leq 0 \quad \text{if} \quad w_{gt} < w_{g't'},$$

$$(2.13b) \quad h_{gt} - \sum_{j=1}^J \psi_j \phi_j(w_{gt}) = 0.$$

J_+ and J_- can be estimated, thereby obtaining \hat{J}_+ and \hat{J}_- , by replacing h_{gt} and w_{gt} in (2.12)-(2.13) with within-group sample averages and adding the constraint (2.4d).

Example 2. Building on the theory of revealed preferences, Ho and Pakes (2014, HP hereafter) develop an inequality estimator for hospital choices using data from privately insured births in California. We consider a simplified setup of HP here.

Following closely the notation used in HP, let $p(c_i, h)$ denote the price an insurer is expected to pay at hospital h for a patient with condition c_i . Let l_i denote patient i 's location, l_h hospital's location, and $d(\cdot, \cdot)$ the distance between the two locations. For hospitals $h \neq h'$, define

$$\begin{aligned}\Delta p(c_i, h, h') &:= p(c_i, h) - p(c_i, h'), \\ \Delta d(l_i, l_h, l_{h'}) &:= d(l_i, l_h) - d(l_i, l_{h'}).\end{aligned}$$

That is, $\Delta p(c_i, h, h')$ is the price difference between hospitals h and h' given patient condition c_i and $\Delta d(l_i, l_h, l_{h'})$ is the distance difference between hospitals h and h' given patient location l_i . Define

$$u_{i,i'}(\psi, h, h') := \psi [\Delta p(c_i, h, h') + \Delta p(c_{i'}, h', h)] - [\Delta d(l_i, l_h, l_{h'}) + \Delta d(l_{i'}, l_{h'}, l_h)],$$

where ψ is a scalar parameter that determines price sensitivity relative to distance. Note that the coefficient for distance is normalized to be -1 . ψ is the key parameter in HP.

Define the four-dimensional vector of instruments based on distance:

$$z_{i,i'}(h, h') := \begin{pmatrix} \max\{\Delta d(l_i, l_h, l_{h'}), 0\} \\ -\min\{\Delta d(l_i, l_h, l_{h'}), 0\} \\ \max\{\Delta d(l_{i'}, l_{h'}, l_h), 0\} \\ -\min\{\Delta d(l_{i'}, l_{h'}, l_h), 0\} \end{pmatrix}.$$

Here, the instruments are based on distance measures and constructed to be positive to preserve the sign of inequalities.

Let $S(h, h', s)$ be the set of patients with severity s who chose hospital h but had hospital h' in their choice set. The identifying assumption in HP is that

$$(2.14) \quad \mathbb{E} \left[u_{i,i'}(\psi, h, h') z_{i,i'}(h, h') \mid i \in S(h, h', s), i' \in S(h', h, s) \right] \geq 0$$

for all s, h, h' such that $h \neq h'$. We can rewrite (2.14) as

$$\psi \mu_p(h, h', s) - \mu_d(h, h', s) \geq 0 \quad \text{for all } (h, h', s) \text{ such that } h \neq h',$$

where

$$\begin{aligned}\mu_p(h, h', s) &:= \mathbb{E} \left[z_{i,i'}(h, h') \{ \Delta p(c_i, h, h') + \Delta p(c_{i'}, h', h) \} \middle| i \in S(h, h', s), i' \in S(h', h, s) \right], \\ \mu_d(h, h', s) &:= \mathbb{E} \left[z_{i,i'}(h, h') \{ \Delta d(l_i, l_h, l_{h'}) + \Delta d(l_{i'}, l_{h'}, l_h) \} \middle| i \in S(h, h', s), i' \in S(h', h, s) \right].\end{aligned}$$

To see the connection between our general framework and HP's inequality estimator, let $f(\psi) = \psi$, $\mu = (\mu_p, \mu_d)$, and g_1 be a collection of inequalities such that

$$g_1(\psi, \mu) = \{ \mu_d(h, h', s) - \psi \mu_p(h, h', s) \text{ for all } (h, h', s) \text{ such that } h \neq h' \}.$$

There is no element in g_2 (no equality constraints here). Since each element in μ can be estimated by a suitable sample mean, our general framework includes HP's estimator as a special case.

Example 3. Kline and Tartari (2016, KT hereafter) study the impact of Connecticut's Jobs First (JF) welfare reform experiment on women's labor supply and welfare participation decisions. KT compare behavior under the JF and federal Aid to Families with Dependent Children (AFDC) regimes. The parameters of interest in KT are the probabilities with which a woman makes certain choices. The choice set under each regime (JF and AFDC) is denoted by $\{0n, 1n, 2n, 0r, 1r, 1u, 2u\}$, where 0 denotes no earnings, 1 denotes earnings below the poverty line, 2 denotes earnings above the poverty line, n denotes non-participation in welfare, r denotes welfare participation with truthful reporting of earnings, and u denotes welfare participation with under-reporting of earnings. Let π_{s^A, s^J} denote the probability of a woman's choosing alternative s^J under JF conditional on her choosing alternative s^A under AFDC. The possible choice probabilities and parameters of interest in KT are

$$\pi_{s^A, s^J} = [\pi_{0n, 1r}, \pi_{0r, 0n}, \pi_{2n, 1r}, \pi_{0r, 2n}, \pi_{0r, 1r}, \pi_{0r, 1n}, \pi_{1n, 1r}, \pi_{0r, 2u}, \pi_{2u, 1r}]'$$

The observable choices are welfare participation status and reported earnings under JF and AFDC. The population probabilities of observable choices are

$$\mathbf{p}^t := [p_{0n}^t, p_{1n}^t, p_{2n}^t, p_{0p}^t, p_{1p}^t, p_{2p}^t]'$$

for $t = A$ or J and the subscript p denotes welfare participation. These probabilities do not identify π_{s^A, s^J} .

To illustrate inference about the partially identified parameter π_{s^A, s^J} in our framework, consider the lower bound on $\pi_{2n, 1r}$. By inequality (15) of KT,

$$\pi_{2n, 1r} \geq \max \left\{ 0, \frac{p_{2n}^A - p_{2n}^J}{p_{2n}^A} \right\}.$$

Therefore,

$$p_{2n}^A - p_{2n}^J - p_{2n}^A \pi_{2n, 1r} \leq 0.$$

Random sampling error is due to estimation of p_{2n}^A and p_{2n}^J , which are population moments. Let \hat{p}_{2n}^A and \hat{p}_{2n}^J be estimates of these moments. Then the estimated lower bound on $\pi_{2n, 1r}$ is

$$(2.15) \quad \hat{J}_- = \min_{\psi, m_A, m_J} \psi$$

subject to

$$(2.16a) \quad m_A - m_J - m_A \psi \geq 0,$$

$$(2.16b) \quad \psi \geq 0,$$

$$(2.16c) \quad \psi \in \Psi,$$

$$(2.16d) \quad [n^{1/2}(\hat{p}_{2n}^A - m_A), n^{1/2}(\hat{p}_{2n}^J - m_J)] \in \mathcal{S}.$$

This example is continued in the empirical application of Section 5, where we specify the set \mathcal{S} and find \hat{J}_- satisfying $\mathbb{P}(\pi_{2n, 1r} \geq \hat{J}_-) \geq 0.95$. Since ψ is a probability, we can set $\Psi = [0, 1]$ in this example.

2.4. Analysis. This section presents the main result of the paper, which is a finite-sample lower bound on

$$\mathbb{P} \left[\hat{J}_-(\bar{X}) \leq J_- \leq f(\psi) \leq J_+ \leq \hat{J}_+(\bar{X}) \right].$$

All proofs are in Appendix A. We begin with the following theorem, which forms the basis of our approach.

Theorem 2.1. *Assume that $g_1(\psi, \mu) \leq 0$ and $g_2(\psi, \mu) = 0$ for some ψ . Then*

$$(2.17) \quad \mathbb{P} \left[\hat{J}_-(\bar{X}) \leq J_- \leq f(\psi) \leq J_+ \leq \hat{J}_+(\bar{X}) \right] \geq \mathbb{P} \left[n^{1/2}(\bar{X} - \mu) \in \mathcal{S} \right].$$

Now define

$$Z_i := X_i - \mu \quad \text{and} \quad \bar{Z} := n^{-1/2} \sum_{i=1}^n Z_i = n^{-1/2} \sum_{i=1}^n (X_i - \mu).$$

Then $\mathbb{E}(\bar{Z}) = 0$. Define $\Sigma := \text{cov}(Z_i) = \text{cov}(\bar{Z}) = \text{cov}(X)$. If Σ is non-singular, let $[\Sigma^{-1/2}(X_i - \mu)]_j$ denote the j 'th component of $\Sigma^{-1/2}(X_i - \mu)$ and Σ_{jk}^{-1} denote the (j, k) component of Σ^{-1} . Let μ_j denote the j 'th component of μ and X_{ij} denote the j 'th component of X_i . Make the following assumptions.

