

OVERIDENTIFICATION IN REGULAR MODELS

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In the unconditional moment restriction model of Hansen (1982), specification tests and more efficient estimators are both available whenever the number of moment restrictions exceeds the number of parameters of interest. We show that a similar relationship between potential refutability of a model and existence of more efficient estimators is present in much broader settings. Specifically, a condition we name *local overidentification* is shown to be equivalent to both the existence of specification tests with nontrivial local power and the existence of more efficient estimators of some “smooth” parameters in general semi/nonparametric models. Under our notion of local overidentification, various locally nontrivial specification tests such as Hausman tests, incremental Sargan tests (or optimally weighted quasi likelihood ratio tests) naturally extend to general semi/nonparametric settings. We further obtain simple characterizations of local overidentification for general models of nonparametric conditional moment restrictions with possibly different conditioning sets. The results are applied to determining when semi/nonparametric models with endogeneity are locally testable, and when nonparametric plug-in and semiparametric two-step GMM estimators are semiparametrically efficient. Examples of empirically relevant semi/nonparametric structural models are presented.

KEYWORDS: Overidentification, semiparametric efficiency, specification testing, nonparametric conditional moment restrictions, semiparametric two step, regular models, non-regular models.

1. INTRODUCTION

IN WORK ORIGINATING WITH Anderson and Rubin (1949) and Sargan (1958), and culminating in Hansen (1982), overidentification in the generalized method of moments (GMM) framework was equated with the number of unconditional moment restrictions exceeding the number of parameters of interest. Under mild regularity conditions, such a surplus of moment restrictions was shown to enable the construction of both more efficient estimators and model specification tests. It is hard to overstate the importance of this result, which has granted practitioners with an intuitive condition characterizing the existence of both efficiency gains and specification tests, and has thus intrinsically linked both phenomena to the notion of overidentification.

Unfortunately, the existence of an analogous simple condition in general semi/nonparametric models is, to the best of our knowledge, unknown. Yet, such a result stands

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to be particularly valuable for these more flexible models, as their richer structure renders their potential refutability harder to evaluate while simultaneously generating a broader set of parameters for which efficiency considerations are of concern.

In this paper, we show that, just as in GMM, efficiency and testability considerations are linked by a single condition we name *local overidentification*. In order to be applicable to general semi/nonparametric models, however, we must abstract from “counting” parameters and moment restrictions as in GMM when defining local overidentification. Instead we employ the *tangent set* $T(P)$ which, given a (data) distribution P and a candidate model \mathbf{P} , consists of the set of scores corresponding to all parametric submodels of \mathbf{P} that contain P (Bickel, Klaassen, Ritov, and Wellner (1993)). Heuristically, $T(P)$ consists of all the paths from which P may be approached from within \mathbf{P} . In particular, whenever the *closure* of $T(P)$ in the mean squared norm equals the set of all possible scores, the model \mathbf{P} is locally consistent with any parametric specification and hence we say P is *locally just identified* by \mathbf{P} . In contrast, whenever there exist scores that do not belong to the *closure* of $T(P)$, the model \mathbf{P} is locally inconsistent with some parametric specification and hence we say P is *locally overidentified* by \mathbf{P} . While these definitions can be generally applied, we mainly focus on models that are regular—in the sense that $T(P)$ is linear—due to the importance of this condition in semiparametric efficiency analysis (van der Vaart (1989)).¹ When specialized to GMM, our notion of local overidentification is equivalent to the standard condition that the number of unconditional moment restrictions exceed the number of parameters of interest.

Our definition of local overidentification arises naturally from embedding estimators of “smooth” (i.e., regular or root- n estimable for n the sample size) parameters and specifications tests in a common limiting experiment of LeCam (1986). This enables us to establish several equivalent characterizations of local overidentification. In particular, we show that if P is locally just identified by \mathbf{P} , then: (i) All asymptotically linear regular estimators of any common “smooth” parameter must be first-order equivalent; and (ii) the local asymptotic power of any local asymptotic level α specification test cannot exceed α along all paths approaching P from outside \mathbf{P} . Moreover, we establish that the local overidentification of P by a regular model \mathbf{P} is equivalent to both: (i) the existence of asymptotically distinct linear regular estimators for any “smooth” parameter that admits one such estimator; and (ii) the existence of a locally unbiased asymptotic level α specification test with nontrivial power against some path approaching P from outside \mathbf{P} .

Our equivalent characterizations of local overidentification are very useful. They offer researchers seemingly different yet equivalent ways to verify whether a data distribution P is locally overidentified by a complicated semi/nonparametric regular model \mathbf{P} . One obvious way is to directly verify the definition by first computing the closure of the tangent set $T(P)$ and then checking whether it is a strict subset of the space of all possible scores. An equivalent but sometimes simpler approach is to examine whether it is possible to obtain two asymptotically distinct regular estimators of a common “smooth” function of $P \in \mathbf{P}$,² such as the cumulative distribution function or a mean parameter $\int f dP$ for a known bounded function f . Given some structure on \mathbf{P} , it is often easy to compute two root- n consistent asymptotically normal estimators of a simple common “smooth” parameter, say as approximate optimizers of weighted criterion functions with different weights, and then verify whether their asymptotic variances differ.

¹See Section 5 for a partial extension of our results for regular models to non-regular models in which $T(P)$ is a convex cone.

²We stress that a “smooth” parameter of $P \in \mathbf{P}$ always exists and does not need to be any structural parameter associated with the model \mathbf{P} ; see Remark 3.1.

The local overidentification condition by itself, however, may not lead to feasible efficient estimators of parameters of interest nor to feasible tests with nontrivial local power in general regular models. Indeed, in parallel to GMM, additional regularity conditions are required to accomplish the latter two objectives. In our general setting, these regularity conditions are imposed by assuming the existence of a score statistic (a stochastic process) \hat{G}_n whose marginals are first-order equivalent to sample means of scores orthogonal to the tangent set $T(P)$. We show that such a score statistic \hat{G}_n can be constructed from two asymptotically distinct regular estimators of a common “smooth” parameter of $P \in \mathbf{P}^3$ —a result that can be exploited to provide low level sufficient conditions for the availability of \hat{G}_n given additional structure on \mathbf{P} . In addition, we show that \hat{G}_n can be used to obtain locally unbiased nontrivial specification tests. The constructed tests encompass, among others, Hausman (1978) type tests, and criterion-based tests such as the J test of Sargan (1958) and Hansen (1982) as special cases. In particular, proceeding in analogy to an incremental J test proposed in Eichenbaum, Hansen, and Singleton (1988) for GMM models, we demonstrate, for general regular models \mathbf{P} and \mathbf{M} satisfying $\mathbf{P} \subset \mathbf{M}$, how to build specification tests that aim their power at deviations from \mathbf{P} that satisfy the maintained larger model \mathbf{M} .

We deduce from the described results that the equivalence between efficiency gains and nontrivial specification tests found in Hansen (1982) is not coincidental, but rather the reflection of a deeper principle applicable to all regular models. Our results on local overidentification in general regular models should be widely applicable. For example, our equivalent characterizations immediately imply that semi/nonparametric models of conditional moment restrictions (with a common conditioning set) containing unknown functions of potentially endogenous variables are locally overidentified because they allow for both inefficient and efficient estimators (Ai and Chen (2003), Chen and Pouzo (2009)). Hence locally unbiased nontrivial specification tests of these models exist. Our results further show that the optimally weighted sieve quasi likelihood ratio tests of Chen and Pouzo (2009, 2015) direct the power at deviations of \mathbf{P} that remain within a larger model \mathbf{M} . We also show that Hausman (1978) type tests that compare estimators efficient under \mathbf{P} to those efficient under a larger model \mathbf{M} aim the power at violations of \mathbf{P} that remain within \mathbf{M} . Therefore, both kinds of tests could be understood as generalized incremental J tests.

In this paper, we focus on a new application to nonparametric conditional moment restriction models with possibly different conditioning sets and potential endogeneity. We derive simple equivalent characterizations of local just identification for this very large class of models so that other researchers do not need to compute the closure of the tangent set $T(P)$ case by case. When specialized to nonparametric conditional moment restrictions with possibly different conditioning sets but without endogeneity, such as nonparametric conditional mean or quantile regressions, our characterization of P being locally just identified reduces to the condition of the nonparametric functions being “exactly identified” in Akerberg, Chen, Hahn, and Liao (2014) for such models. When specialized to semi/nonparametric models using a control function approach for endogeneity (Heckman (1990), Olley and Pakes (1996), Newey, Powell, and Vella (1999), Blundell and Powell (2003)), our characterization implies that P is typically locally overidentified by such models. When specialized to the semi/nonparametric models of sequen-

³We also establish a converse, that is, the availability of such a score statistic \hat{G}_n yields asymptotically distinct regular estimators of any common “smooth” parameter of $P \in \mathbf{P}$.

tial moment restrictions containing unknown functions of potentially endogenous variables, our characterization implies that P is typically locally overidentified, which is consistent with the semiparametric efficiency bound calculation in [Ai and Chen \(2012\)](#) for such models. In Section 4, our results are applied further to determining when nonparametric plug-in and semiparametric two-step GMM estimators are semiparametrically efficient. Empirically relevant examples of semi/nonparametric structural models are also presented.

The rest of the paper is organized as follows. Section 2 formally defines local overidentification, while Section 3 establishes its connections to efficient estimation and testing in regular models. Section 4 applies the general theoretical results to characterize local overidentification in nonparametric conditional moment restriction models with possibly different information sets and potential endogeneity. Section 5 provides a partial extension of the main theoretical results in Section 3 for regular models to non-regular models in which $T(P)$ is a convex cone. Section 6 briefly concludes. Appendix A provides a short discussion of limiting experiments. Appendix B contains the proofs for Sections 2 and 3, while Appendix C contains the proofs for Section 5. The Supplemental Material ([Chen and Santos \(2018\)](#)) contains additional technical lemmas, examples, and the proofs for Section 4.

2. LOCAL OVERIDENTIFICATION

2.1. Main Definition

We let \mathcal{M} denote the collection of all probability measures on a measurable space $(\mathbf{X}, \mathcal{B})$. A model \mathbf{P} is a (not necessarily strict) subset of \mathcal{M} . Typically, a model \mathbf{P} is indexed by (model) parameters that consist of parameters of interest and perhaps additional nuisance parameters. We say a model \mathbf{P} is semiparametric if the parameters of interest are finite dimensional but the nuisance parameters are infinite dimensional (such as the GMM model of [Hansen \(1982\)](#)); semi-nonparametric if the parameters of interest contain both finite- and infinite-dimensional parts; nonparametric if all the parameters are infinite dimensional. We call a model \mathbf{P} fully unrestricted if $\mathbf{P} = \mathcal{M}$.

Throughout, the data $\{X_i\}_{i=1}^n$ are assumed to be an i.i.d. sample from a distribution $P \in \mathbf{P}$ of $X \in \mathbf{X}$. We call P the data distribution, which is always identified from the data, although its associated model parameters might not be. Our analysis is local in nature and hence we introduce suitable perturbations to P . Following the literature on limiting experiments ([LeCam \(1986\)](#)), we consider arbitrary smooth parametric likelihoods, which are defined as follows:

DEFINITION 2.1: A “path” $t \mapsto P_{t,g}$ is a function defined on $[0, 1)$ such that $P_{t,g}$ is a probability measure on $(\mathbf{X}, \mathcal{B})$ for every $t \in [0, 1)$, $P_{0,g} = P$, and

$$\lim_{t \downarrow 0} \int \left[\frac{1}{t} (dP_{t,g}^{1/2} - dP^{1/2}) - \frac{1}{2} g dP^{1/2} \right]^2 = 0. \tag{1}$$

The scalar measurable function $g : \mathbf{X} \rightarrow \mathbf{R}$ is referred to as the “score” of the path $t \mapsto P_{t,g}$.

For any σ -finite positive measure μ_t dominating $(P_t + P)$, the integral in (1) is understood as

$$\int \left[\frac{1}{t} \left(\left(\frac{dP_{t,g}}{d\mu_t} \right)^{1/2} - \left(\frac{dP}{d\mu_t} \right)^{1/2} \right) - \frac{1}{2} g \left(\frac{dP}{d\mu_t} \right)^{1/2} \right]^2 d\mu_t$$

(the choice of μ_t does not affect the value of the integral). Heuristically, a path is a parametric model that passes through P and is smooth in the sense of satisfying (1) or, equivalently, being differentiable in quadratic mean. We note that Definition 2.1 implies any score must have mean zero and be square integrable with respect to P , and therefore belong to the space $L_0^2(P)$ given by

$$L_0^2(P) \equiv \left\{ g : \mathbf{X} \rightarrow \mathbf{R}, \int g dP = 0 \text{ and } \|g\|_{P,2} < \infty \right\}, \quad \|g\|_{P,2}^2 \equiv \int g^2 dP. \quad (2)$$

The restriction $g \in L_0^2(P)$ is solely the result of $P_{t,g} \in \mathcal{M}$ for all t in a neighborhood of zero. If we, in addition, demand that $P_{t,g} \in \mathbf{P}$, then the set of feasible scores reduces to

$$T(P) \equiv \{g \in L_0^2(P) : (1) \text{ holds for some } t \mapsto P_{t,g} \in \mathbf{P}\}, \quad (3)$$

which is called the *tangent set* at P . Finally, we let $\bar{T}(P)$ denote the *closure* of $T(P)$ under $\|\cdot\|_{P,2}$. By definition, $\bar{T}(P)$ is a (not necessarily strict) subset of $L_0^2(P)$. For instance, if $\mathbf{P} = \mathcal{M}$, then $\bar{T}(P) = T(P) = L_0^2(P)$ for any $P \in \mathbf{P}$.

Given the introduced notation, we can now formally define local overidentification.

DEFINITION 2.2: If $\bar{T}(P) = L_0^2(P)$, then we say P is locally just identified by \mathbf{P} . Conversely, if $\bar{T}(P) \neq L_0^2(P)$, then we say P is locally overidentified by \mathbf{P} .

Intuitively, P is locally overidentified by a model \mathbf{P} if \mathbf{P} yields meaningful restrictions on the scores that can be generated by parametric submodels. Conversely, P is locally just identified by \mathbf{P} when the sole imposed restriction is that the scores have mean zero and a finite second moment—a quality common to the scores of all paths regardless of whether they belong to \mathbf{P} or not. It is clear that Definition 2.2 is inherently local in that it concerns only the “shape” of \mathbf{P} at the point P rather than \mathbf{P} in its entirety as would be appropriate for a global notion of overidentification.

REMARK 2.1: [Koopmans and Riersol \(1950\)](#) referred to a model \mathbf{P} as overidentified whenever there is a possibility that P does not belong to \mathbf{P} . Thus, \mathbf{P} is deemed globally overidentified if $\mathbf{P} \neq \mathcal{M}$ (i.e., \mathbf{P} is a strict subset of \mathcal{M}), and globally just identified if $\mathbf{P} = \mathcal{M}$ (i.e., \mathbf{P} is fully unrestricted). Clearly, global just identification implies local just identification, while local overidentification implies global overidentification. Although more demanding, local overidentification will provide a stronger connection to both the testability of \mathbf{P} and the performance of regular estimators.

It is worth emphasizing that local overidentification concerns solely a relationship between the data distribution P and a model \mathbf{P} . As a result, it is possible for P to be locally overidentified (hence globally overidentified) despite underlying structural parameters of the model \mathbf{P} being partially identified—an observation that simply reflects the fact that partially identified models may still be refuted by the data. See, for example, [Koopmans and Riersol \(1950\)](#), [Hansen and Jagannathan \(1997\)](#), [Manski \(2003\)](#), [Haile and Tamer \(2003\)](#), [Hansen \(2014\)](#), and references therein.

2.2. Equivalent Definitions in Regular Models

In many applications, the following condition holds and simplifies our analysis.

ASSUMPTION 2.1: (i) $\{X_i\}_{i=1}^n$ is an i.i.d. sequence with $X_i \in \mathbf{X}$ distributed according to $P \in \mathbf{P}$; (ii) $T(P)$ is linear; that is, if $g, f \in T(P)$, $a, b \in \mathbf{R}$, then $ag + bf \in T(P)$.

The i.i.d. requirement in Assumption 2.1(i) may be relaxed but is imposed to streamline exposition. Assumption 2.1(ii) requires the model \mathbf{P} to be *regular* at P in the sense that its tangent set be linear. This is satisfied by numerous models (such as the GMM model), and is either implicitly or explicitly imposed whenever semiparametric efficiency bounds and efficient estimators are considered (Hájek (1970), Hansen (1985), Chamberlain (1986), Newey (1990), Bickel et al. (1993), Ai and Chen (2003)). We stress, however, that a model \mathbf{P} being regular does not imply that all the parameters underlying the model are regular (i.e., “smooth” or root- n estimable). In fact, some parameters of a regular model \mathbf{P} may only be slower-than-root- n estimable or not even be identified. Nonetheless, Assumption 2.1(ii) does rule out models in which the tangent set $T(P)$ is not linear but a convex cone instead. See Section 5 for a partial extension of the main results for regular models to non-regular models where $T(P)$ is a convex cone.

In the literature, the closed linear span of $T(P)$ under $\|\cdot\|_{P,2}$ is called the *tangent space* at $P \in \mathbf{P}$ (see, e.g., Definition 3.2.2 in Bickel et al. (1993)). Under Assumption 2.1(ii), $\bar{T}(P)$ becomes the tangent space at P , and hence a vector subspace of $L_0^2(P)$. We also define

$$\bar{T}(P)^\perp \equiv \left\{ g \in L_0^2(P) : \int gf \, dP = 0 \text{ for all } f \in \bar{T}(P) \right\}, \tag{4}$$

which is the orthogonal complement of $\bar{T}(P)$. The vector spaces $\bar{T}(P)$ and $\bar{T}(P)^\perp$ then form an orthogonal decomposition of $L_0^2(P)$ (the space of all possible scores)

$$L_0^2(P) = \bar{T}(P) \oplus \bar{T}(P)^\perp, \tag{5}$$

and we let $\Pi_T(\cdot)$ and $\Pi_{T^\perp}(\cdot)$ denote the orthogonal projections under $\|\cdot\|_{P,2}$ onto $\bar{T}(P)$ and $\bar{T}(P)^\perp$, respectively. Every $g \in L_0^2(P)$ then satisfies $g = \Pi_T(g) + \Pi_{T^\perp}(g)$ and $\text{Var}(g) = \text{Var}(\Pi_T(g)) + \text{Var}(\Pi_{T^\perp}(g))$. Intuitively, $\Pi_T(g) \in \bar{T}(P)$ is the component of g that is in accord with model \mathbf{P} , while $\Pi_{T^\perp}(g) \in \bar{T}(P)^\perp$ is the component orthogonal to \mathbf{P} .

The decomposition in (5) implies equivalent characterizations of local overidentification that we summarize in the following simple yet useful lemma.

LEMMA 2.1: *Under Assumption 2.1, the following are equivalent to Definition 2.2:*

- (i) P is locally just identified by \mathbf{P} if and only if $\bar{T}(P)^\perp = \{0\}$, or equivalently, $\text{Var}(\Pi_{T^\perp}(g)) = 0$ for all $g \in L_0^2(P)$.
- (ii) P is locally overidentified by \mathbf{P} if and only if $\bar{T}(P)^\perp \neq \{0\}$, or equivalently, $\text{Var}(\Pi_{T^\perp}(g)) > 0$ for some $g \in L_0^2(P)$.

We next illustrate the introduced concepts in the GMM model.⁴

GMM ILLUSTRATION: Let $\Gamma \subseteq \mathbf{R}^{d_\gamma}$ with $d_\gamma < \infty$ be the parameter space and $\rho : \mathbf{X} \times \mathbf{R}^{d_\gamma} \rightarrow \mathbf{R}^{d_\rho}$ be a known moment function with $d_\rho \geq d_\gamma$. The GMM model \mathbf{P} is

$$\mathbf{P} \equiv \left\{ P \in \mathcal{M} : \int \rho(\cdot, \gamma) \, dP = 0 \text{ for some } \gamma \in \Gamma \right\}, \tag{6}$$

⁴We thank Lars Peter Hansen for sharing with us his notes on GMM and for helping us with the GMM example.

and for any $P \in \mathbf{P}$ we let $\gamma(P)$ solve $\int \rho(\cdot, \gamma(P)) dP = 0$. For simplicity, let ρ be differentiable in γ , and set $D(P) \equiv \int \nabla_{\gamma} \rho(\cdot, \gamma(P)) dP$. For any path $t \mapsto P_{t,g} \in \mathbf{P}$, we then obtain

$$0 = \frac{d}{dt} \int \rho(\cdot, \gamma(P_{t,g})) dP_{t,g} \Big|_{t=0} = \int \rho(\cdot, \gamma(P)) g dP + D(P) \dot{\gamma}(g), \tag{7}$$

where $\dot{\gamma}(g)$ is the derivative of $\gamma(P_{t,g})$ at $t = 0$. If $\int \rho(\cdot, \gamma(P)) \rho(\cdot, \gamma(P))' dP$ is full rank, then the linear functional $g \mapsto \int \rho(\cdot, \gamma(P)) g dP$ maps $L_0^2(P)$ onto $\mathbf{R}^{d_{\rho}}$. On the other hand, $D(P)$ maps $\mathbf{R}^{d_{\gamma}}$ onto a linear subspace of $\mathbf{R}^{d_{\rho}}$ whose dimension equals the rank of $D(P)$. Therefore, (7) imposes restrictions on the possible set of scores g only when the rank of $D(P)$ is smaller than d_{ρ} . When $D(P)$ is full rank, we thus obtain that P is locally just identified by \mathbf{P} if and only if the “standard” GMM just identification condition that $d_{\rho} = d_{\gamma}$ is satisfied.

Our definition of local overidentification extends that in GMM models to general infinite-dimensional models. This will be very useful for nonparametric conditional moment restriction models, where both the number of parameters (of interest) and the number of (unconditional) moments are infinite. Moreover, for general regular models, we will show Definition 2.2 retains the fundamental link to the properties of regular estimators and specification tests present in Hansen (1982). For instance, just as all regular estimators of $\gamma(P)$ in the GMM model are asymptotically equivalent whenever $d_{\rho} = d_{\gamma}$, Theorem 2.1 in Newey (1994) has shown that the asymptotic variance of root- n consistent plug-in estimators is invariant to the choice of first-stage nonparametric estimators whenever $L_0^2(P) = \bar{T}(P)$.

3. GENERAL RESULTS FOR REGULAR MODELS

In this section, we show that, in general regular models, local overidentification is intrinsically linked to the importance of efficiency considerations and the potential refutability of a model.

3.1. The Setup

3.1.1. The Setup: Estimation

Since the data distribution P is always identified, many known functions of P are identified and consistently estimable even if some underlying (structural) parameters of a model \mathbf{P} are not identified. For general regular models, we therefore represent an identifiable parameter as a known mapping $\theta : \mathbf{P} \rightarrow \mathbf{B}$ and the “true” parameter value as $\theta(P) \in \mathbf{B}$, where \mathbf{B} is a Banach space with norm $\|\cdot\|_{\mathbf{B}}$. We further denote the dual space of \mathbf{B} by $\mathbf{B}^* \equiv \{b^* : \mathbf{B} \rightarrow \mathbf{R} : b^* \text{ is linear, } \|b^*\|_{\mathbf{B}^*} < \infty\}$, which is the space of continuous linear functionals with norm $\|b^*\|_{\mathbf{B}^*} \equiv \sup_{\|b\|_{\mathbf{B}} \leq 1} |b^*(b)|$.

An estimator $\hat{\theta}_n : \{X_i\}_{i=1}^n \rightarrow \mathbf{B}$ for $\theta(P) \in \mathbf{B}$ is a map from the data into the space \mathbf{B} . To address the question of whether $\theta(P)$ admits asymptotically distinct estimators (i.e., efficiency “matters”), we focus on *asymptotically linear regular* estimators. In what follows, for any path $t \mapsto P_{t,g} \in \mathcal{M}$, we use the notation $\xrightarrow{L_{n,g}}$ to represent convergence in law under $P_{1/\sqrt{n},g}^n \equiv \otimes_{i=1}^n P_{1/\sqrt{n},g}$, and \xrightarrow{L} for convergence in law under $P^n \equiv \otimes_{i=1}^n P$.

