

MIXED-RATES ASYMPTOTICS (EXTENDED VERSION)

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ABSTRACT. A general method is presented for deriving the limiting behavior of estimators that are defined as the values of parameters optimizing an empirical criterion function. The asymptotic behavior of such estimators is typically deduced from uniform limit theorems for rescaled and reparametrized criterion functions. The new method can handle cases where the standard approach does not yield the complete limiting behavior of the estimator. The asymptotic analysis depends on a decomposition of criterion functions into sums of components with different rescalings. The method is explained by examples from Lasso-type estimation, shorth estimation, and k-means clustering.

1. INTRODUCTION

Consider an estimator (a_n, b_n) that in some sense optimizes a random criterion function $G_n(a, b)$ over an open subset of $\mathbb{R}^{d_1} \times \mathbb{R}^{d_2}$. Two types of *mixed-rates* asymptotic behavior can occur and often occur simultaneously. First, the components a_n and b_n of the estimator may converge at different rates. Second, the criterion function itself may have important components settling down at different rates. The new method presented in this paper can handle both types of mixed-rates behavior.

Deriving the asymptotics of an estimator can be viewed as a three step procedure: proving consistency, establishing the rate of convergence and deriving the limiting distribution.

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This paper concentrates only on the last two steps. The limiting distribution is typically derived via a uniform limit theorem for the rescaled and reparametrized criterion functions. Suppose that the rates of convergence for the two components of the estimator have been established: $q_n^{-1} \|a_n - a_0\| \vee r_n^{-1} \|b_n - b_0\| = O_p(1)$ for some fixed parameter value (a_0, b_0) . Consider *localized criterion functions* of the form

$$H_n(s, t) := G_n(a_0 + q_n s, b_0 + r_n t) - G_n(a_0, b_0).$$

If, after appropriate rescaling, random functions $H_n(s, t)$ settle down to a “nice” stochastic process, the convergence in distribution of vectors $(s_n, t_n) := (q_n^{-1}[a_n - a_0], r_n^{-1}[b_n - b_0])$ to the corresponding optimizer of the limit process may follow from a continuous mapping type of argument. Theorem 3.2.2 of van der Vaart and Wellner (1996) makes this argument precise for estimators defined by maximization. The above approach is standard when the rates r_n and q_n are the same, and it can work in some mixed-rates cases, such as the change-point problem (see, for example, the section on non-regular examples in Kosorok 2006). Other mixed-rates examples where this argument succeeds can be found in Rotnitzky, Cox, Bottai, and Robins (2000), Pollard and Radchenko (2006) and Andrews (1999).

Many mixed-rates problems cannot be completely handled by the above approach. In the examples considered in this paper, the localized criterion function has the form

$$H_n(s, t) = \alpha_n f_n(s) + \beta_n g_n(s, t),$$

where $\beta_n = o(\alpha_n)$, the random function $f_n(s)$ settles down to a stochastic process $f(s)$, and g_n is stochastically bounded. Because the limit of $\alpha_n^{-1} H_n(s, t)$ is a stochastic process indexed only by s , the standard approach fails to establish the limiting distribution of the component t_n . However, if random function $g_n(s, t)$ settles down to a stochastic process $g(s, t)$, a two-step continuous mapping argument can be used to establish the distributional limit of the vector (s_n, t_n) . This general idea is made rigorous by Theorem 1 in Section 2. Corollary 1 handles the special case of the estimators defined by minimization.

Another challenging problem is deriving the *correct* rates of convergence for the two components of the estimator. Standard methods represent the centered criterion function $G_n(a, b) - G_n(a_0, b_0)$ as a sum of a positive deterministic function and a random one, whose rates of growth around the value (a_0, b_0) can be controlled (the deterministic function is typically approximated by a quadratic, and the random function is often approximately linear). Balancing out the two terms produces the rate of convergence: see, for example, Theorem 3.2.5 and Theorem 3.2.16 in van der Vaart and Wellner (1996). When a_n and b_n converge at different rates, this approach yields the “correct” rate only for the slower converging component. A reparametrization of the problem can sometimes be applied beforehand to sidestep this issue (for interesting examples see the references at the end of the paragraph on the standard method for deriving the limiting distribution). Unfortunately, such a trick is not available in general, and a more careful treatment of the criterion function is required. To derive the rate for the faster converging component, say b_n , Theorem 2 in Section 3 balances out the terms in a similar, but typically a more complicated representation for the function $b \mapsto [G_n(a_n, b) - G_n(a_n, b_0)]$.

Section 4 is devoted to mixed-rates problems that arise in M-estimation. Consider a collection of functions $g_\theta(x)$ and an empirical measure P_n , corresponding to independent observations coming from a distribution P . Define the estimator θ_n as the minimizer of the criterion function $G_n(\theta) = P_n g_\theta$, and suppose that function $G(\theta) = \int g_\theta dP$ is minimized by θ_0 . The stochastic bound $\|\theta_n - \theta_0\| = o_p(1)$ usually follows from a uniform law of large numbers, and the central limit theorem for the estimator is typically derived from a quadratic approximation of the form

$$G_n(\theta) - G_n(\theta_0) \approx (\theta - \theta_0)' G''(\theta_0) (\theta - \theta_0) + n^{-1/2} (\theta - \theta_0)' Z_n,$$

under the regularity assumption that matrix $G''(\theta_0)$ is positive definite matrix. If this regularity assumption breaks down and $G''(\theta_0)$ is singular, the approximation has to be carried out to higher order terms, which typically leads to mixed-rates situations that standard

methods cannot handle. Theorem 3 covers exactly such cases. The form of the approximation to function $G(\theta)$ near θ_0 determines the rates of convergence and the main features of the limiting behavior of the components of the estimator. Various remainder terms are handled by simple conditions imposed on functions g_θ .

Mixed-rates behavior naturally arises in the estimation of semiparametric models. Most of the results in this paper do not directly apply to such problems, but, as the example in Section 8 demonstrates, some of the methods and ideas can be carried over.

For the simplicity of the presentation, the estimators and the criterion functions considered in this paper have at most two components converging at different rates. All the results can be easily extended to cover cases of more than two mixed-rates components.

This paper is organized as follows. Sections 2, 3 and 4 contain the general mixed-rates asymptotics results, namely, the limiting distribution theorem, the rates of convergence theorem, and the M-estimation theorem. Proofs of these theorems are confined to Section 9. Sections 5, 6 and 7 contain applications of the general results to particular problems in Lasso-type estimation, shorth estimation and k -means clustering. Section 8 discusses a semiparametric example.

The abbreviation $Qf = \int f dQ$ is used throughout the paper for a given measurable function f and a signed measure Q . In particular, given independent observations X_i coming from a distribution P , let $P_n f$ denote $\sum_{i \leq n} f(X_i)/n$ and define the empirical process ν_n on a class of functions f by

$$f \mapsto \nu_n f = n^{1/2}(P_n - P)f = n^{-1/2} \sum_{i=1}^n [f(X_i) - Pf].$$

Write $\|\cdot\|_2$ for the $L_2(P)$ norm and say that a function f is square-integrable if $\|f\|_2 < \infty$. Interpret $f(\theta) \gtrsim g(\theta)$ to mean that there exists a positive constant c_0 such that $f(\theta) \geq c_0 g(\theta)$ for all θ in a sufficiently small neighborhood of the origin. Analogously, interpret $\alpha_n \gtrsim \beta_n$ to mean $\alpha_n \geq c_0 \beta_n$ for all sufficiently large n .

2. LIMITING DISTRIBUTION

Let the estimator (a_n, b_n) converge in probability to a fixed parameter value (a_0, b_0) . Suppose that the rates of convergence q_n and r_n have been established for the components a_n and b_n , respectively. Points $(s_n, t_n) := (q_n^{-1}[a_n - a_0], r_n^{-1}[b_n - b_0])$ optimize the localized criterion functions $H_n(s, t)$ and satisfy the tightness condition $\|(s_n, t_n)\| = O_p^*(1)$. The results that follow focus on deriving the limiting distribution of these points.

To avoid some measurability issues by allowing non-measurable converging maps, convergence in distribution (denoted by “ \rightsquigarrow ”) is understood in the sense of Hoffmann-Jørgensen. An exposition of this general concept can be found in the monographs of Dudley (1999) and van der Vaart and Wellner (1996). Let $\mathcal{B}_{loc}(\mathbb{R}^d)$ be the space of all locally bounded real functions on \mathbb{R}^d . Convergence of the random processes considered in the examples of this paper is handled by equipping $\mathcal{B}_{loc}(\mathbb{R}^d)$ with the metric ρ for the topology of uniform convergence on compacta:

$$\rho(g, h) = \sum_{k=1}^{\infty} 2^{-k} \min [1, \rho_k(g, h)], \quad \text{where} \quad \rho_k(g, h) = \sup_{\|t\| \leq k} |g(t) - h(t)|.$$

The main result of this section is stated for general optimization estimators that can be viewed as values of the corresponding “optimizing maps”. A continuity condition is included in the definition of an optimizing map (below) because it is required in Theorem 1. The points of continuity in condition (iv) of the definition are continuous functions because the limiting processes in all of the examples of this paper have continuous sample paths.

Write B_r for the closed ball $\{x \in \mathbb{R}^d : \|x\| \leq r\}$. Given a function g on \mathbb{R}^d , write g_r for the restriction of g to B_r .

Definition 1. For each positive r , let \mathcal{D}_r be some set of bounded functions on B_r and let Ψ_r be a map from \mathcal{D}_r to B_r . Let \mathcal{D}_{∞} be a subset of $\mathcal{B}_{loc}(\mathbb{R}^d)$ that is closed under multiplications by positive constants. Call the collection $\Psi = \{\Psi_r\}$ an **optimizing map** on \mathcal{D}_{∞} if the following conditions are satisfied for every function f in \mathcal{D}_{∞} :

- (i) $f_r \in \mathcal{D}_r$ for all r large enough and $\Psi_\infty[f] := \lim_{r \rightarrow \infty} \Psi_r(f_r)$ is a well defined point in \mathbb{R}^d ;
- (ii) if $\Psi_\infty[f] \in B_r$ then $f_r \in \mathcal{D}_r$ and $\Psi_r[f_r] = \Psi_\infty[f]$;
- (iii) $\Psi_\infty[af] = \Psi_\infty[f]$ for all $a > 0$;
- (iv) the map $\Psi_r : \mathcal{D}_r \rightarrow B_r$ is continuous with respect to the uniform metric at each point in the set $\{f_r : f \in \mathcal{D}_\infty, f \text{ is continuous and } \Psi_\infty[f] \in B_r\}$.

