Influence functions and functional derivatives

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1 Motivation

Many statistical problems are naturally formulated using one or more nuisance parameters; that is, we need to introduce some components in the statistical model, which are not of interest in themselves, but which nevertheless are needed to model the question of interest. A good example of this is the ATE ...

introduce briefly

Statistical problems involving a nuisance parameter often lead to an obvious two-step estimation strategy: 1) Estimate the nuisance parameter, and then 2) plug this estimate into the expression for the target parameter.

Example 1.1

Consider the following simple toy example: Given n samples $X_i \in \mathbb{R}$ from some unknown distribution with cumulative distribution function F, we want to estimate $F(x) = P(X \le x)$. Let us say that we are willing to assume that F has a continuous Lebesgue-density f. Then one estimation strategy would be to first use a kernel density estimator to estimate f, and then plug this into the integral operator

$$f \longmapsto \int_{-\infty}^{x} f(z) \, \mathrm{d}z,$$

to obtain an estimate of the cumulative distribution function F(x) at the fixed point $x \in \mathbb{R}$. In this setting our nuisance parameter is f, and our target parameter is $\theta = \Psi(f)$ with

$$\Psi \colon \mathcal{F} \to \mathbb{R}, \quad \Psi(f) = \int_{-\infty}^{x} f(z) \, \mathrm{d}z,$$

where \mathcal{F} is some suitable function space, for instance the collection of continuous functions. This procedure results in the target and nuisance estimators given as

$$\hat{\theta}_n := \int_{-\infty}^x \hat{f}_n(z) \, \mathrm{d}z, \quad \text{and} \quad \hat{f}_n(z) = \hat{\mathbb{P}}_n[k_h(z, \cdot)], \tag{1}$$

for some kernel function k_h with bandwidth h_n . It is well-known that the optimal choice of bandwidth h_n is $h_n \propto n^{-1/5}$ [Wasserman, 2006], so this would also be the natural choice in our case; indeed, the upper panel of figure 1 demonstrates how this choice of bandwidth is superior to the choice $h_n \propto n^{-1/2}$, which instead results in a very rough or undersmoothed estimate. Surprisingly, however, the lower panel shows that for estimation of the target parameter, plugging the undersmoothed estimate into Ψ is superior to using the default, optimal bandwidth estimator. This example is in fact simply enough to allow an exact analytic calculation of the bias and variance of the target parameter, and hence it is fairly straightforward to mathematically prove the behavior suggested by figure 1.

Add explicit bias variance decomposition?

Though the above example is somewhat silly because we already have an obvious estimator of the parameter of interest in this situation, it clearly illustrates that the modus operandi outlined in the two-step estimation procedure can be problematic. With this example in mind, we should not be too confident about our ATE estimator – in fact, with a bit more work it can be showed that similar phenomenon happens if we use the first estimator suggested above.

todo

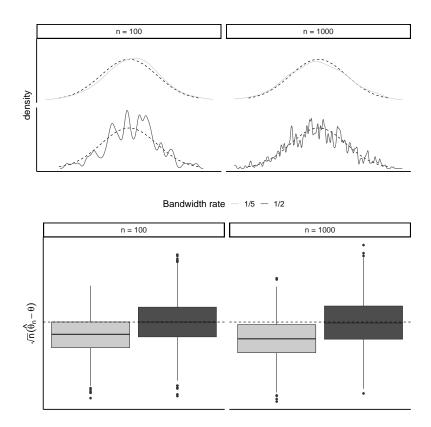


Figure 1: Simulation illustrating the problem for the plug-in approach. The top 2 rows give representative kernel density estimates for two different bandwidths scaling with n; the dashed line is the density of the distribution used to generate the data. The last row gives the distribution of the centralized and \sqrt{n} -scaled corresponding plug-in estimates based on 1000 Monte Carlo samples.

The issue from example 1.1 can be understood more generally by considering the decomposition

$$\begin{split} \sqrt{n} \left(\hat{\theta}_n - \theta \right) &= \sqrt{n} \left(\Psi(\hat{\mathbb{P}}_n, \hat{\nu}_n) - \Psi(\mathbf{P}, \nu) \right) \\ &= \mathbb{G}_n[\varphi(\cdot, \hat{\nu}_n)] + \sqrt{n} \left\{ \Psi(\mathbf{P}, \hat{\nu}_n) - \Psi(\mathbf{P}, \nu) \right\} \\ &= \mathbb{G}_n[\varphi(\cdot, \hat{\nu}_n)] + \sqrt{n} \left\{ \mathbf{D}_{\nu} \Psi(\hat{\nu}_n - \nu) + \mathcal{O}_{\mathbf{P}}(\|\hat{\nu}_n - \nu\|^2) \right\}. \end{split}$$