DEFINITION 3.1: $\hat{\theta}_n : \{X_i\}_{i=1}^n \rightarrow \mathbf{B}$ is a regular estimator of $\theta(P)$ if there is a tight random variable \mathbb{D} such that $\sqrt{n}\{\hat{\theta}_n - \theta(P_{1/\sqrt{n},g})\} \xrightarrow{L_{n,g}} \mathbb{D}$ for any path $t \mapsto P_{t,g} \in \mathbf{P}$.

DEFINITION 3.2: $\hat{\theta}_n : \{X_i\}_{i=1}^n \rightarrow \mathbf{B}$ is an asymptotically linear estimator of $\theta(P)$ if

$$\sqrt{n}\{\hat{\theta}_n - \theta(P)\} = \frac{1}{\sqrt{n}} \sum_{i=1}^n \nu(X_i) + o_p(1) \quad \text{under } P^n, \tag{8}$$

for some $\nu : \mathbf{X} \rightarrow \mathbf{B}$ satisfying $b^*(\nu) \in L_0^2(P)$ for any $b^* \in \mathbf{B}^*$. Here ν is called the influence function of the estimator $\hat{\theta}_n$.

By restricting attention to regular estimators, we focus on root- n consistent estimators whose asymptotic distribution is invariant to local perturbations to P within the model \mathbf{P} . While most commonly employed estimators are regular and asymptotically linear, their existence does impose restrictions on the map $\theta : \mathbf{P} \rightarrow \mathbf{B}$. In fact, the existence of an asymptotically linear regular estimator of $\theta(P)$ in regular models implies $\theta : \mathbf{P} \rightarrow \mathbf{B}$ must be “pathwise differentiable” (or “smooth”) at P relative to $T(P)$ (van der Vaart (1991b)).

REMARK 3.1: Regardless of a model \mathbf{P} being regular or non-regular, there always exists a “smooth” map $\theta : \mathbf{P} \rightarrow \mathbf{B}$ and an asymptotically linear regular estimator $\hat{\theta}_n : \{X_i\}_{i=1}^n \rightarrow \mathbf{B}$ of $\theta(P)$ under i.i.d. data. For example, for any bounded function $f : \mathbf{X} \rightarrow \mathbf{R}$, the sample mean, $n^{-1} \sum_{i=1}^n f(X_i)$, is an asymptotically linear regular estimator of $\theta(P) \equiv \int f dP$ along any path $t \mapsto P_{t,g} \in \mathcal{M}$. Thus, we emphasize that $\theta(P)$ should not be solely thought of as an intrinsic parameter of the model \mathbf{P} , but rather as any “smooth” map of $P \in \mathbf{P}$.

3.1.2. The Setup: Testing

A specification test for a general model \mathbf{P} is a test of the null hypothesis that P belongs to \mathbf{P} against the alternative that it does not; that is, it is a test of the hypotheses

$$H_0 : P \in \mathbf{P} \quad \text{vs.} \quad H_1 : P \in \mathcal{M} \setminus \mathbf{P}. \tag{9}$$

We denote an arbitrary (possibly randomized) test of (9) by $\phi_n : \{X_i\}_{i=1}^n \rightarrow [0, 1]$, which is a function specifying for each realization of the data a corresponding probability of rejecting the null hypothesis. In our analysis, we restrict attention to specification tests ϕ_n that have local asymptotic level α and possess a local asymptotic power function.

DEFINITION 3.3: A specification test $\phi_n : \{X_i\}_{i=1}^n \rightarrow [0, 1]$ for a model \mathbf{P} has local asymptotic level α if, for any path $t \mapsto P_{t,g} \in \mathbf{P}$, it follows that

$$\limsup_{n \rightarrow \infty} \int \phi_n dP^n_{1/\sqrt{n},g} \leq \alpha. \tag{10}$$

DEFINITION 3.4: A specification test $\phi_n : \{X_i\}_{i=1}^n \rightarrow [0, 1]$ for a model \mathbf{P} has a local asymptotic power function $\pi : L_0^2(P) \rightarrow [0, 1]$ if, for any path $t \mapsto P_{t,g} \in \mathcal{M}$, it follows that

$$\lim_{n \rightarrow \infty} \int \phi_n dP^n_{1/\sqrt{n},g} = \pi(g). \tag{11}$$

Finally, a test ϕ_n for (9) with a local asymptotic power function π is said to be locally unbiased if it satisfies: $\pi(g) \leq \alpha$ for all $t \mapsto P_{t,g} \in \mathbf{P}$ and $\pi(g) \geq \alpha$ for all $t \mapsto P_{t,g} \in \mathcal{M} \setminus \mathbf{P}$.

Note that a local asymptotic power function only depends on the score $g \in L_0^2(P)$ and is independent of any other characteristics of the path $t \mapsto P_{t,g} \in \mathcal{M}$. This is because the product measures of any two local paths that share the same score must converge in the Total Variation metric (see Lemma D.1 in the Supplemental Material). Intuitively, a test possesses a local asymptotic power function if the limiting rejection probability of the test is well defined along any local perturbation to P . The existence of a local asymptotic power function is a mild requirement that is typically satisfied; see Remark 3.2.

REMARK 3.2: Tests ϕ_n are often constructed by comparing a test statistic \hat{T}_n to an estimate of the $(1 - \alpha)$ quantile of its asymptotic distribution. By LeCam’s Third Lemma and the Portmanteau Theorem, such tests have a local asymptotic power function provided that: (i) $(\hat{T}_n, \frac{1}{\sqrt{n}} \sum_{i=1}^n g(X_i)) \in \mathbf{R}^2$ converges jointly in distribution under P^n for any $g \in L_0^2(P)$, and (ii) the limiting distribution of \hat{T}_n under P^n is continuous. See Theorem 6.6 in van der Vaart (1998).

3.2. Equivalent Characterizations of Local Overidentification

In Hansen (1982)’s GMM framework, overidentifying restrictions are necessary for both the existence of efficiency gains in estimation and the testability of the model. We now extend this conclusion to general regular models.

THEOREM 3.1: *Let Assumption 2.1 hold and P be locally just identified by \mathbf{P} .*

(i) *Let $\hat{\theta}_n : \{X_i\}_{i=1}^n \rightarrow \mathbf{B}$ and $\tilde{\theta}_n : \{X_i\}_{i=1}^n \rightarrow \mathbf{B}$ be any asymptotically linear regular estimators of any parameter $\theta(P) \in \mathbf{B}$. Then: $\sqrt{n}\{\hat{\theta}_n - \tilde{\theta}_n\} = o_p(1)$ in \mathbf{B} .*

(ii) *Let ϕ_n be any specification test for (9) with local asymptotic level α and a local asymptotic power function π . Then: $\pi(g) \leq \alpha$ for all paths $t \mapsto P_{t,g} \in \mathcal{M}$.*

Theorem 3.1 establishes that the local overidentification of P is a necessary condition for the existence of efficiency gains and nontrivial specification tests. Specifically, Theorem 3.1(i) shows that if P is locally just identified, then all asymptotically linear regular estimators of any “smooth” parameter $\theta(P)$ must be first-order equivalent. This conclusion is a generalization of Newey (1990) who showed scalar (i.e., $\mathbf{B} = \mathbf{R}$) asymptotically linear and regular estimators are first-order equivalent when $\bar{T}(P) = L_0^2(P)$. Theorem 3.1(ii) establishes that if P is locally just identified by \mathbf{P} , then the local asymptotic power of any local asymptotic level α specification test cannot exceed α along any path, including all paths approaching P from outside \mathbf{P} . Heuristically, under local just identification, the set of scores corresponding to paths $t \mapsto P_{t,g} \in \mathbf{P}$ is dense in the set of all possible scores and hence every path $t \mapsto P_{t,g} \in \mathcal{M} \setminus \mathbf{P}$ is locally on the “boundary” of the null hypothesis.

In order to discern how the local overidentification of P can facilitate the existence of efficiency gains and the testability of the model, we consider the asymptotic behavior of sample means of scores. For any $0 \neq \tilde{f} \in L_0^2(P)$, if X_i is distributed according to $P_{1/\sqrt{n},g}$ for a path $t \mapsto P_{t,g} \in \mathcal{M}$, then

$$\mathbb{G}_n(\tilde{f}) \equiv \frac{1}{\sqrt{n}} \sum_{i=1}^n \tilde{f}(X_i) \xrightarrow{L_{n,g}} N\left(\int \tilde{f}g dP, \int \tilde{f}^2 dP\right) \tag{12}$$

by LeCam’s Third Lemma. Recall that, for regular models, local overidentification is equivalent to the existence of at least one score $0 \neq \tilde{f} \in \bar{T}(P)^\perp$. For any such $0 \neq \tilde{f} \in \bar{T}(P)^\perp$ and all path $t \mapsto P_{t,g} \in \mathbf{P}$, we have $\int \tilde{f}g dP = 0$ and hence $\mathbb{G}_n(\tilde{f})$ converges to a centered Gaussian random variable; that is, $\mathbb{G}_n(\tilde{f})$ behaves as “noise” that can alter the efficiency of estimators. On the other hand, for any $0 \neq \tilde{f} \in \bar{T}(P)^\perp$, there is a path $t \mapsto P_{t,g} \in \mathcal{M} \setminus \mathbf{P}$ such that $\int \tilde{f}g dP \neq 0$, and hence $\mathbb{G}_n(\tilde{f})$ can be employed to construct a local asymptotic nontrivial specification test—i.e. $\mathbb{G}_n(\tilde{f})$ is a “signal” that enables the detection of violations of the model \mathbf{P} . Our next result builds on this intuition by using the score statistics $\mathbb{G}_n(\tilde{f})$ to establish a converse to Theorem 3.1.

THEOREM 3.2: *Let Assumption 2.1 hold. Then: the following statements are equivalent:*

- (i) P is locally overidentified by \mathbf{P} .
- (ii) If a parameter $\theta(P) \in \mathbf{B}$ admits an asymptotically linear regular estimator $\hat{\theta}_n$, then there exists another asymptotically linear regular estimator $\tilde{\theta}_n$ of $\theta(P)$ such that $\sqrt{n}\{\hat{\theta}_n - \tilde{\theta}_n\} \xrightarrow{L} \Delta \neq 0$ in \mathbf{B} .
- (iii) There exists a locally unbiased asymptotic level α test ϕ_n for (9) with a local asymptotic power function π such that $\pi(g) > \alpha$ for some path $t \mapsto P_{t,g} \in \mathcal{M} \setminus \mathbf{P}$.

Theorems 3.1 and 3.2 establish that the local overidentification of P is equivalent to the availability of efficiency gains and also to the existence of locally nontrivial specification tests. In addition, Theorems 3.1(i) and 3.2(i)–(ii) imply the following equivalent characterization of local just identification.

COROLLARY 3.1: *Let Assumption 2.1 hold and \mathcal{D} be a set of bounded functions that is dense in $(L^2(P), \|\cdot\|_{P,2})$. For any $f \in \mathcal{D}$, let Ω_f^* be the semiparametric efficient variance bound for estimating $\int f dP$ under \mathbf{P} . Then: $\Omega_f^* = \text{Var}\{f(X)\}$ for all $f \in \mathcal{D}$ if and only if P is locally just identified by \mathbf{P} .*

This corollary is very useful in assessing whether P is locally overidentified by a complicated model \mathbf{P} . For example, in Section 4.1.1, we employ Corollary 3.1 and the semiparametric efficiency bound analysis in Ai and Chen (2012) to characterize local overidentification in nonparametric models defined by sequential moment restrictions.

3.3. Feasible Estimators and Tests

The intuition for Theorem 3.2 suggests that any statistics asymptotically equivalent to the score statistics $\mathbb{G}_n(\tilde{f})$ with some $0 \neq \tilde{f} \in \bar{T}(P)^\perp$ (see (12)) may be employed to obtain distinct regular estimators for arbitrary “smooth” parameters $\theta(P)$ and specification tests with nontrivial local power. To elaborate on this intuition, for any set \mathcal{A} , we let $\ell^\infty(\mathcal{A}) \equiv \{f : \mathcal{A} \rightarrow \mathbf{R} \text{ s.t. } \|f\|_\infty < \infty\}$ where $\|f\|_\infty = \sup_{a \in \mathcal{A}} |f(a)|$, and impose the following condition:

ASSUMPTION 3.1: *For some set \mathbf{T} , there is a statistic $\hat{\mathbb{G}}_n : \{X_i\}_{i=1}^n \rightarrow \ell^\infty(\mathbf{T})$ satisfying:*

- (i) $\hat{\mathbb{G}}_n(\tau) = \frac{1}{\sqrt{n}} \sum_{i=1}^n s_\tau(X_i) + o_p(1)$ uniformly in $\tau \in \mathbf{T}$, where $0 \neq s_\tau \in \bar{T}(P)^\perp$ for all $\tau \in \mathbf{T}$;
- (ii) for some tight nondegenerate centered Gaussian measure $\mathbb{G}_0, \hat{\mathbb{G}}_n \xrightarrow{L} \mathbb{G}_0$ in $\ell^\infty(\mathbf{T})$.

Assumption 3.1 requires the availability of a statistic $n^{-1/2}\hat{G}_n(\tau)$ that is first-order equivalent to the sample mean of some score (or influence function) $s_\tau \in \bar{T}(P)^\perp$. We further let

$$S(P) \equiv \{s_\tau \in \bar{T}(P)^\perp : \tau \in \mathbf{T}\} \tag{13}$$

denote the collection of such scores, which will play an important role in our analysis. As we argue below, statistics \hat{G}_n satisfying Assumption 3.1 are implicitly constructed by various specification tests, such as Hausman tests and criterion-based tests; see Remark 3.5. In order to establish a connection to Hausman tests in particular, we introduce the following assumption:

ASSUMPTION 3.2: *For some parameter $\theta(P) \in \mathbf{B}$, there are asymptotically linear regular estimators $\hat{\theta}_n$ and $\tilde{\theta}_n$ with influence functions ν and $\tilde{\nu}$ such that $\sqrt{n}\{\hat{\theta}_n - \tilde{\theta}_n\} \xrightarrow{L} \Delta \neq 0$.*

Assumption 3.2 simply requires the existence of two distinct estimators of some “smooth” function of $P \in \mathbf{P}$, which need not be structural parameter of the model \mathbf{P} ; see Remark 3.1.

LEMMA 3.1: *Let Assumption 2.1 hold.*

(i) *Let Assumption 3.1 hold. Then: For any parameter $\theta(P) \in \mathbf{B}$ that admits an asymptotically linear regular estimator $\hat{\theta}_n$, Assumption 3.2 is satisfied with $\tilde{\theta}_n = \hat{\theta}_n + \tilde{b} \times n^{-1/2}\hat{G}_n(\tau^*)$ and $\Delta = -\tilde{b} \times G_0(\tau^*)$ for some $0 \neq \tilde{b} \in \mathbf{B}$ and some $\tau^* \in \mathbf{T}$.*

(ii) *Let Assumption 3.2 hold. Then: Assumption 3.1 is satisfied with $\mathbf{T} = \{b^* \in \mathbf{B}^* : \|b^*\|_{\mathbf{B}^*} \leq 1\}$, and $G_0, \hat{G}_n \in \ell^\infty(\mathbf{T})$ given by $G_0(b^*) = b^*(\Delta)$, $\hat{G}_n(b^*) = b^*(\sqrt{n}\{\hat{\theta}_n - \tilde{\theta}_n\})$ where $s_{b^*} = b^*(\nu - \tilde{\nu})$.*

Lemma 3.1 establishes that Assumptions 3.1 and 3.2 are equivalent to each other. In particular, Lemma 3.1(ii) shows that the difference of any two asymptotically distinct linear regular estimators of any common parameter $\theta(P)$ may be employed to construct \hat{G}_n , that is, Assumption 3.2 implies Assumption 3.1. As a result, given the specific structure of a regular model \mathbf{P} , it is straightforward to obtain lower level sufficient conditions for Assumption 3.1. Specifically, we need only ensure the existence of two asymptotically distinct linear regular estimators of some “smooth” parameter, which could be a simple identified reduced form parameter if the structural parameters are not identified. In a large class of semiparametric and nonparametric models, asymptotically distinct estimators may be found as the optimizers of weighted criterion functions with alternative choices of weights. See, for example, Shen (1997) for efficient estimation based on sieve or penalized maximum likelihood in semi/nonparametric likelihood models, and Ai and Chen (2003) for efficient estimation based on optimally weighted sieve minimum distance of semi/nonparametric conditional moment restrictions models. In the Supplemental Material, we apply Lemma 3.1(ii) to verify Assumption 3.1 in general nonparametric conditional moment restriction models (26).

We next employ the fact that \hat{G}_n behaves as a “signal” from a testing perspective (i.e., Theorem 3.2(iii)) to construct nontrivial local specification tests. Let $\bar{S}(P) \equiv \overline{\text{lin}}\{S(P)\}$ be the closed linear span of $S(P)$ in $L_0^2(P)$, and $\Pi_S(g)$ be the metric projection of $g \in L_0^2(P)$ onto $\bar{S}(P)$. We note Assumption 3.1(i) (or (12)) implies that \hat{G}_n exhibits a nonzero asymptotic drift along a path $t \mapsto P_{t,g} \in \mathcal{M}$ if and only if $\Pi_S(g) \neq 0$. Intuitively, $\bar{S}(P)$ therefore

represents the alternatives for which specification tests based on $\hat{\mathbb{G}}_n$ have nontrivial local asymptotic power. To obtain such a test, we employ a map $\Psi : \ell^\infty(\mathbf{T}) \rightarrow \mathbf{R}_+$ to reduce $\hat{\mathbb{G}}_n$ to a scalar test statistic $\hat{T}_n = \Psi(\hat{\mathbb{G}}_n)$.

ASSUMPTION 3.3: (i) $\Psi : \ell^\infty(\mathbf{T}) \rightarrow \mathbf{R}_+$ is continuous, convex, and nonconstant; (ii) $\Psi(0) = 0$, $\Psi(b) = \Psi(-b)$ for all $b \in \ell^\infty(\mathbf{T})$; (iii) $\{b \in \ell^\infty(\mathbf{T}) : \Psi(b) \leq c\}$ is bounded for all $c > 0$.

Finally, we let $c_{1-\alpha} > 0$ be the $(1 - \alpha)$ quantile of $\Psi(\mathbb{G}_0)$ and define the specification test for (9) as

$$\phi_n \equiv 1\{\Psi(\hat{\mathbb{G}}_n) > c_{1-\alpha}\}; \tag{14}$$

that is, we reject correct model specification for large values of $\Psi(\hat{\mathbb{G}}_n)$. Multiple specification tests in the literature are in fact asymptotically equivalent to (14) with different choices of Ψ ; see Theorem 3.3 Part (ii) and Remark 3.5 below. In the rest of the paper we often use $P_{1/\sqrt{n},g}()$ and $E_{P_{1/\sqrt{n},g}}[\cdot]$ for calculations under $P_{1/\sqrt{n},g}^n$.

THEOREM 3.3: Let Assumption 2.1 hold.

(i) Let Assumptions 3.1 and 3.3 hold. Then: ϕ_n defined in (14) with $c_{1-\alpha} > 0$ is a locally unbiased asymptotic level α specification test for (9) with a local asymptotic power function π . Moreover, for any path $t \mapsto P_{t,g} \in \mathcal{M}$ with $\Pi_S(g) \neq 0$, it follows that

$$\pi(g) \equiv \lim_{n \rightarrow \infty} P_{1/\sqrt{n},g}(\Psi(\hat{\mathbb{G}}_n) > c_{1-\alpha}) > \alpha. \tag{15}$$

(ii) Let Assumption 3.2 hold. Then: Assumption 3.3 holds with $\Psi = \|\cdot\|_\infty$, and Part (i) holds with $\Psi(\hat{\mathbb{G}}_n) = \sqrt{n}\|\hat{\theta}_n - \tilde{\theta}_n\|_{\mathbf{B}}$, and $S(P) = \{b^*(\nu - \tilde{\nu}) : b^* \in \mathbf{B}^*, \|b^*\|_{\mathbf{B}^*} \leq 1\}$.

Theorems 3.1, 3.2, 3.3, and Lemma 3.1 link local overidentification to the existence of asymptotically distinct estimators and locally nontrivial specification tests. The latter two concepts were also intrinsically linked by the seminal work of Hausman (1978), who proposed comparing estimators of a common parameter to perform specification tests. Theorem 3.3(ii) shows Hausman tests are a special case of (14) in general regular models.

REMARK 3.3: Whenever $\bar{S}(P) = \bar{T}(P)^\perp$, result (15) holds for any path $t \mapsto P_{t,g}$ with

$$\liminf_{n \rightarrow \infty} \inf_{Q \in \mathbf{P}} n \int [dQ^{1/2} - dP_{1/\sqrt{n},g}^{1/2}]^2 > 0; \tag{16}$$

that is, the proposed test has nontrivial local power against any path that does not approach \mathbf{P} “too fast.” If condition (16) fails, then there is a sequence $Q_n \in \mathbf{P}$ for which

$$\limsup_{n \rightarrow \infty} \left| \int \phi_n(dQ_n^n - dP_{1/\sqrt{n},g}^n) \right| \leq \limsup_{n \rightarrow \infty} \|Q_n^n - P_{1/\sqrt{n},g}^n\|_{\text{TV}} = 0, \tag{17}$$

where $\|\cdot\|_{\text{TV}}$ denotes the total variation distance; see, for example, Theorem 13.1.3 in Lehmann and Romano (2005). Therefore, a violation of (16) implies $P_{1/\sqrt{n},g}$ approaches \mathbf{P} “too fast” in the sense that it is not possible to discriminate the induced distribution on the data $\{X_i\}_{i=1}^n$ from a distribution that is in accord with \mathbf{P} .

REMARK 3.4: Theorem 3.3(ii) states a Hausman test has nontrivial local power against any path $t \mapsto P_{t,g} \in \mathcal{M}$ whose score g is correlated with $b^*(\nu - \tilde{\nu})$ for some $b^* \in \mathbf{B}^*$. When $\hat{\theta}_n$ is a semiparametric efficient estimator, it follows that $\bar{S}(P) = \overline{\text{lin}}\{\Pi_{T^\perp}(b^*(\tilde{\nu})) : b^* \in \mathbf{B}^*\}$ (see Proposition 3.3.1 in Bickel et al. (1993)). In particular, $L_0^2(P) = \overline{\text{lin}}\{b^*(\tilde{\nu}) : b^* \in \mathbf{B}^*\}$ implies $\bar{S}(P) = \bar{T}(P)^\perp$, and hence the corresponding Hausman test for (9) has nontrivial power against all local alternatives.

REMARK 3.5: In addition to Hausman tests, multiple specification tests for (9) are also asymptotically equivalent to test (14). For example, optimally weighted criterion-based tests employ statistics \hat{T}_n that have a chi-squared asymptotic distribution and satisfy

$$\hat{T}_n = \sum_{k=1}^K \left(\frac{1}{\sqrt{n}} \sum_{i=1}^n f_k(X_i) \right)^2 + o_p(1), \tag{18}$$

where K corresponds to the degrees of freedom and $\{f_k\}_{k=1}^K \subset L_0^2(P)$ are orthonormal. A test with this property can only have nominal and local asymptotic level α if $f_k \in \bar{T}(P)^\perp$ for all k . Otherwise, there is a path $t \mapsto P_{t,g} \in \mathbf{P}$ such that $\int g\{\Pi_T(f_k)\} dP \neq 0$ for at least one k , which by (12) leads to a null rejection probability exceeding α . As a result, the structure in test (14) is also present in the J test of Hansen (1982), the semiparametric LR statistic of Murphy and van der Vaart (1997), the sieve QLR statistic in Chen and Pouzo (2009), and the generalized empirical likelihood test in Parente and Smith (2011) among many others.

3.4. Incremental J Tests

In applications, specification tests are often informed by a concern with a particular violation of the model. For instance, in GMM, we may question the validity of a subset of the moment conditions but have confidence in the remaining ones; see, for example, Eichenbaum, Hansen, and Singleton (1988). In such circumstances, a J test, which entertains the possibility of any moment being violated, can be less revealing than the so-called incremental J (Sargan–Hansen) test, which focuses on the specific moments that are of concern (Arellano (2003)).