Assumption 1. (i) $\{X_i : i = 1, \dots, n\}$ is a random sample from the distribution of X . (ii) \mathcal{S} is compact and convex. (iii) Ψ is compact. (iv) $f(\psi)$ is bounded on Ψ .

Assumption 2. (i) Σ is non-singular, and its components are all finite. (ii) There is a constant $\bar{\mu}_3 < \infty$ such that $\mathbb{E} \left(\left| [\Sigma^{-1/2}(X_i - \mu)]_j \right|^3 \right) \leq \bar{\mu}_3$ for all $i = 1, \dots, n$ and $j = 1, \dots, p$. (iii) There is a constant $C_\Sigma < \infty$ such that $|\Sigma_{jk}^{-1}| \leq C_\Sigma$ for each $j, k = 1, \dots, p$.

Assumption 3. There is a finite constant κ_1 such that

$$(2.18) \quad \begin{aligned} \mathbb{E} [|X_{ij} - \mu_j|^r] &\leq \kappa_1^{r-2} \frac{r!}{2}, \\ \mathbb{E} [(X_{ij} - \mu_j)(X_{ik} - \mu_k) - \Sigma_{jk}]^r &\leq \kappa_1^{r-2} \frac{r!}{2} \end{aligned}$$

for every $r = 2, 3, \dots$ and $j, k = 1, \dots, p$.

Assumption 3 requires the distribution of X to be thin-tailed. We use Assumption 3 to apply Bernstein's inequality (see, e.g., Lemma 14.13 Bühlmann and Van De Geer, 2011).

Suppose for the moment that Σ is known. Define the independent random p -vectors $W_i \sim N(0, \Sigma)$ ($i = 1, \dots, n$) and $\bar{W} := n^{-1/2} \sum_{i=1}^n W_i \sim N(0, \Sigma)$. The multivariate generalization of the Lindeberg-Levy central limit theorem shows that \bar{Z} is asymptotically distributed as $N(0, \Sigma)$, so the distribution of \bar{Z} can be approximated by that of W . The following lemma bounds the error of this approximation.

Lemma 2.1. *Let Assumptions 1, 2(i), and 2(ii) hold. Then*

$$|\mathbb{P}(\bar{Z} \in \mathcal{S}) - \mathbb{P}(W \in \mathcal{S})| \leq \frac{(42p^{1/4} + 16)p^{3/2}\bar{\mu}_3}{n^{1/2}}.$$

In applications, Σ is unknown. Let $\hat{\Sigma}$ be the following estimator of Σ :

$$\hat{\Sigma} := n^{-1} \sum_{i=1}^n X_i X_i' - \bar{X} \bar{X}'.$$

Define the random vector $\widehat{W} \sim N(0, \hat{\Sigma})$. We now approximate the distribution of W by the distribution of \widehat{W} with $\hat{\Sigma}$ treated as a non-stochastic matrix. Define $\mathbf{P}(\mathcal{S}, \Sigma) :=$

$\mathbb{P}(W \in \mathcal{S})$ for $W \sim N(0, \Sigma)$ and

$$(2.19) \quad r_n(t) := 8\sqrt{\frac{2t}{n}}, \quad w_n(t) := C_{\Sigma} p^3 2^{p+1} r_n(t), \quad \text{and} \quad \kappa^* := \min \{ \kappa_1^{-1}, 2\kappa_1^{-2}, 1/2 \}.$$

The following lemma gives a finite-sample bound on the error of the approximation.

Lemma 2.2. *Let Assumptions 1-3 hold and that*

$$(2.20) \quad \frac{9}{4} \log(2p) \leq t \leq \kappa^* n.$$

Then,

$$\left| \mathbf{P}(\mathcal{S}, \widehat{\Sigma}) - \mathbb{P}(W \in \mathcal{S}) \right| \leq w_n(t)$$

with probability at least $1 - 2e^{-t}$.

The condition (2.20) is a mild technical condition that can be satisfied easily. The conclusion of Lemma 2.2 holds, that is,

$$\left| \mathbf{P}(\mathcal{S}, \widehat{\Sigma}) - \mathbb{P}(W \in \mathcal{S}) \right| \leq w_n(t)$$

only if $\widehat{\Sigma}$ satisfies certain conditions that are stated in the proof of the lemma in Appendix A. These conditions are satisfied with probability at least $1 - 2e^{-t}$, not with certainty.

Define

$$(2.21) \quad \delta_n^* := \min_t [w_n(t) + 2e^{-t}] \quad \text{subject to (2.20)}.$$

Now combine Lemmas 2.1 and 2.2 to obtain the following theorem.

Theorem 2.2. *Let Assumptions 1-3 and (A.2) hold. Then,*

$$\left| \mathbb{P}(\bar{Z} \in \mathcal{S}) - \mathbf{P}(\mathcal{S}, \widehat{\Sigma}) \right| \leq \frac{(42p^{1/4} + 16)p^{3/2}\bar{\mu}_3}{n^{1/2}} + \delta_n^*.$$

Theorem 2.2 provides a finite-sample upper bound on the error made by approximating $\mathbb{P}[n^{1/2}(\bar{X} - \mu) \in \mathcal{S}]$ by $\mathbf{P}(\mathcal{S}, \widehat{\Sigma})$. Combining Theorem 2.1 and 2.2 yields

Theorem 2.3. *Assume that $g_1(\psi, \mu) \leq 0$ and $g_2(\psi, \mu) = 0$ for some ψ . Further, let Assumptions 1-3 and (A.2) hold. Then,*

$$(2.22) \quad \begin{aligned} & \mathbb{P} \left[\hat{J}_-(\bar{X}) \leq J_- \leq f(\psi) \leq J_+ \leq \hat{J}_+(\bar{X}) \right] \\ & \geq \mathbf{P}(\mathcal{S}, \widehat{\Sigma}) - \left\{ \frac{(42p^{1/4} + 16)p^{3/2}\bar{\mu}_3}{n^{1/2}} + \delta_n^* \right\}. \end{aligned}$$

Theorem 2.3 provides a finite-sample lower bound on $\mathbb{P} \left[\hat{J}_-(\bar{X}) \leq J_- \leq f(\psi) \leq J_+ \leq \hat{J}_+(\bar{X}) \right]$. Theorems 2.2 and 2.3 are the main results of this paper.

Like other large deviation bounds in statistics and the Berry-Esséen bound, the bounds in Theorems 2.2 and 2.3 can be loose unless n is large because they accommodate worst-case distributions of the observed variables. The numerical performance of our method in less extreme cases is illustrated in Section 4.

2.5. Continuous Covariates. In this section, we consider the case in which g_1 and g_2 depend on a continuous covariate ν in addition to (ψ, μ) . This situation occurs, for example, in applications where some observed variables are group averages and others are continuously distributed characteristics of individuals. If ν is discrete, the results of Section 2.4 apply after replacing problem (2.3)-(2.4) with (2.25)-(2.26) below. When there is a continuous covariate, ν , (2.3)-(2.4) become

$$(2.23) \quad \hat{J}_+(\bar{X}) := \max_{\psi, m} f(\psi) \quad \text{and} \quad \hat{J}_-(\bar{X}) := \min_{\psi, m} f(\psi)$$

subject to

$$(2.24a) \quad g_1(\psi, m, \nu) \leq 0 \quad \text{for every } \nu,$$

$$(2.24b) \quad g_2(\psi, m, \nu) = 0 \quad \text{for every } \nu,$$

$$(2.24c) \quad \psi \in \Psi,$$

$$(2.24d) \quad n^{1/2}(\bar{X} - m) \in \mathcal{S}.$$

Thus, there is a continuum of constraints. We form a discrete approximation to (2.24a)-(2.24b) by restricting ν to a discrete grid of points. Let L denote the number of grid points. We give conditions under which the optimal values of the objective functions of the discretized version of (2.23)-(2.24) converge to $\hat{J}_+(\bar{X})$ and $\hat{J}_-(\bar{X})$ as $L \rightarrow \infty$. To minimize the notational complexity of the following discussion we assume that ν is a scalar. The generalization to a vector is straightforward. We also assume that ν is contained in a compact set which, without further loss of generality, we take to be $[0, 1]$.