Define $\mathcal{D}_\infty(cts) = \{f \in \mathcal{D}_\infty, f \text{ is continuous}\}$. The following measurability property holds: if $h : \Omega \rightarrow \mathcal{B}_{loc}(\mathbb{R}^d)$ is a measurable map that takes values in $\mathcal{D}_\infty(cts)$, then $\Psi_\infty[h]$ is measurable. To prove this property it is sufficient to check that for every closed d -dimensional ball B , the set $\mathcal{F} = \{f \in \mathcal{D}_\infty(cts) : \Psi_\infty[f] \in B\}$ is closed in $\mathcal{D}_\infty(cts)$ with respect to the topology of uniform convergence on compacta. But if functions f_n in \mathcal{F} converge to a function f in $\mathcal{D}_\infty(cts)$, then continuity property (iv), applied with an r large enough to satisfy $B \subset B_r$ and $\Psi_\infty[f] \in B_r$, yields the needed result $\Psi_\infty[f_n] \rightarrow \Psi_\infty[f]$.

The following theorem is stated in the cleanest form, thus some of its conditions can be relaxed. For example, condition (iii) does not have to hold exactly, in fact, the criterion functions H_n do not even have to possess unique optimizers. A way to relax this condition is illustrated by Corollary 1, where the estimator is defined as an approximate minimizer. The continuity assumptions on the limit functions f and g can be relaxed at the cost of stronger continuity and measurability assumptions on the optimizing maps. Also, the properties required of sample paths the limit process only need to hold with probability one. The proof of the theorem would remain valid if $\mathcal{B}_{loc}(\mathbb{R}^d)$ were equipped with a different metric d , as long as the continuity assumption on the optimizing maps were formulated in terms of d .

Theorem 1. *Let H_n be random criterion functions on $\mathbb{R}^{d_1} \times \mathbb{R}^{d_2}$ and let Ψ and Φ be optimizing maps on functional classes $\mathcal{D}_\infty^\Psi \subset \mathcal{B}_{loc}(\mathbb{R}^{d_1})$ and $\mathcal{D}_\infty^\Phi \subset \mathcal{B}_{loc}(\mathbb{R}^{d_2})$ respectively. Assume that \mathcal{D}_∞^Φ is closed under translations and $\Phi_\infty[h + c] = \Phi_\infty[h]$ for all $h \in \mathcal{D}_\infty^\Phi$ and $c \in \mathbb{R}$.*

Let (s_n, t_n) be random vectors in $\mathbb{R}^{d_1} \times \mathbb{R}^{d_2}$ for which $H_n(\cdot, t_n) \in \mathcal{D}_\infty^\Psi$ and $H_n(s_n, \cdot) \in \mathcal{D}_\infty^\Phi$. Suppose that the following conditions are satisfied,

- (i) $H_n(s, t) = \alpha_n f_n(s) + \beta_n g_n(s, t)$, where f_n and g_n are random functions on \mathbb{R}^{d_1} and $\mathbb{R}^{d_1} \times \mathbb{R}^{d_2}$ respectively, while α_n and β_n are positive numbers with $\beta_n = o(\alpha_n)$;
- (ii) $(f_n, g_n) \rightsquigarrow (f, g)$ and the limit process has continuous sample paths;
- (iii) $s_n = \Psi_\infty[H_n(\cdot, t_n)]$ and $t_n = \Phi_\infty[H_n(s_n, \cdot)]$;
- (iv) $\|(s_n, t_n)\| = O_p^*(1)$.

Assume that $f \in \mathcal{D}_\infty^\Psi$ and $g(\Psi_\infty[f], \cdot) \in \mathcal{D}_\infty^\Phi$. Define $s^* = \Psi_\infty[f]$ and $t^* = \Phi_\infty[g(s^*, \cdot)]$. Then $(s_n, t_n) \rightsquigarrow (s^*, t^*)$.

The next result covers the special case of the estimators that are defined by minimization or maximization. Note that the sample path properties required of the limit process need to hold only almost surely. See Remark for the alternative to the continuity assumption placed on the sample paths of the stochastic process (f, g) . Note that the unique minimum x^* of a continuous function h is also its *clean* minimum in the sense that the strict inequality $h(x^*) < \inf_{\epsilon \leq |x - x^*| \leq r} h(x)$ is satisfied for all positive r and ϵ . In fact, lower semicontinuity of function h is sufficient.

Corollary 1. Suppose that conditions (i) and (ii) of Theorem 1 are satisfied. Let random vectors (s_n, t_n) be such that

- (i) $H_n(s_n, t_n) \leq \inf_{s, t} H_n(s, t) + o_p^*(\beta_n)$ and
- (ii) $\|(s_n, t_n)\| = O_p^*(1)$.

Assume that the sample paths of $f(\cdot)$ possess a unique minimum at a (random) point s^* and the sample paths of $g(s^*, \cdot)$ possess a unique minimum at t^* . Then $(s_n, t_n) \rightsquigarrow (s^*, t^*)$.

Remark. The assumptions on the sample paths of the limit process (f, g) can be relaxed as follows. Assume that s^* and t^* are measurable random points such that for almost all sample paths of the limit process:

- (a) s^* is the clean minimum of $f(\cdot)$,
- (b) t^* is the clean minimum of $g(s^*, \cdot)$ and

(c) for each ball B , the set of functions $\{g(\cdot, t) : t \in B\}$ is equicontinuous.

Proof. Redefine functions H_n so that points (s_n, t_n) become the unique minima. This can be done by leaving the functions f_n unchanged and decreasing the functions g_n by a $o_p^*(1)$ amount at exactly the points (s_n, t_n) . The assumptions of the corollary remain valid after the change. Take \mathcal{D}_r to be the set of those bounded functions on the d_1 -dimensional ball B_r that possess a unique minimum and take Ψ_r to be the arg min map on \mathcal{D}_r . Let \mathcal{D}_∞^Ψ consist of those locally bounded functions on \mathbb{R}^{d_1} that possess a unique minimum. Observe that $\Psi = \{\Psi_r\}$ satisfies the continuity assumption of Definition 1 and is an optimizing map on the class \mathcal{D}_∞^Ψ . Define Φ and \mathcal{D}_∞^Φ analogously with respect to \mathbb{R}^{d_2} and apply Theorem 1. \square

3. RATES OF CONVERGENCE

Consider two-component estimators (a_n, b_n) that are defined by minimizing random criterion functions $G_n(a, b)$. The following lemma uses an approximation to the criterion function to establish the rate of convergence of the slower converging component a_n and makes an initial guess at the rate of convergence of the component b_n . This guess is not quite correct, but it provides an improvement over existing results, which establish one convergence rate for the whole long vector (a_n, b_n) . Lemma 1 requires a particular representation for the criterion function. In many standard asymptotic problems, this representation is satisfied with the term $M_n(a, b)$ bounded below by a nonsingular quadratic, and the term $N_n(a, b)$ of the order $O_p(n^{-1/2} \|(a, b)\|)$, which yields the usual $n^{-1/2}$ rate of convergence. The lemma handles cases that are more general.

Lemma 1. *Suppose that inequalities $G_n(a_n, b_n) \leq G_n(0, 0)$ hold together with the stochastic bound $\|(a_n, b_n)\| = o_p^*(1)$. Let α and β be positive numbers satisfying $\alpha \geq \beta$, and let $\{\gamma_1, \dots, \gamma_p, \eta_1, \dots, \eta_p\}$ be a collection of nonnegative numbers satisfying $\gamma_i < \alpha$ for all $i \in \{1, \dots, p\}$. Suppose that criterion functions G_n satisfy a representation*

$$G_n(a, b) - G_n(0, 0) = M_n(a, b) - N_n(a, b),$$

such that

$$M_n(a_n, b_n) \gtrsim \|a_n\|^\alpha + \|b_n\|^\beta \quad \text{with inner probability tending to one, and} \\ [N_n(a_n, b_n)]^+ = O_p^*\left(\sum_{i \leq p} n^{-\eta_i} \|(a_n, b_n)\|^{\gamma_i}\right).$$

Define $\tau_a = \min_{i \leq p} \left(\frac{\eta_i}{\alpha - \gamma_i}\right)$. Then $\|a_n\| = O_p^*(n^{-\tau_a})$ and $\|b_n\| = O_p^*(n^{-\alpha\tau_a/\beta})$.

Once the convergence rate of a_n is established, it becomes reasonable to fix $a = a_n$ and consider the function $b \mapsto G_n(a_n, b)$. Existing results do not necessarily yield the convergence rate of the minimizer of this function. The point of difficulty is that the leading terms in the approximation to this function near its minimum are more complex than the ones that appear in the standard asymptotics. The following theorem can handle such cases but it requires a more refined approximation to the criterion function. One may want to use the help of Lemma 1 to obtain such an approximation (see, for example, the proof of Theorem 3), and then apply Theorem 2 to derive the “correct” convergence rate of b_n . Note that Theorem 2 places no assumptions at all on the space containing the a -component.

Theorem 2. *Let $G_n(a, b)$ be a function of two components, where the first component belongs to an abstract set, and the second belongs to a Euclidean space. Suppose that inequalities $G_n(a_n, b_n) \leq G_n(a_n, 0)$ hold together with the stochastic bound $\|b_n\| = o_p^*(1)$. Let β be positive and let $\{\alpha_1, \dots, \alpha_p, \beta_1, \dots, \beta_p\}$ be a collection of nonnegative numbers satisfying $\beta_i < \beta$ for all $i \in \{1, \dots, p\}$. Assume that G_n satisfies a representation*

$$(1) \quad G_n(a, b) - G_n(a, 0) = M_n(a, b) - N_n(a, b),$$

such that

$$M_n(a_n, b_n) \gtrsim \|b_n\|^\beta \quad \text{with inner probability tending to one, and} \\ [N_n(a_n, b_n)]^+ = O_p^*\left(\sum_{i \leq p} n^{-\alpha_i} \|b_n\|^{\beta_i}\right).$$

Then $\|b_n\| = O_p^*(n^{-\tau_b})$ for $\tau_b = \min_{i \leq p} \left\{\frac{\alpha_i}{\beta - \beta_i}\right\}$. If $[N_n]^+ \equiv 0$ then $\mathbb{P}_*\{b_n = 0\} \rightarrow 1$.