The tests in Theorem 3.3 can similarly direct their power at specific violations of the model. To this end, we introduce a set \mathbf{M} satisfying $\mathbf{P} \subseteq \mathbf{M} \subseteq \mathcal{M}$, which represents the characteristics of the model we believe P satisfies even when $P \notin \mathbf{P}$, and consider

$$H_0 : P \in \mathbf{P} \quad \text{vs.} \quad H_1 : P \in \mathbf{M} \setminus \mathbf{P}. \tag{19}$$

The “maintained” model \mathbf{M} generates its own tangent set, which we denote by

$$M(P) \equiv \{g \in L_0^2(P) : (1) \text{ holds for some } t \mapsto P_{t,g} \in \mathbf{M}\}, \tag{20}$$

with $\bar{M}(P)$ being the closure of $M(P)$ in $(L_0^2(P), \|\cdot\|_{p,2})$. If $M(P)$ is linear, then $\bar{M}(P) = \bar{T}(P) \oplus \{\bar{T}(P)^\perp \cap \bar{M}(P)\}$ and the space $L_0^2(P)$ of all possible scores satisfies

$$L_0^2(P) = \bar{M}(P) \oplus \bar{M}(P)^\perp = \bar{T}(P) \oplus \{\bar{T}(P)^\perp \cap \bar{M}(P)\} \oplus \bar{M}(P)^\perp; \tag{21}$$

that is, any score consists of a component that agrees with \mathbf{P} (in $\bar{T}(P)$), a component that disagrees with \mathbf{P} but still agrees with \mathbf{M} (in $\bar{T}(P)^\perp \cap \bar{M}(P)$), and a component that disagrees with \mathbf{M} (in $\bar{M}(P)^\perp$). Intuitively, when testing for the validity of \mathbf{P} while remaining

confident on the correct specification of \mathbf{M} , we should employ tests that direct their power towards the subspace $\bar{T}(P)^\perp \cap \bar{M}(P)$ rather than all of $\bar{T}(P)^\perp$.

In the following, recall that $\Pi_{T^\perp}(\cdot)$ denotes the orthogonal projection under $\|\cdot\|_{P,2}$ onto $\bar{T}(P)^\perp$.

LEMMA 3.2: *Let Assumption 2.1 hold, $\mathbf{P} \subseteq \mathbf{M}$, and $M(P)$ be linear.*

(i) *Let Assumptions 3.1 and 3.3 hold with $S(P) \subseteq \bar{T}(P)^\perp \cap \bar{M}(P)$. Then: Theorem 3.3(i) remains valid for testing (19) for any path $t \mapsto P_{t,g} \in \mathbf{M}$ with $\Pi_S(g) \neq 0$.*

(ii) *Let Assumption 3.1 hold with $\mathbf{T} = \{1, \dots, d\}$, $d < \infty$, and $\{s_\tau\}_{\tau=1}^d$ be an orthonormal basis for $\bar{S}(P) = \bar{T}(P)^\perp \cap \bar{M}(P)$. Then: For any asymptotic level α specification test ϕ_n for (19) with an asymptotic local power function, it follows that*

$$\inf_{g \in \mathcal{G}(\kappa)} \lim_{n \rightarrow \infty} \int \phi_n dP_{1/\sqrt{n},g}^n \leq \inf_{g \in \mathcal{G}(\kappa)} \lim_{n \rightarrow \infty} P_{1/\sqrt{n},g}(\|\hat{G}_n\|^2 > c_{1-\alpha}), \tag{22}$$

where $\mathcal{G}(\kappa) \equiv \{g \in \bar{M}(P) : \|\Pi_{T^\perp}(g)\|_{P,2} \geq \kappa\}$, and $c_{1-\alpha}$ is the $(1 - \alpha)$ quantile of a chi-squared distribution with d degrees of freedom.

(iii) *Let Assumption 3.2 hold with $\hat{\theta}_n$ and $\tilde{\theta}_n$ being efficient estimators of $\theta(P) \in \mathbf{B}$ under \mathbf{P} and \mathbf{M} , respectively. Then: Theorem 3.3(ii) remains valid for testing (19) with $b^*(\nu - \tilde{\nu}) \in \bar{T}(P)^\perp \cap \bar{M}(P)$ for all $b^* \in \mathbf{B}^*$.*

Lemma 3.2(i) revisits the tests examined in Theorem 3.3(i) under the additional requirement that the test focus its power on detecting deviations from \mathbf{P} that remain within \mathbf{M} (i.e., $\bar{S}(P) \subseteq \bar{T}(P)^\perp \cap \bar{M}(P)$) rather than arbitrary deviations from \mathbf{P} (i.e., $\bar{S}(P) \subseteq \bar{T}(P)^\perp$). In order for the resulting test to be able to detect any local deviation of \mathbf{P} that remains within \mathbf{M} , \hat{G}_n must be chosen so that $\bar{S}(P) = \bar{T}(P)^\perp \cap \bar{M}(P)$. When $\bar{T}(P)^\perp \cap \bar{M}(P)$ is finite dimensional, Lemma 3.2(ii) additionally provides a characterization of the optimal specification test in the sense of maximizing local minimum power against alternatives in $\mathbf{M} \setminus \mathbf{P}$ that are a “local distance” of κ away from \mathbf{P} . Specifically, the optimal test corresponds to a quadratic form in \hat{G}_n where \hat{G}_n must be chosen so that it weights every possible local deviation in $\mathbf{M} \setminus \mathbf{P}$ “equally”; that is, $S(P) = \{s_\tau : \tau \in \mathbf{T}\}$ should be an orthonormal basis for $\bar{T}(P)^\perp \cap \bar{M}(P)$.

In parallel to our results in Section 3.3, multiple tests for (19) implicitly possess the structure of the tests described in Lemma 3.2(i)–(ii); see our GMM discussion below. Lemma 3.2(iii), for example, shows that a process \hat{G}_n satisfying the conditions of Lemma 3.2(i) may be obtained by comparing an estimator $\hat{\theta}_n$ that is efficient under \mathbf{P} to an estimator $\tilde{\theta}_n$ that is efficient under the larger model \mathbf{M} . It is again helpful to note that $\hat{\theta}_n$ and $\tilde{\theta}_n$ can be regular estimators of any “smooth” function of $P \in \mathbf{P}$ and need not be of any structural parameter of the model \mathbf{P} . The resulting Hausman type test then satisfies the optimality claim in Lemma 3.2(ii) provided the influence function of $\sqrt{n}\{\hat{\theta}_n - \tilde{\theta}_n\}$ spans $\bar{T}(P)^\perp \cap \bar{M}(P)$; see Remark 3.4. Finally, we emphasize, as in Remark 3.5, that many alternatives to a Hausman type test also satisfy the conditions of Lemmas 3.2(i)–(ii). In fact, the sieve likelihood ratio test of Shen and Shi (2005) and Chen and Liao (2014) for semi/nonparametric likelihood models, and the sieve optimally weighed quasi likelihood ratio test of Chen and Pouzo (2009, 2015) for semi/nonparametric conditional moment restriction models, can be regarded as versions of incremental J tests for (19). These incremental J tests are also applicable to testing hypotheses on structural parameters of a

model \mathbf{M} , in which case \mathbf{P} corresponds to the subset of distributions in \mathbf{M} that satisfy the conjectured null hypothesis on the structural parameters.

GMM ILLUSTRATION—cont.: For \mathbf{P} as defined in (6), we now let $\rho(x, \gamma) = (\rho_1(x, \gamma)', \rho_2(x, \gamma)')$, where $\rho_j : \mathbf{X} \times \mathbf{R}^{d_\gamma} \rightarrow \mathbf{R}^{d_{\rho_j}}$ with $d_{\rho_1} \geq d_\gamma$, and let

$$\mathbf{M} \equiv \left\{ P \in \mathcal{M} : \int \rho_1(\cdot, \gamma) dP = 0 \text{ for some } \gamma \in \Gamma \right\}. \tag{23}$$

Eichenbaum, Hansen, and Singleton (1988) proposed testing \mathbf{P} with \mathbf{M} as a maintained hypothesis by employing an incremental J statistic $J_n(\rho) - J_n(\rho_1)$, where $J_n(\rho)$ and $J_n(\rho_1)$ are the J statistics based on the moments ρ (for \mathbf{P}) and ρ_1 (for \mathbf{M}), respectively. As in Remark 3.5, under their conditions, it can be shown that for $\{p_k\}_{k=1}^{d_\rho - d_\gamma}$ and $\{m_k\}_{k=1}^{d_{\rho_1} - d_\gamma}$ orthonormal bases for $\bar{T}(P)^\perp$ and $\bar{M}(P)^\perp$, we have

$$\begin{aligned} J_n(\rho) - J_n(\rho_1) &= \sum_{k=1}^{d_\rho - d_\gamma} \left(\frac{1}{\sqrt{n}} \sum_{i=1}^n p_k(X_i) \right)^2 - \sum_{k=1}^{d_{\rho_1} - d_\gamma} \left(\frac{1}{\sqrt{n}} \sum_{i=1}^n m_k(X_i) \right)^2 + o_p(1) \\ &= \sum_{k=1}^{d_{\rho_2}} \left(\frac{1}{\sqrt{n}} \sum_{i=1}^n f_k(X_i) \right)^2 + o_p(1), \end{aligned} \tag{24}$$

where the second equality holds for $\{f_k\}_{k=1}^{d_{\rho_2}}$ an orthonormal basis for $\bar{T}(P)^\perp \cap \bar{M}(P)$ since $\bar{M}(P)^\perp \subseteq \bar{T}(P)^\perp$. Therefore, an incremental J test corresponds to a special case of the test discussed in Lemma 3.2(i) for which $\bar{S}(P) = \bar{T}(P)^\perp \cap \bar{M}(P)$. Moreover, by Lemma 3.2(ii), the resulting test is locally maximin optimal. Instead of the statistic $J_n(\rho) - J_n(\rho_1)$, an alternative approach employs the ρ_1 moments for efficient estimation of $\gamma(P)$ (under \mathbf{M}) and the remaining ρ_2 moments for testing; see, for example, Christiano and Eichenbaum (1992), Hansen and Heckman (1996), and Hansen (2010). Such a test corresponds to the Hausman type test in Lemma 3.2(iii). Specifically, $\hat{\theta}_n = 0$ is an efficient estimator of $\theta(P) = \int \rho_2(\cdot, \gamma(P)) dP$ under \mathbf{P} , while an efficient estimator $\hat{\theta}_n$ of $\theta(P)$ under \mathbf{M} equals

$$\frac{1}{n} \sum_{i=1}^n \{ \rho_2(X_i, \hat{\gamma}_n) - \hat{B}'_n \rho_1(X_i, \hat{\gamma}_n) \} \tag{25}$$

for $\hat{\gamma}_n$ an efficient estimator of $\gamma(P)$ using ρ_1 moments (under \mathbf{M}), and \hat{B}_n the OLS coefficients from regressing $\{\rho_2(X_i, \hat{\gamma}_n)\}_{i=1}^n$ on $\{\rho_1(X_i, \hat{\gamma}_n)\}_{i=1}^n$. By Lemma 3.2(ii), an (orthogonalized) quadratic form in (25) leads to a locally maximin optimal test that is asymptotically equivalent to (24).

4. GENERAL NONPARAMETRIC CONDITIONAL MOMENT MODELS

In this section, we apply our previous results to a rich class of models defined by nonparametric conditional moment restrictions with possibly different conditioning sets and potential endogeneity.

4.1. Models and Characterizations

The data distribution P of $X = (Z, W) \in \mathbf{X}$ is assumed to satisfy the following nonparametric conditional moment restrictions:

$$E[\rho_j(Z, h_p)|W_j] = 0 \quad \text{for all } 1 \leq j \leq J \text{ for some } h_p \in \mathbf{H}, \tag{26}$$

for some known measurable mappings $\rho_j : \mathbf{Z} \times \mathbf{H} \rightarrow \mathbf{R}$, where \mathbf{H} is some Banach space (with norm $\|\cdot\|_{\mathbf{H}}$) of measurable functions of $X = (Z, W)$. Here $Z \in \mathbf{Z}$ denotes potentially endogenous random variables, and $W \in \mathbf{W}$ denotes the union of distinct random elements of the conditioning variables (or instruments) (W_1, \dots, W_J) . Note that there are no restrictions imposed on how the conditioning variables relate; for example, W_j and $W_{j'}$ may have all, some, or no elements in common, and some of the W_j could be constants (indicating unconditional moment restrictions).

Model (26) encompasses a very wide array of semiparametric and nonparametric models. It was first studied in Ai and Chen (2007) for root- n consistent estimation of a particular “smooth” linear functional of h_p when the generalized residual functions ρ_j are pointwise differentiable (in h_p) for all $j = 1, \dots, J$. Since Ai and Chen (2007) focused on possibly globally misspecified models, in that P may fail to satisfy (26), they did not characterize the tangent space.

In this section, we characterize the tangent space for model (26) without assuming the differentiability of $\rho_j(Z, \cdot) : \mathbf{H} \rightarrow L^2(P)$ for all j . We assume instead that

$$m_j(W_j, h) \equiv E[\rho_j(Z, h)|W_j] \tag{27}$$

is “smooth” (at h_p) when viewed as a map from \mathbf{H} into $L^2(W_j)$, where $L^2(W_j)$ is the subset of functions $f \in L^2(P)$ depending only on W_j . Specifically, we require Fréchet differentiability of each $m_j(W_j, \cdot) : \mathbf{H} \rightarrow L^2(W_j)$ (at h_p) and denote its derivative by $\nabla m_j(W_j, h_p)$, which could be computed as $\nabla m_j(W_j, h_p)[h] \equiv \frac{\partial}{\partial \tau} m_j(W_j, h_p + \tau h)|_{\tau=0}$ for any $h \in \mathbf{H}$. Employing these derivatives, we may then define the linear map $\nabla m(W, h_p) : \mathbf{H} \rightarrow \bigotimes_{j=1}^J L^2(W_j)$ to be given by

$$\nabla m(W, h_p)[h] \equiv (\nabla m_1(W_1, h_p)[h], \dots, \nabla m_J(W_J, h_p)[h])'. \tag{28}$$

Note that $\bigotimes_{j=1}^J L^2(W_j)$ is itself a Hilbert space when endowed with an inner product (and induced norm) equal to $\langle f, \tilde{f} \rangle \equiv \sum_{j=1}^J E[f_j(W_j)\tilde{f}_j(W_j)]$ for any $f = (f_1, \dots, f_J)$ and $\tilde{f} = (\tilde{f}_1, \dots, \tilde{f}_J)$. The range space \mathcal{R} of the linear map $\nabla m(W, h_p)$ is then defined as

$$\mathcal{R} \equiv \left\{ f \in \bigotimes_{j=1}^J L^2(W_j) : f = \nabla m(W, h_p)[h] \text{ for some } h \in \mathbf{H} \right\}, \tag{29}$$

and we let $\bar{\mathcal{R}}$ be its closure (in $\bigotimes_{j=1}^J L^2(W_j)$), which is a vector space and plays an important role in this section. Finally, we set $\bar{\mathcal{R}}^\perp$ to be the orthocomplement of $\bar{\mathcal{R}}$ (in $\bigotimes_{j=1}^J L^2(W_j)$).

In order to be explicit about the local perturbations we consider, we next introduce a set of conditions on the paths $t \mapsto P_{t,g}$ employed to construct the tangent set $T(P)$.

CONDITION A: (i) Under $X \sim P_{t,g}$, $E[\rho_j(Z, h_t)|W_j] = 0$ for some $h_t \in \mathbf{H}$ and all $1 \leq j \leq J$; (ii) $\|t^{-1}(h_t - h_p) - \Delta\|_{\mathbf{H}} = o(1)$ as $t \downarrow 0$ for some $\Delta \in \mathbf{H}$; (iii) $|\rho_j(z, h_t)| \leq F(z)$ for all $1 \leq j \leq J$ and $F \in L^2(P)$ satisfying $\int F^2 dP_{t,g} = O(1)$ as $t \downarrow 0$.

Thus, a path $t \mapsto P_{t,g}$ satisfies Condition A if whenever X is distributed according to $P_{t,g}$, the conditional moment restrictions in (26) hold for some $h_t \in \mathbf{H}$ and the map $t \mapsto h_t$ is “smooth” in t . These requirements are satisfied, for instance, by the paths considered in semiparametric efficiency calculations, in which distributions are parameterized by h_t and a complementary infinite-dimensional parameter describing aspects of the distribution not characterized by h_t ; see, for example, Begun, Hall, Huang, and Wellner (1983), Hansen (1985), Chamberlain (1986, 1992), Newey (1990), and Ai and Chen (2012). We also introduce a vector space \mathcal{V} given by

$$\mathcal{V} \equiv \left\{ g = \sum_{j=1}^J \rho_j(Z, h_p) \psi_j(W_j) : (\psi_1, \dots, \psi_J) \in \bigotimes_{j=1}^J L^2(W_j) \right\}, \tag{30}$$

which is a subset of $L^2_0(P)$ provided P satisfies (26) and $E[\{\rho_j(Z, h_p)\}^2|W_j]$ is bounded P -a.s. for $1 \leq j \leq J$. Let $\bar{\mathcal{V}}$ be the closure of \mathcal{V} , and $\bar{\mathcal{V}}^\perp$ be the orthocomplement of $\bar{\mathcal{V}}$ (in $L^2_0(P)$).

Finally, we impose the following regularity conditions on the distribution P .

ASSUMPTION 4.1: (i) P satisfies model (26); (ii) $m_j(W_j, \cdot) : \mathbf{H} \rightarrow L^2(W_j)$ is Fréchet differentiable at h_p for $1 \leq j \leq J$; (iii) $\rho_j(Z, \cdot) : \mathbf{H} \rightarrow L^2(P)$ is continuous at h_p for $1 \leq j \leq J$; (iv) there is a $\mathcal{D} \subseteq \mathbf{H}$ such that $\bar{\text{lin}}\{\mathcal{D}\} = \mathbf{H}$ and for every $h \in \mathcal{D}$ there is a $t \mapsto P_{t,g}$ satisfying Condition A with $\Delta = h$; (v) \mathcal{V}^\perp has a dense subset of bounded functions.

ASSUMPTION 4.2: (i) $\sum_{j=1}^J E[\rho_j^2(Z, h_p)|W_j]$ is bounded P -a.s.; (ii) there is $C_0 < \infty$ such that $\sum_{j=1}^J \|\psi_j\|_{P,2} \leq C_0 \|\sum_{j=1}^J \rho_j(\cdot, h_p) \psi_j\|_{P,2}$ for all $(\psi_1, \dots, \psi_J) \in \bigotimes_{j=1}^J L^2(W_j)$.

Assumptions 4.1(i), (ii), (iii), and 4.2(i) are standard. Assumptions 4.1(iv), (v), and 4.2(ii) are sufficient conditions for the simple characterization of the tangent space obtained in Theorem 4.1 below. Assumption 4.1(iv) assumes that \mathbf{H} is the local parameter space for h_p , while Assumption 4.1(v) assumes that any function $g \in \mathcal{V}^\perp$ can be approximated by sequences of bounded functions in \mathcal{V}^\perp —low level sufficient conditions for this requirement are often readily available in specific applications. Assumption 4.2(ii) imposes a linear independence restriction on $\{\rho_j(Z, h_p)\}_{j=1}^J$.

Our next result provides a simple characterization for local overidentification.

THEOREM 4.1: Let P satisfy Assumptions 4.1 and 4.2. Then: $\bar{T}(P)^\perp$ satisfies

$$\bar{T}(P)^\perp = \left\{ g \in L^2_0(P) : g = \sum_{j=1}^J \rho_j(Z, h_p) \psi_j(W_j) \text{ for some } (\psi_1, \dots, \psi_J) \in \bar{\mathcal{R}}^\perp \right\},$$

and moreover, $\bar{T}(P)^\perp = \{0\}$ if and only if $\bar{\mathcal{R}} = \bigotimes_{j=1}^J L^2(W_j)$.

In view of Lemma 2.1, Theorem 4.1 implies that P is locally just identified by a regular model (26) if and only if $\bar{\mathcal{R}} = \bigotimes_{j=1}^J L^2(W_j)$. Heuristically, local overidentification by

model (26) is equivalent to the existence of a nonconstant transformation (ψ_1, \dots, ψ_J) of the conditioning variable that is uncorrelated with the span of the derivative of the conditional expectation with respect to the nonparametric parameter. Theorem 4.1 also has a useful dual representation.

LEMMA 4.1: *Let Assumption 4.1(ii) hold. Let \mathbf{H}^* be the dual space of a Banach space \mathbf{H} , and $\nabla m_j(W_j, h_P)^* : L^2(W_j) \rightarrow \mathbf{H}^*$ be the adjoint of $\nabla m_j(W_j, h_P) : \mathbf{H} \rightarrow L^2(W_j)$ for $j = 1, \dots, J$. Then: $\bar{\mathcal{R}} = \bigotimes_{j=1}^J L^2(W_j)$ if and only if*

$$\left\{ f = (f_1, \dots, f_J) \in \bigotimes_{j=1}^J L^2(W_j) : \sum_{j=1}^J \nabla m_j(W_j, h_P)^*[f_j] = 0 \right\} = \{0\}.$$

Theorem 4.1 and Lemma 4.1 together imply that P is locally just identified by model (26) if and only if the adjoint operator $\nabla m(W, h_P)^* : \bigotimes_{j=1}^J L^2(W_j) \rightarrow \mathbf{H}^*$ is injective. Interestingly, this resembles a necessary condition, the injectivity of $\nabla m(W, h_P) : \mathbf{H} \rightarrow \bigotimes_{j=1}^J L^2(W_j)$, for local identification of the unknown function h_P by model (26) in Chen, Chernozhukov, Lee, and Newey (2014)—when (26) is linear in h_P , this condition is also sufficient (Newey and Powell (2003)).

4.1.1. *Sequential Moment Restrictions*

While the proof of Theorem 4.1 directly computes $\bar{T}(P)^\perp$, we note that in special cases of model (26) for which semiparametric efficiency bounds are known, one could also employ Corollary 3.1 to characterize local just identification. We next follow such an approach by employing the efficiency bound results in Ai and Chen (2012) to characterize the local just identification of P by models defined by sequential moment restrictions.

The data distribution P of $X = (Z, W) \in \mathbf{X}$ is now assumed to satisfy the following nonparametric sequential moment restrictions:

$$\text{model (26) holds with } \sigma(\{W_j\}) \subseteq \sigma(\{W_{j'}\}) \text{ for all } 1 \leq j \leq j' \leq J, \tag{31}$$

where $\sigma(\{W_j\})$ denotes the σ -field generated by W_j for $j = 1, \dots, J$. Note now $W = W_J$, which is assumed to be a non-degenerate random variable.

We will restrict attention to distributions P for which the conditional moments in (31) are suitably linearly independent. To this end, we define

$$s_j^2(W_j) \equiv \inf_{\{a_k\}_{k=j+1}^J} E \left[\left\{ \rho_j(Z, h_P) - \sum_{k=j+1}^J a_k \rho_k(Z, h_P) \right\}^2 \middle| W_j \right] \text{ for } j = 1, \dots, J - 1, \tag{32}$$

and $s_J^2(W_J) \equiv E[\{\rho_J(Z, h_P)\}^2|W_J]$. Since W_j is the most informative conditioning variable, we may interpret $s_j^2(W_j)$ as the residual variance obtained by projecting $\rho_j(Z, h_P)$ on $\{\rho_{j'}(Z, h_P)\}_{j'>j}$ conditionally on all instruments.

The following assumption imposes the basic condition on the distribution P .

ASSUMPTION 4.3: (i) P satisfies (31); (ii) Assumption 4.1(ii) holds; (iii) $\max_j E[\{\rho_j(Z, h)\}^2] < \infty$ for any $h \in \mathbf{H}$ (a Banach space); (iv) $P(\eta \leq E[s_j^2(W_j)|W_j]) = 1$ for some $\eta > 0$ and all $1 \leq j \leq J$; (v) $P(|E[\rho_k(Z, h_P)\rho_j(Z, h_P)|W_j]| \leq M) = 1$ for some $M < \infty$ and all $1 \leq k \leq j \leq J$; (vi) $L^2(W_j)$ is infinite dimensional.