To obtain the grid approximation, let $0 = v_0 < v_1 < v_2 < \dots < v_L = 1$ be a grid of equally spaced points in $[0, 1]$. The distance between grid points is $1/(L - 1)$. Approximate problem (2.23)-(2.24) by

$$(2.25) \quad \tilde{J}_+(\bar{X}) := \max_{\psi, m} f(\psi) \quad \text{and} \quad \tilde{J}_-(\bar{X}) := \min_{\psi, m} f(\psi)$$

subject to the constraints:

$$(2.26a) \quad g_1(\psi, m, \nu_\ell) \leq 0; \ell = 1, \dots, L,$$

$$(2.26b) \quad g_2(\psi, m, \nu_\ell) = 0; \ell = 1, \dots, L,$$

$$(2.26c) \quad \psi \in \Psi,$$

and

$$(2.26d) \quad n^{1/2}(\bar{X} - m) \in \mathcal{S}.$$

We then have

Theorem 2.4. *Assume that f is continuous, $\nu \in [0, 1]$, and m in (2.26) is contained in a compact set \mathcal{M} . Moreover,*

$$\begin{aligned} |g_j(\psi, m; x) - g_j(\psi, m; v_\ell)| &\leq C|x - v_\ell| \\ |g_j(\psi, m; x) - g_j(\psi, m; v_{\ell+1})| &\leq C|x - v_{\ell+1}| \end{aligned}$$

for $j = 1$ or 2 , some $C < \infty$, and all $\psi \in \Psi$, all $m \in \mathcal{M}$, all $x \in [v_\ell, v_{\ell+1}] \in [0, 1]$. Then

$$(2.27) \quad \lim_{L \rightarrow \infty} \tilde{J}_+ = \hat{J}_+ \quad \text{and} \quad \lim_{L \rightarrow \infty} \tilde{J}_- = \hat{J}_-.$$

Theorem 2.4 implies that under weak smoothness assumptions, a sufficiently dense grid provides an arbitrarily accurate approximation to the continuously constrained optimization problem (2.23)-(2.24).

2.6. Alternative Approaches. Two natural choices for \mathcal{S} are an ellipsoid and a hypercube or box. The ellipsoid produces a narrower nominal confidence interval for $f(\psi)$ than the hypercube does. The results of Monte Carlo experiments, presented in Section 4, illustrate this. However, the hypercube makes it possible to replace the multivariate Bentkus (2003) inequality with the one-dimensional Berry-Esséen inequality, which reduces the upper bound on the difference between the true and nominal coverage probabilities of the confidence interval. Horowitz (2020) provides details in a different but related setting.

To provide finite-sample bounds on errors in coverage probabilities, we have focused on the simple case that μ is a vector of unknown population means. In applications, it might be useful to extend to a more general case that μ is a vector of population quantities, including the mean as special case. Our inference method provides asymptotic coverage, provided that a suitable \mathcal{S} is constructed via standard asymptotic

theory. It is an open question how to quantify finite-sample coverage probabilities in such a case.

3. COMPUTATIONAL ALGORITHMS

Recall that our general framework is to obtain the bound

$$[\min_{\psi, m} f(\psi), \max_{\psi, m} f(\psi)]$$

subject to

$$g_1(\psi, m) \leq 0, g_2(\psi, m) = 0, \psi \in \Psi, \text{ and } n^{1/2}(\bar{X} - m) \in \mathcal{S}.$$

3.1. Objective function $f(\psi)$. In many examples, $f(\psi)$ is linear in ψ . For example, ψ is the vector of all the parameters in an econometric model and $f(\psi)$ is just one element of ψ or a linear combination of elements of ψ .

3.2. Restrictions $g_1(\psi, \mu) \leq 0, g_2(\psi, \mu) = 0$, and $\psi \in \Psi$. The restrictions $g_1(\psi, \mu) \leq 0$ include shape restrictions among the elements of ψ . Equality restrictions are imposed via $g_2(\psi, \mu) = 0$. The easiest case is that $g_j(\psi, \mu)$ is linear in (ψ, μ) for each $j = 1, 2$. In some of examples we consider, $g_j(\psi, \mu)$ is linear in ψ , holding μ fixed, and linear in μ , keeping ψ fixed, but not linear in (ψ, μ) jointly. This corresponds to the case of bilinear constraints. For example, $g_j(\psi, \mu)$ may depend on the product between one of elements of ψ and one of elements of μ . In practice, Ψ can always be chosen large enough that the constraint $\psi \in \Psi$ is not binding additionally and can be ignored. For example, suppose that ψ is a probability and the constraints in $g_1(\psi, \mu) \leq 0$ and $g_2(\psi, \mu) = 0$ impose a restriction on ψ such as $[a, b]$ for $0 \leq a < b \leq 1$. Then, it is not necessary to impose $\psi \in \Psi = [0, 1]$ additionally.

3.3. Restrictions $n^{1/2}(\bar{X} - \mu) \in \mathcal{S}$. There are two leading cases of \mathcal{S} : an ellipsoid and a box. We start with the case that \mathcal{S} is a box (that is, the Cartesian product of intervals). Let \hat{D} denote the diagonal matrix consisting of diagonal elements of $\hat{\Sigma}$. Choose $\kappa(1 - \alpha)$ such that

$$\sqrt{n} \max \left\{ \left| \hat{D}_j^{-1/2}(\bar{X}_j - \mu_j) \right| : j = 1, \dots, 2J \right\} \leq \kappa(1 - \alpha)$$

with probability $1 - \alpha$. Here, the subscript j denotes the j -th element of a vector or the (j, j) element of a diagonal matrix. Note that when \mathcal{S} is a box, the critical value can be easily simulated from the $N(0, \hat{\Sigma})$ and the restriction $n^{1/2}(\bar{X} - \mu) \in \mathcal{S}$ can be written as linear constraints.

Consider now the case that \mathcal{S} is an ellipsoid. Choose $\kappa(1 - \alpha)$ such that

$$n(\bar{X} - \mu)' \widehat{\Sigma}^{-1} (\bar{X} - \mu) \leq \kappa(1 - \alpha)$$

with probability $1 - \alpha$.

When \mathcal{S} is an ellipsoid, the critical value $\kappa(1 - \alpha)$ can be obtained from the $\chi^2(J)$ distribution, where J is the dimension of μ . Then, the restriction $n^{1/2}(\bar{X} - \mu) \in \mathcal{S}$ can be written as

$$\mu' \widehat{\Sigma}^{-1} \mu - 2\mu' \widehat{\Sigma}^{-1} \bar{X} \leq n^{-1} \kappa(1 - \alpha) - \bar{X}' \widehat{\Sigma}^{-1} \bar{X}.$$

This is a convex quadratic constraint in μ .

3.4. Mathematical programming for leading cases. Table 1 gives the scheme of mathematical programming we use for leading cases of $f(\psi)$, $g_1(\psi, \mu) \leq 0$, $g_2(\psi, \mu) = 0$, and $n^{1/2}(\bar{X} - \mu) \in \mathcal{S}$. In the table, LP, QP and QCP refer to linear programming, quadratic programming, and quadratically constrained programming, respectively. MILP, MIQP and MIQCP correspond to mixed integer linear programming, mixed integer quadratic programming, and mixed integer quadratically constrained programming, respectively.

TABLE 1. Class of Optimization Problems

| Case | $f(\psi)$ | $g_1(\psi, \mu) \leq 0$ $g_2(\psi, \mu) = 0$ | $n^{1/2}(\bar{X} - \mu) \in \mathcal{S}$ | Programming |
|------|-----------|---|--|-------------|
| 1 | linear | linear | box | LP |
| 2 | linear | linear | ellipsoid | QCP |
| 3 | quadratic | linear | box | QP |
| 4 | quadratic | linear | ellipsoid | QCP |
| 5 | linear | bilinear | box | MILP/LP |
| 6 | linear | bilinear | ellipsoid | MIQCP/QCP |
| 7 | quadratic | bilinear | box | MIQP/QP |
| 8 | quadratic | bilinear | ellipsoid | MIQCP/QCP |

When some of the constraints $g_1(\psi, \mu) \leq 0$ and $g_2(\psi, \mu) = 0$ are bilinear, the resulting problem may not be convex. To deal with non-convexity, we rely on a sequence of convex relaxations to obtain an outer bound for $f(\psi)$ and use a set of restricted inner bounds. When the union of restricted inner bounds matches the best outer bound by convex relaxations, we obtain the exact solution to the problem. Even if they do not match exactly, the best outer and inner bounds will give an approximate solution to the problem. The convex relaxations for bilinear constraints are implemented using

mixed integer optimization (MIO). In Case 5, MILP/LP refers to the use of MILP for the outer bound and that of LP for the inner bound. Cases 6-8 are similar. Appendix B gives a detailed description of dealing with bilinear constraints. By virtue of the developments in MIO solvers and fast computing environments, the MIO has become increasingly used in recent applications. For example, Bertsimas, King, and Mazumder (2016) adopted an MIO approach for obtaining ℓ_0 -constrained estimators in high-dimensional regression models and Reguant (2016) used mixed integer linear programming for computing counterfactual outcomes in game theoretic models.

4. MONTE CARLO EXPERIMENTS

4.1. **Identification Problem.** Suppose that

$$(4.1) \quad Y_i^* = h(Z_i) + e_i,$$

where $h : \mathbb{R} \mapsto \mathbb{R}$ is an unknown function and the error term e_i satisfies $E[e_i|Z_i] = 0$ almost surely. Assume that for each individual i , we do not observe Y_i^* , but only the interval data $[L_i, U_i]$ such that $Y_i^* \in [L_i, U_i]$ along with Z_i . Here, L_i and U_i are random variables.