4. M ESTIMATORS

The following definition introduces notation that is used in the statement of Theorem 3. This notation simplifies the work with polynomials that are homogeneous functions of the elements of vector (a, b) and the absolute values of the elements of vector (a, b) .

Definition 2. Let ψ be a real valued function on \mathbb{R}^d and let γ be a positive constant. Say that $\psi \in H_1^+(\gamma)$ if $\psi(\lambda\theta) = \lambda^\gamma \psi(\theta)$ for all $\lambda \geq 0$ and $\psi(\theta) > 0$ for all $\theta \neq 0$.

Let ϕ be a real valued function on $\mathbb{R}^{d_1} \times \mathbb{R}^{d_2}$ and let α and β be some positive constants. Say that $\phi \in H_2^{(-)}(\alpha, \beta)$ if $\phi(\lambda_1 a, \lambda_2 b) = (\lambda_1)^\alpha (\lambda_2)^\beta \phi(a, b)$ for all nonnegative λ_1 and λ_2 , while function ϕ assumes at least some negative values.

Remark. For each continuous function $\psi(\theta)$ in the class $H_1^+(\gamma)$, there exist positive constants c_1 and c_2 such that $c_1 \|\theta\|^\gamma \leq \psi(\theta) \leq c_2 \|\theta\|^\gamma$.

Suppose that X_1, X_2, \dots, X_n are independent observations in \mathbb{R}^k coming from a distribution P and write P_n for the corresponding empirical distribution. Suppose that A is an open subset of $\mathbb{R}^{d_1} \times \mathbb{R}^{d_2}$ and let $\{g_{a,b}(x) : (a, b) \in A\}$ be a collection of real valued P -integrable functions on \mathbb{R}^k . Assume that this collection of functions is centered to satisfy $g_{0,0} \equiv 0$. Suppose that vectors (a_n, b_n) minimize over A the random criterion functions $G_n(a, b) = P_n g_{a,b}$ and let $(0, 0)$ be the corresponding minimizer of the population analog $G(a, b) = P g_{a,b}$. The following theorem derives the asymptotics of (a_n, b_n) in the challenging case of the singular second derivative matrix $G''(0, 0)$.

Theorem 3. Let $\{\alpha, \beta, \gamma_1, \dots, \gamma_p, \eta_1, \dots, \eta_p\}$ be a collection of positive numbers. Assume that $\alpha > \beta > 1$ and $\beta > \eta_j$ for $1 \leq j \leq p$. Suppose that there exist continuous functions $\psi_1(a) \in H_1^+(\alpha)$, $\psi_2(b) \in H_1^+(\beta)$ and $\phi_i(a, b) \in H_2^{(-)}(\gamma_i, \eta_i)$ for $1 \leq i \leq p$, such that near the origin the population criterion function satisfies the conditions

- (i) $G(a, b) \gtrsim \|a\|^\alpha + \|b\|^\beta$,
- (ii) $G(a, 0) = \psi_1(a) + o(\|a\|^\alpha)$ and
- (iii) $G(a, b) = G(a, 0) + \psi_2(b) [1 + o(1)] + \sum_{i=1}^p \phi_i(a, b) [1 + o(1)] + o\left(\sum_{i=1}^p \|a\|^{\alpha-i} \|b\|^i\right)$.

Let $\tau_a = \frac{1}{2(\alpha-1)}$, $\lambda_0 = \frac{1}{2(\beta-1)}$, $\lambda_i = \frac{\tau_a \gamma_i}{\beta - \gamma_i}$ for $1 \leq i \leq p$, and define $\tau_b = \min_{0 \leq i \leq p} [\lambda_i]$. Suppose there exist on \mathbb{R}^k five square integrable functions, Δ_1 (taking values in \mathbb{R}^{d_1}), Δ_2 (taking values in \mathbb{R}^{d_2}) and real valued $r_{a,b}$, $s_{a,b}$ and $l_{a,b}$, such that:

- (iv) $g_{a,b}(x) = a' \Delta_1(x) + b' \Delta_2(x) + \|(a, b)\| r_{a,b}(x)$;
- (v) $g_{a,b}(x) - g_{a,0}(x) - b' \Delta_2(x) = l_{a,b}(x) + \|b\| s_{a,b}(x)$;
- (vi) $\sup_{\|(a,b)\| \leq \delta_n} |\nu_n r_{a,b}| = o_p(1)$ and $\sup_{\|(a,b)\| \leq \delta_n} |\nu_n s_{a,b}| = o_p(1)$ for all $\delta_n \rightarrow 0$;
- (vii) $\sup_{\|a\| \leq \delta_n, \|b\| \leq \epsilon_n} |\nu_n l_{a,b}| = o_p(n^{-\beta\tau_b+1/2})$ for all $\delta_n = O(n^{-\tau_a})$, $\epsilon_n = O(n^{-\alpha\tau_a/\beta})$.

Assume that $\|(a_n, b_n)\| = o_p(1)$. If $\alpha\tau_a = \beta\tau_b$, then

$$(n^{\tau_a} a_n, n^{\tau_b} b_n) \rightsquigarrow \arg \min_{s,t} [\psi_1(s) + s' Z_1 + \psi_2(t) + 1\{\lambda_0 = \tau_b\} t' Z_2 + \sum_{i=1}^p 1\{\lambda_i = \tau_b\} \phi_i(s, t)];$$

otherwise $(n^{\tau_a} a_n, n^{\tau_b} b_n) \rightsquigarrow (s^*, t^*)$, where

$$\begin{aligned} s^* &= \arg \min_s [\psi_1(s) + s' Z_1], \\ t^* &= \arg \min_t [\psi_2(t) + 1\{\lambda_0 = \tau_b\} t' Z_2 + \sum_{i=1}^p 1\{\lambda_i = \tau_b\} \phi_i(s^*, t)]. \end{aligned}$$

Here (Z_1, Z_2) is a mean zero gaussian vector with covariance matrix $P(\Delta_1, \Delta_2)(\Delta_1, \Delta_2)'$.

Note that a stochastic process $\nu_n f_{a,b}$ necessarily satisfies the uniform stochastic bound required in condition (vi) of the above theorem (cf *asymptotic equicontinuity* defined in van der Vaart 1998) if functions $f_{a,b}$ form a Donsker class and $\|f_{a,b}\|_2 \rightarrow 0$ as $\|(a, b)\| \rightarrow 0$. Simple ways of checking that a class of functions is Donsker are given, for example, in van der Vaart's Theorem 19.5 and Theorem 19.14.

To illustrate the variety of asymptotic results produced by Theorem 3, consider some simple approximations to the function G , which has a singular second derivative at the origin, where its minimum is located. Let $(a, b) \in \mathbb{R}^2$ and consider the case $G(a, b) \approx a^4 + b^2$. Theorem 3 yields $(n^{1/6} a_n, n^{1/2} b_n) \rightsquigarrow (\arg \min_s [s^4 + s Z_1], \arg \min_t [t^2 + t Z_2])$ if the conditions (iv) -(vii) are satisfied. Here (Z_1, Z_2) is a mean zero gaussian vector. Now consider the case $G(a, b) \approx a^4 + b^2 + a^2 b$. Under the same assumptions, the theorem yields $(n^{1/6} a_n, n^{1/3} b_n) \rightsquigarrow (\arg \min_{s,t} [s^4 + s Z_1 + t^2 + s^2 t])$. If the approximation is $G(a, b) \approx a^4 + b^2 + a^3 b$, the corresponding result is $(n^{1/6} a_n, n^{1/2} b_n) \rightsquigarrow (s^*, t^*)$

with $s^* = \arg \min_s [s^4 + sZ_1]$ and $t^* = \arg \min_t [t^2 + (s^*)^3 t + tZ_2]$. Note that Theorem 3 does not attempt to cover every conceivable approximation to $G(a, b)$, as the statement of the result would become too long and complicated, but each such situation can be handled with only minor modifications to the proof of the theorem.

5. EXAMPLE: LASSO-TYPE ESTIMATORS

Assume that the observed variables Y_i satisfy the linear model

$$Y_i = x_i' \beta + \epsilon_i, \quad i = 1, \dots, n.$$

The errors ϵ_i are independent and identically distributed random variables that have mean zero and variance σ^2 . The parameter β is a vector in \mathbb{R}^d that needs to be estimated. The covariates x_i are fixed and centered, and the matrix $C_n = \frac{1}{n} \sum_{i=1}^n x_i x_i'$ is nonsingular.

Suppose λ_n and γ are positive real numbers. Define the "Lasso-type" estimator β_n as the minimizer of the penalized least-squares criterion,

$$W_n(\alpha) = \sum_{i=1}^n (Y_i - x_i' \alpha)^2 + \lambda_n \sum_{j=1}^d |\alpha_j|^\gamma,$$

over all vectors $\alpha = (\alpha_1, \dots, \alpha_d)'$. In the particular cases of $\gamma = 1$ and $\gamma = 2$, this estimator corresponds, respectively, to the "Lasso" of Tibshirani (1996) and the ridge regression. For general γ such estimators were introduced by Frank and Friedman (1993). The limiting behavior of the estimator β_n was described by Knight and Fu (2000) under certain conditions on the growth rate of the weighting sequence $\{\lambda_n\}$.

Assume that the design satisfies the following regularity conditions:

- (i) matrixes C_n converge to a fixed matrix C ;
- (ii) as n tends to infinity, $n^{-1} \max_{i \leq n} (x_i' x_i)$ converges to zero.

In the case of the nonsingular matrix C , Knight and Fu derived the \sqrt{n} -asymptotics for β_n after setting the growth rate for the weighting sequence $\{\lambda_n\}$. They required that for some nonnegative constant λ_0 ,

$$(2) \quad \lambda_n / n^{\min(1/2, \gamma/2)} \rightarrow \lambda_0.$$

Note that when $\lambda_0 = 0$, the penalty contribution is asymptotically negligible and the limiting behavior of the estimator β_n is the same as that of the usual least-squares estimator.

To derive the asymptotics of β_n , Knight and Fu used a standard approach that is based on rescaling the parameters at the same rate and applying a continuous mapping type of argument. When vector β has a zero component, $\gamma < 1$, and λ_n grows faster than the rate given in (2), this approach fails to deliver the complete asymptotics. For concreteness, consider the case $d = 2$, $\beta = (1, 0)'$, $\gamma = 1/2$, and set $\lambda_n = \lambda_0 n^{1/2}$ for some positive constant λ_0 . The standard approach establishes the asymptotics of the first component of β_n , but only yields the $o_p(n^{-1/2})$ stochastic order for the second component of the estimator (see Knight and Fu 2000, page 1361). The techniques developed in Section 3 are applied below to show that the second component is in fact exactly zero with probability tending to one.