Assumptions 4.3(i), (ii), (iii) are standard. Assumption 4.3(iv) restricts the conditional dependence across moments, while Assumption 4.3(v) imposes an almost sure upper bound in the conditional covariance across residuals. When the same instrument is used in all conditioning equations, so that $W_j = W_j$ for all j , Assumptions 4.3(iv), (v) are equivalent to the covariance matrix of the residuals conditional on W_j being nonsingular and finite uniformly in the support of W_j . Finally, Assumption 4.3(vi) ensures that model (31) implies an infinite number of unconditional moment restrictions. If $L^2(W_j)$ is finite dimensional, then model (31) consists of a finite number of unconditional moment restrictions, thus reducing to the well-understood GMM setting.

THEOREM 4.2: *Let Assumption 4.3 hold. Then: P is locally just identified by model (31) if and only if $\bar{\mathcal{R}} = \bigotimes_{j=1}^J L^2(W_j)$.*

It is interesting that the characterization of local just identification in nonparametric sequential moment restrictions (31) coincides with that for the more general model (26) derived in Theorem 4.1. Nevertheless, the additional structure afforded by sequential moment restrictions does allow for the semiparametric efficiency bound calculation in Ai and Chen (2012) and enables us to obtain the local just identification characterization under lower level conditions.

4.1.2. *Models With Triangular Structures*

Numerous nonparametric structural models possess a triangular structure in which the (conditional) moment restrictions depend on a non-decreasing subset of the parameters; see examples in Section 4.2. Lemma 4.2 below focuses on such a setting by assuming the parameter space takes the form $\mathbf{H} = \bigotimes_{j=1}^J \mathbf{H}_j$ and imposing that the moment conditions can be ordered in a manner such that the k th moment condition depends only on the subset $\bigotimes_{j=1}^k \mathbf{H}_j$ of the parameter space. In the lemma, we let $\nabla m_{j,j}(W_j, h_P)^* : L^2(W_j) \rightarrow \mathbf{H}_j^*$ be the adjoint of $\nabla m_{j,j}(W_j, h_P) : \mathbf{H}_j \rightarrow L^2(W_j)$, and $\bar{\mathcal{R}}_j$ be the closure of \mathcal{R}_j (in $L^2(W_j)$), where \mathcal{R}_j is given by

$$\mathcal{R}_j \equiv \{f \in L^2(W_j) : f = \nabla m_{j,j}(W_j, h_P)[h_j] \text{ for } h_j \in \mathbf{H}_j\}.$$

LEMMA 4.2: *Let Assumption 4.1(ii) hold, and $\mathbf{H} = \bigotimes_{j=1}^J \mathbf{H}_j$ with \mathbf{H}_j being Banach spaces for all j . Suppose there are linear maps $\nabla m_{j,k}(W_j, h_P) : \mathbf{H}_k \rightarrow L^2(W_j)$ such that*

$$\nabla m_j(W_j, h_P)[h] = \sum_{k=1}^J \nabla m_{j,k}(W_j, h_P)[h_k] \quad \text{for any } h = (h_1, \dots, h_J) \in \bigotimes_{j=1}^J \mathbf{H}_j, \quad (33)$$

where $\nabla m_{j,k}(W_j, h_P)[h_k] = 0$ for all $k > j$, and there is $0 \leq C < \infty$ such that

$$\|\nabla m_{j,k}(\cdot, h_P)[h_k]\|_{P,2} \leq C \|\nabla m_{k,k}(\cdot, h_P)[h_k]\|_{P,2} \quad \text{for all } k \leq j. \quad (34)$$

Then: $\bar{\mathcal{R}} = \bigotimes_{j=1}^J L^2(W_j)$ if and only if $\bar{\mathcal{R}}_j = L^2(W_j)$ for all j , which also holds if and only if $\{f \in L^2(W_j) : \nabla m_{j,j}(W_j, h_P)^*[f] = 0\} = \{0\}$ for all j .

Lemma 4.2 implies that, under the stated requirements on the partial derivative maps, one may assess whether P is locally overidentified by examining each (conditional) moment restriction separately. This lemma simplifies the verification of local just identification in many nonparametric models. For example, it is directly applicable to the following

class of models:

$$E[\rho_j(Z, h_{P,j})|W_j] = 0 \quad \text{for some } h_{P,j} \in \mathbf{H}_j \quad \text{for all } 1 \leq j \leq J, \tag{35}$$

where $h_P = (h_{P,1}, \dots, h_{P,J}) \in \mathbf{H} = \otimes_{j=1}^J \mathbf{H}_j$, and the unknown functions $h_{P,j} \in \mathbf{H}_j$ could depend on the endogenous variables Z . Our final result applies Lemma 4.2 to special cases of model (35) in which \mathbf{H}_j contains functions of conditioning variables W_j only; that is,

$$E[\rho_j(Z, h_{P,j}(W_j))|W_j] = 0 \quad \text{for some } h_{P,j} \in \mathbf{H}_j \subseteq L^2(W_j) \text{ for all } 1 \leq j \leq J. \tag{36}$$

COROLLARY 4.1: *Let P satisfy model (36) with $h_P = (h_{P,1}, \dots, h_{P,J}) \in \mathbf{H} = \otimes_{j=1}^J \mathbf{H}_j$ and Assumption 4.1(ii) hold. Suppose, for each $1 \leq j \leq J$, there is $d_j \in L^2(W_j)$ that is bounded P -a.s. and $\nabla m_j(W_j, h_P)[h] = d_j(W_j)h_j(W_j)$ for any $h = (h_1, \dots, h_J) \in \mathbf{H}$. Then: $\bar{\mathcal{R}} = \otimes_{j=1}^J L^2(W_j)$ if and only if, for all $1 \leq j \leq J$, \mathbf{H}_j is dense in $L^2(W_j)$ and $P(d_j(W_j) \neq 0) = 1$.*

Corollary 4.1 reduces assessing local just identification of P by model (36) to examining two simple conditions for all $j = 1, \dots, J$: (i) \mathbf{H}_j must be sufficiently “rich” (\mathbf{H}_j is dense in $L^2(W_j)$), and (ii) the derivative of the moment restrictions must be injective ($d_j(W_j) \neq 0$ P -a.s.). It immediately implies, for example, that nonparametric conditional mean and quantile regression models are locally just identified,⁵ and that restricting the parameter space to the space of bounded or differentiable functions is not sufficient for yielding local overidentification as \mathbf{H}_j remains dense in $L^2(W_j)$. On the other hand, Corollary 4.1 does imply that P will be locally overidentified by model (36) as soon as there is one j such that \mathbf{H}_j is not a dense subset of $(L^2(W_j), \|\cdot\|_{P,2})$. Examples in which \mathbf{H}_j is not dense include, among others, the partially linear or additively separable conditional mean specifications of [Robinson \(1988\)](#) and [Stone \(1985\)](#).

REMARK 4.1: Semiparametric two-step GMM models are widely used in applied work. Building on the insights of [Newey \(1994\)](#) and [Newey and Powell \(1999\)](#), [Ackerberg et al. \(2014\)](#) showed that when the unknown function $h_P = (h_{P,1}, \dots, h_{P,J})$ is “exactly identified” by model (36) in the first stage, the second-stage optimally weighted GMM estimator of γ_P identified by unconditional moment restriction $E[g(X, \gamma_P, h_P)] = 0$ is semiparametrically efficient. Our Corollary 4.1 shows that their requirement of nonparametric “exact identification” of h_P is equivalent to our P being locally just identified by model (36) in the first stage. Our Theorem 4.1, Lemmas 4.1 and 4.2 further imply, however, that the second-stage optimally weighted GMM estimator may be inefficient when P is locally overidentified by model (26) in the first stage, such as in the various semiparametric conditional moment restriction models of [Ai and Chen \(2003, 2012\)](#). More generally, our results imply that the asymptotic variance of a plug-in estimator of a regular functional of h_P can depend on the choice of estimator of h_P whenever P is locally overidentified by model (26). See Section 4.2 for examples of P being locally overidentified by nonparametric models.

⁵For $Z = (Y, W)$ with $Y \in \mathbf{R}$, note that a mean regression model corresponds to $\rho_1(Z, h) = Y - h(W)$, so that $d_1(W) = -1$. Instead, in a quantile regression model, $\rho_1(Z, h) = \tau - 1\{Y \leq h(W)\}$, in which case $d_1(W) = -g_{Y|W}(h_P(W)|W)$ for $g_{Y|W}(y|w)$ the conditional density of Y given W .

4.2. Illustrative Examples

This section presents three empirically relevant examples to illustrate the implications of our results; see the Supplemental Material for additional results and discussion.

EXAMPLE 4.1—Differentiated Products Markets: An extensive literature has studied identification of demand and cost functions in differentiated products markets, including the seminal work of [Berry, Levinsohn, and Pakes \(1995\)](#). Here, we follow [Berry and Haile \(2014\)](#) who derived multiple identification results by relying on moment restrictions of the form

$$E[Y_{ij} - h_{j,P}(V_i)|W_{ij}] = 0 \quad \text{for } 1 \leq j \leq J, \tag{37}$$

where $1 \leq i \leq n$ denotes a market and $1 \leq j \leq J$ a good. For instance, in their analysis of demand, $h_{j,P}$ corresponds to the inverse demand function for good j , V_i denotes market shares and prices in market i , Y_{ij} is a “demand shifter,” and W_{ij} is a vector of price instruments and product/market characteristics for good j . Let $h_{j,P} \in \mathbf{H}_j \subseteq L^2(V)$ for all j and $h_P = (h_{1,P}, \dots, h_{J,P}) \in \mathbf{H} = \bigotimes_{j=1}^J \mathbf{H}_j$. We note that this model is a special case of model (35) and we may then apply Lemma 4.2. To this end, we observe that for any $h = (h_1, \dots, h_J) \in \mathbf{H}$, we have

$$\nabla m_j(W_j, h_P)[h] = -E[h_j(V)|W_j]. \tag{38}$$

Hence, following the notation of Lemma 4.2, $\nabla m_{j,j}(W_j, h_P)[h_j] = -E[h_j(V)|W_j]$ and $\nabla m_{j,k}(W_j, h_P)[h_k] = 0$ when $k \neq j$. To find the adjoint $\nabla m_{j,j}(W_j, h_P)^*$, we let $\bar{\mathbf{H}}_j$ be the closure of \mathbf{H}_j under $\|\cdot\|_{P,2}$ and for any $f \in L^2(W_j)$ define

$$\Pi_{\bar{\mathbf{H}}_j} f \equiv \arg \min_{h_j \in \bar{\mathbf{H}}_j} \|(-f) - h_j\|_{P,2}. \tag{39}$$

By orthogonality of projections, the map $\Pi_{\bar{\mathbf{H}}_j} : L^2(W_j) \rightarrow \bar{\mathbf{H}}_j$ equals $\nabla m_{j,j}(W_j, h_P)^*$ and therefore Lemma 4.2 implies that P is locally just identified if and only if

$$\{f \in L^2(W_j) : \Pi_{\bar{\mathbf{H}}_j} f = 0\} = \{0\} \quad \text{for all } 1 \leq j \leq J. \tag{40}$$

For instance, if $\bar{\mathbf{H}}_j = L^2(V)$, then (40) is equivalent to the distribution of (V, W_j) being L^2 -complete with respect to V for all j ([Newey and Powell \(2003\)](#)), which is an untestable condition under endogeneity ([Andrews \(2017\)](#), [Canay, Santos, and Shaikh \(2013\)](#)).⁶ Hence, plug-in estimation of average derivatives may not be efficient when L^2 -completeness fails ([Ai and Chen \(2012\)](#)). We also note that the structure in model (37) is also present in a large literature on consumer demand; see, for example, [Blundell, Duncan, and Pendakur \(1998\)](#), [Blundell, Browning, and Crawford \(2003\)](#). Semiparametric restrictions that are consistent with agents’ optimization behaviors, however, can render P locally overidentified ([Blundell, Chen, and Kristensen \(2007\)](#), [Chen and Pouzo \(2009\)](#)). Finally, we note very similar arguments apply when (37) consists of conditional quantile restrictions instead—simply note that in such a model $\nabla m_{j,j}(W_j, h_P)[h] = -E[g_{Y_j|V,W_j}(h_{j,P}(V))h_j(V)|W_j]$ for $g_{Y_j|V,W_j}$ the conditional density of Y_j given (V, W_j) .

⁶Since there are examples of distributions for which L^2 -completeness fails ([Santos \(2012\)](#)), the model may be locally overidentified even when V and W_j are of equal dimension.

EXAMPLE 4.2—Nonparametric Selection: This example concerns nonparametric versions of the canonical selection model of Heckman (1979) as studied in, for example, Heckman (1990). Suppose that for each individual i , there are latent variables $(Y_{0,i}^*, Y_{1,i}^*)$ satisfying

$$Y_{d,i}^* = g_{d,P}(V_i) + U_{d,i}, \tag{41}$$

where $d \in \{0, 1\}$, V_i is a set of regressors, and $g_{d,P}$ are unknown functions. Instead of $(Y_{0,i}^*, Y_{1,i}^*)$, we observe $Y_i = Y_{0,i}^* + D_i(Y_{1,i}^* - Y_{0,i}^*)$, where $D_i \in \{0, 1\}$ indicates selection into “treatment.” As in Heckman and Vytlacil (2005), we assume there exists a variable R_i excluded from $g_{d,P}$ and impose the index sufficiency requirement

$$E[U_{d,i}|V_i, R_i, D_i = d] = \lambda_{d,P}(P(D_i = 1|V_i, R_i)) \tag{42}$$

for unknown functions $\lambda_{d,P}$. Assuming $E[U_{d,i}|V_i] = 0$ for $d \in \{0, 1\}$, we can then employ equations (41) and (42) to obtain the system of conditional moment restrictions

$$E[D_i - s_P(V_i, R_i)|V_i, R_i] = 0, \tag{43}$$

$$E[Y_i - g_{d,P}(V_i) - \lambda_{d,P}(s_P(V_i, R_i))|V_i, R_i, D_i = d] = 0, \tag{44}$$

which can be used to identify the conditional average treatment effect $g_{1,P}(V_i) - g_{0,P}(V_i)$; see also Newey, Powell, and Vella (1999) and Das, Newey, and Vella (2003) for related models. Hence, in this context, $J = 2$, $h_P = (g_{0,P}, g_{1,P}, \lambda_{0,P}, \lambda_{1,P}, s_P)$, $W_{i1} = (V_i, R_i)$, and $W_{i2} = (V_i, R_i, D_i)$.

We examine a general nonparametric version of this model by only requiring $g_{d,P} \in L^2(V)$ and $\lambda_{d,P}$ be continuously differentiable for $d \in \{0, 1\}$. For any $(g_0, g_1, \lambda_0, \lambda_1, s) \in \mathbf{H}$, restrictions (43) and (44) then possess a sequential moment structure which simplifies applying Theorems 4.1 or 4.2. In particular, Lemma 4.2 implies P is locally just identified if and only if

$$\mathcal{S}_d \equiv \{f \in L^2(V, R) : f(V, R) = g_d(V) + \lambda_d(s_P(V, R)) \text{ for some } g_d, \lambda_d\} \tag{45}$$

is dense in $L^2(V, R)$ for $d \in \{0, 1\}$. However, identification of the functions $g_{d,P}$ and $\lambda_{d,P}$ requires

$$P(\text{Var}\{s_P(V, R)|V\} > 0) > 0, \tag{46}$$

that is, the instrument R must be relevant. When (46) holds, \mathcal{S}_d is not dense in $L^2(V, R)$. Thus, the conditions for the identification of $(g_{d,P}, \lambda_{d,P})$ imply that P is locally overidentified by the model. Hence, the model is testable and efficiency considerations matter when estimating smooth parameters such as the average treatment effects.

EXAMPLE 4.3—Nonparametric Production: This example closely follows the firm’s production structural models proposed by Olley and Pakes (1996), Akerberg, Caves, and Frazer (2015), and others. Econometricians observe a random sample $\{X_i\}_{i=1}^n$ of a panel of firms $i = 1, \dots, n$ from the distribution of $X = \{Y_t, K_t, L_t, I_t\}_{t=1}^T$ for a fixed finite $T \geq 2$, where Y_t, K_t, L_t, I_t respectively denote a firm’s log output, capital, labor, and investment levels at time t . Suppose that

$$Y_{it} = g_P(K_{it}, L_{it}) + \omega_{it} + U_{it}, \quad E[U_{it}|K_{it}, L_{it}, I_{it}] = 0, \tag{47}$$

where g_P is an unknown function, and ω_{it} is a productivity factor observed by the firm but not the econometrician. Olley and Pakes (1996) provided conditions under which the firm’s dynamic optimization problem implies, for some unknown function λ_P , that

$$\omega_{it} = \lambda_P(K_{it}, I_{it}). \tag{48}$$

Let $W = (K_1, L_1, I_1)$, and for simplicity let $T = 2$ and ω_{it} follow an AR(1) process. The literature has employed (47) and (48) to derive the semiparametric conditional moment restrictions

$$E[Y_1 - g_P(K_1, L_1) - \lambda_P(K_1, I_1)|W] = 0, \tag{49}$$

$$E[Y_2 - g_P(K_2, L_2) - \pi_P \lambda_P(K_1, I_1)|W] = 0, \tag{50}$$

where π_P is the coefficient in the AR(1) process for ω_{it} . This model contains multiple overidentifying restrictions that are easily characterized through Theorem 4.1. Specifically, note $h_P = (g_P, \lambda_P, \pi_P)$, and for any $h = (g, \lambda, \pi)$, we have

$$\nabla m_1(W, h_P)[h] = -g(K_1, L_1) - \lambda(K_1, I_1), \tag{51}$$

$$\nabla m_2(W, h_P)[h] = -E[g(K_2, L_2)|W] - \pi_P \lambda(K_1, I_1) - \pi \lambda_P(K_1, I_1). \tag{52}$$

By Theorem 4.1, a necessary condition for local just identification is for the closure of the range of (51) to equal $L^2(W)$. However, this requirement fails since (51) cannot approximate nonseparable functions $f \in L^2((L_1, I_1))$ —a failure reflecting the assumption that labor is not a dynamic variable. Consequently, sequential estimation of average output elasticities, as in Olley and Pakes (1996), can be inefficient. Similarly, we note

$$\begin{aligned} &\nabla m_2(W, h_P)[h] - \pi_P \nabla m_1(W, h_P)[h] \\ &= \pi_P g(K_1, L_1) - E[g(K_2, L_2)|W] - \pi \lambda_P(K_1, I_1), \end{aligned} \tag{53}$$

and local just identification requires the closure of the range of (53) to equal $L^2(W)$. However, such a condition can fail reflecting the empirical content of assuming constancy of g_P through time and additive separability of ω_{it} . As in Section 3.4, the power of specification tests can be directed at violations of these assumptions. See our working paper version, Cowles Foundation Discussion Paper No. 1999R, for a numerical illustration.

5. EXTENSION TO $T(P)$ BEING A CONVEX CONE

Our main theoretical results in Section 3 rely on the requirement that the tangent set $T(P)$ be linear. In this section, we examine whether the main conclusions could be extended to models in which $T(P)$ is a convex cone—a setting that can arise, for example, in mixture models (van der Vaart (1989)) and models where a parameter is on a boundary (Andrews (1999)). To this end, we replace Assumption 2.1 with the following weaker condition:

ASSUMPTION 5.1: (i) Assumption 2.1(i) holds; (ii) the tangent set $T(P)$ is a convex cone; that is, if $g, f \in T(P)$, $a, b \in \mathbf{R}$ with $a \geq 0$ and $b \geq 0$, then $ag + bf \in T(P)$.

We let $\bar{T}(P)$ still denote the closure of $T(P)$ under $\|\cdot\|_{P,2}$ and maintain Definition 2.2. Crucially, Assumption 5.1(ii) implies $\bar{T}(P)$ is a closed convex cone but not necessarily a

closed linear subspace of $L_0^2(P)$ as in regular models. Thus, the alternative characterization of local overidentification in terms of the orthogonal complement of $\bar{T}(P)$ is no longer valid (see Lemma 2.1). However, for any closed convex cone $\bar{T}(P)$ in $L_0^2(P)$, we may define its polar cone, denoted $\bar{T}(P)^-$, which is given by

$$\bar{T}(P)^- \equiv \left\{ g \in L_0^2(P) : \int gf dP \leq 0 \text{ for all } f \in \bar{T}(P) \right\}. \tag{54}$$

Let $\Pi_T(g)$ and $\Pi_{T^-}(g)$ denote the metric projections of a $g \in L_0^2(P)$ onto $\bar{T}(P)$ and $\bar{T}(P)^-$, respectively. For any $g \in L_0^2(P)$, the so-called ‘‘Moreau decomposition’’ (Moreau (1962)) implies

$$g = \Pi_T(g) + \Pi_{T^-}(g), \quad \int \{ \Pi_T(g) \} \{ \Pi_{T^-}(g) \} dP = 0. \tag{55}$$

Unlike the setting in which $\bar{T}(P)$ is a linear subspace, however, there may in fact exist $f \in \bar{T}(P)$ and $g \in \bar{T}(P)^-$ such that $\int fg dP < 0$. Nevertheless, the decomposition in (55) immediately implies the following direct generalization of Lemma 2.1.

LEMMA 5.1: *Under Assumption 5.1, the following are equivalent to Definition 2.2:*

- (i) *P is locally just identified by P if and only if $\bar{T}(P)^- = \{0\}$.*
- (ii) *P is locally overidentified by P if and only if $\bar{T}(P)^- \neq \{0\}$.*

By definition, it is clear that Theorem 3.1 remains valid for the case that P is locally just identified by \mathbf{P} (i.e., $\bar{T}(P) = L_0^2(P)$). Given Lemma 5.1, it should also be possible to establish results similar to Theorem 3.2 for the locally overidentified case. That is, if P is locally overidentified by \mathbf{P} , then the model should be locally testable and ‘‘efficiency’’ should ‘‘matter’’ even when $T(P)$ is a convex cone. To gain some intuition, we can again rely on the sample means of scores $0 \neq \tilde{f} \in L_0^2(P)$. Recall that if $X_i \sim P_{1/\sqrt{n},g}$ for any path $t \mapsto P_{t,g} \in \mathcal{M}$, then $\mathbb{G}_n(\tilde{f}) \equiv \frac{1}{\sqrt{n}} \sum_{i=1}^n \tilde{f}(X_i)$ is asymptotically normally distributed with mean $\int \tilde{f}g dP$ (see equation (12)). By Lemma 5.1, P being locally overidentified by \mathbf{P} is equivalent to the existence of a $0 \neq \tilde{f} \in \bar{T}(P)^-$. For any such \tilde{f} , it follows that $\int \tilde{f}g dP \leq 0$ for all $g \in \bar{T}(P)$. Thus, for the purposes of specification testing, observing a large and positive value for $\mathbb{G}_n(\tilde{f})$ may be viewed as a ‘‘signal’’ that the distribution of $X_i \sim P_{1/\sqrt{n},g}$ is approaching P from outside the model \mathbf{P} . On the other hand, from an estimation perspective, we should be able to employ the knowledge that $\int \tilde{f}g dP \leq 0$ for all $g \in \bar{T}(P)$ to improve on ‘‘inefficient’’ estimators.

The potential lack of orthogonality between $\bar{T}(P)$ and $\bar{T}(P)^-$, however, presents some important complications. For instance, it is no longer natural to restrict attention to regular estimators. We instead focus on a broader class of estimators for parameter $\theta(P) \in \mathbf{B}$ satisfying

$$\sqrt{n} \{ \hat{\theta}_n - \theta(P_{1/\sqrt{n},g}) \} \xrightarrow{L_{n,g}} \mathbb{Z}_g \tag{56}$$

for some tight random variable $\mathbb{Z}_g \in \mathbf{B}$ along any path $t \mapsto P_{t,g} \in \mathbf{P}$. Note that in contrast to regular estimators, the limit \mathbb{Z}_g may depend on g . Focusing on estimators satisfying (56) enables us to easily characterize the local asymptotic risk along any path $t \mapsto P_{t,g} \in \mathbf{P}$.

