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Regularizing Priors for Linear Inverse Problems*

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Abstract

This paper proposes a new Bayesian approach for estimating, nonparametrically, functional parameters in econometric models that are characterized as the solution of a linear inverse problem. By using a Gaussian process prior distribution we propose the posterior mean as an estimator and prove frequentist consistency of the posterior distribution. The latter provides the frequentist validation of our Bayesian procedure. We show that the minimax rate of contraction of the posterior distribution can be obtained provided that either the regularity of the prior matches the regularity of the true parameter or the prior is scaled at an appropriate rate. The scaling parameter of the prior distribution plays the role of a regularization parameter. We propose a new data-driven method for optimally selecting in practice this regularization parameter. We also provide sufficient conditions so that the posterior mean, in a conjugate-Gaussian setting, is equal to a Tikhonov-type estimator in a frequentist setting. Under these conditions our data-driven method is valid for selecting the regularization parameter of the Tikhonov estimator as well. Finally, we apply our general methodology to two leading examples in econometrics: instrumental regression and functional regression estimation.

Key words: nonparametric estimation, Bayesian inverse problems, Gaussian processes, posterior consistency, data-driven method

JEL code: C13, C11, C14

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1 Introduction

In the last decade, econometric theory has shown an increasing interest in the theory of *stochastic inverse problems* as a fundamental tool for functional estimation of structural as well as reduced form models. This paper develops an encompassing Bayesian approach to nonparametrically estimate econometric models based on stochastic linear inverse problems.

We construct a Gaussian process prior for the (functional) parameter of interest and establish sufficient conditions for frequentist consistency of the corresponding posterior mean estimator. We prove that these conditions are also sufficient to guarantee that the posterior mean estimator numerically equals a Tikhonov-type estimator in the frequentist setting. We propose a novel data-driven method, based on an empirical Bayes procedure, for selecting the regularization parameter necessary to implement our Bayes estimator. We show that the value selected by our data-driven method is optimal in a minimax sense if the prior distribution is sufficiently smooth. Due to the equivalence between Bayes and Tikhonov-type estimator our data-driven method has broad applicability and allows to select the regularization parameter necessary for implementing Tikhonov-type estimators.

Stochastic linear inverse problems theory has recently gained importance in many subfields of econometrics to construct new estimation methods. Just to mention some of them, it has been shown to be fundamental in nonparametric estimation of an instrumental regression function, see e.g. Florens (2003), Newey and Powell (2003), Hall and Horowitz (2005), Blundell et~al. (2007), Darolles et~al. (2011), Florens and Simoni (2012a). It has also been used in semiparametric estimation under moment restrictions, see e.g. Carrasco and Florens (2000), Ai and Chen (2003), Chen and Pouzo (2012). In addition, it has been exploited for inference in econometric models with heterogeneity -e.g. Gautier and Kitamura (2012), Hoderlein et~al. (2013) - for inference in auction models -e.g. Florens and Sbaï (2010) - and for frontier estimation for productivity analysis -e.g. Daouia et~al. (2009). Finally, inverse problem theory has been used for functional regression estimation by e.g. Hall and Horowitz (2007) and Johannes (2008). We refer to Carrasco et~al. (2007) and references therein for a general overview of inverse problems in econometrics.

The general framework used in this paper, and that accommodates many functional estimation problems in econometrics just as mentioned above, is the following. Let \mathcal{X} and \mathcal{Y} be infinite dimensional separable Hilbert spaces over \mathbb{R} and denote by $x \in \mathcal{X}$ the functional parameter that we want to estimate. For instance, x can be an Engel curve or the probability density function of the unobserved heterogeneity. The estimating equation characterizes x as the solution of the functional equation

$$y^{\delta} = Kx + U^{\delta}, \quad x \in \mathcal{X}, \quad y^{\delta} \in \mathcal{Y}, \quad \delta > 0$$
 (1)

where y^{δ} is an observable function, $K: \mathcal{X} \to \mathcal{Y}$ is a known, bounded, linear operator and U^{δ} is an error term with values in \mathcal{Y} and covariance operator proportional to a positive scalar δ . Estimating

x is an inverse problem. The function y^{δ} is a transformation of a n-sample of finite dimensional objects and the parameter $\delta^{-1} > 0$ represents the "level of information" (or precision) of the sample, so that $\delta \to 0$ as $n \to \infty$. For $\delta = 0$ (i.e. perfect information) we have $y^0 = Kx$ and $U^0 = 0$. In many econometric models, equation (1) corresponds to a set of moment equations with y^{δ} an empirical conditional moment and $\delta = \frac{1}{n}$. A large class of econometric models write under the form of equation (1) – like moment equality models, consumption based asset pricing models, density estimation of the heterogeneity parameter in structural models, deconvolution in structural models with measurement errors. We illustrate two leading examples that can be estimated with our Bayesian method.

Example 1 (Instrumental variable (IV) regression estimation). Let (Y, Z, W) be an observable real random vector and x(Z) be the IV regression defined through the moment condition $\mathbf{E}(Y|W) = \mathbf{E}(x|W)$. Suppose that the distribution of (Z, W) is confined to the unit square $[0, 1]^2$ and admits a density f_{ZW} . The moment restriction implies that x is a solution to

$$\mathbf{E}_{W}\left[\mathbf{E}(Y|w)a(w,v)\right](v) = \int_{0}^{1} \int_{0}^{1} x(z)a(w,v)f_{ZW}(z,w)dwdz$$

where $a(w,v) \in L^2[0,1]^2$ is a known and symmetric function, \mathbf{E}_W denotes the expectation with respect to the marginal density f_W of W and $L^2[0,1]^2$ denotes the space of square integrable functions on $[0,1]^2$. This transformation of the original moment condition is appealing because in this way its empirical counterpart is asymptotically Gaussian (as required by our Bayesian approach). Assume that $x \in \mathcal{X} \equiv L^2[0,1]$. By replacing the true distribution of (Y,Z,W) with a nonparametric estimator we obtain a problem of the form (1) with $\delta = \frac{1}{n}$,

$$y^{\delta} = \hat{\mathbf{E}}_W \left[\hat{\mathbf{E}}(Y|w) a(w,v) \right] \quad \text{and} \quad Kx = \int_0^1 \int_0^1 x(z) a(w,v) \hat{f}_{ZW}(z,w) dw dz. \tag{2}$$

Example 2 (Functional Linear Regression Estimation). The model is the following:

$$\xi = \int_0^1 x(s)Z(s)ds + \varepsilon, \quad \mathbf{E}(\varepsilon Z(s)) = 0, \quad Z, x \in \mathcal{X} = L^2([0,1]), \quad \mathbf{E} < Z, Z > < \infty$$
 (3)

and $\varepsilon|Z,\tau\sim\mathcal{N}(0,\tau^2)$, with $<\cdot,\cdot>$ the inner product in $L^2[0,1]$. We want to recover the functional regression x. Assuming that Z is a centered random function with covariance operator of traceclass, the most popular approach consists in multiplying both sides of the first equation in (3) by Z(s) and then taking the expectation: $\mathbf{E}(\xi Z(t)) = \int_0^1 x(s)Cov(Z(s),Z(t))ds$, for $t\in[0,1]$. If we dispose of independent and identically distributed data $(\xi_1,Z_1),\ldots,(\xi_n,Z_n)$ we can estimate the unknown moments in the previous equation. Hence, x is solution of an equation of the form of equation (1) with $y^\delta := \frac{1}{n} \sum_i \xi_i Z_i(t)$, $U^\delta = \frac{1}{n} \sum_i \varepsilon_i Z_i(t)$, $\delta = \frac{1}{n}$ and $\forall \varphi \in L^2([0,1]) \mapsto K\varphi :=$

$$\frac{1}{n}\sum_{i} \langle Z_i, \varphi \rangle Z_i(t)$$
. \square

A Bayesian approach to stochastic inverse problems combines the prior and sampling information and proposes the posterior distribution as solution. It allows to deal with two important issues. First, it allows to incorporate in the estimation procedure the prior information about the functional parameter provided by the economic theory or the beliefs of experts. This prior information may be particularly valuable in functional estimation since often the data available are concentrated only in a region of the graph of the functional parameter so that some parts of the function can not be recovered from the data.