Because C is nonsingular and $\lambda_n = o(n)$, the estimator β_n is consistent (see Theorem 1 of Knight and Fu.) The proof is based on the fact that for each fixed α , the penalty part of the criterion function $W_n(\alpha)$ is asymptotically negligible compared to the least-squares part. Focus on vectors α that are near the true parameter β , and write $\alpha = \beta + (a, b)'$. Express the penalized criterion function in terms of a and b . Denote $n^{-1}[W_n(\alpha) - W_n(\beta)]$ by $G_n(a, b)$, and let Z_n stand for $n^{-1/2} \sum_{i=1}^n \epsilon_i x_i$. The regularity conditions on the design guarantee that the sequence of random vectors Z_n has a limiting gaussian distribution with mean zero and covariance $\sigma^2 C$. As a and b tend to zero,

$$G_n(a, b) = (a, b)C_n(a, b)' - 2n^{-1/2}(a, b)Z_n + \frac{\lambda_0}{2}n^{-1/2}a[1 + o(1)] + \lambda_0 n^{-1/2}|b|^{1/2}.$$

The $o(1)$ terms come from the Taylor expansion of $|1 + a|^{1/2}$ near $a = 0$. Function G_n is minimized by the vector $(a_n, b_n)'$ that is defined as the difference between β_n and β .

Define $M_n(a, b)$ to be $(a, b)C_n(a, b)' + \lambda_0 n^{-1/2}|b|^{1/2}$ and let N_n equal $G_n - M_n$. Note that for all n large enough, the eigen values of the sequence of matrixes C_n are bounded away from zero. Apply Lemma 1 from Section 3 and conclude that $\|(a, b)\| = O_p(n^{-1/2})$.

Let v_n denote the bottom right element of the matrix C_n . Observe that

$$G_n(a_n, b_n) - G_n(a_n, 0) = v_n b_n^2 + O_p(n^{-1/2}|b_n|) + \lambda_0 n^{-1/2}|b_n|^{1/2}.$$

Note that λ_0 and v_n are positive and v_n is bounded away from zero for all sufficiently large n . Deduce that with probability tending to one, the right hand side of the above display is bounded below by cb_n^2 for some positive c . Apply Theorem 2 with $N_n \equiv 0$ and conclude that $\mathbb{P}\{b_n = 0\} \rightarrow 1$.

More examples of mixed-rates behavior in Lasso-type estimation can be found in Radchenko (2005).

6. EXAMPLE: SHORTH

Assume that the observations are independently sampled from a distribution P on the real line and let $[m_n - r_n, m_n + r_n]$ be the shortest interval that contains at least half of the first n observations. The shorth estimator is defined as the average over such an interval, but the goal of this section is the limiting behavior of m_n and r_n . Grübel (1988) derived the root- n asymptotics for r_n and Kim and Pollard (1990) derived the cube root asymptotics for m_n . The methods of the present paper allow one to establish the joint limiting behavior of (m_n, r_n) using a simple approximation to the criterion function.

Denote by μ and ρ the population solution, in other words let $[\mu - \rho, \mu + \rho]$ be the shortest interval to which P assigns at least half the probability. Assume that the population solution is unique and let P have a bounded density f that is differentiable at the endpoints $\mu \pm \rho$. Define the criterion function V_n by $V_n(\epsilon, \delta) = P_n[(\mu + \epsilon) - (\rho + \delta), (\mu + \epsilon) + (\rho + \delta)] - 1/2$, and let $V(\epsilon, \delta)$ denote the population analog obtained by replacing P_n with P . Observe that $V(0, 0) = 0$ and write out the Taylor expansion for function V near the origin:

$$(3) \quad V(\epsilon, \delta) = c_1\delta + c_2\epsilon^2 + c_3\epsilon\delta + c_4\delta^2 + o(\epsilon^2 + \delta^2),$$

where the coefficients are $c_1 = f(\mu - \rho) + f(\mu + \rho)$, $c_2 = c_4 = [f'(\mu + \rho) - f'(\mu - \rho)]/2$ and $c_3 = f'(\mu + \rho) + f'(\mu - \rho)$. The coefficient of the linear term in ϵ equals zero because the function $V(\epsilon, 0)$ is maximized at $\epsilon = 0$. This forces the equality $f(\mu + \rho) = f(\mu - \rho)$. By the same reasoning, coefficient c_2 must be non-positive. Assume $c_2 < 0$ and $c_1 > 0$ for regularity.

Recall the bound $\sup_{m,r} |P_n[m-r, m+r] - P[m-r, m+r]| = O_p(n^{-1/2})$ from the standard empirical process theory. Denote this supremum by Δ_n . Uniqueness of the population solution and regularity assumptions on the coefficients of Taylor expansion (3) guarantee that there exists a positive constant c such that for all small enough positive δ , inequality $\sup_m P[m - (\rho - \delta), m + (\rho - \delta)] < 1/2 - c\delta$ holds. Consequently,

$$\sup_m P_n[m - (\rho - \Delta_n/c), m + (\rho - \Delta_n/c)] < \Delta_n + 1/2 - c\Delta_n/c = 1/2,$$

and hence $\delta_n \geq -\Delta_n/c$. Expansion (3) also implies existence of a positive constant b such that $P[\mu - (\rho + \delta), \mu + (\rho + \delta)] \geq 1/2 + b\delta$ for all small enough positive δ . Take $\delta = \Delta_n/b$ and deduce that $P_n[\mu - (\rho + \Delta_n/b), \mu + (\rho + \Delta_n/b)] \geq 1/2$. Conclude that $\delta_n \leq \Delta_n/b$ and hence $\delta_n = O_p(n^{-1/2})$.

Note that function $V_n(\cdot, \delta_n)$ is maximized by ϵ_n and function $V(\cdot, 0)$ has a clean maximum at zero. Uniform convergence in probability of $V_n(\cdot, \delta_n)$ to $V(\cdot, 0)$ implies $\epsilon_n = o_p(1)$. Introduce functions $M_n(\epsilon, \delta) = |c_2|\epsilon^2/4$ and define functions N_n by equalities

$$- [V_n(\epsilon, \delta) - V_n(0, \delta)] = M_n(\epsilon, \delta) - N_n(\epsilon, \delta).$$

Note that when $\delta = \delta_n$, the expression on left hand side is minimized by $\epsilon = \epsilon_n$. Denote the difference between the indicator functions of intervals $[(\mu + \epsilon) - (\rho + \delta), (\mu + \epsilon) + (\rho + \delta)]$ and $[\mu - (\rho + \delta), \mu + (\rho + \delta)]$ by $J(\epsilon, \delta)$. Observe that

$$V_n(\epsilon, \delta) - V_n(0, \delta) = V(\epsilon, \delta) - V(0, \delta) + (P_n - P)J(\epsilon, \delta).$$

Recall that $c_2 < 0$ by the regularity assumptions placed on the coefficients of expansion (3), and use Taylor expansion (3) to deduce a stochastic bound

$$(4) \quad N_n(\epsilon_n, \delta_n) \leq (P_n - P)J(\epsilon_n, \delta_n) - |c_2|\epsilon_n^2/2 + O_p(n^{-1/2}|\epsilon_n|) + O_p(n^{-1}).$$

Note that the collection of functions $J(\epsilon, \delta)$ is a Vapnik-Cervonenkis class. For R near zero, the envelope function $G_R = \sup_{\{\epsilon^2 + \delta^2 \leq R^2\}} |J(\epsilon, \delta)|$ is the indicator of the two interval of total length bounded above by $4R$; boundedness of the density implies $PG_R^2 = O(R)$. Hence

the conditions of Lemma 4.1 of Kim and Pollard (1990) are satisfied and, consequently, the bound $|(P_n - P)J(\epsilon_n, \delta_n)| - c\epsilon_n^2 \leq O_p(n^{-2/3})$ is valid for each positive c . It follows that

$$[(P_n - P)J(\epsilon_n, \delta_n) - |c_2|\epsilon_n^2/2]^+ = O_p(n^{-2/3}).$$

Combine this stochastic bound with bound (4) and deduce that

$$[N_n(\epsilon_n, \delta_n)]^+ = O_p(n^{-1/2}|\epsilon_n|) + O_p(n^{-2/3}).$$

An application of Theorem 2 yields $\epsilon_n = O_p(n^{-1/3})$.

Set $I_{s,t} = 1[(\mu + n^{-1/3}t) - (\rho + n^{-1/2}s), (\mu + n^{-1/3}t) + (\rho + n^{-1/2}s)] - 1[\mu - \rho, \mu + \rho]$ and define the localized criterion functions $H_n(s, t) = V_n(n^{-1/3}t, n^{-1/2}s)$. Use the empirical process notation to write an approximation to $H_n(s, t)$ that holds uniformly on compacta:

$$(5) \quad H_n(s, t) = n^{-1/2} \left\{ c_1 s + \nu_n [\mu - \rho, \mu + \rho] \right\} + n^{-2/3} \left\{ c_2 t + n^{1/6} \nu_n I_{s,t} + o(1) \right\}.$$

On each compact set, the stochastic processes $X_n(s, t) = n^{1/6} \nu_n I_{n^{1/6}s, t}$ converges in distribution to a tight Gaussian process by Theorem 19.28 of van der Vaart (1998). The conditions of the theorem are checked in van der Vaart's Example 19.29 for essentially the same process as X_n . Consequently, X_n satisfies the asymptotic equicontinuity condition of van der Vaart's Theorem 18.14, and approximation $n^{1/6} \nu_n [I_{s,t} - \nu_n I_{0,t}] = o_p(1)$ holds uniformly over s and t in each given compact set. Note that

$$\begin{aligned} \lim_{n \rightarrow \infty} n^{1/3} P I_{0,t} I_{0,t'} &= c_1 \min(|t|, |t'|) \quad \text{and} \\ \lim_{n \rightarrow \infty} n^{1/6} P I_{0,t} 1[\mu - \rho, \mu + \rho] &= 0. \end{aligned}$$