Assume that the support of Z_i is finite, that is, $\mathcal{Z} \in \{z_1, \dots, z_J\}$. Denote the values of $h(\cdot)$ on \mathcal{Z} by $\{\psi_1, \dots, \psi_J\}$. That is, $h(u) = \sum_{j=1}^J \psi_j 1(u = z_j)$ for $u \in \mathcal{Z}$.

Suppose that the object of interest is the value of $\psi^* \equiv h(z^*)$, where z^* is not in the support of Z_i but $z_{j-1} < z^* < z_j$ for some j . This type of extrapolation problem is given as a motivating example in Manski (2007a, pp. 4-5).

To partially identify ψ^* , assume that $h(\cdot)$ is monotone non-decreasing. Specifically, we impose the monotonicity on $\mathcal{Z} \cup \{z^*\}$. That is, $h(z_1) \leq h(z_2)$ whenever $z_1 \leq z_2$ for any $z_1, z_2 \in \mathcal{Z} \cup \{z^*\}$. In addition, we have the following inequality constraints:

$$(4.2) \quad E[L_i|Z_i = z_j] \leq \psi_j \leq E[U_i|Z_i = z_j]$$

for any $1 \leq j \leq J$. Note that (4.2) alone does not provide a bounded interval on ψ^* since z^* is not in \mathcal{Z} . The monotonicity assumption combined with (4.2) provides an informative bound on ψ^* .

To write the optimization problem in our canonical form, let μ denote the population moments of the following \bar{X} :

$$n\bar{X} = \begin{pmatrix} \sum_{i=1}^n L_i 1(Z_i = z_1) \\ \vdots \\ \sum_{i=1}^n L_i 1(Z_i = z_J) \\ \sum_{i=1}^n U_i 1(Z_i = z_1) \\ \vdots \\ \sum_{i=1}^n U_i 1(Z_i = z_J) \\ \sum_{i=1}^n 1(Z_i = z_1) \\ \vdots \\ \sum_{i=1}^n 1(Z_i = z_J) \end{pmatrix}.$$

Then, we can rewrite the constraints (4.2) in a bilinear form:

$$(4.3) \quad E[L_i 1(Z_i = z_j)] \leq \psi_j E[1(Z_i = z_j)] \leq E[U_i 1(Z_i = z_j)]$$

for any $1 \leq j \leq J$. To deal with the bilinear constraints, we rely on a method called piecewise McCormick relaxation, which is given in Appendix B.

4.2. Results of a Monte Carlo Experiment. Suppose that (4.1) holds with $h(z) = 2z$, the covariate Z_i is uniformly distributed on

$$\mathcal{Z} = \{-3/2, -1, -1/2, 1/2, 1, 3/2\},$$

and $e_i \sim \text{Unif}[-1/2, 1/2]$. The interval data are generated from $L_i = Y_i^* + V_i 1(V_i < 0)$ and $U_i = Y_i^* + V_i 1(V_i \geq 0)$, where $V_i \sim N(0, 1)$. Here, V_i and e_i in (4.1) are independent of each other. The parameter of interest is $\psi^* = h(0)$. Note that zero is not included in \mathcal{Z} . The monotonicity constraint is imposed as

$$\begin{pmatrix} 1 & 0 & 0 & 0 & -1 & 0 & 0 & 0 \\ -1 & 0 & 0 & 1 & 0 & 0 & 0 & 0 \\ 0 & 1 & -1 & 0 & 0 & 0 & 0 & 0 \\ 0 & 0 & 1 & -1 & 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 & 1 & -1 & 0 & 0 \\ 0 & 0 & 0 & 0 & 0 & 1 & -1 & 0 \end{pmatrix} \begin{pmatrix} h(0) \\ h(-3/2) \\ h(-1) \\ h(-1/2) \\ h(1/2) \\ h(1) \\ h(3/2) \end{pmatrix} \leq \mathbf{0}_{6 \times 1}.$$

The simulation design here is similar to that of Bontemps, Magnac, and Maurin (2012) except that the support of X_i is discrete and the linearity of $h(\cdot)$ is not used in

estimation. The sample size was 200, 500, 1000 and 2000. There were 100 repetitions for each Monte Carlo experiment. The identified set containing $h(0)$ is $-1.40 \leq h(0) \leq 1.40$. The reported coverage probability is the frequency that the estimated lower bound is smaller than or equal to the true lower bound (-1.40) and the estimated upper bound is greater than or equal to the true upper bound (1.40). The nominal coverage probability was 0.95.

We consider both cases that \mathcal{S} is a box and an ellipsoid. We take $\Psi = \{\psi_j \in [-5, 5] \text{ for each } j = 1, \dots, J\}$. The outer bounds were computed with piecewise linear relaxations with $K = 10$ that is described in Appendix B. To describe how to obtain inner bounds, first partition μ and \bar{X} into $\mu \equiv (\mu_1, \mu_2)$ and $\bar{X} \equiv (\bar{X}_1, \bar{X}_2)$, where $\mu_1 \equiv (E[L_i 1(Z_i = z_1)], \dots, E[L_i 1(Z_i = z_J)], E[U_i 1(Z_i = z_1)], \dots, E[U_i 1(Z_i = z_J)])$, $\mu_2 \equiv (E[1(Z_i = z_1)], \dots, E[1(Z_i = z_J)])$, and \bar{X}_1 and \bar{X}_2 are corresponding sample moments. In other words, only components of μ_2 appear as bilinear terms in (4.3) and those of μ_1 are linearly separable. Then, to obtain a lower bound, we fix μ_2 at its feasible value and optimize with respect to (ψ, μ_1) . When \mathcal{S} is a box, it is straightforward to obtain a lower bound since the feasible value of μ_1 does not depend on that of μ_2 . When \mathcal{S} is an ellipsoid, recall that the restriction $n^{1/2}(\bar{X} - \mu) \in \mathcal{S}$ can be written as

$$(4.4) \quad \mu' \hat{\Sigma}^{-1} \mu - 2\mu' \hat{\Sigma}^{-1} \bar{X} \leq n^{-1} \kappa (1 - \alpha) - \bar{X}' \hat{\Sigma}^{-1} \bar{X}.$$

Following μ and \bar{X} , partition $\hat{\Sigma}^{-1}$ into blocks such that $\hat{\Sigma}^{-1} \equiv \{\hat{\Sigma}_{(k,\ell)}^{-1}, k = 1, 2, \ell = 1, 2\}$. Now rewrite (4.4) as

$$(4.5) \quad \begin{aligned} & \mu'_1 \hat{\Sigma}_{(1,1)}^{-1} \mu_1 + 2\mu'_1 \left[\hat{\Sigma}_{(1,2)}^{-1} \mu_2 - \hat{\Sigma}_{(1,1)}^{-1} \bar{X}_1 - \hat{\Sigma}_{(1,2)}^{-1} \bar{X}_2 \right] \\ & \leq n^{-1} \kappa (1 - \alpha) - \bar{X}' \hat{\Sigma}^{-1} \bar{X} - \mu'_2 \hat{\Sigma}_{(2,2)}^{-1} \mu_2 + 2\mu'_2 \left[\hat{\Sigma}_{(2,1)}^{-1} \bar{X}_1 + \hat{\Sigma}_{(2,2)}^{-1} \bar{X}_2 \right]. \end{aligned}$$

Given μ_2 , this is a convex, quadratic constraint. First, we generate a random grid of μ_2 using the box version of \mathcal{S} . Then we optimize with respect to (ψ, μ_1) under the restrictions $g(\psi, \mu) \leq 0$ and (4.5). In both cases, the inner bounds were obtained with a random grid search with $G = 1000$. As an alternative to the outer and inner bounds, we also consider the bounds when μ_2 is fixed at \bar{X}_2 . These correspond to the bounds when the observed covariates Z_1, \dots, Z_n are regarded as non-stochastic. These bounds will be tighter than those constructed under the random design. In each Monte Carlo repetition when \mathcal{S} is a box, the number of simulations to draw $N(0, \hat{\Sigma})$ was 1,000. The χ^2 critical value is used when \mathcal{S} is an ellipsoid.