Here we use $D_{\nu}\Psi$ to denote some kind of derivate (so be defined shortly) of the map $\nu \mapsto \Psi(P,\nu)$, where P is held fixed, and $\|\cdot\|$ is a norm defined on the space in which the nuisance estimator $\hat{\nu}_n$ takes it values – typically a function space. Using empirical process theory or sample spitting, it can in many cases be shown that $\mathbb{G}_n[\varphi(\cdot,\hat{\nu}_n)] = \mathbb{G}_n[\varphi(\cdot,\nu)] + \mathcal{O}_P(1)$ [van der Vaart and Wellner, 1996, van der Vaart, 2000, Chernozhukov et al., 2018]. Furthermore, if we assume that $\hat{\nu}_n$ can be estimated at rate $n^{-1/4}$, meaning that $\|\hat{\nu}_n - \nu\| = \mathcal{O}_P(n^{-1/4})$, the decomposition becomes

$$\sqrt{n}\left(\hat{\theta}_n - \theta\right) = \mathbb{G}_n[\varphi(\cdot, \nu)] + D_{\nu}\Psi\left(\sqrt{n}(\hat{\nu}_n - \nu)\right) + \mathcal{O}_{P}(1), \tag{2}$$

where we use that a derivative $\dot{\Psi}_{\nu}$ should be linear. The first term of this expression is controlled by reference to the central limit theorem, but unless we are able to estimate ν at the improved parametric rate of $n^{-1/2}$, we cannot hope to control the second term. Indeed, the fact that the default kernel estimator from example 1.1 converges at rate $n^{-2/5}$ is what

check

ruins the asymptotic behavior of the target estimator. In the example we showed that the target estimator could be improved by using a suitably undersmoothed density estimator. However, we would not in general know how to choose a properly undersmoothed nuisance estimator, so this strategy is not applicable in practice (though see ...). If we want to use flexible machine learning estimators (which cannot be expected to converge at parametric rate) a better strategy is to design the map to Ψ such that its partial derivate $D_{\nu}\Psi$ equals 0 at ν . This would make the second term in (2) vanish, and we would be able to estimate the target parameter θ at parametric rate, while still using a nuisance estimator converging only at rate $n^{-1/4}$.

ref to undersmoothed

In the next section we formally consider how to define a functional derivative, and then indicate how this can be used to design Ψ with partial derivate equal to 0. In section 3 we show how these results in addition allow us to analyze and construct *efficient* estimators of θ , that is, estimators with lowest possible asymptotic variance.

2 Functional derivatives and geometry

The most straightforward form of functional differentiability is the generalization of the directional derivative of multivariate calculus. This is known as Gâteaux differentiability and defined as follows.

Definition 2.1 (Gâteaux derivate). Let \mathcal{M} and \mathcal{Y} be a normed real vector spaces, $\Psi \colon \mathcal{M} \to \mathcal{Y}$ a map, and $x \in \mathcal{M}$ a point in the domain. If there exists a linear, continuous operator $\dot{\Psi}_x \colon \mathcal{X} \to \mathcal{Y}$ such that for all $h \in \mathcal{X}$

$$\left\|\Psi(x+\varepsilon h)-\Psi(x)-\dot{\Psi}_x(\varepsilon h)\right\|_{\mathcal{Y}}=\mathcal{O}(\varepsilon),\quad\text{when}\quad\varepsilon\longrightarrow0\in\mathbb{R},$$

we say that Ψ is Gâteaux differentiable at x. We call the operator $\dot{\Psi}_x$ the Gâteaux derivative of Ψ at x.

If the map Ψ is Gâteaux differentiable at x then it follows from the definition that

$$\left\| \frac{\Psi(x + \varepsilon h) - \Psi(x)}{\varepsilon} - \dot{\Psi}_x(h) \right\|_{\mathcal{V}} \longrightarrow 0, \quad \text{when} \quad \varepsilon \longrightarrow 0.$$
 (3)

In particular, when $\mathcal{Y} = \mathbb{R}$ the Gâteaux derivative at x in the direction h can be derived as the ordinary derivative of the real-valued function $\varepsilon \mapsto \Psi(x + \varepsilon h)$ evaluated at 0, i.e.,

$$\dot{\Psi}_x(h) = \partial_0 \Psi(x + \varepsilon h) := \frac{\partial}{\partial \varepsilon} \bigg|_{\varepsilon = 0} \Psi(x + \varepsilon h). \tag{4}$$