Write $f_n(s)$ and $g_n(s, t)$ for the two expressions in curly brackets that appear in representation (5). Let $B(t)$ be a two-sided Brownian motion and let Z be an independent $N(0, 1/2)$ random variable. Conclude that

$$(f_n(s), g_n(s, t)) \rightsquigarrow (c_1 s - Z, c_2 t + \sqrt{c_1} B(t)),$$

and, consequently, conditions (i) and (ii) of Theorem 1 are satisfied. Recall that the rescaled solution $(s_n, t_n) = (n^{1/2}\delta_n, n^{1/3}\epsilon_n)$ is stochastically bounded and note the relationship

$$s_n = \inf \{s : H_n(s, t_n) \geq 0\} \quad \text{and} \quad t_n = \arg \max [H_n(s_n, \cdot)].$$

Consider functionals $\Psi_r : h \mapsto \inf\{s \in [-r, r] : h(s) \geq 0\}$ and let \mathcal{D}_∞^Ψ consist of those locally bounded functions h on the real line for which the value $\inf\{s : h(s) \geq 0\}$ is finite. Note that $\Psi = \{\Psi_r\}$ is an optimizing map on \mathcal{D}_∞^Ψ , and functions $H_n(\cdot, t)$ belong to \mathcal{D}_∞^Ψ for each value of t . Define \mathcal{D}_∞^Φ to be the set of those locally bounded functions on the real line that possess a unique maximum. Let $\{\Phi_r\}$ be the arg max functionals corresponding to functions defined on intervals $[-r, r]$. Note that $\Phi = \{\Phi_r\}$ is an optimizing map on \mathcal{D}_∞^Φ invariant to translations, and functions $g_n(s_n, \cdot)$ belong to \mathcal{D}_∞^Φ with probability one. Apply Theorem 1 and express the result in terms of the original variables:

$$(n^{1/2}[r_n - r], n^{1/3}[m_n - \mu]) \rightsquigarrow (Z/c_1, \arg \max_t [c_2 t + \sqrt{c_1} B(t)]).$$

Standard techniques fail to extract the joint limiting behavior of the estimators from approximation (5) because the first component of the approximation dominates the essential second component as n tends to infinity.

7. EXAMPLE: k -MEANS

The k -means procedure divides observations X_1, \dots, X_n in \mathbb{R}^d into k sets by locating the cluster centers and then assigning each observation to the closest center. The set of cluster centers $C_n = \{c_{1n}, \dots, c_{kn}\}$ is chosen to minimize

$$(6) \quad W_n(C) = n^{-1} \sum_{i \leq n} \min_{1 \leq j \leq k} \|x_i - c_j\|^2$$

as a function of sets $C = \{c_1, \dots, c_k\}$ of k not necessarily distinct points in \mathbb{R}^d . Assume that the observations are independent and come from a distribution P on \mathbb{R}^d . Let $W(C)$ be the population counterpart to the criterion function $W_n(C)$,

$$(7) \quad W(C) = P W_n(C) = P \min_{1 \leq j \leq k} \|X_1 - c_j\|^2,$$

and let C_0 be a set that minimizes the function W . Call C_n a set of *optimal sample centers*, and call C_0 a set of *optimal population centers*.

Note that if P has a finite second moment and is not concentrated on fewer than k points, then each set of optimal population centers has to contain exactly k points. Under these conditions, and given that the set C_0 of optimal population centers is uniquely defined, Pollard (1981) showed that the sets C_n of optimal empirical centers are strongly consistent with respect to the Hausdorff metric.

In the examples that follow, condition (vi) of Theorem 3 needs to be verified for classes of functions that possess the following simple property.

Property 1. *The class of functions $f_\theta(x)$ satisfies the following conditions:*

- (i) *the envelope function $F(x)$ is square integrable with respect to P ;*
- (ii) *there exist positive integers N and d such that each f_θ can be represented as a sum of at most N functions of the form LQ , where L is a linear function and Q is the indicator function of the intersection of at most N half-spaces in \mathbb{R}^d ;*
- (iii) $\|f_\theta\|_2 \rightarrow 0$ as $\theta \rightarrow 0$.

The first two conditions imply that the class of functions f_θ is Donsker. This fact is proved on page 921 of Pollard (1982), but it can also be easily deduced from the standard results on pages 274-276 of van der Vaart (1998). The third condition together with the Donsker property yield the required $\sup_{\|\theta\| \leq \delta_n} |\nu_n f_\theta| \rightarrow 0$ for each $\delta_n \rightarrow 0$.

7.1. 2-means for the double-exponential distribution. Let P be the double-exponential distribution on the real line and consider the case $k = 2$. Serinko and Babu (1992) derived the $n^{-1/4}$ rate asymptotics for the optimal sample centers and showed that the distance between the centers settles down to a constant faster than the rate $n^{-1/4}$. A simple application of Theorem 3 will refine their result and derive the $n^{-1/2}$ rate asymptotics for the distance between the optimal sample centers.

The optimal population centers for this problem are located at -1 and 1. Introduce new variables for convenience: for each set C of potential centers $c_1 \leq c_2$, define $a = (c_1 + c_2)/2$

and $b = 1 + (c_1 - c_2)/2$. Let $g_{a,b}(x)$ be the squared distance from x to the closest center, written in terms of a and b and centered at $(a, b) = (0, 0)$:

$$\begin{aligned} g_{a,b}(x) &= (x - c_1)^2 \wedge (x - c_2)^2 - (x + 1)^2 \wedge (x - 1)^2 \\ &= (x + 1 - a - b)^2 \{x \leq a\} + (x - 1 - a + b)^2 \{x > a\} - (x + 1)^2 \wedge (x - 1)^2. \end{aligned}$$

Define criterion functions $G_n(a, b) = P_n g_{a,b}$ and $G(a, b) = P g_{a,b}$, and note that they are just the functions $W_n(C)$ and $W(C)$ centered at the set of optimal population centers and written in terms of the new variables. Thus, if (a_n, b_n) minimizes G_n , then a_n is the midpoint between the optimal sample centers, and $1 - b_n$ equals the half distance between the centers. Also note that $G(a, b)$ is minimized at zero and $\|(a_n, b_n)\| = o_p(1)$ because of consistency.

The following approximation holds for the criterion function G near zero (see, for example, the derivation in subsection 2.2.1 of Radchenko 2004),

$$G(a, b) = |a|^3/3 + b^2 + a^2b + O(\|(a, b)\|^4).$$

Note that conditions (i), (ii) and (iii) of Theorem 3 are satisfied with $\alpha = 3, \beta = 2, p = 1, \gamma_1 = 2$ and $\eta_1 = 1$. Take functions $\Delta_1(z)$ and $\Delta_2(z)$ in condition (iv) as the L_2 partial derivatives of $g_{a,b}(z)$ with respect to a and b at $(0, 0)$. Because $g_{a,b}(0)$ is not differentiable at zero, define $\Delta_1(0)$ and $\Delta_2(0)$ arbitrarily. For example, set

$$\begin{aligned} \Delta_1(x) &= -2(x + 1)\{x \leq 0\} - 2(x - 1)\{x > 0\} \\ \Delta_2(x) &= -2(x + 1)\{x \leq 0\} + 2(x - 1)\{x > 0\}. \end{aligned}$$

Condition (vi) for the remainder functions $r_{a,b}(x)$ follows from a general result on k -means (see Pollard 1982, Lemma B). The proof essentially consists of checking that Property 1 holds for the class $\{r_{a,b}\}$. Note that $P\Delta_1^2 = P\Delta_2^2 = 4$ and $P\Delta_1\Delta_2 = 0$.

Let $l_{a,b} \equiv 0$ and define $s_{a,b}(x) = [g_{a,b}(x) - g_{a,0}(x) - b\Delta_2(x)]/|b|$. Then

$$\begin{aligned} s_{a,b}(x) &= [-b(2x + 2 - 2a - b)\{x \leq a\} + b(2x - 2 - 2a + b)\{x > a\} \\ &\quad + b(2x + 2)\{x \leq 0\} - b(2x - 2)\{x > 0\}]/|b| \end{aligned}$$

Note that Property 1 holds for the class of functions $s_{a,b}$. Hence, $\sup_{\|(a,b)\| < \delta_n} |\nu_n s_{a,b}| \rightarrow 0$ for each $\delta_n \rightarrow 0$.

All the conditions of Theorem 3 are now satisfied. Let (Z_1, Z_2) be a mean zero gaussian vector with covariance matrix $4I$ and conclude that $(n^{1/4}a_n, n^{1/2}b_n)$ converges in distribution to the random vector (s^*, t^*) that satisfies $s^* = \arg \min_s [|v|^3/3 + s'Z_1]$ and $t^* = \arg \min_t [t^2 + t'Z_2 + (s^*)^2 t]$. A closed-form expression for (s^*, t^*) can be found in subsection 2.2.3 of Radchenko (2004).

7.2. 2-means for a distribution in \mathbb{R}^2 . The following example is a two-dimensional extension of the one given above. Consider a distribution Q on the plane (x, y) that concentrates on the lines $\{x = 1\}$ and $\{x = -1\}$. Let Q put probability one half on each line, and let the conditional distribution on each line be the double exponential. Write Q as $P \times \mu$, where P is the double exponential distribution and $\mu\{-1\} = \mu\{1\} = 1/2$.

There are two pairs of optimal population centers, $\{(-1, 0), (1, 0)\}$ and $\{(0, -1), (0, 1)\}$; denote them by $C^v = \{c_1^v, c_2^v\}$ and $C^h = \{c_1^h, c_2^h\}$ respectively. The superscripts reflect either the vertical or the horizontal direction of the *split line*, which is defined as the common boundary for the two Voronoi half-planes generated by a given pair of centers. Let C_n^v and C_n^h minimize the criterion function (6) over two fixed non-overlapping Hausdorff neighborhoods of the sets C^v and C^h respectively, and let C_n be a global minimizer. A slight extension of Pollard's consistency result yields

$$\begin{aligned} C_n &\in \{C_n^v, C_n^h\} \quad \text{with probability tending to one,} \\ C_n^h &\rightarrow C^h \quad \text{and} \quad C_n^v \rightarrow C^v \quad \text{almost surely,} \end{aligned}$$

where the set convergence is understood with respect to the Hausdorff metric. In fact, the probability with which C_n chooses between the two configurations converges to a half. Near the set C^h , the population criterion function W is approximated by a nonsingular quadratic. As a result, the solution C_n^h settles down at the standard $n^{-1/2}$ rate and satisfies a

central limit theorem. The remainder of the section is concerned with deriving the asymptotics of C_n^v , which is a challenging problem because the quadratic approximation to the population criterion function near the set C^v is singular.