Concretely, for a loss function $\Psi : \mathbf{B} \rightarrow \mathbf{R}_+$, the local asymptotic risk of $\hat{\theta}_n$ is given by

$$\limsup_{n \rightarrow \infty} E_{P_{1/\sqrt{n},g}} [\Psi(\sqrt{n}\{\hat{\theta}_n - \theta(P_{1/\sqrt{n},g})\})], \tag{57}$$

which represents the expected loss of employing $\hat{\theta}_n$ to estimate $\theta(P)$ when the data generating process is locally perturbed within \mathbf{P} . For simplicity, we consider Ψ -loss functions, defined as follows:

DEFINITION 5.1: Ψ -loss is a map from \mathbf{B} to \mathbf{R}_+ such that: (i) $\{b \in \mathbf{B} : \Psi(b) \leq t\}$ is convex for all $t \in \mathbf{R}$; (ii) $\Psi(0) = 0$ and $\Psi(b) = \Psi(-b)$; (iii) Ψ is bounded, continuous, and nonconstant.

A minimal requirement on an estimator is that its local asymptotic risk not be dominated by that of an alternative estimator; that is, a sensible estimator should be “asymptotically locally admissible.”

DEFINITION 5.2: $\hat{\theta}_n : \{X_i\}_{i=1}^n \rightarrow \mathbf{B}$ is “asymptotically locally admissible” for $\theta(P)$ under Ψ -loss if it satisfies (56) and there is no estimator $\tilde{\theta}_n : \{X_i\}_{i=1}^n \rightarrow \mathbf{B}$ satisfying (56) and

$$\limsup_{n \rightarrow \infty} E_{P_{1/\sqrt{n},g}} [\Psi(\sqrt{n}\{\tilde{\theta}_n - \theta(P_{1/\sqrt{n},g})\})] \leq \limsup_{n \rightarrow \infty} E_{P_{1/\sqrt{n},g}} [\Psi(\sqrt{n}\{\hat{\theta}_n - \theta(P_{1/\sqrt{n},g})\})]$$

for all paths $t \mapsto P_{t,g} \in \mathbf{P}$, and with the inequality holding strictly for some path $t \mapsto P_{t,g} \in \mathbf{P}$.

Given the introduced concepts, we can document an equivalence result between the local overidentification of P , the importance of “efficiency” in estimation, and the potential refutability of a model.

THEOREM 5.1: *Let Assumption 5.1 hold. Then: the following statements are equivalent:*

- (i) P is locally overidentified by \mathbf{P} .
- (ii) There exists a bounded function $f : \mathbf{X} \rightarrow \mathbf{R}$ such that $\sum_{i=1}^n f(X_i)/n$ is not an asymptotically locally admissible estimator for $\theta(P) = \int f dP$ under any Ψ -loss.
- (iii) There exists a local asymptotic level α test ϕ_n for (9) with a local asymptotic power function π satisfying $\pi(g) > \alpha$ for some path $t \mapsto P_{t,g} \in \mathcal{M} \setminus \mathbf{P}$.

Theorem 3.2 and Theorem 5.1 reflect both the similarities and the differences between regular and non-regular models. With regard to estimation, for example, Theorems 3.2(i), (ii) and 5.1(i), (ii) both show that local overidentification of P is equivalent to the availability of “efficiency” gains in estimation. However, since in non-regular models we need to consider a broader class of estimators than just regular estimators, Theorem 5.1(i), (ii) links the availability of “efficiency” gains to the local overidentification of P through the estimation of simple “smooth” maps $\theta(P) = \int f dP$ (population means) for bounded functions f . In particular, while sample means are always locally admissible when P is locally just identified (see Lemma C.1 in Appendix C), Theorem 5.1(ii) shows this fails to be the case when P is locally overidentified.

With regard to specification testing, Theorem 3.2(i), (iii) and Theorem 5.1(i), (iii) both show that local overidentification of P is equivalent to the potential refutability of the model. However, important differences also exist in the properties of local specification

tests for regular and non-regular models. Notably, our next result shows that for any non-regular model whose convex cone $\bar{T}(P)$ contains at least two linearly independent elements, any asymptotically locally unbiased specification test for (9) will have local power no larger than its level against an important class of alternatives.

THEOREM 5.2: *Let Assumption 5.1 hold and let there be linearly independent $f_1, f_2 \in \bar{T}(P)^-$ with $\lambda f_1, \lambda f_2 \in \bar{T}(P)$ for any $\lambda \leq 0$. Let ϕ_n be any specification test for (9) with a local asymptotic power function π such that $\pi(g) \leq \alpha$ for all $g \in \bar{T}(P)$ and $\pi(g) \geq \alpha$ for all $g \notin \bar{T}(P)$. Then: $\pi(g) = \alpha$ for any path $t \mapsto P_{t,g} \in \mathcal{M} \setminus \mathbf{P}$ with $\lambda \Pi_{\bar{T}^-}(g) \in \bar{T}(P)$ for any $\lambda \leq 0$.*

Given Theorem 5.1(iii), Theorem 5.2 does not preclude the existence of asymptotically nontrivial specification tests, but rather implies such tests can necessarily be asymptotically locally biased for non-regular models whose convex cone $\bar{T}(P)$ is not a ray. We next examine in more detail the construction of both such specification tests and of “better” estimators than the sample mean. To this end, we impose the following:

ASSUMPTION 5.2: *For some set \mathbf{T} , there is a statistic $\hat{\mathbb{G}}_n : \{X_i\}_{i=1}^n \rightarrow \ell^\infty(\mathbf{T})$ satisfying:*

- (i) $\hat{\mathbb{G}}_n(\tau) = \frac{1}{\sqrt{n}} \sum_{i=1}^n s_\tau(X_i) + o_p(1)$ uniformly in $\tau \in \mathbf{T}$, where $0 \neq s_\tau \in \bar{T}(P)^-$ for all $\tau \in \mathbf{T}$;
- (ii) Assumption 3.1(ii) holds.

Assumption 5.2(i) is identical to Assumption 3.1(i) except that s_τ is required to belong to $\bar{T}(P)^-$ instead of $\bar{T}(P)^+$. As in Theorem 3.3(i), $\hat{\mathbb{G}}_n$ can be employed to construct a specification test for (9). For any $0 \leq \omega \in \ell^\infty(\mathbf{T})$, we define $\hat{\mathbb{G}}_n^\omega(\tau) \equiv \omega(\tau) \times \hat{\mathbb{G}}_n(\tau)$ and $\mathbb{G}_0^\omega(\tau) \equiv \omega(\tau) \times \mathbb{G}_0(\tau)$ for $\tau \in \mathbf{T}$. Let $c_{1-\alpha}^\omega$ be the $1 - \alpha$ quantile of $\|\max\{\mathbb{G}_0^\omega, 0\}\|_\infty$. We then define the test

$$\phi_n^\omega \equiv 1 \{ \|\max\{\hat{\mathbb{G}}_n^\omega, 0\}\|_\infty > c_{1-\alpha}^\omega \}. \tag{58}$$

Intuitively, $0 \leq \omega \in \ell^\infty(\mathbf{T})$ is a weight function that determines the local alternatives against which ϕ_n^ω has nontrivial power. In parallel to Theorem 3.3(i), the power properties of ϕ_n^ω also depend on the set $C(P) \equiv \{s_\tau \in \bar{T}(P)^- : \tau \in \mathbf{T}\}$ being sufficiently “rich.” Let $\bar{C}(P)$ denote the closed convex cone generated by $C(P)$ (in $L_0^2(P)$), that is, $\bar{C}(P)$ parallels $\bar{S}(P)$ in Theorem 3.3(i). For any $g \in L_0^2(P)$, we let $\Pi_C(g)$ denote the metric projection of g onto $\bar{C}(P)$.

Our next result shows that for any path $t \mapsto P_{t,g} \in \mathcal{M}$ with $\Pi_C(g) \neq 0$, it is possible to select an $\omega^* \in \ell^\infty(\mathbf{T})$ such that the corresponding specification test $\phi_n^{\omega^*}$ has nontrivial local power against that alternative. Given Theorem 5.2, $\phi_n^{\omega^*}$ can be asymptotically locally biased, however.

THEOREM 5.3: *Let Assumptions 5.1, 5.2 hold, and $0 \leq \omega \in \ell^\infty(\mathbf{T})$ satisfy $c_{1-\alpha}^\omega > 0$. Then: ϕ_n^ω is a local asymptotic level α test for (9) with a local asymptotic power function. Moreover, for any $t \mapsto P_{t,g} \in \mathcal{M}$ with $\Pi_C(g) \neq 0$, there is $0 \leq \omega^* \in \ell^\infty(\mathbf{T})$ with $c_{1-\alpha}^{\omega^*} > 0$ for $\alpha \in (0, \frac{1}{2})$, and*

$$\lim_{n \rightarrow \infty} P_{1/\sqrt{n},g}(\|\max\{\hat{\mathbb{G}}_n^{\omega^*}, 0\}\|_\infty > c_{1-\alpha}^{\omega^*}) > \alpha. \tag{59}$$

Turning to estimation, we note that when restricting attention to paths $t \mapsto P_{t,g} \in \mathbf{P}$, knowledge that $\int s_\tau g dP \leq 0$ should be useful. Specifically, for any bounded function $f : \mathbf{X} \rightarrow \mathbf{R}$, we define

$$\hat{\mu}_n(f, \tau) \equiv \frac{1}{n} \sum_{i=1}^n f(X_i) - \beta(f, \tau) \times n^{-1/2} \max\{\hat{G}_n(\tau), 0\} \tag{60}$$

for $\beta(f, \tau) \equiv \max\{\int f s_\tau dP, 0\} / \|s_\tau\|_{p,2}^2$. The function $\beta(f, \tau)s_\tau$ is the projection of f onto the cone generated by $s_\tau \in \bar{T}(P)^-$ in $L_0^2(P)$. Our final theorem shows that when P is locally overidentified, $\hat{\mu}_n(f, \tau)$ can be viewed as a more “efficient” estimator for $\theta(P) = \int f dP$ than the sample mean whenever $f \notin \bar{T}(P)$. It is analogous to Lemma 3.1(i).

THEOREM 5.4: *Let Assumptions 5.1 and 5.2 hold. Then: for any bounded $f : \mathbf{X} \rightarrow \mathbf{R}$ and $\tau^* \in \mathbf{T}$ such that $\int f s_{\tau^*} dP > 0$ and f and s_{τ^*} are linearly independent, we have: $\hat{\mu}_n(f, \tau^*)$ defined in (60) satisfies (56) and*

$$\begin{aligned} & \limsup_{n \rightarrow \infty} E_{P_{1/\sqrt{n},g}} \left[\Psi \left(\sqrt{n} \left\{ \hat{\mu}_n(f, \tau^*) - \int f dP_{1/\sqrt{n},g} \right\} \right) \right] \\ & < \limsup_{n \rightarrow \infty} E_{P_{1/\sqrt{n},g}} \left[\Psi \left(\frac{1}{\sqrt{n}} \sum_{i=1}^n \left\{ f(X_i) - \int f dP_{1/\sqrt{n},g} \right\} \right) \right] \end{aligned} \tag{61}$$

for any path $t \mapsto P_{t,g} \in \mathbf{P}$ and any Ψ -loss.

Thus, when P is locally overidentified, there is information in the model that can be employed to both render the model testable (Theorem 5.3) and to obtain “efficiency” gains (Theorem 5.4). As a result, the local testability of a model and “efficiency” considerations remain intrinsically linked to P being locally overidentified by \mathbf{P} even when $\bar{T}(P)$ is a convex cone. We emphasize that many important issues, such as optimality in estimation and specification testing, and analog of incremental J test for (19), remain open when $\bar{T}(P)$ is a convex cone. We leave these questions for future research.

6. CONCLUSION

This paper reinterprets the common practice of counting the numbers of restrictions and parameters of interest in GMM to determine overidentification as an approach that examines whether the tangent space is a strict subset of $L_0^2(P)$. This abstraction naturally leads to a notion of local overidentification, which we show is responsible for an intrinsic link between efficiency considerations in estimation and the local testability of a model. While we have relied on an i.i.d. assumption for simplicity, there are ample works deriving efficiency bounds in time series settings (Hansen (1985, 1993)) and characterizing limit experiments under nonstationary, strongly dependent data (Ploberger and Phillips (2012)). We conjecture the results in this paper could similarly be extended to allow for dependence, but leave such extensions for future work.

APPENDIX A: LIMITING EXPERIMENT

In this appendix, we embed specification tests and regular estimators in a common statistical experiment that highlights their connection to each other and to the local overi-

identification of P . The main result in this appendix, Theorem A.1 below, plays an important role in the proofs of our main results in Section 3, and is therefore presented here for completeness. The proof of Theorem A.1 can be found in the Supplemental Material.

Heuristically, in an asymptotic framework that is local to P , our parameter uncertainty is over what “direction” P is being approached from. We may intuitively interpret such a direction as the score g of $P_{1/\sqrt{n},g}$ and represent our parameter uncertainty as possessing only a “noisy” measure of g . Let $d_T \equiv \dim\{\bar{T}(P)\}$ and $d_{T^\perp} \equiv \dim\{\bar{T}(P)^\perp\}$ denote the (possibly infinite) dimensions of the tangent space and its orthogonal complement. Both $\bar{T}(P)$ and $\bar{T}(P)^\perp$ are Hilbert spaces with norm $\|\cdot\|_{P,2}$ and hence there exist orthonormal bases $\{\psi_k^T\}_{k=1}^{d_T}$ and $\{\psi_k^{T^\perp}\}_{k=1}^{d_{T^\perp}}$ for $\bar{T}(P)$ and $\bar{T}(P)^\perp$, respectively. We then consider a random variable $(\mathbb{Y}^T, \mathbb{Y}^{T^\perp}) \in \mathbf{R}^{d_T} \times \mathbf{R}^{d_{T^\perp}}$ whose law is such that the vectors $\mathbb{Y}^T \equiv (\mathbb{Y}_1^T, \dots, \mathbb{Y}_{d_T}^T)'$ and $\mathbb{Y}^{T^\perp} \equiv (\mathbb{Y}_1^{T^\perp}, \dots, \mathbb{Y}_{d_{T^\perp}}^{T^\perp})'$ have mutually independent coordinates and satisfy, for some (unknown) $g_0 \in L_0^2(P)$, the relation

$$\begin{aligned} \mathbb{Y}_k^T &\sim N\left(\int g_0 \psi_k^T dP, 1\right) \quad \text{for } 1 \leq k \leq d_T, \\ \mathbb{Y}_k^{T^\perp} &\sim N\left(\int g_0 \psi_k^{T^\perp} dP, 1\right) \quad \text{for } 1 \leq k \leq d_{T^\perp}. \end{aligned} \tag{A.1}$$

Here, if $d_{T^\perp} = 0$, then we interpret \mathbb{Y}^{T^\perp} as being equal to zero with probability 1. Finally, we let Q_g denote the distribution of $(\mathbb{Y}^T, \mathbb{Y}^{T^\perp}) \in \mathbf{R}^{d_T} \times \mathbf{R}^{d_{T^\perp}}$ when (A.1) holds with $g_0 = g \in L_0^2(P)$. Thus, by definition, we know that the (unknown) distribution Q_{g_0} of $(\mathbb{Y}^T, \mathbb{Y}^{T^\perp})$ belongs to the nonparametric family $\{Q_g : g \in L_0^2(P)\}$.

The following theorem formalizes the connection between specification tests, regular estimators, and the tangent space through the introduced limiting experiment.⁷ Recall that \xrightarrow{L} means convergence in law under $P^n \equiv \otimes_{i=1}^n P$.

THEOREM A.1: *Under Assumption 2.1, the following two propositions hold:*

(i) *Let ϕ_n be any local asymptotic level α specification test for (9) with a local asymptotic power function π . Then: there is a level α test $\phi : (\mathbb{Y}^T, \mathbb{Y}^{T^\perp}) \rightarrow [0, 1]$ of*

$$H_0 : \Pi_{T^\perp}(g_0) = 0, \quad H_1 : \Pi_{T^\perp}(g_0) \neq 0 \tag{A.2}$$

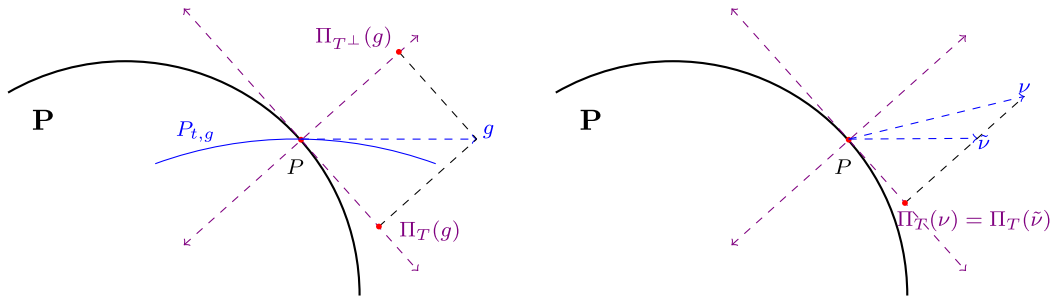
based on one observation $(\mathbb{Y}^T, \mathbb{Y}^{T^\perp})$ such that $\pi(g_0) = \int \phi dQ_{g_0}$ for all $g_0 \in L_0^2(P)$.

(ii) *(Convolution Theorem) Let $\hat{\theta}_n : \{X_i\}_{i=1}^n \rightarrow \mathbf{B}$ be any asymptotically linear regular estimator of any parameter $\theta(P) \in \mathbf{B}$. Then: for any $b^* \in \mathbf{B}^*$, there exist linear maps $F^T : \mathbf{R}^{d_T} \rightarrow \mathbf{R}$ and $F^{T^\perp} : \mathbf{R}^{d_{T^\perp}} \rightarrow \mathbf{R}$ such that under the law P^n , it follows that*

$$\sqrt{n}\{b^*(\hat{\theta}_n) - b^*(\theta(P))\} \xrightarrow{L} F^T(\mathbb{Y}^T) + F^{T^\perp}(\mathbb{Y}^{T^\perp}), \tag{A.3}$$

where $(\mathbb{Y}^T, \mathbb{Y}^{T^\perp}) \sim Q_{g_0}$ with $g_0 = 0$, and the map $F^T : \mathbf{R}^{d_T} \rightarrow \mathbf{R}$ depends on $\theta : \mathbf{P} \rightarrow \mathbf{B}$ and $b^ \in \mathbf{B}^*$ but not on the estimator $\hat{\theta}_n$.*

⁷See Choi, Hall, and Schick (1996) and Hirano and Porter (2009) for other applications of the limit experiment.



(a) Intuition for Theorem A.1(i)

(b) Intuition for Theorem A.1(ii)

FIGURE 1.—The tangent space, specification tests, and regular estimators.

Theorem A.1(i) relates the local properties of any specification test for \mathbf{P} (as in (9)) to a testing problem concerning the unknown distribution Q_{g_0} of $(\mathbb{Y}^T, \mathbb{Y}^{T^\perp})$ through the asymptotic representation theorem (van der Vaart (1991a)). Intuitively, any path $t \mapsto P_{t,g}$ that approaches P from outside the model \mathbf{P} should be such that its score does not belong to the tangent set or, equivalently, $\Pi_{T^\perp}(g) \neq 0$; see Figure 1(a). In contrast, the score g of any submodel $t \mapsto P_{t,g} \in \mathbf{P}$ must belong to the tangent set, implying $\Pi_{T^\perp}(g) = 0$. Thus, any specification test for \mathbf{P} should behave locally as a test of the null hypothesis in (A.2). Theorem A.1(i) formalizes these heuristics by showing that if π is the local asymptotic power function of a specification test for \mathbf{P} (as in (9)), then π must also be the power function of a test of (A.2) based on a single observation $(\mathbb{Y}^T, \mathbb{Y}^{T^\perp})$ whose (unknown) law Q_{g_0} is known to belong to the nonparametric family $\{Q_g : g \in L_0^2(P)\}$.

Theorem A.1(ii) is essentially the convolution theorem of Hájek (1970), stated here in a manner that facilitates a connection to Theorem A.1(i). To gain intuition on this result, we focus on the scalar case ($\mathbf{B} = \mathbf{R}$) and suppose there are two asymptotically linear regular estimators $\hat{\theta}_n$ and $\tilde{\theta}_n$ of a common parameter with influence functions ν and $\tilde{\nu}$, respectively. Regularity constrains the projection of ν and $\tilde{\nu}$ onto the tangent space to be equal, which originates a term in the asymptotic distribution that is independent of the choice of estimator ($F^T(\mathbb{Y}^T)$); see Figure 1(b). The estimators, however, may differ on a component that is extraneous to the model ($\Pi_{T^\perp}(\nu) \neq \Pi_{T^\perp}(\tilde{\nu})$), contributing a “noise” term to the asymptotic distribution that depends on the choice of estimator ($F^{T^\perp}(\mathbb{Y}^{T^\perp})$). An efficient estimator is the one for which the “noise” component is zero.

Crucial for our purposes is the observation that \mathbb{Y}^{T^\perp} plays fundamental yet distinct roles in the asymptotic behavior of both specification tests and regular estimators. From a specification testing perspective, \mathbb{Y}^{T^\perp} is a partially sufficient statistic for $\Pi_{T^\perp}(g)$ and is needed to construct any nontrivial test of (A.2). In contrast, from a regular estimation perspective, \mathbb{Y}^{T^\perp} is an ancillary statistic that can only contribute “noise” to estimators.⁸ Thus, the limit experiment requires P to be locally overidentified ($\bar{T}(P)^\perp \neq \{0\}$) in order to allow for both nontrivial tests and asymptotically distinct estimators.

⁸In regular estimation, only paths within the model are considered; see Definition 3.1. The resulting limiting experiment is then indexed by $\{Q_g : g \in \bar{T}(P)\}$, in which \mathbb{Y}^{T^\perp} is ancillary.

APPENDIX B: PROOFS FOR SECTIONS 2 AND 3

In this appendix, we present proofs of the theoretical results in Sections 2 and 3. All of the additional technical lemmas used in this appendix can be found in the Supplemental Material.

PROOF OF LEMMA 2.1: Since $T(P)$ is linear by Assumption 2.1(ii), $\bar{T}(P)$ is a vector subspace of $L_0^2(P)$, and therefore $L_0^2(P) = \bar{T}(P) \oplus \bar{T}(P)^\perp$; see, for example, Theorem 3.4.1 in Luenberger (1969). The claims of the lemma then immediately follow from $\bar{T}(P) = L_0^2(P)$ if and only if $\bar{T}(P)^\perp = \{0\}$. Q.E.D.

PROOF OF THEOREM 3.1: To establish part (i) of the theorem, we let ν and $\tilde{\nu}$ denote the influence functions of $\hat{\theta}_n$ and $\tilde{\theta}_n$, respectively. Then note, for any $b^* \in \mathbf{B}^*$ and $\lambda \in \mathbf{R}$, that

$$\begin{aligned} & \sqrt{n}\{b^*(\lambda\hat{\theta}_n + (1-\lambda)\tilde{\theta}_n) - b^*(\theta(P))\} \\ &= \frac{1}{\sqrt{n}} \sum_{i=1}^n \{ \lambda b^*(\nu(X_i)) + (1-\lambda)b^*(\tilde{\nu}(X_i)) \} + o_p(1) \xrightarrow{L} N(0, \sigma_\lambda^2) \end{aligned} \tag{B.1}$$

for $\sigma_\lambda^2 = \|b^*(\lambda\nu + (1-\lambda)\tilde{\nu})\|_{P,2}^2$ by asymptotic linearity and the central limit theorem. Further note that if P is locally just identified, then Theorem A.1(ii) implies σ_λ^2 does not depend on λ . However, since $\|b^*(\lambda\nu + (1-\lambda)\tilde{\nu})\|_{P,2}^2$ being constant in λ implies that $\|b^*(\nu - \tilde{\nu})\|_{P,2} = 0$, and $b^* \in \mathbf{B}^*$ was arbitrary, we can conclude that

$$b^*(\sqrt{n}\{\hat{\theta}_n - \tilde{\theta}_n\}) = \frac{1}{\sqrt{n}} \sum_{i=1}^n b^*(\nu(X_i) - \tilde{\nu}(X_i)) + o_p(1) = o_p(1) \tag{B.2}$$

for any $b^* \in \mathbf{B}^*$. Since $\sqrt{n}\{\hat{\theta}_n - \tilde{\theta}_n\}$ is asymptotically tight and measurable by Lemmas 1.4.3 and 1.4.4 in van der Vaart and Wellner (1996), result (B.2) and Lemma E.1 (in the Supplemental Material) imply $\sqrt{n}\{\hat{\theta}_n - \tilde{\theta}_n\} = o_p(1)$ in \mathbf{B} , which establishes part (i) of the theorem.