Second, since the quality of the estimation of the solution of the inverse problem (1) relies on the value of a *regularization parameter*, it is particularly important to choose such a parameter in an accurate way. This can be done with our Bayesian approach.

The majority of the Bayesian approaches to stochastic inverse problems proposed so far are based on a finite approximation of (1) and so, cannot be applied to functional estimation in econometrics, see e.g. Chapter 5 in Kaipio and Somersalo (2004) and Helin (2009), Lassas et al. (2009), Hofinger and Pikkarainen (2007), Hofinger and Pikkarainen (2007), Neubauer and Pikkarainen (2008). These papers consider a finite dimensional projection of (1) and recover x only on a finite grid of points. Hence, they do not work for econometric models that consider functional observations and parameters. Knapik et al. (2011) study Bayesian inverse problems with functional observations and parameter. However, their framework does not allow to accommodate the econometric models of interest because of the different definition of the error term U^{δ} . Indeed, the analysis of Knapik et al. (2011) works under the assumption that U^{δ} is an isonormal Gaussian process, which implies that U^{δ} , and by consequence y^{δ} , is not realizable as a random element in \mathcal{Y} . This assumption greatly simplifies the analysis but unfortunately does not hold, in general, in econometric models because real (functional) data cannot be generated by an isonormal Gaussian process. On the contrary, our paper works in the more realistic situation where U^{δ} , and therefore y^{δ} , are random elements with realizations in \mathcal{Y} . Thus, our approach works for econometric models.

For the special case where model (1) results from a conditional moment restricted model, Liao and Jiang (2011) proposed a quasi-Bayesian procedure. Their approach, which works also in the nonlinear case, is based on limited-information-likelihood and a sieve approximation technique. It is essentially different from our approach since we work with Gaussian process priors and do not use a finite-dimensional approximation.

Working with Gaussian process priors is computationally convenient since in many cases the sampling distribution is (asymptotically) Gaussian and thus the posterior is also Gaussian (conjugate-Gaussian setting). The current paper gives sufficient conditions under which the posterior mean of x, in a conjugate-Gaussian setting, exists in a closed-form and thus can be used to estimate x. Existence of such a closed-form is not verified in general, as we explain in the next

paragraph and after Theorem 1. Agapiou et al. (2013) propose an alternative approach to deal with this problem. Florens and Simoni (2012b) overcome this problem by constructing a regularized posterior distribution which works well in practice but the regularization of the posterior distribution is ad hoc and cannot be justified by any prior-to-posterior transformation. In comparison to Agapiou et al. (2013) and Florens and Simoni (2012b), the current paper also provides an adaptive method for choosing the regularization parameter as we explain in the next section.

Our contribution. Our estimation procedure is based on a conjugate-Gaussian setting which is suggested by the linearity of problem (1). On one hand, such a setting is appealing because the corresponding posterior distribution can be computed analytically without using any Markov Chain Monte Carlo algorithm which, even if very powerful, slows down the estimate computation. On the other hand, a conjugate-Gaussian Bayesian inverse problem has the drawback that the posterior mean is, in general, not defined as a linear estimator but as a measurable linear transformation (mlt, hereafter) which is a weaker notion, see Mandelbaum (1984). In particular, there is no explicit form for the mlt estimator and so it is unclear how we can compute the posterior mean estimator of x in practice. Moreover, whether consistency of the mlt estimator holds or not is still an open question.

The first contribution of our paper is to provide a sufficient condition under which the posterior mean, in a conjugate-Gaussian setting, is defined as a linear estimator, has a closed-form and thus can be easily computed and used as an estimator for x (as it is justified for a broad class of loss functions). We assume a Gaussian process prior distribution for x, with mean function $x_0 \in \mathcal{X}$ and covariance operator $\Omega_0: \mathcal{X} \to \mathcal{X}$. In the case where \mathcal{X} and \mathcal{Y} are finite-dimensional and (x, y^{δ}) are jointly Gaussian, the posterior mean of x is the linear estimator $[x_0 + \Omega_0 K^* Var(y^{\delta})^{-1} (y^{\delta} - Kx_0)]$ provided $Var(y^{\delta})$ is invertible, where $Var(y^{\delta})$ denotes the marginal covariance operator of y^{δ} . Unfortunately, when the dimension of \mathcal{X} and \mathcal{Y} is infinite, the linear operator $\Omega_0 K^* Var(y^{\delta})^{-1}$ is no longer defined on \mathcal{Y} but only on a dense subspace of \mathcal{Y} of measure zero and is typically non-continuous (i.e. unbounded). This paper gives a sufficient condition that guarantees that $\Omega_0 K^* Var(y^{\delta})^{-1}$ is continuous (and defined) on the whole $\mathcal Y$ and shows that this condition is in general satisfied in many econometric models. Then, we derive the closed form for the posterior mean of x that is implementable in practice and prove that it is a continuous and linear (thus consistent) estimator defined on \mathcal{Y} . Under this condition, the prior-to-posterior transformation can be interpreted as a regularization scheme so that no ad hoc regularization schemes are required as e.g. in Florens and Simoni (2012b).

Our second contribution consists in the study of frequentist asymptotic properties of the conjugate-Gaussian Bayesian estimation of equation (1). For that, we admit the existence of a true x, say x_* , that generates the data. We establish that the posterior mean estimator and pos-

terior distribution have good frequentist asymptotic properties for $\delta \to 0$. More precisely, we show that the posterior distribution converges towards a Dirac mass at x_* almost surely with respect to the sampling distribution (frequentist posterior consistency, see e.g. Diaconis and Freedman (1986, 1998)). This property provides the frequentist validation of our Bayesian procedure.

We also recover the rate of contraction of the risk associated with the posterior mean and of the posterior distribution. This rate depends on the smoothness and the scale of the prior as well as on the smoothness of x_* . Depending on the specification of the prior this rate may be *minimax* over a Sobolev ellipsoid. In particular, (i) when the regularity of the prior matches the regularity of x_* , the minimax rate of convergence is obtained with a fixed prior covariance; (ii) when the prior is rougher or smoother at any degree than the truth, the minimax rate can still be obtained if the prior is scaled at an appropriate rate depending on the unknown regularity of x_* .

Our third contribution consists in proposing a new data-driven method for optimally selecting the regularization parameter. This parameter enters the prior distribution as a scaling hyperparameter of the prior covariance and is needed to compute the posterior mean of x. Our adaptive data-driven method is based on an empirical Bayes (EB) approach. Because the posterior mean is, under our assumptions, equal to a Tikhonov-type estimator for problem (1), our EB approach for selecting the regularization parameter is valid, and can be used, also for computing frequentist estimators based on Tikhonov regularization.¹ Finally, the EB-selected regularization parameter is plugged into the prior distribution of x and for the corresponding EB-posterior distribution we prove frequentist posterior consistency.