Suppose that $C = \{c_1, c_2\}$ is a candidate to minimize the criterion function (6) over a small Hausdorff neighborhood of the set C^v . Let $c_1 = (c_{1x}, c_{1y})$ be the point lying close to $c_2^v = (-1, 0)$ and let $c_2 = (c_{2x}, c_{2y})$ lie close to $c_2^v = (1, 0)$. Write z to denote a point on the plane and let (x, y) be the coordinate form of z . Introduce new variables by

$$\delta_s = \frac{1}{2}(c_{1x} + c_{2x}), \quad \delta_d = 1 + \frac{1}{2}(c_{1x} - c_{2x}), \quad \epsilon_s = \frac{1}{2}(c_{1y} - c_{2y}) \quad \text{and} \quad \epsilon_d = \frac{1}{2}(c_{1y} + c_{2y}).$$

These variables contain the information on how far the centers in the set C lie from the corresponding centers in C^v . Define $a = (\delta_s, \epsilon_d)$ and $b = (\delta_d, \epsilon_s)$. Let $g_{a,b}(z)$ be the squared distance from z to the closest center in C , written in terms of (a, b) and centered:

$$\begin{aligned} g_{a,b}(z) &= [(x + 1 - \delta_s - \delta_d)^2 + (y - \epsilon_s - \epsilon_d)^2] \wedge [(x - 1 - \delta_s + \delta_d)^2 + (y - \epsilon_s + \epsilon_d)^2] \\ &\quad - \|z + (1, 0)\|^2 \wedge \|z - (1, 0)\|^2. \end{aligned}$$

Define criterion functions $G_n(a, b) = P_n g_{a,b}$ and $G(a, b) = P g_{a,b}$, and note that they are just the functions $W_n(C)$ and $W(C)$ centered at the set C^v and written in terms of the new variables. Note that $G(a, b)$ is minimized at zero and the points (a_n, b_n) that minimize $G_n(a, b)$ are of order $o_p(1)$ because of consistency.

The following approximation holds for G near zero (see Section 2.3 of Radchenko 2004),

$$G(a, b) = \frac{1}{6}(|\delta_s| + |\epsilon_d|)^3 + \frac{1}{6}||\delta_s| - |\epsilon_d||^3 + \delta_d^2 + \epsilon_s^2 + \delta_s^2 \delta_d + 2\delta_s \epsilon_d \epsilon_s - \epsilon_d^2 \delta_d + O(\|(a, b)\|^4).$$

Note that conditions (i), (ii) and (iii) of Theorem 3 are satisfied with $\alpha = 3, \beta = 2, p = 3$ and $(\gamma_i, \eta_i) = (2, 1)$ for $1 \leq i \leq 3$. The corresponding homogeneous functions are

$$\begin{aligned} \psi_1(a) &= \frac{1}{6}(|\delta_s| + |\epsilon_d|)^3 + \frac{1}{6}||\delta_s| - |\epsilon_d||^3, \quad \psi_2(b) = \delta_d^2 + \epsilon_s^2, \quad \text{and} \\ \phi_1(a, b) &= \delta_s^2 \delta_d, \quad \phi_2(a, b) = 2\delta_s \epsilon_d \epsilon_s, \quad \phi_3(a, b) = -\epsilon_d^2 \delta_d. \end{aligned}$$

Take functions $\Delta_1(z)$ and $\Delta_2(z)$ in condition (iv) as the L_2 partial derivatives of $g_{a,b}(z)$ with respect to a and b at $(0, 0)$. For example, let

$$b' \Delta_2(z) = -[2\delta_d(x+1) + 2\epsilon_s y] H_-(z) + [2\delta_d(x-1) - 2\epsilon_s y] H_+(z),$$

where $H_-(z)$ is the indicator function of the half-plane $\{z : x \leq 0\}$ and $H_+(z)$ is the indicator functions of the half-plane $\{z : x > 0\}$. Condition (vi) for the remainder functions $r_{a,b}(x)$ follows from the fact that Property 1 holds for the class $\{r_{a,b}\}$.

The expression for $g_{a,b}(z)$ depends on on the sign of x and on which of the centers in the set C lies closer to the point z . Let D and U be the x -coordinates of the crossing points of the split line corresponding to C with the lines $\{y = -1\}$ and $\{y = 1\}$ respectively. Note that when $b = 0$, the values D and U are simply $\delta_s - \epsilon_d$ and $\delta_s + \epsilon_d$. Introduce functions

$$\begin{aligned} A(z) &= 1\{|x| \leq |D|, xD > 0, y = -1\} + 1\{|x| \leq |U|, xU > 0, y = 1\} \quad \text{and} \\ A^0(z) &= A(z) - 1\{|x| \leq |\delta_s - \epsilon_d|, x(\delta_s - \epsilon_d) > 0, y = -1\} \\ &\quad - 1\{|x| \leq |\delta_s + \epsilon_d|, x(\delta_s + \epsilon_d) > 0, y = 1\}. \end{aligned}$$

Simplify the notation for products of indicator functions by writing, for example, $AH_+(z)$ for $A(z)H_+(z)$, and derive that

$$\begin{aligned} (8) \quad g_{a,b}(z) - g_{a,0}(z) - b' \Delta_2(z) &= \delta_d^2 + \epsilon_s^2 + 2(\delta_s \delta_d + \epsilon_s \epsilon_d) [H_-(z) - H_+(z)] \\ &\quad + 4(x\delta_d - \delta_s \delta_d - \epsilon_s \epsilon_d) [AH_-(z) - AH_+(z)] \\ &\quad + 4(\delta_s - x + y\epsilon_d) [A^0 H_-(z) - A^0 H_+(z)]. \end{aligned}$$

Define the remainder functions $s_{a,b}(z)$ by equalities

$$\begin{aligned} \|b\| s_{a,b}(z) &= \delta_d^2 + \epsilon_s^2 + 2(\delta_s \delta_d + \epsilon_s \epsilon_d) [H_-(z) - H_+(z)] \\ &\quad + 4(x\delta_d - \delta_s \delta_d - \epsilon_s \epsilon_d) [AH_-(z) - AH_+(z)], \end{aligned}$$

and observe that Property 1 holds for the class $\{s_{a,b}\}$. Thus conditions (v) and (vi) of Theorem 3 are satisfied if the functions $l_{a,b}(z)$ are defined as the remaining part of expression (8), namely $4(\delta_s - x + y\epsilon_d) [A^0 H_-(z) - A^0 H_+(z)]$. Define $\tau_a = 1/4$ and $\tau_b = 1/2$ as

Theorem 3 prescribes. It is only left to check condition (vii) of Theorem 3 by establishing

$$(9) \quad \sup_{(a,b) \in \mathcal{N}_n} |\nu_n l_{a,b}| = o_p(n^{-1/2})$$

for all sequences of rectangles \mathcal{N}_n of the order $O(n^{-1/4}) \times O(n^{-3/8})$ that are centered at the origin. Write out the Taylor approximations $U = \delta_s + \epsilon_d + \delta_d \epsilon_d - \epsilon_s \epsilon_d + o(\|(a, b)\|^2)$ and $D = \delta_s - \epsilon_d - \delta_d \epsilon_d - \epsilon_s \epsilon_d + o(\|(a, b)\|^2)$, and conclude that quantities $|D - (\delta_s - \epsilon_d)|$ and $|U - (\delta_s + \epsilon_d)|$ are of order $O(n^{-5/8})$ uniformly over (a, b) in the neighborhoods \mathcal{N}_n . Use the oscillation properties of the empirical process established on page 765 in Shorack and Wellner (1986) to conclude that

$$\sup_{(a,b) \in \mathcal{N}_n} |\nu_n AH_-(z) - \nu_n AH_+(z)| = o_p(n^{-1/4}).$$

Stochastic bound (9) follows directly.

Apply Theorem 3 and deduce that $(n^{1/4}a_n, n^{1/2}b_n) \rightsquigarrow (s^*, t^*)$, where

$$s^* = \arg \min_s [\psi_1(s) + s'Z_1] \quad \text{and} \quad t^* = \arg \min_t [\psi_2(t) + t'Z_2 + \sum_{i=1}^3 \phi_i(s^*, t)].$$

A closed form expression for (s^*, t^*) is given in Section 2.3 of Radchenko (2004).

8. EXAMPLE: PARTIAL SPLINES

The following semiparametric example is discussed in Van de Geer (1999, chapter 11), where a CLT is established for the parametric component using its characterization as a zero of the derivative of the criterion function. Below, the same result is derived by working directly with the definition of the estimator as a minimizer, using the approach introduced in Sections 2 and 3 for mixed-rates parametric problems.

Let $(Y_1, Z_1), \dots, (Y_n, Z_n), \dots$ be independent copies of (Y, Z) , where Y is a real-valued response variable and Z is a covariate. Suppose, for simplicity, that Z takes values in $[0, 1]^2$, write $Z = (U, V)$ and assume that the model

$$Y = g(U, V) + W,$$

is satisfied with $E(W|Z) = 0$ and $g(U, V) = \theta U + \gamma(V)$. Here $\theta \in \mathbb{R}$ is an unknown parameter, and γ is an unknown member of the functional class

$$\mathcal{S} = \left\{ \eta : [0, 1] \rightarrow \mathbb{R}, \int_0^1 |\eta^{(m)}(v)|^2 dv < \infty \right\},$$

defined for a fixed positive integer m . Assume that the tails of the error distribution decrease exponentially fast: there exist positive constants K and σ_0^2 , such that for all $z \in [0, 1]^2$,

$$2K^2 E\left(e^{|W|/K} - 1 - |W|/K \mid Z = z\right) \leq \sigma_0^2.$$

Denote the distribution of (U, V) by Q and write $\|f\|_2$ for the $L_2(Q)$ -norm of a function f . Define functions $e(v) = E(U \mid V = v)$ and $h(u, v) = u - e(v)$. Assume that $\|h\|_2 > 0$.