To simplify notation in the following we will use the notation $\partial_0 f(\varepsilon)$ for maps f with a real domain as just defined in (4). Gâteaux differentiability is used by Chernozhukov et al. [2018] to define Neyman orthogonality: A function $\varphi \colon \mathcal{Z} \times \mathcal{V} \to \mathbb{R}$ fulfills the Neyman orthogonality condition (wrt. \mathcal{V}) if the map (as defined in section 1)

define the sample space \mathcal{Z} proper def?

$$F \colon \mathcal{V} \to \mathbb{R}, \quad \nu \longmapsto F(\nu) := P[\varphi(Z, \nu)]$$

is Gâteaux differentiable at ν with vanishing derivative, i.e., $\dot{F}_{\nu}(h) = 0$ for all $h \in \mathcal{V}$.

Neyman orthogonality ensures that the second component of the decomposition in (2) vanishes, and as the condition can be checked using (4) it is rather straightforward to verify that a given function fulfills the condition. Still, Gâteaux is a weak form of differentiability; for instance, as it is equivalent to the directional derivative for ordinary multivariate

functions, Gâteaux differentiability is not enough to guarantee (ordinary) differentiability of such functions. Similarly, to get a richer theory in the functional setting a stronger notion of differentiability is needed. For our setting, a particular useful concept is *Hadamard* differentiability.

Definition 2.2 (Hadamard derivate). Let \mathcal{M} and \mathcal{Y} be a normed real vector spaces, $\Psi \colon \mathcal{M} \to \mathcal{Y}$ a map, and $x \in \mathcal{M}$ a point in the domain. If there exists a linear, continuous operator $\dot{\Psi}_x \colon \mathcal{X} \to \mathcal{Y}$ such that

$$\left\| \frac{\Psi(x + \varepsilon_n h_n) - \Psi(x)}{\varepsilon_n} - \dot{\Psi}_x(h) \right\|_{\mathcal{Y}} \longrightarrow 0,$$

for any $\varepsilon_n \to 0$ and $\{h_n\}_{n \in \mathbb{N}} \subset \mathcal{X}$ with $h_n \to h \in \mathcal{X}$, we say that Ψ is Hadamard differentiable at x and call the operator $\dot{\Psi}_x$ the Hadamard derivative of Ψ at x.

Comparing this definition with (3) we see that the only difference is that the linear approximation provided by $\dot{\Psi}_x$ should holds along any converging sequence h_n and not merely in a fixed direction h. For this reason Hadamard differentiability is also known a path-wise differentiability. A still stronger condition is to demand that the approximation should hold for any bounded sequence h_n ; this gives the concept of Fréchet differentiability, which we will not use in this text. One can show that if a map is Hadamard (or Fréchet) differentiable, then the Hadamard (or Fréchet) derivative is equal to the Gâteaux derivative. Hence, to find the Hadamard derivative of a given functional Ψ , the common strategy is to first use high school math tools to calculate (4) and then verify that the obtained candidate fulfills the requirements of definition 2.2.

So far we have assumed the domain to be a linear space. When working with probability measures, this is often not the case, and hence we need one final definition, which allows the operator Ψ and its derivative $\dot{\Psi}$ to be defined only on subsets of the normed vector space \mathcal{M} . We should think of the following as mirroring differentiation of a multivariate function defined on a manifold embedded in a higher dimensional Euclidean space (for instance, a surface embedded in \mathbb{R}^3 , see picture).

Make the pictur

Definition 2.3 (Tangential Hadamard derivative). Let \mathcal{P} and $\dot{\mathcal{P}}_x$ be subsets of the normed real vector space \mathcal{M} , with $x \in \mathcal{P} \subset \mathcal{M}$. For a map $\Psi \colon \mathcal{P} \to \mathcal{Y}$, we say that Ψ is Hadamard differentiable (at x) tangential to $\dot{\mathcal{P}}_x$ if there exists a continuous, linear operator $\dot{\Psi}_x \colon \dot{\mathcal{P}}_x \to \mathcal{Y}$ such that

$$\left\| \frac{\Psi(x + \varepsilon_n h_n) - \Psi(x)}{\varepsilon_n} - \dot{\Psi}_x(h) \right\|_{\mathcal{Y}} \longrightarrow 0,$$

for any $\{h_n\} \subset \mathcal{M}$ and $\{\varepsilon_n\} \subset \mathbb{R}$ with $h_n \to h \in \dot{\mathcal{P}}_x$, $\varepsilon_n \to 0$, and $x + \varepsilon_n h_n \in \mathcal{P}$.