In the following, we present the Bayesian approach and the asymptotic results for general models of the form (1); then, we develop further results that apply to the specific examples 1 and 2. In section 2 we set the Bayesian model associated with (1) and the main assumptions. In section 3 the posterior distribution of x is computed and its frequentist asymptotic properties are analyzed. Section 4 focuses on the mildly ill-posed case. The EB method is developed in section 5. Section 6 shows numerical implementations and section 7 concludes. All the proofs are in the Appendix.

2 The Model

Let \mathcal{X} and \mathcal{Y} be infinite dimensional separable Hilbert spaces over \mathbb{R} with norm $||\cdot||$ induced by the inner product $\langle \cdot, \cdot \rangle$. Let $\mathfrak{B}(\mathcal{X})$ and $\mathfrak{B}(\mathcal{Y})$ be the Borel σ -fields generated by the open sets of \mathcal{X} and \mathcal{Y} , respectively. We consider the inverse problem of estimating the function $x \in \mathcal{X}$ which

¹Notice that in general the posterior mean in a conjugate-Gaussian problem stated in infinite-dimensional Hilbert spaces cannot be equal to the Tikhonov solution of (1). This is due to the particular structure of the covariance operator of the error term U^{δ} and it will be detailed in section 2.

is linked to the data y^{δ} through the linear relation

$$y^{\delta} = Kx + U^{\delta}, \quad x \in \mathcal{X}, \quad y^{\delta} \in \mathcal{Y}, \quad \delta > 0$$
 (4)

where y^{δ} is an observable function and $K: \mathcal{X} \to \mathcal{Y}$ is a known, bounded, linear operator (we refer to Carrasco *et al.* (2007) for definition of terminology from functional analysis). The elements y^{δ} and U^{δ} are Hilbert space-valued random variables (H-r.v.), that is, for a complete probability space $(S, \mathcal{S}, \mathbb{P}), U^{\delta}$ (resp. y^{δ}) defines a measurable map $U^{\delta}: (S, \mathcal{S}, \mathbb{P}) \to (\mathcal{Y}, \mathfrak{B}(\mathcal{Y}))$, see *e.g.* Kuo (1975). Realizations of y^{δ} are functional transformations of the observed data and U^{δ} is an error term. In the following, we simply write U, instead of U^{δ} , in order to lighten the notation. The true value of x that generates the data is denoted by x_* .

We assume a mean-zero Gaussian distribution on $\mathfrak{B}(\mathcal{Y})$ for $U: U \sim \mathcal{N}(0, \delta\Sigma)$ where $\delta > 0$ is the noise level and $\Sigma: \mathcal{Y} \to \mathcal{Y}$ is a covariance operator, that is, Σ is such that $< \delta\Sigma\phi_1, \phi_2 > = \mathbf{E}(< U, \phi_1 > < U, \phi_2 >)$ for all $\phi_1, \phi_2 \in \mathcal{Y}$. Therefore, Σ is a one-to-one, linear, positive definite, self-adjoint and trace-class operator. Because Σ is one-to-one the support of U is all \mathcal{Y} , see Kuo (1975) and Ito (1970). A trace-class operator is a compact operator with eigenvalues that are summable. This property rules out a covariance Σ proportional to the identity operator I and this is a key difference between our model and the model used in a large part of the statistical inverse problem literature, see e.g. Cavalier and Tsybakov (2002), Bissantz et al. (2007) and Knapik et al. (2011). The reason why we focus on a covariance Σ different from I is that most of the functional estimation problems in econometrics that writes as equation (4) does not allow for an identity (or proportional to identity) covariance operator, see examples 1 and 2.

Remark 2.1. We stress that $\Sigma \neq I$ is not an assumption but a feature of the econometric models we want to treat and many difficulties arise because $\Sigma \neq I$. One of the main contributions of our paper is to deal with the important case where Σ is trace-class, i.e. $\Sigma \neq I$. In addition, the identity operator cannot be the covariance operator of a well-defined Gaussian process in a Hilbert space because a Gaussian process with an identity covariance operator cannot generate trajectories that are proper elements in infinite dimensional Hilbert spaces and is only defined in a weaker sense as a weak distribution, see e.g. Kuo (1975). A weak distribution is a very complicated object and, in general, real (functional) data cannot be generated by a Gaussian process with an identity covariance operator. Thus, the model considered by Knapik et al. (2011) is not appropriate for many econometric models of interest like the IV and the functional linear regression estimation. \square

Under model (4) with $U \sim \mathcal{N}(0, \delta \Sigma)$, the sampling distribution P^x , *i.e.* the conditional distribution of y^{δ} given x, is a Gaussian distribution on $\mathfrak{B}(\mathcal{Y})$:

$$y^{\delta}|x \sim P^x = \mathcal{N}(Kx, \delta\Sigma).$$
 (5)

The parameter $\delta > 0$ is assumed to converge to 0 when the information from the observed data increases. Hereafter, $\mathbf{E}_x(\cdot)$ will denote the expectation taken with respect to P^x .

Remark 2.2. The assumption of Gaussianity of the error term U in the econometric model (4) is not necessary and only made in order to construct (and give a Bayesian interpretation to) the estimator. The proofs of our results of frequency consistency do not rely on the normality of U. In particular, asymptotic normality of $y^{\delta}|x$, which is verified e.g. in example 1, is enough for our estimation procedure and also for our EB data-driven method for choosing the regularization parameter. \square

Remark 2.3. All the results in the paper are given for the general case where K and Σ do not depend on δ . This choice is made in order to keep our presentation as simple as possible. We discuss how our results apply to the case with unknown K and Σ through examples 1 and 2. \square

2.1 Notation

We set up some notational convention used in the paper. We simply write U to denote the error term U^{δ} . For positive quantities M_{δ} and N_{δ} depending on a discrete or continuous index δ , we write $M_{\delta} \approx N_{\delta}$ to mean that the ratio M_{δ}/N_{δ} is bounded away from zero and infinity. We write $M_{\delta} = \mathcal{O}(N_{\delta})$ if M_{δ} is at most of the same order as N_{δ} . For an H-r.v. W we write $W \sim \mathcal{N}$ for denoting that W is a Gaussian process. We denote by $\mathcal{R}(\cdot)$ the range of an operator and by $\mathcal{D}(\cdot)$ its domain. For an operator $B: \mathcal{X} \to \mathcal{Y}$, $\overline{\mathcal{R}(B)}$ denotes the closure in \mathcal{Y} of the range of B. For a bounded operator $A: \mathcal{Y} \to \mathcal{X}$, we denote by A^* its adjoint, i.e. $A^*: \mathcal{X} \to \mathcal{Y}$ is such that $\langle A\psi, \varphi \rangle = \langle \psi, A^*\varphi \rangle$, $\forall \varphi \in \mathcal{X}$, $\psi \in \mathcal{Y}$. The operator norm is defined as $||A|| := \sup_{||\phi||=1} ||A\phi|| = \min\{C \geq 0 \; ; \; ||A\phi|| \leq C||\phi||$ for all $\phi \in \mathcal{Y}$. The operator I denotes the identity operator on both spaces \mathcal{X} and \mathcal{Y} , i.e. $\forall \psi \in \mathcal{X}$, $\varphi \in \mathcal{Y}$, $I\psi = \psi$ and $I\varphi = \varphi$.