TABLE 2. Results of Monte Carlo Experiments

| Type of \mathcal{S} | Type of bounds | sample size | Avg. of estimated lower bound | Avg. of estimated upper bound | Coverage probability |
|-----------------------|--------------------------------|-------------|-------------------------------|-------------------------------|----------------------|
| Box | Outer bounds ($K = 10$) | 200 | -3.87 | 3.78 | 1 |
| | | 500 | -2.59 | 2.54 | 1 |
| | | 1000 | -2.14 | 2.14 | 1 |
| | | 2000 | -1.89 | 1.89 | 1 |
| Box | Inner bounds ($G = 1000$) | 200 | -3.89 | 3.75 | 1 |
| | | 500 | -2.59 | 2.54 | 1 |
| | | 1000 | -2.14 | 2.14 | 1 |
| | | 2000 | -1.89 | 1.89 | 1 |
| Box | Fixing Z_1, \dots, Z_n | 200 | -2.07 | 2.12 | 1 |
| | | 500 | -1.83 | 1.84 | 1 |
| | | 1000 | -1.71 | 1.69 | 1 |
| | | 2000 | -1.62 | 1.61 | 1 |
| Ellipsoid | Outer bounds ($K = 10$) | 200 | -2.34 | 2.35 | 1 |
| | | 500 | -1.90 | 1.91 | 1 |
| | | 1000 | -1.78 | 1.76 | 1 |
| | | 2000 | -1.67 | 1.66 | 1 |
| Ellipsoid | Inner bounds ($G = 1000$) | 200 | -2.25 | 2.19 | 1 |
| | | 500 | -1.84 | 1.81 | 1 |
| | | 1000 | -1.68 | 1.68 | 1 |
| | | 2000 | -1.59 | 1.59 | 1 |
| Ellipsoid | Fixing Z_1, \dots, Z_n | 200 | -1.92 | 1.93 | 0.99 |
| | | 500 | -1.72 | 1.73 | 1 |
| | | 1000 | -1.64 | 1.62 | 1 |
| | | 2000 | -1.57 | 1.56 | 1 |

Note: “Box” and “Ellipsoid” are our proposed methods with a box S and an ellipsoidal S , respectively.

Table 2 presents the simulation results. First, we comment on the results with a box. When $n = 200$, there are minor discrepancies between the outer and inner bounds. For all other large sample sizes, averages of the bounds are identical. As the sample size increases, the length of the estimated bounds decreases rapidly. We now look at the results with an ellipsoid. The estimated bounds are much tighter with an ellipsoid than with a box. There are more noticeable differences between the average values of the outer and inner bounds when \mathcal{S} is an ellipsoid. However, these differences shrink as the sample size gets larger. The bounds with fixed Z_i 's are tighter especially when n is relatively small or if a box is used for \mathcal{S} . The empirical coverage

probabilities are all larger than the nominal probability of 0.95. This is consistent with Theorem 2.3, which provides a lower bound on the coverage probability, not a point probability.

4.3. Comparison with Minsker’s Method. In this subsection, we provide the results of a small Monte Carlo experiment that is designed to compare our main proposal with Minsker’s method described in Appendix C. We make two changes in the experimental design of Section 4.2. First, in addition to the standard normal distribution, the errors V_i are generated from the t -distribution with degrees of freedom equal to 3 to consider a fat-tailed distribution. Note that the fat-tailed distribution is adopted here since the median-of-means approach is robust to outliers. Second, we set $n = 10000$ or 20000 because Minsker’s method requires a relatively large sample size. The estimated bounds obtained with Minsker’s method were uninformative when $n = 200, 500, 1000, 2000$. Specifically, they were $[-5, 5]$ for all of these sample sizes. We consider coordinate-wise medians for the median of means since μ is relatively low-dimensional. For simplicity, we consider only the outer bounds with 10 repetitions for each experiment.

TABLE 3. Comparison with Minsker’s Method

| Distribution of V_i | Type of \mathcal{S} | Type of bounds | sample size | Avg. of estimated lower bound | Avg. of estimated upper bound | Coverage probability |
|-----------------------|-----------------------|----------------|-------------|-------------------------------|-------------------------------|----------------------|
| $N(0, 1)$ | Ellipsoid | Outer bounds | 10000 | -1.52 | 1.53 | 1 |
| | | | 20000 | -1.49 | 1.48 | 1 |
| $N(0, 1)$ | Minsker | Outer bounds | 10000 | -3.87 | 3.86 | 1 |
| | | | 20000 | -3.20 | 3.22 | 1 |
| $t(3)$ | Ellipsoid | Outer bounds | 10000 | -1.71 | 1.73 | 1 |
| | | | 20000 | -1.67 | 1.67 | 1 |
| $t(3)$ | Minsker | Outer bounds | 10000 | -4.19 | 4.17 | 1 |
| | | | 20000 | -3.48 | 3.46 | 1 |

Note: “Ellipsoid” is our proposed method with an ellipsoidal S .

Table 3 reports the experimental results. The true lower and upper bounds are -1.40 and 1.40 for the $N(0, 1)$ errors and -1.37 and 1.37 for the $t(3)$ errors, respectively. The estimated bounds for Minsker’s method are much wider than the bounds estimated by our proposed method, although they shrink from $n = 10000$ to

$n = 20000$. Moreover, Minsker’s method does not produce a better result for the t -distribution. Even with the t -distribution, Minsker’s method gives wider bounds than our method does. The bounds from the ellipsoid, which is our proposed method, provides much tighter bounds but they also seem conservative, as noted in the previous section.

5. EMPIRICAL EXAMPLES

5.1. Bounding the Average of Log Weekly Wages. The example in this subsection mimics the Monte Carlo Experiments. Suppose that Y_i^* is the log weekly wage. In the 1960 and 1970 US censuses, hours of work during a week are measured in brackets. As a result, weekly wages are only available in terms of interval data. Let Z_i be the years of schooling.

We use a sample extract from the US 2000 census to illustrate our methods. The sample is restricted to males who were 40 years old with educational attainment at least grade 10, positive wages and positive hours of work. In the 2000 census, variable `WKSWORK1` is the hours of work in integer value, whereas `WKSWORK2` is weeks worked in intervals. As a parameter of interest, we focus on the average log wage for one year of college. In our estimation procedure, we drop all observations with those with one year of college. The resulting sample size is $n = 15,647$.

TABLE 4. 95% Confidence Interval for the Average of Log Weekly Wages

| method | lower bound | upper bound |
|---------------------------------|-------------|-------------|
| outer box | 6.10 | 7.66 |
| inner box | 6.11 | 7.62 |
| inner box covariate fixed | 6.24 | 7.17 |
| outer ellipsoid | 6.32 | 6.84 |
| inner ellipsoid | 6.32 | 6.83 |
| inner ellipsoid covariate fixed | 6.34 | 6.81 |
| oracle | 6.56 | 6.61 |

In estimation, the latent Y_i^* is $\log(\text{INCWAGE}/\text{WKSWORK1})$, where `INCWAGE` is total pre-tax wage and salary income and `WKSWORK1` is weeks worked in integer value. We observe brackets $[L_i, U_i]$ of $\log(\text{INCWAGE}/\text{WKSWORK2})$, where `WKSWORK2` is weeks worked in intervals: $[1,13]$, $[14,26]$, $[27,39]$, $[40,47]$, $[48,49]$, $[50,52]$. The brackets $[L_i, U_i]$ are random because the numerator `INCWAGE` is random. The set \mathcal{Z} includes 7 different values: 10th grade, 11th grade, 12th grade, 1 year of college, 2 years of college, 4 years

of college and 5 or more years of college. Since we drop observations with those with 1 year of college, we have a missing data problem as in the Monte Carlo experiment.

Table 4 shows estimation results. The nominal level of coverage is 0.95. The oracle refers to the confidence interval of the average log wage for one year of college using the actual value of weekly wages for those with one year of college. As in the Monte Carlo experiments, ellipsoid-based confidence intervals are tighter than box-type confidence intervals and fixing covariates tightens the confidence intervals as well. All the ellipsoid intervals are only slightly wider than the oracle interval.

5.2. Kline and Tartari (2016) Revisited. This section provides an empirical example based on the study of KT that is described in Section 2.3. Specifically, we use the information in Table 4 of KT to obtain the set \mathcal{S} in (2.16d) and the lower endpoint of a 95% confidence bound for $\pi_{2n,1r}$. We consider only the lower endpoint because KT found the upper endpoint to be 1 and, therefore, uninformative.

KT used the JF welfare reform experimental data and pooled all person-quarter observations in the seven quarters following randomization of participants. They treated each person-quarter observation as a potentially separate decision, allowing time-varying behaviors. Because assignment of individuals to the JF treatment and AFDC control groups was random, we assume that observations in each regime are independent of the observations in the other. We further assume that observations within the JF and AFDC regimes are independently and identically distributed (iid). The set \mathcal{S} can be expressed as a confidence region for p_{2n}^A and p_{2n}^J . We used the Clopper–Pearson (1934) procedure to construct the rectangular 95% confidence region

$$\begin{aligned} L_A &\leq p_{2n}^A \leq U_A, \\ L_J &\leq p_{2n}^J \leq U_J, \end{aligned}$$

where L_A, U_A, L_J , and U_J are random lower and upper bounds chosen so that

$$\mathbb{P}(L_A \leq p_{2n}^A \leq U_A; L_J \leq p_{2n}^J \leq U_J) \geq 0.95.$$

The sample sizes in KT are $n_A = 16,268$ and $n_J = 16,226$ for the AFDC and JF regimes, respectively. The estimated values of p_{2n}^A and p_{2n}^J are $\hat{p}_{2n}^A = 0.099$ and $\hat{p}_{2n}^J = 0.068$. The resulting confidence region for p_{2n}^A and p_{2n}^J is $0.092 \leq p_{2n}^A \leq 0.106$ and $0.062 \leq p_{2n}^J \leq 0.074$. Solving (2.15)-(2.16) yielded 0.195 as the lower endpoint of a 95% confidence region for $\pi_{2n,1r}$. KT obtained a lower endpoint of 0.198, which is similar to ours. However, there are important differences between our method and that of KT. KT used a block bootstrap procedure that resamples a woman's

entire profile of choices for the first seven quarters after randomization, whereas we used the non-asymptotic inference method for iid data. The method of KT relies on asymptotic arguments but allows serial dependence of the observations of the same woman. Our method is valid for any sample size but does not take account of any serial dependence. It turns out that the conservative nature of our inference method is offset by our assumption of independence, thereby yielding a confidence interval of approximately the same size as that of KT.