To establish part (ii) of the theorem, we note that by Theorem A.1(i), there exists a level α test ϕ of (A.2) such that, for any $g \in L_0^2(P)$ and path $t \mapsto P_{t,g}$,

$$\lim_{n \rightarrow \infty} \int \phi_n dP_{1/\sqrt{n},g}^n = \int \phi dQ_g. \tag{B.3}$$

However, if P is locally just identified by \mathbf{P} , then $\bar{T}(P) = L_0^2(P)$, or equivalently, $\bar{T}(P)^\perp = \{0\}$. Therefore, the null hypothesis in (A.2) holds for all $g \in L_0^2(P)$, which implies $\int \phi dQ_g \leq \alpha$ for all $g \in L_0^2(P)$, and part (ii) of the theorem holds by (B.3). Q.E.D.

PROOF OF THEOREM 3.2: First, by Theorem 3.1, it follows that (ii) implies (i) and that (iii) implies (i). Therefore, it suffices to show that (i) (i.e., P being locally overidentified by \mathbf{P}) implies that both (ii) and (iii) hold. To this end, we observe that if P is locally overidentified by \mathbf{P} , then Lemma 2.1 implies there exists a $0 \neq \tilde{f} \in \bar{T}(P)^\perp$, which without loss of generality we assume satisfies $\|\tilde{f}\|_{P,2} = 1$. We next aim to employ such a \tilde{f} to verify that (ii) and (iii) indeed hold.

To establish that (i) implies (ii), we note that $\mathbb{G}_n(\tilde{f}) \equiv \sum_{i=1}^n \tilde{f}(X_i)/\sqrt{n}$ trivially satisfies Assumption 3.1(i) and Assumption 3.1(ii) with $\mathbb{G}_0 \sim N(0, 1)$ since $\|\tilde{f}\|_{P,2} = 1$ and $\tilde{f} \in L_0^2(P)$. Thus, part (i) implying part (ii) is a special case of Lemma 3.1(i).

To establish that (i) implies (iii), we let $\mathbb{Z} \sim N(0, 1)$ and note that Theorem 3.10.12 in van der Vaart and Wellner (1996) implies that, for any path $t \mapsto P_{t,g} \in \mathcal{M}$,

$$\mathbb{G}_n(\tilde{f}) \xrightarrow{L_{n,g}} \mathbb{Z} + \int \tilde{f}g dP \tag{B.4}$$

since $\|\tilde{f}\|_{P,2} = 1$. For $z_{1-\alpha/2}$ the $(1 - \alpha/2)$ quantile of a standard normal distribution, we define the test $\phi_n \equiv 1\{|\mathbb{G}_n(\tilde{f})| > z_{1-\alpha/2}\}$. Then:(B.4) and the Portmanteau theorem imply that

$$\pi(g) \equiv \lim_{n \rightarrow \infty} \int \phi_n dP_{1/\sqrt{n},g} = P\left(\left|\mathbb{Z} + \int \tilde{f}g dP\right| > z_{1-\alpha/2}\right) \tag{B.5}$$

for any path $t \mapsto P_{t,g} \in \mathcal{M}$. Hence, (B.5) implies ϕ_n indeed has a local asymptotic power function. Moreover, since $\tilde{f} \in \bar{T}(P)^\perp$, result (B.5) implies $\pi(g) = \alpha$ whenever $g \in \bar{T}(P)$, which establishes ϕ_n is a local asymptotic level α specification test. In addition, for any $g \in \bar{T}(P)^\perp$, we have either $\int \tilde{f}g dP = 0$ (and hence $\pi(g) = \alpha$ by (B.5)), or $\int \tilde{f}g dP \neq 0$ (and hence $\pi(\tilde{f}) > \alpha$ by (B.5)). Thus, this test is locally unbiased. Finally, there exists a path $t \mapsto P_{t,\tilde{f}} \in \mathcal{M}$ with score $\tilde{f} \in \bar{T}(P)^\perp$, in which case (B.5) implies $\pi(\tilde{f}) > \alpha$ and hence (i) implies (iii). Q.E.D.

PROOF OF COROLLARY 3.1: First note that since every $f \in \mathcal{D}$ is bounded, $\theta_f(P) \equiv \int f dP$ is pathwise differentiable at P relative to $T(P)$ with derivative $\dot{\theta}_f(g) \equiv \int \Pi_T(f)g dP$; see Lemma F.1 (in the Supplemental Material). Therefore, by Theorem 5.2.1 in Bickel et al. (1993), its efficiency bound is given by $\Omega_f^* = \|\Pi_T(f)\|_{P,2}^2$. For any $f \in L^2(P)$, let $\Pi_{L_0^2(P)}(f)$ denote its projection onto $L_0^2(P)$ and note that $\Pi_{L_0^2(P)}(f) = \{f - \int f dP\}$, and hence $\text{Var}\{f(X)\} = \|\Pi_{L_0^2(P)}(f)\|_{P,2}^2$. By orthogonality of $\bar{T}(P)$ and $\bar{T}(P)^\perp$, then

$$\begin{aligned} \text{Var}\{f(X)\} &= \|\Pi_{L_0^2(P)}(f)\|_{P,2}^2 = \|\Pi_T(\Pi_{L_0^2(P)}(f)) + \Pi_{T^\perp}(\Pi_{L_0^2(P)}(f))\|_{P,2}^2 \\ &= \|\Pi_T(\Pi_{L_0^2(P)}(f))\|_{P,2}^2 + \|\Pi_{T^\perp}(\Pi_{L_0^2(P)}(f))\|_{P,2}^2 = \Omega_f^* + \|\Pi_{T^\perp}(f)\|_{P,2}^2, \end{aligned} \tag{B.6}$$

where in the final equality we used $\Pi_T(\Pi_{L_0^2(P)}(f)) = \Pi_T(f)$ and $\Pi_{T^\perp}(\Pi_{L_0^2(P)}(f)) = \Pi_{T^\perp}(f)$ for any $f \in L^2(P)$ due to $\bar{T}(P)$ and $\bar{T}(P)^\perp$ being subspaces of $L_0^2(P)$. Thus, by (B.6), $\text{Var}\{f(X)\} = \Omega_f^*$ for all $f \in \mathcal{D}$ if and only if $\Pi_{T^\perp}(f) = 0$ for all $f \in \mathcal{D}$, which by denseness of \mathcal{D} is equivalent to $\bar{T}(P)^\perp = \{0\}$. Q.E.D.

PROOF OF LEMMA 3.1: For part (i) of the lemma, note that since \mathbb{G}_0 is non-degenerate, Assumption 3.1(i) implies $s_{\tau^*} \neq 0$ for some $\tau^* \in \mathbf{T}$, and for a $0 \neq \tilde{b} \in \mathbf{B}$, we set

$$\tilde{\theta}_n \equiv \hat{\theta}_n + \tilde{b} \times n^{-1/2} \hat{\mathbb{G}}_n(\tau^*). \tag{B.7}$$

Notice $\hat{\theta}_n$ is asymptotically linear by hypothesis and denote its influence function by ν . Assumption 3.1(i), definition (B.7), and the continuous mapping theorem then yield

$$\sqrt{n}\{\tilde{\theta}_n - \theta(P)\} = \frac{1}{\sqrt{n}} \sum_{i=1}^n \{\nu(X_i) + \tilde{b} \times s_{\tau^*}(X_i)\} + o_p(1). \tag{B.8}$$

Setting $\tilde{\nu}(X_i) \equiv \nu(X_i) + \tilde{b} \times s_{\tau^*}(X_i)$, we obtain for any $b^* \in \mathbf{B}^*$ that $b^*(\tilde{\nu}) = \{b^*(\nu) + b^*(\tilde{b}) \times s_{\tau^*}\} \in L_0^2(P)$ since $b^*(\nu) \in L_0^2(P)$ due to $\hat{\theta}_n$ being asymptotically linear and $s_{\tau^*} \in \bar{T}(P)^\perp \subseteq L_0^2(P)$ by Assumption 3.1(i). Hence, (B.8) implies $\tilde{\theta}_n$ is indeed asymptotically linear and its influence function equals $\tilde{\nu}$. Moreover, by Lemma D.4 (in the Supplemental Material), $(\sqrt{n}\{\hat{\theta}_n - \theta(P)\}, \frac{1}{\sqrt{n}} \sum_{i=1}^n s_{\tau^*}(X_i))$ converge jointly in distribution in $\mathbf{B} \times \mathbf{R}$ under P^n , and hence the continuous mapping theorem implies

$$\sqrt{n}\{\tilde{\theta}_n - \theta(P)\} = \sqrt{n}\{\hat{\theta}_n - \theta(P)\} + \tilde{b} \times \left\{ \frac{1}{\sqrt{n}} \sum_{i=1}^n s_{\tau^*}(X_i) \right\} \xrightarrow{L} \mathbb{Z} \tag{B.9}$$

on \mathbf{B} under P^n for some tight Borel random variable \mathbb{Z} . In addition, we have that

$$\sqrt{n}\{\hat{\theta}_n - \tilde{\theta}_n\} = -\tilde{b} \times \left\{ \frac{1}{\sqrt{n}} \sum_{i=1}^n s_{\tau^*}(X_i) \right\} \xrightarrow{L} \Delta \tag{B.10}$$

by the central limit and continuous mapping theorems. Further note that since $\tilde{b} \neq 0$, we trivially have $\Delta \neq 0$ in \mathbf{B} because $b^*(\Delta) \sim N(0, \|b^*(\tilde{b})s_{\tau^*}\|_{P,2}^2)$ and $\|b^*(\tilde{b})s_{\tau^*}\|_{P,2} > 0$ for some $b^* \in \mathbf{B}^*$ since $\tilde{b} \neq 0$. Thus, to conclude the proof of part (i), it only remains to show that $\tilde{\theta}_n$ is regular. To this end, let $t \mapsto P_{t,g} \in \mathbf{P}$, and note Lemma 25.14 in van der Vaart (1998) yields

$$\sum_{i=1}^n \log \left(\frac{dP_{1/\sqrt{n},g}}{dP}(X_i) \right) = \frac{1}{\sqrt{n}} \sum_{i=1}^n g(X_i) - \frac{1}{2} \int g^2 dP + o_p(1) \tag{B.11}$$

under P^n , and thus Example 3.10.6 in van der Vaart and Wellner (1996) implies P^n and $P_{1/\sqrt{n},g}^n$ are mutually contiguous. Since $\tilde{\theta}_n$ is asymptotically linear, $(\sqrt{n}\{\tilde{\theta}_n - \theta(P)\}, \frac{1}{\sqrt{n}} \sum_{i=1}^n g(X_i))$ converge jointly in $\mathbf{B} \times \mathbf{R}$ by Lemma D.4. Hence, by (B.11) and Lemma A.8.6 in Bickel et al. (1993), we obtain that

$$\sqrt{n}\{\tilde{\theta}_n - \theta(P)\} \xrightarrow{L_{n,g}} \mathbb{Z}_g \tag{B.12}$$

for some tight Borel \mathbb{Z}_g on \mathbf{B} . Furthermore, since $T(P)$ is linear by Assumption 2.1(ii), and $\hat{\theta}_n$ is regular by hypothesis, Lemma D.4 and Theorem 5.2.3 in Bickel et al. (1993) imply there is a bounded linear map $\dot{\theta}: \bar{T}(P) \rightarrow \mathbf{B}$ such that, for any $t \mapsto P_{t,g} \in \mathbf{P}$,

$$\lim_{t \downarrow 0} \|t^{-1}\{\theta(P_{t,g}) - \theta(P)\} - \dot{\theta}(g)\|_{\mathbf{B}} = 0. \tag{B.13}$$

Therefore, combining (B.12) and (B.13) and the continuous mapping theorem yields

$$\sqrt{n}\{\tilde{\theta}_n - \theta(P_{1/\sqrt{n},g})\} \xrightarrow{L_{n,g}} \mathbb{Z}_g + \dot{\theta}(g). \tag{B.14}$$

Next, we note that for any $b^* \in \mathbf{B}^*$, (B.9), (B.11), and the central limit theorem imply

$$\left(\sqrt{n} \left\{ b^*(\tilde{\theta}_n) - b^*(\theta(P)) \right\} \right. \\ \left. \sum_{i=1}^n \log \left(\frac{dP_{1/\sqrt{n},g}}{dP} (X_i) \right) \right) \xrightarrow{L} N \left(\begin{bmatrix} 0 \\ -\frac{1}{2} \int g^2 dP \end{bmatrix}, \Sigma \right) \tag{B.15}$$

under P^n , where since $\int g s_{\tau^*} dP = 0$ due to $g \in T(P)$ and $s_{\tau^*} \in \bar{T}(P)^\perp$, we have

$$\Sigma = \begin{bmatrix} \int (b^*(\nu) + b^*(\tilde{b})s_{\tau^*})^2 dP & \int b^*(\nu)g dP \\ \int b^*(\nu)g dP & \int g^2 dP \end{bmatrix}. \tag{B.16}$$

In addition, since $b^*(\hat{\theta}_n)$ is an asymptotically linear regular estimator of $b^*(\theta(P))$, Proposition 3.3.1 in Bickel et al. (1993) and $g \in \bar{T}(P)$ imply $\int b^*(\nu)g dP = b^*(\dot{\theta}(g))$. Hence, results (B.15) and (B.16) and Lemma A.9.3 in Bickel et al. (1993) establish

$$\sqrt{n} \{ b^*(\tilde{\theta}_n) - b^*(\theta(P_{1/\sqrt{n},g})) \} \xrightarrow{L^{n,g}} N \left(0, \int (b^*(\nu) + b^*(\tilde{b})s_{\tau^*})^2 dP \right). \tag{B.17}$$

Define $\zeta_{b^*}(X_i) \equiv \{ b^*(\nu(X_i)) + b^*(\tilde{b})s_{\tau^*}(X_i) \}$, and for any finite collection $\{b_k^*\}_{k=1}^K \subset \mathbf{B}^*$ let $(\mathbb{W}_{b_1^*}, \dots, \mathbb{W}_{b_K^*})$ denote a multivariate normal vector with $E[\mathbb{W}_{b_k^*}] = 0$ for all $1 \leq k \leq K$ and $E[\mathbb{W}_{b_k^*} \mathbb{W}_{b_j^*}] = E[\zeta_{b_k^*}(X_i) \zeta_{b_j^*}(X_i)]$ for any $1 \leq j \leq k \leq K$. Letting $C_b(\mathbf{R}^K)$ denote the set of continuous and bounded functions on \mathbf{R}^K , we then obtain from (B.14), (B.17), the Cramer–Wold device, and the continuous mapping theorem that

$$E[f(b_1^*(\mathbb{Z}_g + \dot{\theta}(g)), \dots, b_K^*(\mathbb{Z}_g + \dot{\theta}(g)))] = E[f(b_1^*(\mathbb{W}_{b_1^*}), \dots, b_K^*(\mathbb{W}_{b_K^*}))], \tag{B.18}$$

for any $f \in C_b(\mathbf{R}^K)$. Since $\mathcal{G} \equiv \{ f \circ (b_1^*, \dots, b_K^*) : f \in C_b(\mathbf{R}^K), \{b_k^*\}_{k=1}^K \subset \mathbf{B}^*, 1 \leq K < \infty \}$ is a vector lattice that separates points in \mathbf{B} , it follows from Lemma 1.3.12 in van der Vaart and Wellner (1996) that there is a unique tight Borel measure \mathbb{W} on \mathbf{B} satisfying (B.18). In particular, since the right-hand side of (B.18) does not depend on g , we conclude that the law of $\mathbb{Z}_g + \dot{\theta}(g)$ is constant in g , establishing the regularity of $\tilde{\theta}_n$.

For part (ii) of the lemma, we let ν and $\tilde{\nu}$ denote the influence functions of $\hat{\theta}_n$ and $\tilde{\theta}_n$, respectively, and note that since $\|b^*\|_{\mathbf{B}^*} \leq 1$ for all $b^* \in \mathbf{T}$, it follows that

$$\sup_{b^* \in \mathbf{T}} \left| \hat{\mathbb{G}}_n(b^*) - \frac{1}{\sqrt{n}} \sum_{i=1}^n b^*(\nu(X_i) - \tilde{\nu}(X_i)) \right| \\ \leq \sup_{b^* \in \mathbf{B}^*} \|b^*\|_{\mathbf{B}^*} \times \left\| \sqrt{n} \{ \hat{\theta}_n - \tilde{\theta}_n \} - \frac{1}{\sqrt{n}} \sum_{i=1}^n \{ \nu(X_i) - \tilde{\nu}(X_i) \} \right\|_{\mathbf{B}} = o_p(1). \tag{B.19}$$

Moreover, note that since $b^*(\hat{\theta}_n)$ and $b^*(\tilde{\theta}_n)$ are both asymptotically linear regular estimators of the parameter $b^*(\theta(P)) \in \mathbf{R}$, Proposition 3.3.1 in Bickel et al. (1993) implies

$$\Pi_T(b^*(\nu)) = \Pi_T(b^*(\tilde{\nu})). \tag{B.20}$$

In particular, since $b^*(\nu) \in L_0^2(P)$, we may decompose $b^*(\nu) = \Pi_T(b^*(\nu)) + \Pi_{T^\perp}(b^*(\nu))$, and therefore, applying an identical argument to $b^*(\tilde{\nu})$, we can conclude that

$$b^*(\nu - \tilde{\nu}) = \Pi_{T^\perp}(b^*(\nu)) - \Pi_{T^\perp}(b^*(\tilde{\nu})) \tag{B.21}$$

by result (B.20). It follows that $b^*(\nu - \tilde{\nu}) \in \tilde{T}(P)^\perp$ for any $b^* \in \mathbf{T}$, which together with (B.19) verifies that Assumption 3.1(i) holds. Next, define $F : \mathbf{B} \rightarrow \ell^\infty(\mathbf{T})$ to be given by $F(b)(b^*) = b(b^*)$ for any $b \in \mathbf{B}$, and note F is linear and in addition

$$\|F(b)\|_\infty = \sup_{\|b^*\|_{\mathbf{B}^*} \leq 1} |b(b^*)| = \|b\|_{\mathbf{B}}, \tag{B.22}$$

due to the definition of \mathbf{T} and Lemma 6.10 in Aliprantis and Border (2006). In particular, (B.22) implies F is continuous, and by the continuous mapping theorem, we obtain

$$\hat{\mathbb{G}}_n = F(\sqrt{n}(\hat{\theta}_n - \tilde{\theta}_n)) \xrightarrow{L} F(\Delta) \text{ in } \ell^\infty(\mathbf{T}). \tag{B.23}$$

Let $\mathbb{G}_0 \equiv F(\Delta)$ and note Gaussianity of \mathbb{G}_0 follows by (B.19). Moreover, we note there must exist a $b^* \in \mathbf{B}^*$ such that $\|b^*(\nu - \tilde{\nu})\|_{P,2} > 0$, for otherwise Lemma E.1 (in the Supplemental Material) would imply $\Delta = 0$, contradicting Assumption 3.2. Hence, \mathbb{G}_0 is in addition non-degenerate, which verifies Assumption 3.1(ii). *Q.E.D.*

PROOF OF THEOREM 3.3: For part (i) of the theorem, we note that Lemma E.2 (in the Supplemental Material), Assumption 3.3(i), and the continuous mapping theorem imply that for any path $t \mapsto P_{t,g} \in \mathcal{M}$,

$$\Psi(\hat{\mathbb{G}}_n) \xrightarrow{L_{n,g}} \Psi(\mathbb{G}_0 + \Delta_g), \tag{B.24}$$

where $\Delta_g : \mathbf{T} \rightarrow \mathbf{R}$ satisfies $\Delta_g(\tau) = \int s_\tau g dP$ for any $\tau \in \mathbf{T}$. Further note that by direct calculation, $\Delta_{-g} = -\Delta_g$, and hence Lemma E.3 (in the Supplemental Material) implies $-\Delta_g$ belongs to the support of \mathbb{G}_0 for any $g \in L_0^2(P)$. In particular, since $\Psi(0) = 0$ and $\Psi(b) \geq 0$ for all $b \in \ell^\infty(\mathbf{T})$, it follows that, for any $c > 0$, there exists an open neighborhood N_c of $-\Delta_g \in \ell^\infty(\mathbf{T})$ such that $0 \leq \Psi(b + \Delta_g) \leq c$ for all $b \in N_c$. Thus, we can conclude, for any $c > 0$, that

$$P(\Psi(\mathbb{G}_0 + \Delta_g) \leq c) \geq P(\mathbb{G}_0 \in N_c) > 0, \tag{B.25}$$

where the final inequality follows from $-\Delta_g$ belonging to the support of \mathbb{G}_0 . Next, note that Theorem 7.1.7 in Bogachev (2007) implies \mathbb{G}_0 is a regular measure, and hence, since it is tight by Assumption 3.1(ii), it follows that it is also a Radon measure. Together with the convexity of the map $\Psi(\cdot + \Delta_g) : \ell^\infty(\mathbf{T}) \rightarrow \mathbf{R}$, \mathbb{G}_0 being Radon allows us to apply Theorem 11.1 in Davydov, Lifshits, and Smorodina (1998) to conclude that the point

$$c_0 \equiv \inf\{c : P(\Psi(\mathbb{G}_0 + \Delta_g) \leq c) > 0\} \tag{B.26}$$

is the only possible discontinuity point of the c.d.f. of $\Psi(\mathbb{G}_0 + \Delta_g)$. However, since $\Psi(b) \geq 0$ for all $b \in \ell^\infty(\mathbf{T})$, result (B.25) holding for any $c > 0$ implies that $c_0 = 0$. In particular, $c_{1-\alpha} > 0$ by hypothesis implies that $c_{1-\alpha}$ is a continuity point of the c.d.f. of $\Psi(\mathbb{G}_0 + \Delta_g)$ for any $g \in L_0^2(P)$. Therefore, result (B.24) allows us to conclude: for any path $t \mapsto P_{t,g} \in \mathcal{M}$,

$$\pi(g) \equiv \lim_{n \rightarrow \infty} P_{1/\sqrt{n},g}(\Psi(\hat{\mathbb{G}}_n) > c_{1-\alpha}) = P(\Psi(\mathbb{G}_0 + \Delta_g) > c_{1-\alpha}), \tag{B.27}$$

which establishes that the test ϕ_n indeed has an asymptotic local power function. Moreover, if $t \mapsto P_{t,g} \in \mathbf{P}$, then Lemma E.2 (in the Supplemental Material) implies $\Delta_g = 0$ and hence result (B.27) yields

$$\lim_{n \rightarrow \infty} P_{1/\sqrt{n},g}(\Psi(\hat{\mathbb{G}}_n) > c_{1-\alpha}) = P(\Psi(\mathbb{G}_0) > c_{1-\alpha}) = \alpha, \tag{B.28}$$

where we exploited that $c_{1-\alpha}$ is the $1 - \alpha$ quantile of $\Psi(\mathbb{G}_0)$ and that the c.d.f. of $\Psi(\mathbb{G}_0)$ is continuous at $c_{1-\alpha}$. Thus, we conclude from (B.28) that ϕ_n is also an asymptotic level α specification test. On the other hand, we note that $\Delta_g \neq 0$ whenever $\Pi_S(g) \neq 0$ since $\Delta_g(\tau) = \int s_\tau g dP$ and $\bar{S}(P) = \overline{\text{lin}}\{s_\tau : \tau \in \mathbf{T}\}$. In addition, Theorem 3.6.1 in Bogachev (1998) implies the support of \mathbb{G}_0 is a separable vector subspace of $\ell^\infty(\mathbf{T})$, and hence $\Delta_g \neq 0$ belonging to the support of \mathbb{G}_0 and Lemma E.5 (in the Supplemental Material) establish

$$P(\Psi(\mathbb{G}_0 + \Delta_g) < c_{1-\alpha}) < P(\Psi(\mathbb{G}_0) < c_{1-\alpha}) = 1 - \alpha. \tag{B.29}$$

We can now exploit that $c_{1-\alpha} > 0$ is a continuity point of the c.d.f. of $\Psi(\mathbb{G}_0 + \Delta_g)$ together with results (B.27) and (B.29) to conclude that, for any path $t \mapsto P_{t,g} \in \mathcal{M}$ with $\Pi_S(g) \neq 0$,

$$\lim_{n \rightarrow \infty} P_{1/\sqrt{n},g}(\Psi(\hat{\mathbb{G}}_n) > c_{1-\alpha}) = 1 - P(\Psi(\mathbb{G}_0 + \Delta_g) \leq c_{1-\alpha}) > \alpha, \tag{B.30}$$

which satisfies (15). Finally, for any path $t \mapsto P_{t,g} \in \mathcal{M}$ with $\Pi_S(g) = 0$, we have $\Delta_g = 0$ and hence $\pi(g) = \alpha$ by equations (B.27) and (B.28). Thus, the test ϕ_n is also locally unbiased, and we establish part (i) of the theorem.