Let $\{\varphi_j\}_j$ denote an orthonormal basis of \mathcal{Y} . The trace of a bounded linear operator $A: \mathcal{Y} \to \mathcal{X}$ is defined as $tr(A) := \sum_{j=1}^{\infty} <(A^*A)^{\frac{1}{2}}\varphi_j, \varphi_j >$ independently of the basis $\{\varphi_j\}_j$. If A is compact then its trace writes $tr(A) = \sum_{j=1}^{\infty} \lambda_j$, where $\{\lambda_j\}$ are the singular values of A. The Hilbert-Schmidt norm of a bounded linear operator $A: \mathcal{Y} \to \mathcal{X}$ is denoted by $||A||_{HS}$ and defined as $||A||_{HS}^2 = tr(A^*A)$, see Kato (1995).

2.2 Prior measure and main assumptions

In this section we introduce the prior distribution and two sets of assumptions. (i) The first set (Assumptions A.2 and B below) will be used for establishing the rate of contraction of the posterior distribution and concerns the smoothness of the operator $\Sigma^{-1/2}K$ and of the true value x_* . (ii) The assumptions in the second set (A.1 and A.3 below) are new in the literature and guarantee continuity of the posterior mean of x as a linear operator on \mathcal{Y} . The detection of the latter assumptions is

an important contribution because, as remarked in Luschgy (1995) and Mandelbaum (1984), in the Gaussian infinite-dimensional model the posterior mean is in general only defined as a *mlt* which is a weaker notion than that one of a continuous linear operator. Therefore, in general the posterior mean has not an explicit form and may be an inconsistent estimator in the frequentist sense while our Assumptions A.1 and A.3 ensure closed-form (easy to compute) and consistency for the posterior mean.

Assumption A.1. $\mathcal{R}(K) \subset \mathcal{D}(\Sigma^{-1/2})$.

Since K and Σ are integral operators, $\Sigma^{-1/2}$ is a differential operator and Assumption A.1 demands that the functions in $\mathcal{R}(K)$ are at least as smooth as the functions in $\mathcal{R}(\Sigma^{1/2})$. A.1 ensures that $\Sigma^{-1/2}$ is defined on $\mathcal{R}(K)$ so that $\Sigma^{-1/2}K$, which is used in Assumption A.2, exists.

Assumption A.2. There exists an unbounded, self-adjoint, densely defined operator L in the Hilbert space \mathcal{X} for which $\exists \eta > 0$ such that $\langle L\psi, \psi \rangle \geq \eta ||\psi||^2$, $\forall \psi \in \mathcal{D}(L)$, and that satisfies

$$\underline{m}||L^{-a}x|| \le ||\Sigma^{-1/2}Kx|| \le \overline{m}||L^{-a}x|| \tag{6}$$

on \mathcal{X} for some a>0 and $0<\underline{m}\leq\overline{m}<\infty$. Moreover, L^{-2s} is trace-class for some s>a.

Assumption A.2 means that $\Sigma^{-1/2}K$ regularizes at least as much as L^{-a} . Because $\Sigma^{-1/2}K$ must satisfy (6) it is necessarily an injective operator.

We turn now to the construction of the prior distribution of x. Our proposal is to use the operator L^{-2s} to construct the prior covariance operator. We assume a Gaussian *prior distribution* μ on $\mathfrak{B}(\mathcal{X})$:

$$x|\alpha, s \sim \mu = \mathcal{N}\left(x_0, \frac{\delta}{\alpha}\Omega_0\right), \qquad x_0 \in \mathcal{X}, \quad \Omega_0 := L^{-2s}, \ s > a$$
 (7)

with $\alpha > 0$ such that $\alpha \to 0$ as $\delta \to 0$. The parameter α describes a class of prior distributions and it may be viewed as an hyperparameter. Section 5 provides an EB approach for selecting it.

By definition of L, the operator $\Omega_0: \mathcal{X} \to \mathcal{X}$ is linear, bounded, positive-definite, self-adjoint, compact and trace-class. It results evident from Assumption A.2 that such a choice for the prior covariance is aimed at linking the prior distribution to the sampling model. A similar idea was proposed by Zellner (1986) for linear regression models for which he constructed a class of prior called g-prior. Our prior (7) is an extension of the Zellner's g-prior and we call it extended g-prior.

The distribution μ (resp. P^x) is realizable as a proper random element in \mathcal{X} (resp. \mathcal{Y}) if and only if Ω_0 (resp. Σ) is trace-class. Thus, neither Σ nor Ω_0 can be proportional to I so that, in general, in infinite-dimensional inverse problems, the posterior mean cannot be equal to the Tikhonov regularized estimator $x_{\alpha}^T := (\alpha I + K^*K)^{-1}K^*y^{\delta}$. However, we show in this paper that under A.1, A.2 and A.3, the posterior mean equals the Tikhonov regularized solution in the *Hilbert Scale* generated by L. We give later the definition of Hilbert Scale.

The following assumption ties further the prior to the sampling distribution by linking the smoothing properties of Σ , K and $\Omega_0^{\frac{1}{2}}$.

Assumption A.3. $\mathcal{R}(K\Omega_0^{\frac{1}{2}}) \subset \mathcal{D}(\Sigma^{-1})$.

Hereafter, we denote $B = \Sigma^{-1/2} K \Omega_0^{\frac{1}{2}}$. Assumption A.3 guarantees that B and $\Sigma^{-1/2} B$ exist.

We now discuss the relationship between Assumption A.2, which quantifies the smoothness of $\Sigma^{-1/2}K$, and Assumption B below, which quantifies the smoothness of the true value x_* . In order to explain these assumptions and their link we will: (i) introduce the definition of Hilbert scale, (ii) explain the meaning of the parameter a in (6), (iii) discuss the smoothness conditions of $\Sigma^{-1/2}K$ and of the true value of x.

- (i) The operator L in Assumption A.2 is a generating operator of the Hilbert scale $(\mathcal{X}_t)_{t\in\mathbb{R}}$ where $\forall t\in\mathbb{R}$, \mathcal{X}_t is the completion of $\bigcap_{k\in\mathbb{R}} \mathcal{D}(L^k)$ with respect to the norm $||x||_t := ||L^t x||$ and is a Hilbert space, see Definition 8.18 in Engl *et al.* (2000), Goldenshluger and Pereverzev (2003) or Krein and Petunin (1966). For t>0 the space $\mathcal{X}_t\subset\mathcal{X}$ is the domain of definition of L^t : $\mathcal{X}_t=\mathcal{D}(L^t)$. Typical examples of \mathcal{X}_t are Sobolev spaces of various kinds.
- (ii) We refer to the parameter a in A.2 as the "degree of ill-posedness" of the estimation problem under study and a is determined by the rate of decreasing of the spectrum of $\Sigma^{-1/2}K$ (and not only by that one of K as it would be in a classical inverse problems framework for (1)). Since the spectrum of $\Sigma^{-1/2}K$ is decreasing slower than that one of K we have to control for less ill-posedness than if we used the classical approach.
- (iii) In inverse problems theory it is natural to impose conditions on the regularity of x_* by relating it to the regularity of the operator that characterizes the inverse problem (that is, the operator $\Sigma^{-1/2}K$ in our case). A possible implementation of this consists in introducing a Hilbert Scale and expressing the regularity of both x_* and $\Sigma^{-1/2}K$ with respect to this common Hilbert Scale. This is the meaning of and the link between Assumptions A.2 and B where we use the Hilbert Scale $(\mathcal{X}_t)_{t\in\mathbb{R}}$ generated by L. We refer to Chen and Reiss (2011) and Johannes et al. (2011) for an explanation of the relationship between Hilbert Scale and regularity conditions. The following assumption expresses the regularity of x_* according to \mathcal{X}_t .