Fix a positive λ_0 and take $\lambda_n = \lambda_0 n^{-m/(2m+1)}$. Consider a class \mathcal{F} of all regression functions f of the form $f(u, v) = \alpha u + \eta(v)$ with $\alpha \in \mathbb{R}$ and $\eta \in \mathcal{S}$. Denote the roughness of such a function by $I^2(f) = I^2(\eta) = \int_0^1 |\eta^{(m)}(v)|^2 dv$. Define

$$g_n = \arg \min_{f \in \mathcal{F}} \left\{ \frac{1}{n} \sum_{i=1}^n [Y_i - f(U_i, V_i)]^2 + \lambda_n^2 I^2(f) \right\},$$

the penalized least squares estimator of function g over the class \mathcal{F} . Assume that the regression of U on V is sufficiently smooth by requiring $I(e) < \infty$. Given a function τ from the class \mathcal{S} and a real δ , define $f_{\tau, \delta}(u, v) = [\theta + \delta]u + [\gamma(v) + \tau(v) - \delta e(v)]$ and note that function $f_{\tau, \delta}$ is a member of the class \mathcal{F} . Introduce criterion functions

$$G_n(\tau, \delta) = \frac{1}{n} \sum_{i=1}^n [Y_i - f_{\tau, \delta}(U_i, V_i)]^2 + \lambda_n^2 I^2(f_{\tau, \delta}).$$

Write $g_n(u, v)$ as $[\theta + \delta_n]u + [\gamma(v) + \tau_n(v) - \delta_n e(v)]$ and observe that the pair (τ_n, δ_n) minimizes G_n over the class $\{(\tau, \delta) : \tau \in \mathcal{S}, \delta \in \mathbb{R}\}$.

Methods from penalized least-squares estimation establish the common rate of convergence for the two components of the estimator (τ_n, δ_n) . Define $r = m/(2m+1)$. Stochastic bound $\|g_n - g\|_2 = O_p(n^{-r})$ is derived in Lemma 11.1 of Van de Geer (1999). Note that $\|g_n - g\|_2^2 = \delta_n^2 \|h\|_2^2 + \|\tau_n\|_2^2$, because conditional expectation $E(h(U, V) \mid V = v)$ is zero for each v in $[0, 1]$. Conclude that $\delta_n = O_p(n^{-r})$ and $\|\tau_n\|_2 = O_p(n^{-r})$.

Apply the approach of Section 3 to improve the convergence rate of δ_n . Write X_n for the standardized sum $n^{-1/2} \sum_{i=1}^n h(U_i, V_i)W_i$ and deduce that

$$(10) \quad G_n(\tau, \delta) - G_n(\tau, 0) = \delta^2 Q_n h^2 - 2\delta[n^{-1/2}X_n - Q_n h\tau] + \lambda_n^2[I^2(f_{\tau, \delta}) - I^2(f_{\tau, 0})].$$

Equality $Eh(U, V)\tau_n(V) = 0$ implies $Q_n h\tau_n = n^{-1/2}\nu_n h\tau_n$, and asymptotic equicontinuity of the empirical process indexed by functions $\{h\tau : \tau \in \mathcal{S}\}$ yields $Q_n h\tau_n = o_p(n^{-1/2})$. Use the definition of the roughness to derive $|I^2(f_{\tau_n, \delta}) - I^2(f_{\tau_n, 0})| \leq I(\delta e) I(2\gamma + 2\tau_n - \delta e)$. Note the stochastic bound $I(\tau_n) = O_p(1)$, implied by Van de Geer's Lemma 11.1, and conclude that $\lambda_n^2[I^2(f_{\tau_n, \delta}) - I^2(f_{\tau_n, 0})] = o_p(n^{-1/2}\delta)$. Expression (10) evaluated at $\tau = \tau_n$ and $\delta = \delta_n$ simplifies to

$$G_n(\tau_n, \delta_n) - G_n(\tau_n, 0) = \delta_n^2 Q_n h^2 - 2\delta_n n^{-1/2}[X_n + o_p(1)].$$

The law of large numbers yields $Q_n h^2 \rightarrow \|h\|_2^2$, and the limit is positive by assumption. Observe that $X_n = O_p(1)$ and apply Theorem 2, with $\delta^2 Q_n h^2$ playing the role of $M_n(\tau, \delta)$, to derive the correct $n^{-1/2}$ convergence rate of δ_n .

Note that δ_n minimizes the criterion function $G_n(\tau_n, \delta) - G_n(\tau_n, 0)$ over δ . Localize this function by writing $\delta = n^{-1/2}t$, and use the results of the previous paragraph to derive a quadratic approximation that holds uniformly on compacta,

$$(11) \quad G_n(\tau_n, n^{-1/2}t) - G_n(\tau_n, 0) = n^{-1} \left[t^2 Q_n h^2 - 2tX_n + o_p(1) \right].$$

Define $\sigma^2(z) = E(W^2|Z = z)$ and note that $X_n \rightsquigarrow X$, where $X \sim N(0, \|\sigma h\|_2^2)$. Minimization of the random quadratic function in (11) yields $n^{1/2}\delta_n = X_n/Q_n h^2 + o_p(1)$, and a CLT for δ_n follows directly. Note that because the criterion functions $G_n(\tau_n, \cdot)$ are convex, the formal derivation of the bound $\delta_n = O_p(n^{-1/2})$ could have been sidestepped.

9. PROOFS

9.1. Proof of Theorem 1. The next result is a version of the continuous mapping theorem.

Lemma 2 (Modified continuous mapping). *Consider a metric space (\mathcal{X}, d) . Let random maps $X_n : A_n \rightarrow \mathcal{X}$ be defined on some sets $A_n \subset \Omega$ and consider a function $g : \mathcal{X} \rightarrow \mathbb{R}^d$*

that is continuous at every point of a set $\mathcal{X}_0 \subset \mathcal{X}$. Suppose that $X : \Omega \rightarrow \mathcal{X}$ is a Borel measurable map for which there exists a Borel measurable set A , containing each of the sets A_n , such that $X \in \mathcal{X}_0$ on A . Suppose that $\mathbb{P}^*\{d(X_n, X) > \epsilon\} \cap A_n \rightarrow 0$ for all $\epsilon > 0$. Then $\mathbb{P}^*\{\|g(X_n) - g(X)\| > \delta\} \cap A_n \rightarrow 0$ for all $\delta > 0$.

Proof. Apply a standard device for proving continuous mapping theorems (see, for example, the proof of Theorem 1.9.5 in van der Vaart and Wellner 1996). Fix a positive ϵ . Let D_k be the set of all x in \mathcal{X} for which there exist y and z within the open ball of radius $1/k$ around x with $\|g(y) - g(z)\| > \delta$. Note that D_k is open and the sequence D_k is decreasing. Also note that $\mathbb{P}\{X \in D_k\} \cap A \downarrow 0$, because every point in $\cap_{k=1}^{\infty} D_k$ is a point in \mathcal{X}_0^c . Observe that for every fixed k ,

$$\mathbb{P}^*\{\|g(X_n) - g(X)\| > \delta\} \cap A_n \leq \mathbb{P}\{X \in D_k\} \cap A + \mathbb{P}^*\{d(X_n, X) > 1/k\} \cap A_n.$$

The first term on the right hand side can be made arbitrarily small by choosing k large enough. For a given choice of k , the second term tends to zero as n goes to infinity. \square

Dudley (1985) proved a representation theorem for the convergence in distribution in the sense of Hoffmann-Jørgensen. The following argument uses Dudley's result in the convenient form of Theorem 9.4 in Pollard (1990), referred to as Representation Theorem.

Proof of Theorem 1. It is enough to show that $\mathbb{P}^*h(s_n, t_n) \rightarrow \mathbb{P}^*h(s^*, t^*)$ for all bounded, uniformly continuous, real functions h on $\mathbb{R}^{d_1} \times \mathbb{R}^{d_2}$. Invoke Representation Theorem for the convergence $(f_n, g_n) \rightsquigarrow (f, g)$, denote the corresponding perfect maps by ϕ_n and write $\tilde{\omega}$ for the elements of the new probability space. Simplify the notation by replacing the composition $f_n(\phi_n(\tilde{\omega}), s)$ with $\tilde{f}_n(s)$, writing \tilde{s}_n for $s_n(\phi_n(\tilde{\omega}))$, and so on, omitting the $\tilde{\omega}$. Perfectness of ϕ_n implies $|\mathbb{P}^*h(s_n, t_n) - \mathbb{P}^*h(s^*, t^*)| \leq \tilde{\mathbb{P}}^*|h(\tilde{s}_n, \tilde{t}_n) - h(\tilde{s}^*, \tilde{t}^*)|$, hence it is enough to show that random points $(\tilde{s}_n, \tilde{t}_n)$ converge to $(\tilde{s}^*, \tilde{t}^*)$ in outer probability (see, for example, Theorem 1.9.5 in van der Vaart and Wellner). Write A_n^r for the set $\{\tilde{\omega} : \|\tilde{s}_n\| \vee \|\tilde{t}_n\| \vee \|\tilde{s}^*\| \vee \|\tilde{t}^*\| \leq r\}$. Because the quantities $\|\tilde{s}_n\|$, $\|\tilde{t}_n\|$, $\|\tilde{s}^*\|$ and $\|\tilde{t}^*\|$

are stochastically bounded with respect to \tilde{P}^* , it is sufficient to check that for each fixed r ,

$$(12) \quad \tilde{\mathbb{P}}^* \{ \|\tilde{s}_n - \tilde{s}^*\| > \delta \} \cap A_n^r \rightarrow 0 \quad \text{and} \quad \tilde{\mathbb{P}}^* \{ \|\tilde{t}_n - \tilde{t}^*\| > \delta \} \cap A_n^r \rightarrow 0 \quad \text{for all } \delta.$$

Fix a positive r and simplify the notation by agreeing that all functions are now restricted to the closed balls of radius r that are centered at the origin. On the set A_n^r , the points \tilde{s}_n and \tilde{s}^* are two values of the same map: $\tilde{s}_n = \Psi_r[\alpha_n^{-1}\tilde{H}_n(\cdot, \tilde{t}_n)]$ and $\tilde{s}^* = \Psi_r[\tilde{f}]$. The two corresponding arguments are close when n is large. Indeed, on the set A_n^r ,

$$\sup_{\|s\| \leq r} |\alpha_n^{-1}\tilde{H}_n(s, \tilde{t}_n) - \tilde{f}(s)| \leq \sup_{\|s\| \leq r} |\tilde{f}_n(s) - \tilde{f}(s)| + \beta_n/\alpha_n \sup_{\|s\| \wedge \|t\| \leq r} |\tilde{g}_n(s, t)|.$$

The right hand side of the above inequality goes to zero in outer probability because of the bound from Representation Theorem and the boundedness of $g(s, t)$. Let A^r stand for the Borel measurable set $\{\tilde{\omega} : \|\tilde{s}^*\| \leq r\}$ and observe that on the set A^r the map Ψ_r is continuous at \tilde{f} . Apply the modified continuous mapping Lemma with A_n^r playing the role of A_n and A^r playing the role of A , and deduce the first convergence in display (12).