For part (ii) of the theorem, we proceed as in Lemma 3.1(ii) and set $\mathbf{T} = \{b^* \in \mathbf{B}^* : \|b^*\|_{\mathbf{B}^*} \leq 1\}$ and $\hat{\mathbb{G}}_n(b^*) = \sqrt{n}b^*(\hat{\theta}_n - \tilde{\theta}_n)$ for any $b^* \in \mathbf{B}^*$. Since $s_{b^*} = b^*(\nu - \tilde{\nu})$ by Lemma 3.1(ii), we obtain by definition that $S(P) = \{b^*(\nu - \tilde{\nu}) : b^* \in \mathbf{T}\} = \{b^*(\nu - \tilde{\nu}) : b^* \in \mathbf{B}^*, \|b^*\|_{\mathbf{B}^*} \leq 1\}$. Moreover, if $\Psi = \|\cdot\|_\infty$, then

$$\Psi(\hat{\mathbb{G}}_n) = \sup_{\|b^*\|_{\mathbf{B}^*} \leq 1} |b^*(\sqrt{n}\{\hat{\theta}_n - \tilde{\theta}_n\})| = \sqrt{n}\|\hat{\theta}_n - \tilde{\theta}_n\|_{\mathbf{B}}, \tag{B.31}$$

where the final equality follows by Lemma 6.10 in Aliprantis and Border (2006). Since $\Psi = \|\cdot\|_\infty$ satisfies Assumption 3.3, the second claim of the theorem follows. *Q.E.D.*

PROOF OF LEMMA 3.2: Part (i) of the lemma is immediate since $\bar{T}(P)^\perp \cap \bar{M}(P) \subseteq \bar{T}(P)^\perp$.

For part (ii) of the lemma, we will exploit Theorem A.1(i) and its notation. Note that Theorem A.1(i) implies there exists a level α test $\phi : (\mathbb{Y}^T, \mathbb{Y}^{T^\perp}) \rightarrow [0, 1]$ of the hypothesis in (A.2), and such that, for any path $t \mapsto P_{t,g} \in \mathcal{M}$,

$$\lim_{n \rightarrow \infty} \int \phi_n dP_{1/\sqrt{n},g}^n = \int \phi dQ_g, \tag{B.32}$$

where Q_g denotes the (unknown) distribution of $(\mathbb{Y}^T, \mathbb{Y}^{T^\perp})$ as defined in (A.1). Recall $\{s_\tau\}_{\tau=1}^d$ with $d < \infty$ is an orthonormal basis for $\bar{T}(P)^\perp \cap \bar{M}(P)$; we let d_r denote the (possibly infinite) dimension of $\bar{M}(P)^\perp$ and $\{r_k\}_{k=1}^{d_r}$ be an orthonormal basis for $\bar{M}(P)^\perp$. By (21), $\{s_\tau\}_{\tau=1}^d \cup \{r_k\}_{k=1}^{d_r}$ is then an orthonormal basis for $\bar{T}(P)^\perp$. Thus, in Theorem A.1(i), we may set $\{\psi_k^{T^\perp}\}_{k=1}^{d_r} = \{s_\tau\}_{\tau=1}^d \cup \{r_k\}_{k=1}^{d_r}$, which implies we may write $\mathbb{Y}^{T^\perp} = (\mathbb{M}, \mathbb{R}) \in \mathbf{R}^d \times \mathbf{R}^{d_r}$,

where the vectors $\mathbb{M} \equiv (\mathbb{M}_1, \dots, \mathbb{M}_d)'$ and $\mathbb{R} = (\mathbb{R}_1, \dots, \mathbb{R}_{d_r})'$ have mutually independent coordinates, and whenever $(\mathbb{Y}^T, \mathbb{Y}^{T^\perp})$ are distributed according to Q_g , the induced distribution on (\mathbb{M}, \mathbb{R}) is

$$\begin{aligned} \mathbb{M}_\tau &\sim N\left(\int g s_\tau dP, 1\right) \quad \text{for } 1 \leq \tau \leq d, \\ \mathbb{R}_k &\sim N\left(\int g r_k dP, 1\right) \quad \text{for } 1 \leq k \leq d_r. \end{aligned} \tag{B.33}$$

Let Φ denote the standard normal measure on \mathbf{R} . Note that $Q_0 = \bigotimes_{k=1}^{d_r} \Phi \times \bigotimes_{k=1}^d \Phi \times \bigotimes_{k=1}^{d_r} \Phi$, and define a test $\bar{\phi} : \mathbb{M} \rightarrow [0, 1]$ to be given by

$$\bar{\phi}(\mathbb{M}) \equiv E_{Q_0}[\phi(\mathbb{Y}^T, \mathbb{M}, \mathbb{R}) | \mathbb{M}], \tag{B.34}$$

where the expectation is taken over $(\mathbb{Y}^T, \mathbb{Y}^{T^\perp}) \sim Q_0$. Since $\{g \in \bar{T}(P)^\perp \cap \bar{M}(P) : \|g\|_{p,2} \geq \kappa\} \subseteq \mathcal{G}(\kappa)$, we can conclude from result (B.32) that

$$\begin{aligned} \inf_{g \in \mathcal{G}(\kappa)} \lim_{n \rightarrow \infty} \int \phi_n dP^n_{1/\sqrt{n},g} &\leq \inf_{g \in \bar{T}(P)^\perp \cap \bar{M}(P) : \|g\|_{p,2} \geq \kappa} \int \phi dQ_g \\ &= \inf_{g \in \bar{T}(P)^\perp \cap \bar{M}(P) : \|g\|_{p,2} \geq \kappa} \int \bar{\phi} d\left\{ \bigotimes_{\tau=1}^d \Phi\left(\cdot - \int g s_\tau dP\right) \right\}, \end{aligned} \tag{B.35}$$

where in the equality we exploited (B.34), the independence of $(\mathbb{Y}^T, \mathbb{R})$ and \mathbb{M} , and that for any $g \in \bar{T}(P)^\perp \cap \bar{M}(P)$ it follows that $(\mathbb{Y}^T, \mathbb{R}) \sim \bigotimes_{k=1}^{d_r} \Phi \times \bigotimes_{k=1}^d \Phi$ under Q_g as a result of g being orthogonal to $\{\psi_k^T\}_{k=1}^{d_r} \cup \{r_k\}_{k=1}^{d_r}$. Finally, note that

$$\begin{aligned} \inf_{g \in \bar{T}(P)^\perp \cap \bar{M}(P) : \|g\|_{p,2} \geq \kappa} \int \bar{\phi} d\left\{ \bigotimes_{\tau=1}^d \Phi\left(\cdot - \int g s_\tau dP\right) \right\} \\ = \inf_{h \in \mathbf{R}^d : \|h\| \geq \kappa} \int \bar{\phi} d\left\{ \bigotimes_{\tau=1}^d \Phi(\cdot - h_k) \right\} \end{aligned} \tag{B.36}$$

by Parseval's equality and where $h = (h_1, \dots, h_d)$. Let $\chi_d^2(\kappa)$ denote a chi-squared random variable with d degrees of freedom and non-centrality parameter κ . It then follows from $\int \bar{\phi} d\{\bigotimes_{\tau=1}^d \Phi\} \leq \alpha$ due to (B.34) and ϕ being a level α test of (A.2), results (B.35) and (B.36), and Problem 8.29 in Lehmann and Romano (2005), that

$$\inf_{g \in \mathcal{G}(\kappa)} \lim_{n \rightarrow \infty} \int \phi_n dP^n_{1/\sqrt{n},g} \leq P(\chi_d^2(\kappa) > q_{d,1-\alpha}), \tag{B.37}$$

where $q_{d,1-\alpha}$ denotes the $(1 - \alpha)$ quantile of a chi-squared random variable with d degrees of freedom. However, note that since $\{s_\tau\}_{\tau=1}^d$ is orthonormal by hypothesis, Assumption 3.1(i) implies $\|\hat{\mathbb{G}}_n\|^2 \xrightarrow{L} \chi_d^2(0)$ under P^n and therefore $c_{1-\alpha} = q_{d,1-\alpha}$. Furthermore, Lemma E.2 (in the Supplemental Material) implies that, for $\mathbb{G}_0 \sim N(0, I_d)$ with I_d the $d \times d$ identity matrix and $\Delta_g \in \mathbf{R}^d$ given by $\Delta_g = (\int g s_1 dP, \dots, \int g s_d dP)'$, we must have,

for any path $t \mapsto P_{t,g} \in \mathcal{M}$,

$$\hat{\mathbb{G}}_n \xrightarrow{L_{n,g}} \mathbb{G}_0 + \Delta_g. \tag{B.38}$$

In particular, since $\|\Delta_g\| = \|\Pi_{T^\perp}(g)\|_{P,2}$ for any $g \in \bar{M}(P)$, we obtain from (B.38) that

$$\begin{aligned} & \inf_{g \in \mathcal{G}(\kappa)} \lim_{n \rightarrow \infty} P_{1/\sqrt{n},g}(\|\hat{\mathbb{G}}_n\|^2 > c_{1-\alpha}) \\ &= \inf_{g \in \mathcal{G}(\kappa)} P(\|\mathbb{G}_0 + \Delta_g\|^2 > q_{d,1-\alpha}) = P(\chi_d^2(\kappa) > q_{d,1-\alpha}). \end{aligned} \tag{B.39}$$

Therefore, part (ii) of the lemma follows from (B.37) and (B.39).

For part (iii) of the lemma, it suffices to verify that $b^*(\nu - \tilde{\nu}) \in \bar{T}(P)^\perp \cap \bar{M}(P)$ for all $b^* \in \mathbf{B}^*$. To this end, we note that $b^*(\hat{\theta}_n)$ and $b^*(\tilde{\theta}_n)$ are both asymptotically linear regular (with respect to \mathbf{P}) estimators of $b^*(\theta(P))$ with influence functions $b^*(\nu)$ and $b^*(\tilde{\nu})$, respectively. We also have that

$$b^*(\tilde{\nu}) - b^*(\nu) = \Pi_{T^\perp}(b^*(\tilde{\nu})), \tag{B.40}$$

since by Proposition 3.3.1 in Bickel et al. (1993), $b^*(\nu) \in \bar{T}(P)$ due to $b^*(\hat{\theta}_n)$ being efficient (with respect to \mathbf{P}), and $\Pi_T(b^*(\tilde{\nu})) = b^*(\nu)$ due to $b^*(\tilde{\theta}_n)$ being regular (with respect to \mathbf{M} and \mathbf{P}). However, $b^*(\tilde{\nu})$ being efficient with respect to \mathbf{M} and Proposition 3.3.1 in Bickel et al. (1993) imply $b^*(\tilde{\nu}) \in \bar{M}(P)$. Since $\bar{M}(P)$ is a vector subspace and $b^*(\nu) \in \bar{T}(P) \subseteq \bar{M}(P)$, result (B.40) additionally implies $\Pi_{T^\perp}(b^*(\tilde{\nu})) \in \bar{M}(P)$, and thus $b^*(\nu - \tilde{\nu}) \in \bar{T}(P)^\perp \cap \bar{M}(P)$ as claimed. Q.E.D.

APPENDIX C: PROOFS FOR $T(P)$ A CONVEX CONE

PROOF OF LEMMA 5.1: Since $\bar{T}(P)$ is a convex cone by Assumption 5.1, Proposition 46.5.4 in Zeidler (1984) implies $L_0^2(P) = \bar{T}(P) \oplus \bar{T}(P)^\perp$. The lemma then follows since $\bar{T}(P)^\perp = \{0\}$ if and only if $\bar{T}(P) = L_0^2(P)$. Q.E.D.

LEMMA C.1: *Let Assumption 5.1 hold and let P be locally just identified by \mathbf{P} . Then: for all bounded function $f : \mathbf{X} \rightarrow \mathbf{R}$, the sample mean, $n^{-1} \sum_{i=1}^n f(X_i)$, is an asymptotically locally admissible estimator of $\int f dP$ under any Ψ -loss.*

PROOF: We aim to show that if P is locally just identified by \mathbf{P} , then $n^{-1} \sum_{i=1}^n f(X_i)$ is an asymptotically locally admissible estimator of $\int f dP$. To this end, we note that for any bounded $f : \mathbf{X} \rightarrow \mathbf{R}$, Theorem 3.10.12 in van der Vaart and Wellner (1996) implies that, for any path $t \mapsto P_{t,g} \in \mathcal{M}$,

$$\frac{1}{\sqrt{n}} \sum_{i=1}^n \left\{ f(X_i) - \int f dP_{1/\sqrt{n},g} \right\} \xrightarrow{L_{n,g}} \mathbb{G}_0, \tag{C.1}$$

where $\mathbb{G}_0 \sim N(0, \text{Var}\{f(X_i)\})$. Therefore, since Ψ is bounded and continuous, we obtain from (C.1) that, for any path $t \mapsto P_{t,g} \in \mathcal{P}$,

$$\limsup_{n \rightarrow \infty} E_{P_{1/\sqrt{n},g}} \left[\Psi \left(\frac{1}{\sqrt{n}} \sum_{i=1}^n \left\{ f(X_i) - \int f dP_{1/\sqrt{n},g} \right\} \right) \right] = E[\Psi(\mathbb{G}_0)]. \tag{C.2}$$

By way of contradiction, next suppose that $\frac{1}{n} \sum_{i=1}^n f(X_i)$ is not an asymptotically locally admissible estimator of $\int f dP$ under Ψ -loss. It then follows that there must exist another estimator $\hat{\theta}_n : \{X_i\}_{i=1}^n \rightarrow \mathbf{R}$ of $\int f dP$ satisfying, for any path $t \mapsto P_{t,g} \in \mathbf{P}$ and some tight law \mathbb{Z}_g ,

$$\sqrt{n} \left\{ \hat{\theta}_n - \int f dP_{1/\sqrt{n},g} \right\} \xrightarrow{L_{n,g}} \mathbb{Z}_g, \tag{C.3}$$

and moreover, by result (C.2), for any $t \mapsto P_{t,g} \in \mathbf{P}$, $\hat{\theta}_n$ must additionally be such that

$$\limsup_{n \rightarrow \infty} E_{P_{1/\sqrt{n},g}} \left[\Psi \left(\sqrt{n} \left\{ \hat{\theta}_n - \int f dP_{1/\sqrt{n},g} \right\} \right) \right] \leq E[\Psi(\mathbb{G}_0)] \tag{C.4}$$

with strict inequality holding for some $t \mapsto P_{t,g} \in \mathbf{P}$. In particular, since Ψ is bounded and continuous, results (C.3) and (C.4) imply that

$$E[\Psi(\mathbb{G}_0)] \geq \sup_{g \in T(P)} E[\Psi(\mathbb{Z}_g)] = \sup_{g \in \bar{T}(P)} E[\Psi(\mathbb{Z}_g)], \tag{C.5}$$

where the equality follows from Lemma E.6 (in the Supplemental Material), which establishes both that \mathbb{Z}_g is well defined for $g \in \bar{T}(P)$ and that the supremums over $T(P)$ and $\bar{T}(P)$ must be equal. Since P is just identified by \mathbf{P} , however, we have $\bar{T}(P) = L_0^2(P)$, which implies $f - \int f dP \in \bar{T}(P)$. Therefore, result (C.5), Theorem 2.6 in van der Vaart (1989), and Proposition 8.6 in van der Vaart (1998) together establish that under $\otimes_{i=1}^n P$, we must have

$$\sqrt{n} \left\{ \hat{\theta}_n - \int f dP \right\} = \frac{1}{\sqrt{n}} \sum_{i=1}^n \left(f(X_i) - \int f dP \right) + o_p(1). \tag{C.6}$$

Equivalently, $\sqrt{n} \{ \hat{\theta}_n - \frac{1}{n} \sum_{i=1}^n f(X_i) \} = o_p(1)$ under $\otimes_{i=1}^n P$ and, by contiguity, also under $\otimes_{i=1}^n P_{1/\sqrt{n},g}$ for any path $t \mapsto P_{t,g}$. However, by results (C.1) and (C.3), we can then conclude that \mathbb{Z}_g must equal \mathbb{G}_0 in distribution, thus establishing the desired contradiction since, as a result, (C.4) cannot hold strictly for any path $t \mapsto P_{t,g} \in \mathbf{P}$. *Q.E.D.*

PROOF OF THEOREM 5.1: We note that if P is locally just identified, then by Lemma C.1, it follows that (ii) implies (i). Similarly, we also note that by Theorem 3.1(ii), it follows that (iii) implies (i). Thus, to conclude the proof, we need only show that (i) implies (ii) and (iii). To this end, note that if P is locally overidentified by \mathbf{P} , then Lemma 5.1 implies there exists a $0 \neq \tilde{f} \in \bar{T}(P)^-$, which without loss of generality we assume satisfies $\|\tilde{f}\|_{P,2} = 1$. For $\mathbb{G}_0 \sim N(0, 1)$, then note that Theorem 3.10.12 in van der Vaart and Wellner (1996) implies that

$$\mathbb{G}_n(\tilde{f}) \equiv \frac{1}{\sqrt{n}} \sum_{i=1}^n \tilde{f}(X_i) \xrightarrow{L_{n,g}} \mathbb{G}_0 + \int \tilde{f} g dP \quad \text{for any path } t \mapsto P_{t,g} \in \mathcal{M}. \tag{C.7}$$

In order to establish that (i) implies (ii), we consider two cases. First, note that if \tilde{f} is unbounded, then we may set f to equal $f(x) = \tilde{f}(x)1\{|\tilde{f}(x)| \leq M\}$, which satisfies $\int \tilde{f} f dP > 0$ for M large enough. Moreover, for any finite M , unboundedness of \tilde{f} implies f and \tilde{f} are linearly independent, and therefore (ii) follows by applying Theorem 5.4

with $s_{r^*} = \tilde{f}$. On the other hand, if \tilde{f} is bounded, then we may set $f = \tilde{f}$ and define the estimator

$$\hat{\theta}_n \equiv \min \left\{ \frac{1}{n} \sum_{i=1}^n f(X_i), 0 \right\}. \tag{C.8}$$

Next note that since $\int f dP = 0$ due to $f = \tilde{f} \in L_0^2(P)$, f being bounded and Lemma F.1 (in the Supplemental Material) imply, for any path $t \mapsto P_{t,g} \in \mathcal{M}$, that

$$\lim_{n \rightarrow \infty} \int \sqrt{n} f dP_{1/\sqrt{n},g} = \int f g dP. \tag{C.9}$$

Therefore, results (C.7), (C.8), (C.9), and the continuous mapping theorem establish that

$$\sqrt{n} \left\{ \hat{\theta}_n - \int f dP_{1/\sqrt{n},g} \right\} \xrightarrow{L_{n,g}} \min \left\{ \mathbb{G}_0, - \int f g dP \right\} \equiv \mathbb{Z}_g \tag{C.10}$$

for any path $t \mapsto P_{t,g} \in \mathcal{M}$, which verifies that $\hat{\theta}_n$ indeed satisfies (56). Similarly, results (C.7) and (C.9) together imply, for any path $t \mapsto P_{t,g} \in \mathcal{M}$, that

$$\sqrt{n} \left\{ \frac{1}{n} \sum_{i=1}^n f(X_i) - \int f dP_{1/\sqrt{n},g} \right\} \xrightarrow{L_{n,g}} \mathbb{G}_0. \tag{C.11}$$

Since $0 \neq f \in T(P)^-$, however, it follows that $\int f g dP \leq 0$ whenever $t \mapsto P_{t,g} \in \mathbf{P}$ and therefore

$$P(|\mathbb{Z}_g| \leq t) \geq P(|\mathbb{G}_0| \leq t), \tag{C.12}$$

with strict inequality holding whenever $t > - \int f g dP$. By definition of Ψ -loss, however, $\Psi(b) = \Psi(|b|)$, and $\Psi(b) \geq \Psi(b')$ whenever $|b| \geq |b'|$, and therefore result (C.12) implies that

$$E[\Psi(\mathbb{Z}_g)] \leq E[\Psi(\mathbb{G}_0)]. \tag{C.13}$$

The path $t \mapsto P_{t,g}$ satisfying $P_{t,g} = P$ for all t , however, trivially satisfies $t \mapsto P_{t,g} \in \mathbf{P}$ and has score $g = 0$. Therefore, the fact that Ψ is not constant together with results (C.10) and (C.11) imply $E[\Psi(\mathbb{Z}_0)] < E[\Psi(\mathbb{G}_0)]$. Since Ψ being bounded and continuous implies

$$\lim_{n \rightarrow \infty} E_{P_{1/\sqrt{n},g}} \left[\Psi \left(\sqrt{n} \left\{ \hat{\theta}_n - \int f dP_{1/\sqrt{n},g} \right\} \right) \right] = E[\Psi(\mathbb{Z}_g)], \tag{C.14}$$

we conclude from (C.2) that $\sum_{i=1}^n f(X_i)/n$ is indeed not asymptotically locally admissible since it is dominated by $\hat{\theta}_n$.

Finally, to establish that (i) implies (iii), we define the test $\phi_n \equiv 1\{\mathbb{G}_n(\tilde{f}) > z_{1-\alpha}\}$ for $z_{1-\alpha}$ the $1 - \alpha$ quantile of \mathbb{G}_0 . Then equation (C.7) implies that

$$\pi(g) \equiv \lim_{n \rightarrow \infty} \int \phi_n dP_{1/\sqrt{n},g} = P \left(\mathbb{G}_0 + \int \tilde{f} g dP > z_{1-\alpha} \right) \tag{C.15}$$

for any path $t \mapsto P_{t,g} \in \mathcal{M}$. Thus, whenever the path $t \mapsto P_{t,g} \in \mathbf{P}$, it follows from $g \in \bar{T}(P)$ and $\tilde{f} \in \bar{T}(P)^-$ that $\int g \tilde{f} dP \leq 0$ and hence by (C.15) that $\pi(g) \leq \alpha$, that is, ϕ_n has

asymptotic local level α . On the other hand, there exists a path $t \mapsto P_{t,\tilde{f}} \in \mathcal{M}$ with score $\tilde{f} \in \bar{T}(P)^-$. This observation and (C.15) together imply $\pi(\tilde{f}) > \alpha$; hence we conclude (i) implies (iii). Q.E.D.