The only change from the previous definition is that the "path" $x + \varepsilon_n h_n$ is restricted to lie in the subset \mathcal{P} , and that the "direction" is h is restricted to lie in $\dot{\mathcal{P}}_x$.

Finally, to be able to talk about differentiability of the statistical problem (\mathcal{P}, Ψ) , we need to embed the model \mathcal{P} into a suitable normed vector space. To do so we now assume for simplicity that the family \mathcal{P} is dominated by a single σ -finite measure μ (for our running example of the ATE this measure would be a product of Lebesgue and counting measures), and then think of \mathcal{P} as lying inside the Banach space \mathcal{M}_{μ} of finite signed measures dominated by μ equipped with the variational norm

define this earlier

$$||M||_{\mathcal{M}_{\mu}} = \int |m| \,\mathrm{d}\mu, \quad \text{ for } \quad M = m \cdot \mu.$$

Definition 2.4 (Tangent space for \mathcal{P}). Let $\mathcal{P} \subset \mathcal{M}_{\mu}$ be a collection of probability measures and let $P \in \mathcal{P}$. For any one-dimensional path $\varepsilon \mapsto P_{\varepsilon} \in \mathcal{P}$ with $P_0 = P$, which is Hadamard differentiable at 0, let $\partial_0 P_{\varepsilon}$ be the derivative, and let $\{\partial_0 P_{\varepsilon}\}$ be the collection of derivatives of all such paths. We call the closed linear span of this collection the tangent space of \mathcal{P} at P and denote it by \mathcal{P}_P . Formally,

 $\dot{\mathcal{P}}_{P} := \overline{\operatorname{span}} \left\{ \partial_{0} P_{\varepsilon} \mid \varepsilon \mapsto P_{\varepsilon} \text{ is Hadamard differentiable and } P_{0} = P \right\}.$

should there be added some note about this – always sensible thing, or need to assume something?

The definition of a tangent space for a collection of probability measures simple mimics the definition of a tangent space for a surface embedded in \mathbb{R}^3 : We move along differentiable paths through a given point on the surface, and the tangent space is then the span of the derivatives of all such paths. With a differentiable structure on the collection \mathcal{P} we can talk about a gradient of a functional defined on this set.

Definition 2.5 (Canonical gradient). Let (\mathcal{P}, Ψ) be a statistical problem, with $\mathcal{P} \subset \mathcal{M}_{\mu}$, and $\dot{\mathcal{P}}_{P}$ the tangent space of \mathcal{P} at $P \in \mathcal{P}$. If $\Psi \colon \mathcal{P} \to \mathbb{R}$ is Hadamard differentiable at P tangential to $\dot{\mathcal{P}}_{P}$, we refer to the Hadamard derivative $\dot{\Psi}_{P}$ as the canonical gradient of the statistical problem.

The definitions of the tangent space $\dot{\mathcal{P}}_{P}$ and the gradient $\dot{\Psi}_{P}$ above capture the intuitive meaning of these concepts as generalizations of well-known concepts from multivariate calculus: The statistical model \mathcal{P} is viewed as a subset of the unit sphere in \mathcal{M}_{μ} and the canonical gradient $\dot{\Psi}_{P}$ is simply the infinite-dimensional version of the gradient of a map defined on a surface in Euclidean space. We shall see in a moment why this object $\dot{\Psi}_{P}$ is of interest for the statistician, but firstly we consider a different representation of \mathcal{P} which allows us to represent $\dot{\mathcal{P}}_{P}$ as a subset of \mathcal{L}_{P}^{2} ; this is particularly useful as we then have the rich Hilbert space structure of \mathcal{L}_{P}^{2} at our disposal.

For a fixed element $P \in \mathcal{P}$ with μ -density p, consider a one-dimensional parametric submodel of μ -densities p_{ε} with $p_0 = p$ and $p_{\varepsilon} \cdot \mu \in \mathcal{P}$ for all ε . Restricting attention to such submodels for which the function

$$\dot{\ell}_0 = \partial_0 \log(p_{\varepsilon})$$

exists as an element in $\mathcal{L}_{\mathrm{P}}^2$, we call $\dot{\ell}_0$ the *score* of the submodel p_{ε} . This reflects the terminology used in likelihood inference for ordinary parametric models.