Note that $\tilde{t}_n = \Phi_r[\tilde{g}_n(\tilde{s}_n, \cdot)]$ and $\tilde{t}^* = \Phi_r[\tilde{g}(\tilde{s}^*, \cdot)]$ on the set A_n^r . Also on this set,

$$\sup_{\|t\| \leq r} |\tilde{g}_n(\tilde{s}_n, t) - \tilde{g}(\tilde{s}^*, t)| \leq \sup_{\|s\| \wedge \|t\| \leq r} |\tilde{g}_n(s, t) - \tilde{g}(s, t)| + \sup_{\|t\| \leq r} |\tilde{g}(\tilde{s}_n, t) - \tilde{g}(\tilde{s}^*, t)|.$$

The first term on the right hand side tends to zero in outer probability because of the bound from Representation Theorem; the second term tends to zero in outer probability by the standard continuous mapping theorem. Deduce the second convergence in display (12) using an argument analogous to the one concluding the previous paragraph. \square

9.2. Proofs of the results in Section 3. The following lemma simplifies the work with random polynomial functions.

Lemma 3. *Let α be positive and let $\{\gamma_1, \dots, \gamma_p, \eta_1, \dots, \eta_p\}$ be a collection of nonnegative numbers satisfying $\gamma_i < \alpha$ for all $i \in \{1, \dots, p\}$. Define $\tau = \min_{i \leq p} (\frac{\eta_i}{\alpha - \gamma_i})$. For each positive δ and each $O_p^*(1)$ sequence of random variables L_n , there exist a $O_p^*(1)$ sequence*

of random variables M_n , such that the following upper bound holds for all positive u :

$$L_n \left(\sum_{i \leq p} n^{-\eta_i} u^{\gamma_i} \right) \leq \delta u^\alpha + M_n n^{-\alpha\tau}.$$

Proof. It is enough to establish the bound for $\delta = 1$. Let $M_n(\omega)$ be the smallest real number satisfying the inequality $\sup_{u \geq 0} (L_n(\omega) \sum_{i \leq p} n^{-\eta_i} u^{\gamma_i} - u^\alpha) \leq M_n(\omega) n^{-\alpha\tau}$. Given a positive ϵ , select a large enough L to ensure that $P^*\{L_n > L\} < \epsilon$. Note that

$$P^*\{M_n > M\} \leq P^*\left\{ \sup_{u \geq 0} \left(L \sum_{i \leq p} n^{-\eta_i} u^{\gamma_i} - u^\alpha \right) > M n^{-\alpha\tau} \right\} + \epsilon.$$

To see that the first term on the right-hand side of the above inequality is zero for all M large enough, combine the upper bound

$$\sup_{u \geq 0} \left(L \sum_{i \leq p} n^{-\eta_i} u^{\gamma_i} - u^\alpha \right) \leq \max_{i \leq k} \sup_{u \geq 0} (p L n^{-\eta_i} u^{\gamma_i} - u^\alpha)$$

with the inequalities

$$\sup_{u \geq 0} (p L n^{-\eta_i} u^{\gamma_i} - u^\alpha) = c_i n^{-\alpha\eta_i/(\alpha-\gamma_i)} \leq c_i n^{-\alpha\tau}, \quad i = 1, \dots, p.$$

Conclude that $M_n = O_p^*(1)$. □

Proof of Theorem 2 is omitted because it is similar to following argument.

Proof of Lemma 1. Deduce $\|a_n\|^\alpha + \|b_n\|^\beta = O_p^*\left(\sum_{i \leq p} n^{-\eta_i} \|(a_n, b_n)\|^{\gamma_i}\right)$ from inequality $G_n(\alpha_n, b_n) - G_n(0, 0) \leq 0$. For each positive δ , use Lemma 3 to establish

$$\|a_n\|^\alpha + \|b_n\|^\beta \leq \delta \|(a_n, b_n)\|^\alpha + O_p^*(n^{-\alpha\tau_a}).$$

Take a small enough δ and use inequality $\alpha \geq \beta$ to derive $\|a_n\|^\alpha + \|b_n\|^\beta = O_p^*(n^{-\alpha\tau_a})$.

Conclude that $\|a_n\| = O_p^*(n^{-\tau_a})$ and $\|b_n\| = O_p^*(n^{-\alpha\tau_a/\beta})$. □

9.3. Proof of Theorem 3.

Proof. According to condition (iv),

$$G_n(a, b) = G(a, b) + n^{-1/2} a' \nu_n \Delta_1 + n^{-1/2} b' \nu_n \Delta_2 + n^{-1/2} \|(a, b)\| \nu_n r_{a,b}.$$

Combine this representation with the stochastic bound $\|(a_n, b_n)\| = o_p(1)$ and deduce the approximation $G_n(a_n, b_n) = G(a_n, b_n) + O_p(n^{-1/2} \|(a_n, b_n)\|)$. Apply Lemma 1 with function G playing the role of M_n and derive the stochastic bounds $\|a_n\| = O_p(n^{-\tau_a})$ and $\|b_n\| = O_p(n^{-\alpha\tau_a/\beta})$. It follows from condition (v) that

$$G_n(a, b) - G_n(a, 0) = G(a, b) - G(a, 0) + n^{-1/2} b' \nu_n \Delta_2 + n^{-1/2} \|b\| \nu_n s_{a,b} + n^{-1/2} \nu_n l_{a,b}.$$

Conditions (vi) and (vii) yield $\nu_n s_{a_n, b_n} = o_p(1)$ and $\nu_n l_{a_n, b_n} = o_p(n^{-\beta\tau_b+1/2})$. Thus,

$$(13) \quad G_n(a_n, b_n) - G_n(a_n, 0) = G(a_n, b_n) - G(a_n, 0) + O_p(n^{-1/2} \|b_n\|) + o_p(n^{-\beta\tau_b}).$$

Observe that $\sum_{i=1}^{\alpha} \|a_n\|^{\alpha-i} \|b_n\|^i = O_p(n^{-1/2} \|b_n\|)$. Consequently,

$$G(a_n, b_n) - G(a_n, 0) = \psi_2(b_n) [1 + o_p(1)] + O_p\left(\sum_{i=1}^p n^{-\tau_a \gamma_i} \|b_n\|^{\eta_i}\right) + o_p(n^{-1/2} \|b_n\|).$$

Combine this approximation with approximation (13) and let the term $\psi_2(b_n) [1 + o_p(1)]$ play the role of $M_n(a_n, b_n)$ in Theorem 2. Conclude that $\|b_n\| = O_p(n^{-\tau_b})$.

Introduce new variables s and t by $s = n^{\tau_a} a$ and $t = n^{\tau_b} b$. Observe that

$$\begin{aligned} \sum_{i=1}^p \phi_i(n^{-\tau_a} s, n^{-\tau_b} t) &= n^{-\beta\tau_b} \sum_{i=1}^p n^{-(\beta-\eta_i)[\lambda_i-\tau_b]} \phi_i(s, t) \quad \text{and} \\ \sum_{i=1}^{\alpha} \|n^{-\tau_a} s\|^{\alpha-i} \|n^{-\tau_b} t\|^i &\leq n^{-\tau_a(\alpha-1)-\tau_b} \sum_{i=1}^{\alpha} \|s\|^{\alpha-i} \|t\|^i \\ &= n^{-\beta\tau_b} n^{-(\beta-1)[\lambda_0-\tau_b]} \sum_{i=1}^{\alpha} \|s\|^{\alpha-i} \|t\|^i. \end{aligned}$$

Combine the last two displays with the approximations in conditions (ii) and (iii) to deduce

$$(14) \quad \begin{aligned} G(n^{-\tau_a} s, n^{-\tau_b} t) &= n^{-\alpha\tau_a} [\psi_1(s) + q_n(s)] + \\ &\quad n^{-\beta\tau_b} [\psi_2(t) + \phi_i(s, t) 1\{\lambda_i = \tau_b\} + w_n(s, t)], \end{aligned}$$

where $\sup_{K_1} |q_n(s)| = o(1)$ for every compact set K_1 in \mathbb{R}^{d_1} , and $\sup_{K_2} |w_n(s, t)| = o(1)$ for every compact set K_2 in $\mathbb{R}^{d_1} \times \mathbb{R}^{d_2}$. Conditions (iv) through (vii) yield

$$(15) \quad \begin{aligned} G_n(n^{-\tau_a} s, n^{-\tau_b} t) &= G(n^{-\tau_a} s, n^{-\tau_b} t) + n^{-\alpha\tau_a} [s' \nu_n \Delta_1 + q'_n(s)] \\ &\quad + n^{-\beta\tau_b} [n^{-(\beta-1)[\lambda_0-\tau_b]} t' \nu_n \Delta_2 + w'_n(s, t)], \end{aligned}$$

where $\sup_{K_1} |q'_n(s)| = o_p(1)$ for every compact set K_1 in \mathbb{R}^{d_1} , and $\sup_{K_2} |w'_n(s, t)| = o_p(1)$ for every compact set K_2 in $\mathbb{R}^{d_1} \times \mathbb{R}^{d_2}$.

Denote $G_n(n^{-\tau_a} s, n^{-\tau_b} t)$ by $H_n(s, t)$. Combine approximations (14) and (15) and conclude that uniformly on compacta in $\mathbb{R}^{d_1} \times \mathbb{R}^{d_2}$,

$$\begin{aligned} H_n(s, t) &= n^{-\alpha\tau_a} [\psi_1(s) + s' \nu_n \Delta_1 + o_p(1)] + \\ &\quad n^{-\beta\tau_b} [\psi_2(t) + 1\{\lambda_0 = \tau_b\} t' \nu_n \Delta_2 + \sum_{i=1}^p 1\{\lambda_i = \tau_b\} \phi_i(s, t) + o_p(1)]. \end{aligned}$$

Note that $\alpha\tau_a \leq \beta\tau_b$. Indeed, this inequality is valid in the case $\lambda_0 = \tau_b$; in the case $\lambda_0 \neq \tau_b$, it follows from approximation (14) and the fact that function G assumes only nonnegative values near the origin. If $\alpha\tau_a < \beta\tau_b$, apply Corollary 1 to complete the proof. If $\alpha\tau_a = \beta\tau_b$, the standard arg min theorem will suffice. \square

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