PROOF OF THEOREM 5.2: Fix a path $t \mapsto P_{t,g}$ such that its score $g \in L_0^2(P)$ satisfies $\lambda \Pi_{T^-}(g) \in \bar{T}(P)$ for any $\lambda \leq 0$, and note that by Proposition 46.5.4 in Zeidler (1984),

$$g = \Pi_T(g) + \Pi_{T^-}(g). \tag{C.16}$$

Moreover, we note that if $\Pi_{T^-}(g) = 0$, then $g \in \bar{T}(P)$ and thus $\pi(g) \leq \alpha$ since $\pi(g) \leq \alpha$ for all $g \in \bar{T}(P)$ by hypothesis. We therefore assume without loss of generality that $\Pi_{T^-}(g) \neq 0$ and observe that by the hypotheses of the theorem, there exists an $f^* \in \{f_1, f_2\}$ such that $f^* \in \bar{T}(P)^-$, f^* is linearly independent of $\Pi_{T^-}(g)$, and f^* satisfies $\lambda f^* \in \bar{T}(P)$ for all $\lambda \leq 0$. Defining

$$H \equiv \{h \in L_0^2(P) : h = \Pi_T(g) + \gamma_1 \Pi_{T^-}(g) + \gamma_2 f^* \text{ for some } (\gamma_1, \gamma_2) \in \mathbf{R}^2\}, \tag{C.17}$$

we may then construct for any $h \in H$ a path $t \mapsto \bar{P}_{t,h}$ whose score is h and such that $\bar{P}_{t,h} \ll P \ll \bar{P}_{t,h}$; see, for example, Example 3.2.1 in Bickel et al. (1993). Recall that \mathcal{B} is the σ -algebra on \mathbf{X} , and consider the sequence of experiments \mathcal{E}_n given by

$$\mathcal{E}_n \equiv \left(\mathbf{X}^n, \mathcal{B}^n, \bigotimes_{i=1}^n \bar{P}_{1/\sqrt{n},h} : h \in H \right). \tag{C.18}$$

Setting $h_0 \equiv \Pi_T(g)$, then observe that Lemma 25.14 in van der Vaart (1998) implies

$$\sum_{i=1}^n \log \left(\frac{d\bar{P}_{1/\sqrt{n},h}}{d\bar{P}_{1/\sqrt{n},h_0}}(X_i) \right) = \frac{1}{\sqrt{n}} \sum_{i=1}^n (h(X_i) - h_0(X_i)) - \frac{1}{2} \int (h^2 - h_0^2) dP + o_p(1) \tag{C.19}$$

under $P^n \equiv \bigotimes_{i=1}^n P$, and where we exploited that $\bar{P}_{t,h} \ll P \ll \bar{P}_{t,h_0}$. Since similarly

$$\sum_{i=1}^n \log \left(\frac{d\bar{P}_{1/\sqrt{n},h_0}}{dP}(X_i) \right) = \frac{1}{\sqrt{n}} \sum_{i=1}^n h_0(X_i) - \frac{1}{2} \int h_0^2 dP + o_p(1) \tag{C.20}$$

under P^n by Lemma 25.14 in van der Vaart (1998), it follows by LeCam’s Third Lemma (see, e.g., Lemma A.8.6 in Bickel et al. (1993)) that, for an arbitrary finite subset $\{h_j\}_{j=1}^J \equiv I \subseteq H$ and L_{n,h_0} denoting the law under $\bigotimes_{i=1}^n \bar{P}_{1/\sqrt{n},h_0}$, we have

$$\left(\sum_{i=1}^n \log \left(\frac{d\bar{P}_{1/\sqrt{n},h_1}}{d\bar{P}_{1/\sqrt{n},h_0}}(X_i) \right), \dots, \sum_{i=1}^n \log \left(\frac{d\bar{P}_{1/\sqrt{n},h_J}}{d\bar{P}_{1/\sqrt{n},h_0}}(X_i) \right) \right)' \xrightarrow{L_{n,h_0}} N(-\mu_I, \Sigma_I), \tag{C.21}$$

where $\Sigma_I \equiv \int (h_1 - h_0, \dots, h_J - h_0)(h_1 - h_0, \dots, h_J - h_0)' dP$ and the mean is given by $\mu_I \equiv \frac{1}{2} (\int (h_1 - h_0)^2 dP, \dots, \int (h_J - h_0)^2 dP)'$. Next, define $v_h \in \mathbf{R}^2$ and $\Omega \in \mathbf{R}^{2 \times 2}$ by

$$v_h \equiv \begin{pmatrix} \int \{\Pi_{T^-}(g)\} h dP \\ \int f^* h dP \end{pmatrix}, \quad \Omega \equiv \begin{pmatrix} \int \{\Pi_{T^-}(g)\}^2 dP & \int \{\Pi_{T^-}(g)\} f^* dP \\ \int \{\Pi_{T^-}(g)\} f^* dP & \int \{f^*\}^2 dP \end{pmatrix}, \tag{C.22}$$

and note that the linear independence of f^* and $\Pi_{T^-}(g)$ in $L_0^2(P)$ imply Ω is invertible. For any $h \in H$, then let Q_h be the bivariate normal law on \mathbf{R}^2 satisfying

$$Q_h \stackrel{L}{=} N(\Omega^{-1}\{v_h - 2v_{h_0}\}, \Omega^{-1}). \tag{C.23}$$

Further observe that for any $h \in H$ and $U \in \mathbf{R}^2$, we can obtain by direct calculation

$$\log\left(\frac{dQ_h}{dQ_{h_0}}(U)\right) = U'(v_h - v_{h_0}) + \frac{1}{2}\{v'_{h_0}\Omega^{-1}v_{h_0} - (v_h - 2v_{h_0})'\Omega^{-1}(v_h - 2v_{h_0})\}, \tag{C.24}$$

and therefore, exploiting (C.24) and $(v_{h_i} - v_{h_0})'\Omega^{-1}(v_{h_k} - v_{h_0}) = \int (h_i - h_0)(h_k - h_0) dP$ for any $h_i, h_k \in H$ implies that, for any finite subset $\{h_j\}_{j=1}^I \equiv I \subseteq H$, we have

$$\left(\log\left(\frac{dQ_{h_1}}{dQ_{h_0}}\right), \dots, \log\left(\frac{dQ_{h_I}}{dQ_{h_0}}\right)\right) \sim N(-\mu_I, \Sigma_I) \tag{C.25}$$

under Q_{h_0} . Since (C.21) and Corollary 12.3.1 in Lehmann and Romano (2005) imply $\{P_{1/\sqrt{n}, h}\}$ and $\{P_{1/\sqrt{n}, h_0}\}$ are mutually contiguous for any $h \in H$, results (C.21) and (C.25) together with Lemma 10.2.1 in LeCam (1986) establish \mathcal{E}_n converges weakly to

$$\mathcal{E} \equiv (\mathbf{R}^2, \mathcal{A}^2, Q_h : h \in H), \tag{C.26}$$

where \mathcal{A} denotes the Borel σ -algebra on \mathbf{R} .

By the asymptotic representation theorem (see, e.g., Theorem 7.1 in van der Vaart (1991a)), it then follows from ϕ_n having a local asymptotic power function π that there exists a test ϕ based on a single observation of $U \sim Q_h$ such that, for all $h \in H$,

$$\pi(h) \equiv \lim_{n \rightarrow \infty} \int \phi_n d\bar{P}_{1/\sqrt{n}, h}^n = \int \phi dQ_h. \tag{C.27}$$

Further note that any $h \in H$ can be written as $h = \Pi_T(g) + \gamma_1(h)\Pi_{T^-}(g) + \gamma_2(h)f^*$ for some $\gamma(h) = (\gamma_1(h), \gamma_2(h))' \in \mathbf{R}^2$ and that, moreover, $\gamma(h) = \Omega^{-1}\{v_h - v_{h_0}\}$. In addition, we observe that $\lambda f^*, \lambda \Pi_{T^-}(g) \in \bar{T}(P)^-$ whenever $\lambda \geq 0$ and $\lambda f^*, \lambda \Pi_{T^-}(g) \in \bar{T}(P)$ whenever $\lambda \leq 0$ together with the linear independence of $\Pi_{T^-}(g)$ imply that $h \in \bar{T}(P)$ if and only if $\gamma_1(h) \leq 0$ and $\gamma_2(h) \leq 0$. Thus, the hypothesis on π and result (C.27) yield

$$\begin{aligned} \int \phi dQ_h &\leq \alpha \quad \text{if } \gamma_1(h) \leq 0 \text{ and } \gamma_2(h) \leq 0, \\ \int \phi dQ_h &\geq \alpha \quad \text{if } \gamma_1(h) > 0 \text{ or } \gamma_2(h) > 0. \end{aligned} \tag{C.28}$$

However, (C.28), $h \mapsto \gamma(h)$ being bijective between H and \mathbf{R}^2 , and Lehmann (1952, p. 542) imply that $\int \phi dQ_h = \alpha$ for all $h \in H$. In particular, since $g \in H$, the claim of the theorem finally follows from $\int \phi dQ_g = \alpha$, result (C.27), and Lemma D.1 (in the Supplemental Material). Q.E.D.

PROOF OF THEOREM 5.3: For any path $t \mapsto P_{t,g} \in \mathcal{M}$, we first note that applying Lemma E.2 (in the Supplemental Material) with Assumption 5.2 in place of Assumption 3.1 implies

$$\hat{\mathbb{G}}_n \xrightarrow{L_{n,g}} \mathbb{G}_0 + \Delta_g, \tag{C.29}$$

for $\Delta_g \in \ell^\infty(\mathbf{T})$ given by $\Delta_g(\tau) \equiv \int s_\tau g dP$. Let $x \vee y \equiv \max\{x, y\}$. Defining $\Delta_g^\omega \in \ell^\infty(\mathbf{T})$ to be $\Delta_g^\omega(\tau) \equiv \omega(\tau) \times \Delta_g(\tau)$, we then obtain from (C.29) and the continuous mapping theorem

$$\|\hat{\mathbb{G}}_n^\omega \vee 0\|_\infty \xrightarrow{L_{n,g}} \|(\mathbb{G}_0^\omega + \Delta_g^\omega) \vee 0\|_\infty. \tag{C.30}$$

Moreover, note that \mathbb{G}_0^ω is a regular measure by Theorem 7.1.7 in Bogachev (2007), and hence since \mathbb{G}_0^ω is also tight due to $\omega \in \ell^\infty(\mathbf{T})$ and \mathbb{G}_0 being tight by Assumption 5.2(ii), we conclude \mathbb{G}_0^ω is Radon. Together with the convexity of the map $\|\cdot \vee 0\|_\infty : \ell^\infty(\mathbf{T}) \rightarrow \mathbf{R}$, \mathbb{G}_0^ω being Radon allows us to apply Theorem 11.1 in Davydov, Lifshits, and Smorodina (1998) to obtain that

$$c_0 \equiv \inf\{c : P(\|(\mathbb{G}_0^\omega + \Delta_g^\omega) \vee 0\|_\infty \leq c) > 0\} \tag{C.31}$$

is the only possible discontinuity point of the c.d.f. of $\|(\mathbb{G}_0^\omega + \Delta_g^\omega) \vee 0\|_\infty$. However, note that since $c_{1-\alpha}^\omega > 0$ by hypothesis, we must have $\|\omega\|_\infty > 0$, and therefore, for any $c > 0$,

$$P(\|(\mathbb{G}_0^\omega + \Delta_g^\omega) \vee 0\|_\infty \leq c) \geq P\left(\|\mathbb{G}_0 + \Delta_g\|_\infty \leq \frac{c}{\|\omega\|_\infty}\right) > 0, \tag{C.32}$$

where we exploited that $\omega \geq 0$, and the final inequality follows from Proposition 12.1 in Davydov, Lifshits, and Smorodina (1998) and $-\Delta_g = \Delta_{-g}$ belonging to the support of \mathbb{G}_0 by Lemma E.3 (in the Supplemental Material).⁹ Since $c_{1-\alpha}^\omega > 0$ by hypothesis, it follows from (C.31) and (C.32) that $c_{1-\alpha}^\omega$ is a continuity point of the c.d.f. of $\|(\mathbb{G}_0^\omega + \Delta_g^\omega) \vee 0\|_\infty$. Therefore, we obtain from (C.29) that

$$\lim_{n \rightarrow \infty} P_{1/\sqrt{n},g}(\|\hat{\mathbb{G}}_n^\omega \vee 0\|_\infty > c_{1-\alpha}^\omega) = P(\|(\mathbb{G}_0^\omega + \Delta_g^\omega) \vee 0\|_\infty > c_{1-\alpha}^\omega), \tag{C.33}$$

which verifies that ϕ_n^ω indeed has an asymptotic local power function. Moreover, note that if $t \mapsto P_{t,g} \in \mathbf{P}$, then $g \in \bar{T}(P)$ by definition and hence $\int s_\tau g dP \leq 0$ for all $\tau \in \mathbf{T}$ since $s_\tau \in \bar{T}(P)^-$. Thus, $\omega \geq 0$ implies $\Delta_g^\omega \leq 0$, and therefore (C.33) yields

$$\lim_{n \rightarrow \infty} P_{1/\sqrt{n},g}(\|\hat{\mathbb{G}}_n^\omega \vee 0\|_\infty > c_{1-\alpha}^\omega) \leq P(\|\mathbb{G}_0^\omega \vee 0\|_\infty > c_{1-\alpha}^\omega) = \alpha, \tag{C.34}$$

where we exploited that $c_{1-\alpha}^\omega$ is the $1 - \alpha$ quantile of $\|\mathbb{G}_0^\omega \vee 0\|_\infty$ and the c.d.f. of $\|\mathbb{G}_0^\omega \vee 0\|_\infty$ is continuous at $c_{1-\alpha}^\omega$. Since (C.34) holds for any path $t \mapsto P_{t,g} \in \mathbf{P}$, we conclude ϕ_n^ω is indeed an asymptotic level α specification test.

To establish (59), note that $\bar{C}(P) \subseteq L_0^2(P)$ is a closed convex cone by definition, and let $\bar{C}(P)^- \equiv \{g \in L_0^2(P) : \int gf dP \leq 0 \text{ for all } f \in \bar{C}(P)\}$. For any $g \in L_0^2(P)$, Proposition 46.5.4 in Zeidler (1984) implies $g = \Pi_C(g) + \Pi_{C^-}(g)$ and $\int \{\Pi_C(g)\}\{\Pi_{C^-}(g)\} dP = 0$. In particular, if a path $t \mapsto P_{t,g} \in \mathcal{M}$ is such that $\Pi_C(g) \neq 0$, then

$$\int g\{\Pi_C(g)\} dP = \int \{\Pi_C(g) + \Pi_{C^-}(g)\}\{\Pi_C(g)\} dP = \int \{\Pi_C(g)\}^2 dP > 0. \tag{C.35}$$

⁹Lemma E.3 requires Assumption 3.1 in place of Assumption 5.2, but the proof of Lemma E.3 also holds under the latter assumption with no modifications.

Since $\Pi_C(g) \in \bar{C}(P)$ and $\bar{C}(P)$ is the closed convex cone generated by $\{s_\tau\}_{\tau \in \mathbf{T}}$, there exist an integer $K < \infty$, positive scalars $\{\alpha_k\}_{k=1}^K$, and $\{\tau_k\}_{k=1}^K \subseteq \mathbf{T}$ such that

$$\left\| \Pi_C(g) - \sum_{k=1}^K \alpha_k s_{\tau_k} \right\|_{P,2} < \frac{1}{2} \frac{\|\Pi_C(g)\|_{P,2}^2}{\|g\|_{P,2}}. \tag{C.36}$$

Therefore, results (C.35) and (C.36) together with the Cauchy–Schwarz inequality yield

$$\begin{aligned} \int g \left\{ \sum_{k=1}^K \alpha_k s_{\tau_k} \right\} dP &\geq \int g \{ \Pi_C(g) \} dP - \left| \int g \left\{ \Pi_C(g) - \sum_{k=1}^K \alpha_k s_{\tau_k} \right\} dP \right| \\ &\geq \frac{1}{2} \|\Pi_C(g)\|_{P,2}^2 > 0. \end{aligned} \tag{C.37}$$

Since $\alpha_k \geq 0$ for all $1 \leq k \leq K$, result (C.37) implies that $\int g s_{\tau^*} dP > 0$ for some $\tau^* \in \mathbf{T}$. To conclude, we then let $\omega^*(\tau) \equiv 1\{\tau = \tau^*\}$ and note $\mathbb{G}_0^{\omega^*}(\tau^*) \sim N(0, \int s_{\tau^*}^2 dP)$. Furthermore, since $\int s_{\tau^*}^2 dP > 0$ because $\int g s_{\tau^*} dP > 0$, and $\|\mathbb{G}_0^{\omega^*} \vee 0\|_\infty = \max\{\mathbb{G}_0(\tau^*), 0\}$ almost surely, it follows that $c_{1-\alpha}^{\omega^*} > 0$ provided $\alpha \in (0, \frac{1}{2})$. We may then exploit result (C.33) since $c_{1-\alpha}^{\omega^*} > 0$, and employ $\int g s_{\tau^*} dP > 0$ to obtain

$$\lim_{n \rightarrow \infty} P_{1/\sqrt{n}, g}(\|\hat{\mathbb{G}}_n^{\omega^*} \vee 0\|_\infty > c_{1-\alpha}^{\omega^*}) = P\left(\mathbb{G}_0(\tau^*) + \int g s_{\tau^*} dP > c_{1-\alpha}^{\omega^*}\right) > \alpha, \tag{C.38}$$

which establishes the second claim of the theorem.

Q.E.D.

PROOF OF THEOREM 5.4: Let $R(\tau^*) \equiv \{\lambda s_{\tau^*} : \lambda \geq 0\}$ which is a closed convex cone and set $R(\tau^*)^-$ to be the polar cone of $R(\tau^*)$, which satisfies

$$R(\tau^*)^- = \left\{ g \in L_0^2(P) : \int g s_{\tau^*} dP \leq 0 \right\}. \tag{C.39}$$

In addition, for any $g \in L_0^2(P)$, we let $\Pi_R(g)$ and $\Pi_{R^-}(g)$ denote the metric projections of g onto $R(\tau^*)$ and $R(\tau^*)^-$, respectively, and we note by direct calculation that

$$\Pi_R\left(f - \int f dP\right) = \beta(f, \tau^*) \times s_{\tau^*} \tag{C.40}$$

for any $f \in L^2(P)$. Moreover, by Proposition 46.5.4 in [Zeidler \(1984\)](#), it also follows that

$$\begin{aligned} &\frac{1}{\sqrt{n}} \sum_{i=1}^n \left\{ f(X_i) - \int f dP \right\} \\ &= \frac{1}{\sqrt{n}} \sum_{i=1}^n \left\{ \left\{ \Pi_{R^-}\left(f - \int f dP\right) \right\}(X_i) + \beta(f, \tau^*) s_{\tau^*}(X_i) \right\}, \end{aligned} \tag{C.41}$$

where $\int \Pi_{R^-}(f - \int f dP) \beta(f, \tau^*) s_{\tau^*} dP = 0$. Let $\Delta_g \equiv \int \Pi_R(f - \int f dP) g dP$. Then we obtain from results (C.40) and (C.41), Assumption 5.2(i), and Theorem 3.10.12 in [van der](#)

Vaart and Wellner (1996) that, for any path $t \mapsto P_{t,g} \in \mathcal{M}$, we have

$$\left(\frac{1}{\sqrt{n}} \sum_{i=1}^n \left\{ f(X_i) - \int f dP \right\}, \beta(f, \tau^*) \hat{\mathbb{G}}_n(\tau^*) \right) \xrightarrow{L_{n,g}} \left(\mathbb{G}_R + \mathbb{G}_{R^-} + \int fg dP, \mathbb{G}_R + \Delta_g \right), \tag{C.42}$$

where $(\mathbb{G}_R, \mathbb{G}_{R^-})$ are independent normals with $\text{Var}\{\mathbb{G}_R + \mathbb{G}_{R^-}\} = \text{Var}\{f(X_i)\}$ and $\text{Var}\{\mathbb{G}_R\} = \|\beta(f, \tau^*)s_{\tau^*}\|_{P,2}^2$. Moreover, for any bounded $f : \mathbf{X} \rightarrow \mathbf{R}$, Lemma F.1 (in the Supplemental Material) implies

$$\left| \sqrt{n} \int f(dP_{1/\sqrt{n},g} - dP) - \int fg dP \right| = o(1) \tag{C.43}$$

for any path $t \mapsto P_{t,g} \in \mathcal{M}$. Therefore, results (C.42) and (C.43), the definition of $\hat{\mu}_n(f, \tau^*)$, and the continuous mapping theorem allow us to conclude that

$$\sqrt{n} \left\{ \hat{\mu}_n(f, \tau^*) - \int f dP_{1/\sqrt{n},g} \right\} \xrightarrow{L_{n,g}} \mathbb{G}_{R^-} + \min\{\mathbb{G}_R, -\Delta_g\} \equiv \mathbb{Z}_g, \tag{C.44}$$

which implies $\hat{\mu}_n(f, \tau^*)$ indeed satisfies (56). We next aim to show that (61) holds for any path $t \mapsto P_{t,g} \in \mathbf{P}$ provided $\int fs_{\tau^*} dP > 0$, which implies $\beta(f, \tau^*) > 0$. We thus assume $\beta(f, \tau^*) \neq 0$, and note this implies $\text{Var}\{\mathbb{G}_R\} > 0$. Hence, since $\mathbb{G}_0 = \mathbb{G}_R + \mathbb{G}_{R^-}$ by results (C.1) and (C.42), we can exploit the definition of \mathbb{Z}_g in (C.44) to obtain, for any $t > 0$, that

$$\begin{aligned} P(|\mathbb{Z}_g| \leq t) &= P(|\mathbb{G}_0| \leq t, \mathbb{G}_R \leq -\Delta_g) + P(|\mathbb{G}_{R^-} - \Delta_g| \leq t, \mathbb{G}_R > -\Delta_g) \\ &= P(|\mathbb{G}_0| \leq t) + P(|\mathbb{G}_{R^-} - \Delta_g| \leq t, \mathbb{G}_R > -\Delta_g) \\ &\quad - P(|\mathbb{G}_0| \leq t, \mathbb{G}_R > -\Delta_g). \end{aligned} \tag{C.45}$$

Let $\sigma_{R^-}^2 \equiv \text{Var}\{\mathbb{G}_{R^-}\}$, and note that $\sigma_{R^-}^2 > 0$ since f and s_{τ^*} are linearly independent. For Φ the c.d.f. of a standard normal random variable, we can then conclude from $\mathbb{G}_0 = \mathbb{G}_R + \mathbb{G}_{R^-}$ and the independence of \mathbb{G}_R and \mathbb{G}_{R^-} that

$$\begin{aligned} P(|\mathbb{G}_{R^-} - \Delta_g| \leq t | \mathbb{G}_R + \Delta_g > 0) &= \Phi\left(\frac{t + \Delta_g}{\sigma_{R^-}}\right) - \Phi\left(\frac{-t + \Delta_g}{\sigma_{R^-}}\right), \\ P(|\mathbb{G}_0| \leq t | \mathbb{G}_R + \Delta_g > 0) &= E\left[\Phi\left(\frac{t - \mathbb{G}_R}{\sigma_{R^-}}\right) - \Phi\left(\frac{-t - \mathbb{G}_R}{\sigma_{R^-}}\right) | \mathbb{G}_R > -\Delta_g\right]. \end{aligned} \tag{C.46}$$

We note that the function $F_t(a) \equiv \Phi((t - a)/\sigma_{R^-}) - \Phi((-t - a)/\sigma_{R^-})$ is decreasing in $a \in [0, \infty)$ whenever $t \geq 0$. Since $s_{\tau^*} \in \tilde{T}(P)^-$, we have $\Delta_g \equiv \int \{\Pi_R(f - \int f dP)\}g dP \leq 0$ whenever $g \in \tilde{T}(P)$. It follows from (C.46) that, for any $g \in \tilde{T}(P)$, we have

$$P(|\mathbb{Z}_g| \leq t) > P(|\mathbb{G}_0| \leq t) \quad \text{for all } t > 0. \tag{C.47}$$

Thus, since $\Psi(b) = \Psi(|b|)$, $\Psi(b) \geq \Psi(b')$ whenever $|b| \geq |b'|$, and Ψ is nonconstant, result (C.47) implies

$$E[\Psi(\mathbb{Z}_g)] < E[\Psi(\mathbb{G}_0)]. \tag{C.48}$$

Since (C.48) holds for any $g \in \bar{T}(P)$, we then obtain from (C.44) and Definition 5.1(iii) that

$$\limsup_{n \rightarrow \infty} E_{P_{1/\sqrt{n},g}} \left[\Psi \left(\sqrt{n} \left\{ \hat{\mu}_n(f, \tau^*) - \int f dP_{1/\sqrt{n},g} \right\} \right) \right] = E[\Psi(\mathbb{Z}_g)] < E[\Psi(\mathbb{G}_0)] \quad (\text{C.49})$$

for any path $t \mapsto P_{t,g} \in \mathbf{P}$, which together with (C.2) establishes (61). Q.E.D.

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