Proposition 2.6. The tangent space $\dot{\mathcal{P}}_P$ can be represented as the closed linear span of all score functions as defined above, considered as a subset of \mathcal{L}_P^2 , which we denote as $\Gamma_P := \overline{\operatorname{span}}\{\dot{\ell}_0\} \subset \mathcal{L}_P^2$.

Proof. Something like this: We can also represent \mathcal{P} as a subset of \mathcal{L}^2_{μ} through the map $P \mapsto \sqrt{p}$, and the topology on \mathcal{P} induced in this way is the same as that induced by $\|\cdot\|_{\mathcal{M}_{\mu}}$ [Bickel et al., 1993]. This implies that convergence of $(P_{\varepsilon h} - P_0)\varepsilon^{-1}$ in \mathcal{M}_{μ} is equivalent to convergence of convergence of $(\sqrt{p_{\varepsilon h}} - \sqrt{p_0})\varepsilon^{-1}$ in \mathcal{L}^2_{μ} . Then Gâteaux + map with $/p_0$...

Proposition 2.6 in essence states that we can think of the tangent space as the (closed linear span of) the collection of score functions for all parametric submodels $\mathcal{P}_{\varepsilon}$ passing through P; hence we shall often simply identify $\dot{\mathcal{P}}_{P}$ with Γ_{P} . This representation in turn implies the following useful representation of the canonical gradient.

interval...

right words to use? More precise...?

todo/finish.
Pretty sure this
results is
correct...
check – also the
generalized
versions on the
whole space?
finish

 $^{^{1}}$ When the domain is \mathbb{R} , all the previously considered types of differentiability are equivalent, so any type could be used here.

Proposition 2.7. Let (\mathcal{P}, Ψ) be a statistical problem with canonical gradient $\dot{\Psi}_{P}$ at $P \in \mathcal{P}$. There exists a unique element $\varphi_{P} \in \Gamma_{P}$ such that

$$\partial_0 \Psi(\mathbf{P}_{\varepsilon}) = \langle \varphi_{\mathbf{P}}, \dot{\ell}_0 \rangle_{\mathbf{P}} \tag{5}$$

holds for any differentiable submodel P_{ε} with score function $\dot{\ell}_0$.

Proof. The chain rule for Hadamard derivative implies that $\partial_0 \Psi(P_{\varepsilon}) = \dot{\Psi}_P(\partial_0 P_{\varepsilon})$, and then the representation given by proposition 2.6 implies that this expression equals $\Phi_P(\dot{\ell}_0)$ for some continuous linear functional $\Phi_P \colon \Gamma_P \to \mathbb{R}$. As Γ_P is a closed subspace of a Hilbert space it is itself a Hilbert space, and hence Riesz representation theorem (for Hilbert spaces) gives the existence of a unique element $\varphi_P \in \Gamma_P$ such that $\Phi_P(\dot{\ell}_0) = \langle \varphi_P, \dot{\ell}_0 \rangle$ for all elements $\dot{\ell}_0 \in \Gamma_P$.

ntroduce or

By the chain rule of ordinary multivariate calculus we have for a differentiable function $f: \mathbb{R}^d \to \mathbb{R}$ and any smooth curve $s: \mathbb{R} \to \mathbb{R}^d$ with $s(0) = x_0 \in \mathbb{R}^d$ that

$$\partial_0 (f \circ s) = \nabla f(s(0)) \cdot (\partial_0 s)^\top = \langle \nabla f(x_0), \partial_t s \rangle,$$

and as \mathbb{R}^d can be spanned by smooth 1-dimensional curves, this property characterizes the gradient locally. Proposition 2.7 above states that this characterization extends to the Hilbert space setting with Hadamard differentiability, and in fact this characterization can be used to define Hadamard differentiability of the map Ψ tangential to Γ_P [Bickel et al., 1993, A.5].

I think this is

As with the identification $\dot{\mathcal{P}}_{P} = \Gamma_{P}$ we also identify the function φ_{P} with the canonical gradient $\dot{\Psi}_{P}$. Note that if Γ_{P} is a proper subset of \mathcal{L}_{P}^{2} there can be several functions $\tilde{\varphi} \in \mathcal{L}_{P}^{2}$ fulfilling condition (5); we refer to such functions as *gradients*. It follows from standard Hilbert space theory that the unique canonical gradient φ_{P} can be derived from any gradient $\tilde{\varphi}$ as the projection onto Γ_{P} , i.e., $\varphi_{P} = \Pi(\tilde{\varphi} \mid \Gamma_{P})$. This geometric representation of the tangent space and the gradient is very useful.