6. CONCLUSIONS

This paper has described a method for carrying out inference on partially identified parameters that are solutions to a class of optimization problems. The parameters arise, for example, in applications in which grouped data are used for estimation of a model's structural parameters. Inference consists of obtaining confidence intervals for the structural parameters. The method of this paper provides a finite-sample upper bound on the difference between the true and nominal probabilities with which a confidence interval based on an asymptotic approximation contains the true but unknown value of a parameter. The paper has described computational algorithms for implementing the method. The results of Monte Carlo experiments and empirical examples illustrate the method's usefulness.

APPENDIX A. PROOFS OF THEOREMS

Proof of Theorem 2.1. Suppose $n^{1/2}(\bar{X} - \mu)$ is in \mathcal{S} . Any feasible solution of (2.3)–(2.4) is also a feasible solution of (2.1)–(2.2). Therefore, the feasible region of (2.1)–(2.2) contains the feasible region of (2.3)–(2.4). Consequently,

$$\hat{J}_-(\bar{X}) \leq J_- \leq J_+ \leq \hat{J}_+(\bar{X}),$$

which in turn proves (2.17). □

Proof of Lemma 2.1. Define the random p -vector $\bar{V} := \Sigma^{-1/2}\bar{Z}$. Then $\mathbb{E}(\bar{V}) = 0$ and $\text{cov}(\bar{V}) = I_{p \times p}$. Define the set

$$(A.1) \quad \mathcal{S}_\Sigma := \{\xi : \xi = \Sigma^{-1/2}\zeta; \zeta \in \mathcal{S}\}.$$

Then \mathcal{S}_Σ is a convex set. Define the random vectors $U_i \sim N(0, I_{p \times p})$ and $\bar{U} := n^{1/2} \sum_{i=1}^n U_i$. It follows from Assumption 2(ii) and the generalized Minkowski inequality that

$$\mathbb{E} \left[\left(\sum_{j=1}^p [\Sigma^{-1/2}(X_i - \mu)]_j^2 \right)^{3/2} \right] \leq p^{3/2} \bar{\mu}_3.$$

In addition, it follows from Theorem 1.1 of Raič (2019) that

$$|\mathbb{P}(\bar{V} \in \mathcal{S}_\Sigma) - \mathbb{P}(\bar{U} \in \mathcal{S}_\Sigma)| \leq \frac{(42p^{1/4} + 16)p^{3/2}\bar{\mu}_3}{n^{1/2}},$$

which proves the lemma. \square

Proof of Lemma 2.2. Define

$$\Delta_n := \sup_{\mathcal{S}} \left| \mathbf{P}(\mathcal{S}, \widehat{\Sigma}) - \mathbb{P}(W \in \mathcal{S}) \right|.$$

Now

$$\mathbf{P}(\mathcal{S}, \widehat{\Sigma}) - \mathbb{P}(W \in \mathcal{S}) = \mathbb{P}(\Sigma^{-1/2}\widehat{W} \in \mathcal{S}_\Sigma) - \mathbb{P}(\xi \in \mathcal{S}_\Sigma),$$

where $\xi \sim N(0, I_{p \times p})$ and \mathcal{S}_Σ is defined in (A.1). Therefore,

$$\begin{aligned} \Delta_n &= \sup_{\mathcal{S}} \left| \mathbb{P}(\Sigma^{-1/2}\widehat{W} \in \mathcal{S}_\Sigma) - \mathbb{P}(\xi \in \mathcal{S}_\Sigma) \right| \\ &\leq \text{TV} \left[N(0, I_{p \times p}), N(0, \Sigma^{-1}\widehat{\Sigma}) \right], \end{aligned}$$

where $\text{TV}(P_1, P_2)$ is the total variation distance between distributions P_1 and P_2 . By Example 2.3 of Dasgupta (2008),

$$\text{TV} \left[N(0, I_{p \times p}), N(0, \Sigma^{-1}\widehat{\Sigma}) \right] \leq p2^{p+1} \left\| \Sigma^{-1}\widehat{\Sigma} - I_{p \times p} \right\|_{\text{F}}$$

where for any matrix A ,

$$\|A\|_{\text{F}}^2 := \sum_{j=1}^p \sum_{k=1}^p a_{jk}^2.$$

Define $\omega := \widehat{\Sigma} - \Sigma$. Then,

$$\Sigma^{-1}\widehat{\Sigma} - I_{p \times p} = \Sigma^{-1}(\widehat{\Sigma} - \Sigma) = \Sigma^{-1}\omega,$$

$$|(\Sigma^{-1}\omega)_{jk}| \leq \sum_{\ell=1}^p |\Sigma_{j\ell}^{-1}\omega_{\ell k}| \leq C_\Sigma \sum_{\ell=1}^p |\omega_{\ell k}|$$

and

$$\left\| \Sigma^{-1}\widehat{\Sigma} - I_{p \times p} \right\|_{\text{F}} \leq C_\Sigma p^{1/2} \left[\sum_{k=1}^p \left(\sum_{\ell=1}^p |\omega_{\ell k}| \right)^2 \right]^{1/2}.$$

To obtain the conclusion of the lemma, it remains to show that $|\omega_{jk}| \leq r_n(t)$ with probability at least $1 - 2e^{-t}$. We prove this claim below.

Write

$$\omega = n^{-1} \sum_{i=1}^n [(X_i - \mu)(X_i - \mu)' - \Sigma] - (\bar{X} - \mu)(\bar{X} - \mu)'.$$

Then, we have that

$$\begin{aligned} & \mathbb{P} \left[\max_{j,k} |\omega_{jk}| \geq r_n(t) \right] \\ & \leq \mathbb{P} \left[\max_{1 \leq j,k \leq p} \left| \sum_{i=1}^n [(X_{ij} - \mu)(X_{ik} - \mu)' - \Sigma_{jk}] \right| \geq \frac{r_n(t)}{2} \right] + \mathbb{P} \left[\left\{ \max_{1 \leq j \leq p} |\bar{X}_j - \mu_j| \right\}^2 \geq \frac{r_n(t)}{2} \right]. \end{aligned}$$

Define

$$\lambda(\kappa_1, n, p) := \sqrt{\frac{2 \log(2p)}{n}} + \frac{\kappa_1 \log(2p)}{n}.$$

By Bernstein's inequality for the maximum of p averages (see, e.g., Lemma 14.13 Bühlmann and Van De Geer, 2011),

$$\mathbb{P} \left[\max_{1 \leq j \leq p} |\bar{X}_j - \mu_j| \geq \frac{\kappa_1 t}{n} + \sqrt{\frac{2t}{n}} + \lambda(\kappa_1, n, p) \right] \leq \exp(-t)$$

and

$$\mathbb{P} \left[\max_{1 \leq j,k \leq p} \left| \sum_{i=1}^n [(X_{ij} - \mu)(X_{ik} - \mu)' - \Sigma_{jk}] \right| \geq \frac{\kappa_1 t}{n} + \sqrt{\frac{2t}{n}} + \lambda(\kappa_1, n, p^2) \right] \leq \exp(-t).$$

Suppose that

$$(A.2) \quad \max\{t, \log(2p)\} \leq \min \left\{ \frac{n}{\kappa_1}, \frac{n}{2} \right\}.$$

Under (A.2),

$$\begin{aligned} & \mathbb{P} \left[\left\{ \max_{1 \leq j \leq p} |\bar{X}_j - \mu_j| \right\}^2 \geq \frac{\kappa_1 t}{n} + \sqrt{\frac{2t}{n}} + \lambda(\kappa_1, n, p) \right] \\ & \leq \mathbb{P} \left[\left\{ \max_{1 \leq j \leq p} |\bar{X}_j - \mu_j| \right\}^2 \geq \left\{ \frac{\kappa_1 t}{n} + \sqrt{\frac{2t}{n}} + \lambda(\kappa_1, n, p) \right\}^2 \right] \\ & \leq \exp(-t). \end{aligned}$$

Therefore, if we can choose $r_n(t)$ by

$$(A.3) \quad r_n(t) := 2\sqrt{\frac{2t}{n}} + 3\sqrt{\frac{2\log(2p)}{n}} + \frac{2\kappa_1 t}{n} + \frac{3\kappa_1 \log(2p)}{n}$$

for $t > 0$, we have that

$$\mathbb{P} \left[\max_{j,k} |\omega_{jk}| \leq r_n(t) \right] \geq 1 - 2e^{-t}.$$