Assumption B. For some $0 \leq \beta$, $(x_* - x_0) \in \mathcal{X}_{\beta}$, that is, there exists a $\rho_* \in \mathcal{X}$ such that $(x_* - x_0) = L^{-\beta} \rho_* \ (\equiv \Omega_0^{\frac{\beta}{2s}} \rho_*)$.

The parameter β characterizes the "regularity" of the centered true function $(x_* - x_0)$ and is generally unknown. Assumption B is satisfied by regular functions x_* . In principle, it could be satisfied also by irregular x_* if we were able to decompose x_* in the sum of a regular part plus an irregular part and to choose x_0 such that it takes all the irregularity of x_* . This is clearly infeasible in practice as x_* is unknown. On the contrary, we could choose a very smooth function x_0 so that Assumption B would be less demanding about the regularity of x_* . When \mathcal{X}_{β} is the

scale of Sobolev spaces, Assumption B is equivalent to assume that $(x_* - x_0)$ has at least β square integrable derivatives.

Assumption B is classical in inverse problems literature, see e.g. Chen and Reiss (2011) and Nair et al. (2005), and is closely related to the so-called source condition which expresses the regularity of x_* according to the spectral representation of the operator K^*K defining the inverse problem, see Engl et al. (2000) and Carrasco et al. (2007). In our case, the regularity of $(x_* - x_0)$ is expressed according to the spectral representation of L.

Remark 2.4. Assumption A.2 covers not only the mildly ill-posed but also the severely ill-posed case if $(x_* - x_0)$ in Assumption B is infinitely smooth. In the mildly ill-posed case the singular values of $\Sigma^{-1/2}K$ decay slowly to zero (typically at a geometric rate) which means that the kernel of $\Sigma^{-1/2}K$ is finitely smooth. In this case the operator L is generally some differential operator so that L^{-1} is finitely smooth. In the severely ill-posed case the singular values of $\Sigma^{-1/2}K$ decay very rapidly (typically at an exponential rate). Assumption A.2 covers also this case if $(x_* - x_0)$ is very smooth. This is because when the singular values of $\Sigma^{-1/2}K$ decay exponentially, Assumption A.2 is satisfied if L^{-1} has an exponentially decreasing spectrum too. On the other hand, L^{-1} is used to describe the regularity of $(x_* - x_0)$, so that in the severely ill-posed case, Assumption B can be satisfied only if $(x_* - x_0)$ is infinitely smooth. In this case we could for instance take $L = (K^*\Sigma^{-1}K)^{-\frac{1}{2}}$ which implies a = 1. We could make Assumption A.2 more general, as in Chen and Reiss (2011), in order to cover the severely ill-posed case even when $(x_* - x_0)$ is not infinitely smooth. Since computations to find the rate would become more cumbersome (even if still possible) we do not pursue this direction here. \square

Remark 2.5. The specification of the prior covariance operator can be generalized as $\frac{\delta}{\alpha}\Omega_0 = \frac{\delta}{\alpha}QL^{-2s}Q^*$, for some bounded operator Q not necessarily compact. Then, the previous case is a particular case of this one for Q=I. In this setting, Assumptions A.1 and A.3 are replaced by the weaker assumptions $\mathcal{R}(KQ) \subset \mathcal{D}(\Sigma^{-1/2})$ and $\mathcal{R}(KQL^{-s}) \subset \mathcal{D}(\Sigma^{-1})$, respectively. In Assumption A.2 the operator $\Sigma^{-1/2}K$ must be replaced by $\Sigma^{-1/2}KQ$ and Assumption B becomes: there exists $\tilde{\rho}_* \in \mathcal{X}$ such that $(x_* - x_0) = QL^{-\beta}\tilde{\rho}_*$. \square

Example 1 (Instrumental variable (IV) regression estimation (continued)). Let us consider the integral equation (4), with y^{δ} and K defined as in (2), that characterizes the IV regression x. Suppose to use the kernel smoothing approach to estimate f_{YW} and f_{ZW} , where f_{YW} denotes the density of the distribution of (Y,W) with respect to the Lebesgue measure. For simplicity we assume that (Z,W) is a bivariate random vector. Let $K_{Z,h}$ and $K_{W,h}$ denote two univariate kernel functions in $L^2[0,1]$, h be the bandwidth and $(y_i,w_i,z_i)_{i=1}^n$ be the n-observed random sample. Denote $\Lambda:L^2[0,1]\to L^2[0,1]$ the operator $\Lambda\varphi=\int a(w,v)\varphi(w)dw$, with a(w,v) a known function, and $\tilde{K}:L^2[0,1]\to L^2[0,1]$ the operator $\tilde{K}\phi=\frac{1}{n}\sum_{i=1}^n\frac{K_{W,h}(w_i-w)}{h}<\phi(z),\frac{K_{Z,h}(z_i-z)}{h}>$. Therefore,

 $K = \Lambda \tilde{K}$ so that the quantities in (2) can be rewritten as

$$y^{\delta} = \Lambda \left[\hat{\mathbf{E}}(Y|W=w) \hat{f}_W \right](v) = \int a(w,v) \frac{1}{nh} \sum_{i=1}^n y_i K_{W,h}(w_i - w) dw$$
and
$$Kx = \int a(w,v) \frac{1}{n} \sum_{i=1}^n \frac{K_{W,h}(w_i - w)}{h} \int x(z) \frac{K_{Z,h}(z_i - z)}{h} dz dw$$
(8)

Remark that $\lim_{n\to\infty} \tilde{K}\phi = f_W(w)\mathbf{E}(\phi|w) = M_f\mathbf{E}(\phi|w)$ where M_f denotes the multiplication operator by f_W . If $a = f_{WZ}$ then $\Lambda \lim_{n\to\infty} \tilde{K}$ is the same integral operator in Hall and Horowitz (2005).

In this example, the assumption that $U \sim \mathcal{N}(0, \delta \Sigma)$ (where $U = y^{\delta} - Kx$) holds asymptotically and the transformation of the model through Λ is necessary in order to guarantee such a convergence of U towards a zero-mean Gaussian process. We explain this fact by extending Ruymgaart (1998). It is possible to show that the covariance operator $\tilde{\Sigma}_h$ of $\frac{\sqrt{n}}{n} \sum_{i=1}^n \left(y_i - \langle x, \frac{K_{Z,h}(z_i - z)}{h} \rangle \right) \frac{K_{W,h}(w_i - w)}{h}$ satisfies

$$<\phi_1, \tilde{\Sigma}_h \phi_2> \longrightarrow <\phi_1, \tilde{\Sigma}\phi_2>, \text{ as } h\to 0, \forall \phi_1, \phi_2\in L^2[0,1]$$