Returning to the decomposition in (2) we have the following results.

Proposition 2.8. Let (\mathcal{P}, Ψ) be a statistical problem with \mathcal{P} dominated by μ and such that $\Psi(P) = P[\varphi(Z, \nu(P))]$ for some function $\varphi \colon \mathcal{Z} \times \mathcal{V} \to \mathbb{R}$. If $\varphi(\cdot, \nu(P))$ is the canonical gradient of (\mathcal{P}, Ψ) then, under regularity conditions, φ fulfills the Neyman orthogonality condition (wrt. $\dot{\nu}_P(\dot{\mathcal{P}}_P)$), i.e., $\partial_0 P[\varphi(\cdot, \nu(P) + \varepsilon h)] = 0$ for all $h \in \dot{\nu}_P(\dot{\mathcal{P}}_P)$.

can these be made precise? ok? Or should it actually be $\nu(\mathcal{P})$? some unaligned

Proof. Under regularity conditions that allow us to interchange integration and differentation in the following, we have for any differentiable path $\{P_{\varepsilon}\}_{\varepsilon} \subset \mathcal{P}$ with $P_0 = P$ that

$$\partial_{0} \Psi(\mathbf{P}_{\varepsilon}) = \partial_{0} \int p_{\varepsilon}(z) \varphi(z, \nu(\mathbf{P}_{\varepsilon})) \, \mathrm{d}\mu(z)$$

$$= \int p\{\partial_{0} \varphi(z, \nu(\mathbf{P}_{\varepsilon}))\} + \{\partial_{0} p_{\varepsilon}(z)\} \varphi(z, \nu(\mathbf{P})) \, \mathrm{d}\mu(z)$$

$$= \partial_{0} \mathbf{P}[\varphi(\cdot, \nu(\mathbf{P}_{\varepsilon}))] + \langle \dot{\ell}_{0}, \varphi(\cdot, \nu(\mathbf{P})) \rangle_{\mathbf{P}},$$
(6)

where ℓ_0 is the score function of the parametric submodel P_{ε} and we used the relation $\partial_0 \log p_{\varepsilon} = (\partial_0 p_{\varepsilon}) p_0^{-1}$. As $\varphi(\cdot, \nu(P))$ is the canonical gradient at P, proposition 2.7 and (6) imply that $\partial_0 P[\varphi(\cdot, \nu(P_{\varepsilon}))] = 0$. Assuming Hadamard differentiability of $P \mapsto \nu(P)$ and $\nu \mapsto F(\nu) = P[\varphi(\cdot, \nu)]$ we can write $\partial_0 P[\varphi(\cdot, \nu(P) + \varepsilon h)] = \partial_0 \dot{F}_{\nu(P)}(h)$, for $h \in \dot{\nu}_P(\dot{\mathcal{P}}_P)$; hence (by definition of $\dot{\mathcal{P}}_P$) we can find path P_{ε} such that $\dot{\nu}_P(\partial_0 P_{\varepsilon}) = h$, and then picking this path in the expression above we get $\partial_0 \dot{F}_{\nu(P)}(h) = 0$.

finish the argument with something like this

The result shows that estimators based on the canonical gradient has a "build-in" so-called debiasing mechanism because the first order bias (the second expression in (2)), due to estimation of the nuisance parameter $\hat{\nu}_n$, vanishes. This debiasing mechanism is crucial for $n^{-1/2}$ -rate inference of a target parameter in a statistical model with nuisance parameters that are not estimable at this rate themselves.

Note that the fact that a function φ fulfills the Neyman orthogonality condition does not necessarily imply that φ is the canonical gradient (see, for instance, Chernozhukov et al. [2016] for an example). This can be seen from the fact that we in the proof above only used that $\varphi(\cdot, \nu(P))$ satisfies (5), and not that $\varphi(\cdot, \nu(P)) \in \Gamma_P$; hence, any gradient for the statistical problem will fulfill the Neyman orthogonality condition (under regularity conditions).

What is the intuitive explanation of this result? In standard calculus, the derivatives are orthogonal – can this intuition explain the results? I think this must be correct

3 Estimators and information bounds

The last result of the previous section demonstrated that estimators based on the canonical gradient provide target parameter estimators with no first order bias. In this section we show that they also provide *efficient* estimators.

briefly to this

- RAL estimators and IFs
- Information bounds
- Connection to canonical gradient

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