To simplify the upper bound $r_n(t)$, we now strengthen (A.2) to (2.20) stated in the lemma. Then, we can take

$$r_n(t) := 8\sqrt{\frac{2t}{n}},$$

which proves the claim. \square

Proof of Theorem 2.2. Write

$$\begin{aligned} \left| \mathbb{P}(\bar{Z} \in \mathcal{S}) - \mathbf{P}(\mathcal{S}, \widehat{\Sigma}) \right| &= \left| [\mathbb{P}(\bar{Z} \in \mathcal{S}) - \mathbb{P}(W \in \mathcal{S})] - [\mathbf{P}(\mathcal{S}, \widehat{\Sigma}) - \mathbb{P}(W \in \mathcal{S})] \right| \\ &\leq \left| \mathbb{P}(\bar{Z} \in \mathcal{S}) - \mathbb{P}(W \in \mathcal{S}) \right| + \left| \mathbf{P}(\mathcal{S}, \widehat{\Sigma}) - \mathbb{P}(W \in \mathcal{S}) \right|. \end{aligned}$$

Thus, the theorem follows immediately by combining Lemmas 2.1 and 2.2. \square

Proof of Theorem 2.3. Combining Theorem 2.1 and 2.2 yields Theorem 2.3. \square

Proof of Theorem 2.4. We focus on the maximization problem since the minimization problem can be analyzed analogously.

Let (ψ_L, μ_L) denote the optimal solution to the maximization version of (2.25)-(2.26). Define $g(\psi, \mu; \nu) = [g_1(\psi, \mu; \nu), g_2(\psi, \mu; \nu), -g_2(\psi, \mu; \nu)]$, so that $g(\psi, \mu; \nu) \leq 0$ componentwise. Define $\ell(\nu) := \arg \min_{\ell} |\nu - \nu_{\ell}|$. Then

$$\sup_{\nu \in [0,1]} |g(\psi, \mu; \nu) - g(\psi, \mu; \ell(\nu))| \leq C/(L-1)$$

and $g(\psi, \mu; \nu_{\ell}) \leq 0$ implies that

$$g(\psi, \mu; \nu) \leq C/(L-1)$$

componentwise uniformly over $\nu \in [0, 1]$. Therefore, (ψ_L, μ_L) is a feasible solution to

$$(A.4) \quad J_+^* := \max_{\psi, \mu} f(\psi)$$

subject to the new constraint:

$$g(\psi, \mu; \nu) \leq C/(L-1) \text{ for all } \nu \in [0, 1], \nu = \text{rational, and } n^{1/2}(\bar{X} - \mu) \in \mathcal{S}.$$

Consequently, $J_+^* \geq \tilde{J}_+ \geq J_+$, where $J_+ = \max f(\psi)$ subject to $g_1(\psi, \mu, x) \leq 0$, $g_2(\psi, \mu, x) = 0$, and $\psi \in \Psi$. Define

$$\Pi := \left\{ \xi \geq 1, \eta : \text{there is } (\psi, \mu) \text{ such that } n^{1/2}(\bar{X} - \mu) \in \mathcal{S}, f(\psi) \leq \eta, \right. \\ \left. \text{and } g(\varphi, \mu; \nu) \leq C/\xi \text{ for all } \nu \in [0, 1] \right\}.$$

Note that Π is a closed set. Therefore, by Proposition 3.3 of Jeyakumar and Wolkowicz (1990), $J_+^* \rightarrow J_+$ as $L \rightarrow \infty$ if the constraints are restricted to rational values of $\nu \in [0, 1]$. It follows from continuity of g as a function of ν that the constraints hold for all $\nu \in [0, 1]$. \square

APPENDIX B. DETAILS ABOUT COMPUTATION WITH BILINEAR CONSTRAINTS

To explain how to deal with the constraints in a bilinear form, suppose that we have a cross product term $\mu_j \psi_\ell$ in $g(\psi, \mu) \leq 0$ for some j and ℓ , where $\mu = (\mu_1, \dots, \mu_J)'$ and $\psi = (\psi_1, \dots, \psi_L)'$.

The existence of the bilinear term $\mu_j \psi_\ell$ can make the corresponding optimization problem non-convex. As mentioned in the main text, we rely on a sequence of convex relaxations to obtain an outer bound for $f(\psi)$. Specifically, we use piecewise-linear relaxations that are called piecewise McCormick relaxation in the operation research and engineering literature.

There exist a number of different formulations for piecewise McCormick relaxations. For instance, Gounaris, Misener, and Floudas (2009) applied 15 different formulations. We follow the formulation called ‘nf4l’ in Gounaris, Misener, and Floudas (2009). This formulation was one of recommended formulations in Gounaris, Misener, and Floudas (2009). To simplify the notation, we will drop dependence on the subscripts and write $\mu_j \psi_\ell$ as $\mu\psi$. In practice, one has to apply piecewise McCormick relaxation to each bilinear term.

For any two positive terms $a \in [0, \bar{a}]$ and $b \in [0, \bar{b}]$, McCormick relaxation of $c \equiv ab$ consists of the following four inequalities:

$$(B.1) \quad c \geq 0, \quad c \geq a\bar{b} + \bar{a}b - \bar{a}\bar{b}, \quad c \leq a\bar{b}, \quad c \leq \bar{a}b,$$

which is known as the tightest possible convex relaxation.

To explain how to apply McCormick relaxation to $\mu\psi$, we introduce a new variable φ and replace $\mu\psi$ with φ . Then instead of imposing the bilinear constraint that $\varphi = \mu\psi$, we relax this in a piecewise fashion.

Suppose that ψ belongs to a known interval $[\underline{\psi}, \overline{\psi}]$. Assume that $\mu \in [\underline{\mu}, \overline{\mu}]$ with known end points. In practice, they can be deduced from \mathcal{S} since \mathcal{S} will be imposed simultaneously.

We now partition the space $[\underline{\psi}, \overline{\psi}]$ for ψ by a grid of $(K + 1)$ points $\{m_k : k = 0, \dots, K, m_0 = \underline{\psi}, m_K = \overline{\psi}\}$. Define λ_k to be a set of binary variables such that

$$\lambda_k = \begin{cases} 1 & \text{if } m_{k-1} \leq \psi \leq m_k \\ 0 & \text{otherwise} \end{cases}$$

for $k = 1, \dots, K$. Since we would like to ensure that ψ belongs to only one of segments $[m_{k-1}, m_k]$, we impose the summing up constraint such that

$$(B.2) \quad \sum_{k=1}^K \lambda_k = 1.$$

To reflect that $[\underline{\psi}, \overline{\psi}]$ is partitioned as described above, we introduce a set of continuous variables $\delta_k, k = 1, \dots, K$, where $0 \leq \delta_k \leq (m_k - m_{k-1})$. Then we impose the following set of restrictions

$$(B.3) \quad \begin{aligned} \psi &= \sum_{k=1}^K \{m_{k-1}\lambda_k + \delta_k\}, \\ 0 &\leq \delta_k \leq (m_k - m_{k-1})\lambda_k \quad \forall k. \end{aligned}$$

It can be seen that $\delta_k = 0$ if $\lambda_k = 0$ and $\delta_k = \psi - m_{k-1}$ for the index k such that $\lambda_k = 1$. For μ , we also introduce a set of continuous variables $\eta_k, k = 1, \dots, K$, where $0 \leq \eta_k \leq (\overline{\mu} - \underline{\mu})$. Impose the following restrictions

$$(B.4) \quad \begin{aligned} \mu &= \underline{\mu} + \sum_{k=1}^K \eta_k, \\ 0 &\leq \eta_k \leq (\overline{\mu} - \underline{\mu})\lambda_k \quad \forall k. \end{aligned}$$

As before, $\eta_k = 0$ if $\lambda_k = 0$ and $\eta_k = \mu - \underline{\mu}$ for the index k such that $\lambda_k = 1$.