where $\tilde{\Sigma}\phi_2=\sigma^2 f_W(v)\phi_2(v)=\sigma^2 M_f\phi_2(v)$ under the assumption $\mathbf{E}[(Y-x(Z))^2|W]=\sigma^2<\infty$. Unfortunately, because $\tilde{\Sigma}$ has not finite trace, it is incompatible with the covariance structure of a Gaussian limiting probability measure. The result is even worst, since Ruymgaart (1998) shows that there are no scaling factors n^{-r} , for 0< r<1, such that $n^{-r}\sum_{i=1}^n\left(y_i-\langle x,\frac{K_{Z,h}(z_i-z)}{h}\rangle\right)\times\frac{K_{W,h}(w_i-w)}{h}$ converges weakly in $L^2[0,1]$ to a Gaussian distribution (unless this distribution is degenerate at the zero function). However, if we choose a(w,v) appropriately so that Λ is a compact operator and $\Lambda^*\Lambda$ has finite trace, then $\sqrt{n}\left(y^\delta-Kx\right)\Rightarrow \mathcal{N}(0,\Lambda\tilde{\Sigma}\Lambda^*)$, in $L^2[0,1]$ as $n\to\infty$, where ' \Rightarrow ' denotes weak convergence and $\Lambda\tilde{\Sigma}\Lambda^*=:\Sigma$. The adjoint operator $\Lambda^*:L^2[0,1]\to L^2[0,1]$ is defined as: $\forall \varphi\in L^2[0,1]$, $\Lambda^*\varphi=\int a(w,v)\varphi(v)dv$ and $\Lambda=\Lambda^*$ if a(w,v) is symmetric in w and v. The operator Σ is unknown and can be estimated by $\hat{\Sigma}=\Lambda\frac{\hat{\sigma}^2}{n}\sum_{i=1}^n\frac{K_{W,h}(w_i-w)}{h}\Lambda^*$.

We now discuss Assumptions A.1, A.2 and A.3. While A.1 and A.3 need to hold both in finite sample and for $n \to \infty$, A.2 only has to hold for $n \to \infty$. We start by checking Assumption A.1. In large sample, the operator K converges to $\Lambda M_f \mathbf{E}(\phi|w)$ and it is trivial to verify that $\mathcal{D}(\Sigma^{-1/2}) = \mathcal{R}(\Lambda M_f^{\frac{1}{2}}) \supset \mathcal{R}(\Lambda M_f) \supset \mathcal{R}(\Lambda M_f \mathbf{E}(\cdot|w))$. In finite sample, the same holds with M_f and $\mathbf{E}(\cdot|w)$ replaced by their empirical counterparts. Next, we check the validity of Assumption A.2 for $n \to \infty$. Remark that the operator $\Sigma^{1/2}$ may be equivalently defined in two ways: (1) as a self-adjoint operator, that is $\Sigma^{1/2} = (\Sigma^{1/2})^* = (\Lambda M_f \Lambda^*)^{1/2}$, so that $\Sigma = \Sigma^{1/2} \Sigma^{1/2}$ or (2) as $\Sigma^{1/2} = \Lambda M_f^{1/2}$ so that $\Sigma = \Sigma^{1/2} (\Sigma^{1/2})^*$ where $(\Sigma^{1/2})^* = M_f^{\frac{1}{2}} \Lambda^*$. By using the second definition, we obtain that $\Sigma^{-1/2}K = (\Lambda M_f^{\frac{1}{2}})^{-1}\Lambda \tilde{K} = f_W^{-\frac{1}{2}}\Lambda^{-1}\Lambda \tilde{K} = f_W^{-\frac{1}{2}}\tilde{K}$ and $\forall x \in \mathcal{X}$, $||\lim_{n\to\infty} \Sigma^{-1/2}Kx|| = ||\mathbf{E}(x|w)||_W^2$, where $||\phi||_W^2 := \int [\phi(w)]^2 f_W(w) dw$. This shows that, in the IV case, Assumption A.2 is a particular case of Assumptions 2.2 and 4.2 in Chen and Reiss (2011).

Finally, we check Assumption A.3 for both $n \to \infty$ and finite sample. In finite sample this assumption holds trivially since $\mathcal{R}(\hat{\Sigma})$ ($\equiv \mathcal{D}(\hat{\Sigma}^{-1})$) and $\mathcal{R}(K\Omega_0^{\frac{1}{2}})$ have finite ranks. Suppose that the conditions for the application of the Dominated Convergence Theorem hold, then Assumption A.3 is satisfied asymptotically if and only if $||\Sigma^{-1}\Lambda\lim_{n\to\infty}\tilde{K}\Omega_0^{\frac{1}{2}}||^2 < \infty$. This holds if Ω_0 is appropriately chosen. One possibility could be to set $\Omega_0 = T^*\Lambda^*\Lambda T$, where $T: L^2[0,1] \to L^2[0,1]$ is a trace-class integral operator $T\phi = \int \omega(w,z)\phi(z)dz$ for a known function ω and T^* is its adjoint. Define $\Omega_0^{\frac{1}{2}} = T^*\Lambda^*$. Then,²

$$\begin{split} ||\Sigma^{-1}\Lambda \lim_{n\to\infty} \tilde{K}\Omega_0^{\frac{1}{2}}||^2 &= ||\Sigma^{-1}\Lambda \lim_{n\to\infty} \tilde{K}T^*\Lambda^*||^2 \leq ||\Sigma^{-1}\Lambda \lim_{n\to\infty} \tilde{K}T^*\Lambda^*||^2_{HS} \leq tr(\Sigma^{-1}\Lambda \lim_{n\to\infty} \tilde{K}T^*\Lambda^*) \\ &= tr(\Lambda^*\Sigma^{-1}\Lambda \lim_{n\to\infty} \tilde{K}T^*) = tr(\mathbf{E}(T^*\cdot|w)) \leq tr(T^*)||\mathbf{E}(\cdot|W)|| < \infty. \quad \Box \end{split}$$

2.2.1 Covariance operators proportional to K

A particular situation often encountered in applications is the case where the sampling covariance operator has the form $\Sigma = (KK^*)^r$, for some $r \in \mathbb{R}_+$, and is related to the classical example of g-priors given in Zellner (1986). In this situation it is convenient to choose $L = (K^*K)^{-\frac{1}{2}}$ so that $\Omega_0 = (K^*K)^s$, for $s \in \mathbb{R}_+$. Because $(KK^*)^r$ and $(K^*K)^s$ are proper covariance operators only if they are trace-class then K is necessarily compact. Assumptions A.1 and A.2 hold for $r \leq 1$ and a = 1 - r, respectively. Assumption A.3 holds for $s \geq 2r - 1$.

Example 2 (Functional Linear Regression Estimation (continued)). Let us consider model (3) and the associated integral equation $\mathbf{E}(\xi Z(t)) = \int_0^1 x(s) Cov(Z(s), Z(t)) ds$, for $t \in [0, 1]$. If we dispose of i.i.d. data $(\xi_1, Z_1), \ldots, (\xi_n, Z_n)$ the unknown moments in this equation can be estimated. Thus, x is solution of (4) with $y^{\delta} := \frac{1}{n} \sum_i \xi_i Z_i(t)$, $U = \frac{1}{n} \sum_i \varepsilon_i Z_i(t)$ and $\forall \varphi \in L^2([0, 1]) \mapsto K\varphi := \frac{1}{n} \sum_i \langle Z_i, \varphi \rangle Z_i(t)$. The operator $K : L^2([0, 1]) \to L^2([0, 1])$ is self-adjoint, i.e. $K = K^*$. Moreover, conditional on Z, the error term U is exactly a Gaussian process with covariance operator $\delta \Sigma = \delta \tau^2 K$ with $\delta = \frac{1}{n}$ which is trace-class since its range has finite dimension. Thus, we can write $\delta \Sigma = \frac{1}{n} \tau^2 (KK^*)^r$ with $r = \frac{1}{2}$. Assumption A.1 is trivially satisfied in finite sample as well as for $n \to \infty$. We discuss later on how to choose L in order to satisfy Assumptions A.2 and A.3. \square

2.2.2 Imposing shape restrictions through the Gaussian process prior

We briefly discuss how restrictions such as monotonicity, convexity, or more generally shape restrictions can be incorporated into a Gaussian process prior, see Berger and Wang (2011) for more details. First, to be sure that the realizations of a Gaussian process with mean \tilde{x}_0 and covariance operator $\tilde{\Omega}_0$ have derivatives of all orders one should take as \tilde{x}_0 an infinitely differentiable function and as kernel of $\tilde{\Omega}_0$ a squared exponential covariance function, that is: $\tilde{\omega}_0(t,\tilde{t}) = b_1 \exp\{-b_2(t-\tilde{t})^2\}$,

This is because for a compact operator $A: L^2[0,1] \to L^2[0,1]$ and by denoting $|A| = (A^*A)^{1/2}$, we have: $||A||^2 \le ||A||^2_{HS} = ||A||^2_{HS} \le tr(|A|) = tr(A)$.

 $b_1 > 0, b_2 \ge 0$ where $\tilde{\Omega}_0 \varphi = \int \tilde{\omega}_0(\tilde{t}, t) \varphi(\tilde{t}) d\tilde{t}, \forall \varphi \in \mathcal{X}.$

Second, since differentiation is a linear operator, derivatives of a Gaussian process remain a Gaussian process. Thus, if $x^{(k)}$ denotes the k-th derivative of x, then $x^{(k)}$ is a Gaussian process with mean $x_0^{(k)}$ and covariance operator $\tilde{\Omega}_{kk}$ with kernel $\sigma_{kk}(\tilde{t},t) = \frac{\partial^{2k}}{\partial t^k \partial t^k} \tilde{\omega}_0(\tilde{t},t)$. Moreover, the covariance operator between x and $x^{(k)}$ is denoted by $\tilde{\Omega}_{0k}$ and has kernel $\tilde{\sigma}_{kk}(\tilde{t},t) = \frac{\partial^k}{\partial t^k} \tilde{\omega}_0(\tilde{t},t)$. To incorporate the shape restrictions one should constrain the derivative $x^{(k)}(t)$ at a discrete set of points t_1, \ldots, t_m (dense enough to guarantee that the realizations of x will satisfy the constraint) and then fix the prior μ , for x, equal to the conditional distribution of $x|\{x^{(k)}(t_i)\}_{i=1}^m$, which is Gaussian. That is, by abuse of notation: $x|\alpha,t\sim\mu=\mathcal{N}(x_0,\Omega_0)$ with

$$x_0 = \tilde{x}_0 + \tilde{\Omega}_{0k}\tilde{\Omega}_{kk}^{-1} \{x^{(k)}(t_i) - x_0^{(k)}(t_i)\}_{i=1}^m, \qquad \Omega_0 = \tilde{\Omega}_0 - \tilde{\Omega}_{0k}\tilde{\Omega}_{kk}^{-1}\tilde{\Omega}_{k0}$$

where $\{x^{(k)}(t_i) - x_0^{(k)}(t_i)\}_{i=1}^m$ denotes the *m*-vector with *i*-th element $x^{(k)}(t_i) - x_0^{(k)}(t_i)$.

3 Main results

The posterior distribution of x, denoted by μ_{δ}^{Y} , is the Bayesian solution of the inverse problem (4). Because a separable Hilbert space is Polish, there exists a regular version of the posterior distribution μ_{δ}^{Y} , that is, a conditional probability characterizing μ_{δ}^{Y} . In many applications \mathcal{X} and \mathcal{Y} are L^{2} spaces and L^{2} spaces are Polish if they are defined on a separable metric space. In the next theorem we characterize the joint distribution of (x, y^{δ}) and the posterior distribution μ_{δ}^{Y} of x. The notation $\mathfrak{B}(\mathcal{X}) \otimes \mathfrak{B}(\mathcal{Y})$ means the Borel σ -field generated by the product topology.

Theorem 1. Consider two separable infinite dimensional Hilbert spaces \mathcal{X} and \mathcal{Y} . Let $x|\alpha$, s and $y^{\delta}|x$ be two Gaussian H-r.v. on \mathcal{X} and \mathcal{Y} as in (7) and (5), respectively. Then,

(i) (x, y^{δ}) is a measurable map from $(S, \mathcal{S}, \mathbb{P})$ to $(\mathcal{X} \times \mathcal{Y}, \mathfrak{B}(\mathcal{X}) \otimes \mathfrak{B}(\mathcal{Y}))$ and has a Gaussian distribution: $(x, y^{\delta})|_{\alpha}, s \sim \mathcal{N}((x_0, Kx_0), \Upsilon)$, where Υ is a trace-class covariance operator defined as $\Upsilon(\varphi, \psi) = \left(\frac{\delta}{\alpha}\Omega_0\varphi + \frac{\delta}{\alpha}\Omega_0K^*\psi, \frac{\delta}{\alpha}K\Omega_0\varphi + (\delta\Sigma + \frac{\delta}{\alpha}K\Omega_0K^*)\psi\right)$ for all $(\varphi, \psi) \in \mathcal{X} \times \mathcal{Y}$. The marginal sampling distribution of $y^{\delta}|_{\alpha}$, s is $P_{\alpha} \sim \mathcal{N}(Kx_0, (\delta\Sigma + \frac{\delta}{\alpha}K\Omega_0K^*))$.

Moreover, let $A: \mathcal{Y} \to \mathcal{X}$ be a P_{α} -measurable linear transformation (P_{α} -mlt), that is, $\forall \phi \in \mathcal{Y}$, $A\phi$ is a P_{α} -almost sure limit, in the norm topology, of $A_k\phi$ as $k \to \infty$, where $A_k: \mathcal{Y} \to \mathcal{X}$ is a sequence of continuous linear operators. Then,

(ii) the conditional distribution μ_{δ}^{Y} of x given y^{δ} exists, is regular and almost surely unique. It is Gaussian with mean $\mathbf{E}(x|y^{\delta},\alpha,s)=A(y^{\delta}-Kx_{0})+x_{0}$ and trace-class covariance operator $Var(x|y^{\delta},\alpha,s)=\frac{\delta}{\alpha}[\Omega_{0}-AK\Omega_{0}]:\mathcal{X}\to\mathcal{X}$. Furthermore, $A=\Omega_{0}K^{*}(\alpha\Sigma+K\Omega_{0}K^{*})^{-1}$ on $\mathcal{R}((\delta\Sigma+\frac{\delta}{\alpha}K\Omega_{0}K^{*})^{\frac{1}{2}})$.

(iii) Under Assumptions A.1 and A.3, the operator A characterizing μ_{δ}^{Y} is a continuous linear operator on \mathcal{Y} and can be written as

$$A = \Omega_0^{\frac{1}{2}} (\alpha I + B^* B)^{-1} (\Sigma^{-1/2} B)^* : \mathcal{Y} \to \mathcal{X}$$
 (9)

with
$$B = \Sigma^{-1/2} K \Omega_0^{\frac{1}{2}}$$
.

Point (ii) in the theorem is an immediate application of the results of Mandelbaum (1984), we refer to this paper for the proof and for a rigorous definition of P_{α} -mlt. As stated above, the quantity Ay^{δ} is defined as a P_{α} -mlt, which is a weaker notion than that of a linear and continuous operator and A is in general not continuous. In fact, since Ay^{δ} is a P_{α} -almost sure limit of A_ky^{δ} , for $k \to \infty$, the null set where this convergence is not satisfied depends on y^{δ} and we do not have an almost sure convergence of A_k to A. Moreover, in general, A takes the form $A = \Omega_0 K^* (\alpha \Sigma + K \Omega_0 K^*)^{-1}$ only on a dense subspace of \mathcal{Y} of P_{α} -probability measure zero. Outside of this subspace, A is defined as the unique extension of $\Omega_0 K^* (\alpha \Sigma + K \Omega_0 K^*)^{-1}$ to \mathcal{Y} for which we do not have an explicit expression. This means that in general it is not possible to construct a feasible estimator for x.