Using newly defined variables δ_k and η_k , we now write

$$(B.5) \quad \varphi = \underline{\mu}\psi + \sum_{k=1}^K m_{k-1}\eta_k + \sum_{k=1}^K \delta_k\eta_k.$$

The first two terms on the right-hand side of (B.5) are linear in ψ and η_k ; whereas, the third term involves K bilinear terms of $\delta_k\eta_k$. Applying McCormick relaxation

(B.1) to $\delta_k \eta_k$ gives four inequalities for each k :

$$\begin{aligned}
 & \delta_k \eta_k \geq 0, \\
 (B.6) \quad & \delta_k \eta_k \geq (m_k - m_{k-1})\eta_k + (\underline{\mu} - \mu)\delta_k - (\underline{\mu} - \mu)(m_k - m_{k-1}), \\
 & \delta_k \eta_k \leq (\underline{\mu} - \mu)\delta_k, \\
 & \delta_k \eta_k \leq (m_k - m_{k-1})\eta_k.
 \end{aligned}$$

Instead of introducing a k -specific variable for each $\delta_k \eta_k$, define a single continuous variable Δ , where $0 \leq \Delta \leq \max_{k=1, \dots, K} (m_k - m_{k-1})(\bar{\mu} - \underline{\mu})$. Then rewrite (B.5) as

$$(B.7) \quad \varphi = \underline{\mu}\psi + \sum_{k=1}^K m_{k-1}\eta_k + \Delta$$

and aggregate equations in (B.6) over k to yield the following restrictions

$$\begin{aligned}
 (B.8) \quad & \Delta \geq \sum_{k=1}^K (m_k - m_{k-1})\eta_k + (\bar{\mu} - \underline{\mu}) \left(\sum_{k=1}^K [\delta_k - (m_k - m_{k-1})\lambda_k] \right), \\
 & \Delta \leq (\bar{\mu} - \underline{\mu}) \sum_{k=1}^K \delta_k, \\
 & \Delta \leq \sum_{k=1}^K (m_k - m_{k-1})\eta_k.
 \end{aligned}$$

In summary, the formulation of piecewise McCormick relaxation consists of (B.2), (B.3), (B.4), (B.7), and (B.8). The variables of optimization are μ , ψ , φ , Δ , $\lambda_k \in \{0, 1\}$, $\delta_k \in [0, (m_k - m_{k-1})]$, $\eta_k \in [0, (\bar{\mu} - \underline{\mu})]$, where $k = 1, \dots, K$. The total number of variables for optimization has increased from 2 to $4 + 3K$, but a bilinear constraint is relaxed to mixed integer linear constraints. A modern optimization solver (e.g. Gurobi) can handle efficiently mixed integer linear constraints.

We now describe how to construct inner bounds. Recall that $(\psi, \mu) \in [\underline{\psi}, \bar{\psi}] \times [\underline{\mu}, \bar{\mu}]$. When the bilinear term $\mu\psi$ exists in the optimization problem and we fix μ at one of values on its feasible set, the corresponding constrained optimization problem becomes convex but sup-optimal. Hence, solving the constrained optimization problem yields an inner bound. To obtain a tighter inner bound, we can create a grid of points for possible values of μ with size G and solve a constrained problem at each value of the grid. Taking the union of all these inner bounds gives a tight inner bound.

Note that K and G are tuning parameters to choose in implementation. To implement the method described above, we can start with small K and G and increase K

and G gradually up to the point that the set difference between the resulting outer and inner bounds is negligible up to some tolerance level. Even if the algorithm does not converge in a fixed time, we can compute the gap between the outer and inner bounds. This optimality gap is useful for evaluating the quality of the solution.

We state the proposed algorithm as follows.

Algorithm 1: Algorithm for outer and inner bounds

1. Select the type of \mathcal{S} and choose tuning parameters K and G .
 2. Obtain the outer bounds by solving $[\min_{\psi, \mu} f(\psi), \max_{\psi, \mu} f(\psi)]$ subject to

$$(B.9) \quad g(\psi, \mu) \leq 0, \psi \in \Psi, \text{ and } n^{1/2}(\bar{X} - \mu) \in \mathcal{S},$$
 while replacing each incidence of a bilinear term with the formulation of K -piecewise McCormick relaxation consisting of (B.2), (B.3), (B.4), (B.7), and (B.8).
 3. Construct a G -dimensional grid for components of μ , say μ_2 , appearing in the problem as bilinear terms. Obtain the lower bounds by solving $[\min_{\psi, \mu} f(\psi), \max_{\psi, \mu} f(\psi)]$ subject to (B.9), while fixing μ_2 at a fixed value of the grid points. Take the union of all G inner bounds to construct the best inner bounds.
 4. If the gap between outer and inner bounds is small, terminate. If not, increase K and G to see whether the gap can decrease further. Repeat the last step only fixed number of times.
 5. Report the resulting outer and inner bounds.
-

APPENDIX C. MINSKER'S (2015) MEDIAN OF MEANS METHOD

In this appendix, we describe how to carry out inference based on Minsker (2015). In particular, we consider two versions of the median of means: the one based on geometric median and the other using coordinate-wise medians. Lugosi and Mendelson (2019) propose a different version of the median of means estimator that has theoretically better properties but is more difficult to compute.

First, for the case of geometric median, let $\alpha_* := 7/18$ and $p_* := 0.1$. Define

$$\psi(\alpha_*; p_*) = (1 - \alpha_*) \log \frac{1 - \alpha_*}{1 - p_*} + \alpha_* \log \frac{\alpha_*}{p_*}.$$

Let $0 < \delta < 1$ be the level of the confidence set and set

$$(C.1) \quad k := \left\lceil \frac{\log(1/\delta)}{\psi(\alpha_*; p_*)} \right\rceil + 1.$$

Assume that δ is small enough that $k \leq n/2$. Divide the sample X_1, \dots, X_n into k disjoint groups G_1, \dots, G_k of size $\lfloor \frac{n}{k} \rfloor$ each, and define

$$\hat{\mu}_j := \frac{1}{|G_j|} \sum_{i \in G_j} X_i, \quad j = 1, \dots, k,$$

$$\hat{\mu} := \text{G.med}(\hat{\mu}_1, \dots, \hat{\mu}_k),$$

where G.med refers to the geometric median. See Minsker (2015) and references therein for details on the geometric median. The intuition behind $\hat{\mu}$ is that it is a robust measure of the population mean vector μ since each subsample mean vector $\hat{\mu}_j$ is an unbiased estimator for μ and the aggregation method via the geometric median is robust to outliers. Because of this feature, it turns out that the finite sample bound for the Euclidean norm distance between $\hat{\mu}$ and μ depends only on $\text{tr}(\Sigma)$, but not on the higher moments (see Corollary 4.1 of Minsker, 2015). This is the main selling point of the median of means since the finite sample probability bound for the usual sample mean assumes the existence of a higher moment (e.g. the third absolute moment in Bentkus (2003) and Lemma 2.1 in Section 2.4).

Second, Minsker (2015) also considered using coordinate-wise medians instead of using the geometric median. In this case, let $\alpha_* = 1/2$ and $p_* = 0.12$. Then k is redefined via (C.1). Let $\hat{\mu}_*$ denote the vector of coordinate-wise medians.

To estimate $\text{tr}(\Sigma)$, Minsker (2015) proposed the following:

$$\hat{T}_j := \frac{1}{|G_j|} \sum_{i \in G_j} \|X_i - \hat{\mu}_j\|^2, \quad j = 1, \dots, k,$$

$$\hat{T} := \text{med}(\hat{T}_1, \dots, \hat{T}_k).$$

where $\|a\|$ is the Euclidean norm of a vector a . Let $B(h, r)$ denote the ball of radius r centered at h and let

$$r_n := 11\sqrt{2} \sqrt{\hat{T} \frac{\log(1.4/\delta)}{n}},$$

$$r_{n,*} := 4.4\sqrt{2} \sqrt{\hat{T} \frac{\log(1.6d_\mu/\delta)}{n - 2.4 \log(1.6d_\mu/\delta)}},$$

where d_μ is the dimension of μ .

Lemma C.1 (Minsker (2015)). *Assume that*

$$(C.2) \quad 15.2 \sqrt{\frac{\mathbb{E} \|X - \mu\|^4 - (\text{tr}(\Sigma))^2}{(\text{tr}(\Sigma))^2}} \leq \left(\frac{1}{2} - 178 \frac{\log(1.4/\delta)}{n} \right) \sqrt{\frac{n}{\log(1.4/\delta)}}.$$

Then

$$(C.3) \quad \mathbb{P}[\mu \in B(\hat{\mu}, r_n)] \geq 1 - 2\delta \quad \text{and} \quad \mathbb{P}[\mu \in B(\hat{\mu}_*, r_{n,*})] \geq 1 - 2\delta.$$

Proof of Lemma C.1. The result on the geometric median is the exactly the same as Corollary 4.2 of Minsker (2015). The case for the vector of coordinate-wise medians follows from combining equation (4.4) in Minsker (2015) with Proposition 4.1 of Minsker (2015). \square

Lemma C.1 indicates that \mathcal{S} in our setup can be chosen as

$$(\hat{\mu} - \mu)'(\hat{\mu} - \mu) \leq r_n^2,$$

or

$$(\hat{\mu}_* - \mu)'(\hat{\mu}_* - \mu) \leq r_{n,*}^2,$$

either of which gives the bound with probability at least $1 - 2\delta$. The former produces a tighter bound than the latter only when the dimension of μ is sufficiently high. Note that (C.2) requires the existence of fourth moments due to the fact that $\text{tr}(\Sigma)$ is estimated by the median of means as well. The inequality in (C.2) is a relatively mild condition when n is large. In Section 4, we provide a numerical comparison between our main proposal and Minsker’s method.

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