On the contrary, point (iii) of the theorem shows that, under A.1 and A.3, A is defined as a continuous linear operator on the whole \mathcal{Y} . This is the first important contribution of our paper since A.1 and A.3 permit to construct a linear estimator for x – equal to the posterior mean – that is implementable in practice. Thus, our result (iii) makes operational the Bayesian approach for linear statistical inverse problems in econometrics. When Assumptions A.1 and A.3 do not hold then we can use a quasi-Bayesian approach as proposed in Florens and Simoni (2012a,b). Summarizing, under a quadratic loss function, the Bayes estimator for a functional parameter x characterized by (4) is the posterior mean

$$\hat{x}_{\alpha} = \Omega_0^{\frac{1}{2}} (\alpha I + B^* B)^{-1} (\Sigma^{-1/2} B)^* (y^{\delta} - K x_0), \text{ with } B = \Sigma^{-1/2} K \Omega^{1/2}.$$

Remark 3.1. Under our assumptions it is possible to show the existence of a close relationship between Bayesian and frequentist approach to statistical inverse problems. In fact, the posterior mean \hat{x}_{α} is equal to the Tikhonov regularized solution in the Hilbert scale $(\mathcal{X}_s)_{s\in\mathbb{R}}$ generated by L of the equation $\Sigma^{-1/2}y^{\delta} = \Sigma^{-1/2}Kx + \Sigma^{-1/2}U$. The existence of $\Sigma^{-1/2}K$ is guaranteed by A.1. Since $\mathbf{E}(x|y^{\delta},\alpha,s) = A(y^{\delta} - Kx_0) + x_0$, we have, under A.1 and A.3:

$$\mathbf{E}(x|y^{\delta}, \alpha, s) = L^{-s}(\alpha I + L^{-s}K^*\Sigma^{-1}KL^{-s})^{-1}L^{-s}K^*\Sigma^{-1/2}\Sigma^{-1/2}(y^{\delta} - Kx_0) + x_0$$

$$= (\alpha L^{2s} + T^*T)^{-1}T^*(\hat{y}^{\delta} - Tx_0) + x_0, \text{ where } T = \Sigma^{-1/2}K, \quad \tilde{y}^{\delta} = \Sigma^{-1/2}y^{\delta}$$

and it is equal to the minimizer, with respect to x, of the Tikhonov functional

$$||\tilde{y}^{\delta} - Tx||^2 + \alpha ||x - x_0||_s^2$$

The quantities $\Sigma^{-1/2}y^{\delta}$ and $\Sigma^{-1/2}U$ have to be interpreted in the sense of weak distributions, see Kuo (1975) and Bissantz *et al.* (2007). In the IV regression estimation (see Example 1 for the notation), $\Sigma^{-1/2}y^{\delta} = \hat{\mathbf{E}}(Y|W)\hat{f}_W^{\frac{1}{2}}$, $Tx = \hat{\mathbf{E}}(x|W)\hat{f}_W^{\frac{1}{2}}$ and the posterior mean estimator writes

$$\hat{x}_{\alpha} = \left(\alpha L^{2s} + \tilde{K}^* \frac{1}{\hat{f_W}} \tilde{K}\right)^{-1} \tilde{K}^* \left(\hat{E}(Y - x_0|W)\right) = \left(\alpha L^{2s} + \hat{\mathbf{E}} \left[\hat{\mathbf{E}}(\cdot|W)|Z\right]\right)^{-1} \hat{\mathbf{E}} \left[\hat{\mathbf{E}}(Y - x_0|W)|Z\right].$$

For $x_0 = 0$ this is, in the framework of Darolles *et al.* (2011) or Florens *et al.* (2011), the Tikhonov estimator in the Hilbert scale generated by $\frac{1}{f_Z}L^{2s}$. \square

3.1 Asymptotic analysis

We analyze now frequentist asymptotic properties of the posterior distribution μ_{δ}^{Y} of x. The asymptotic analysis is for $\delta \to 0$. Let P^{x_*} denote the sampling distribution (5) with $x = x_*$, we remind the definition of *posterior consistency*, see Diaconis and Freedman (1986) or Ghosh and Ramamoorthi (2003):

Definition 1. The posterior distribution is consistent at x_* with respect to P^{x_*} if it weakly converges towards the Dirac measure δ_{x_*} at the point x_* , i.e. if, for every neighborhood \mathcal{U} of x_* , $\mu_{\delta}^{Y}(\mathcal{U}|y^{\delta},\alpha,s) \to 1$ in P^{x_*} -probability or P^{x_*} -a.s. as $\delta \to 0$.

Posterior consistency provides the basic frequentist validation of our Bayesian procedure because it ensures that with a sufficiently large amount of data, it is almost possible to recover the truth accurately. Lack of consistency is extremely undesirable, and one should not use a Bayesian procedure if the corresponding posterior distribution is inconsistent.

Our asymptotic analysis is organized as follows. First, we consider the posterior mean \hat{x}_{α} as an estimator for the solution of (4) and study the rate of convergence of the associated risk. Second, we state posterior consistency and recover the rate of contraction of the posterior distribution. We denote the risk (MISE) associated with \hat{x}_{α} by $\mathbf{E}_{x_*}||\hat{x}_{\alpha}-x_*||^2$ where \mathbf{E}_{x_*} denotes the expectation taken with respect to P^{x_*} . Let $\{\lambda_{jL}\}$ denote the eigenvalues of L^{-1} , we define:

$$\gamma := \inf\{\tilde{\gamma} \in (0,1] \; ; \; \sum_{j=1}^{\infty} \lambda_{jL}^{2\tilde{\gamma}(a+s)} < \infty\} \equiv \inf\{\tilde{\gamma} \in (0,1] \; ; \; tr(L^{-2\tilde{\gamma}(a+s)}) < \infty\}.$$
 (10)

We point out that γ is known since it depends on L. For instance, $\gamma = \frac{1}{2(a+s)}$ means that either L^{-1} is trace-class but $L^{-(1-\omega)}$ is not trace-class for every $\omega \in \mathbb{R}_+$ or that $tr\left(L^{-(1+\omega)}\right) < \infty$, $\forall \omega \in \mathbb{R}_+$ but L^{-1} is not trace-class. Since under A.2 the operator L^{-2s} is trace-class, the parameter γ cannot be larger than $\frac{s}{a+s}$. Thus, if $\gamma = \frac{1}{2(a+s)}$ this implies that $s \geq \frac{1}{2}$ since $\frac{1}{2(a+s)}$ must be less than or equal to $\frac{s}{(a+s)}$. Remark that the smaller the γ is and the smaller the eigenvalues of L^{-1}

are. Furthermore, we denote by $\mathcal{X}_{\beta}(\Gamma)$ the ellipsoid of the type

$$\mathcal{X}_{\beta}(\Gamma) := \{ \varphi \in \mathcal{X} \; ; \; ||\varphi||_{\beta}^{2} \leq \Gamma \}, \qquad 0 < \Gamma < \infty$$
 (11)

where $||\varphi||_{\beta} := ||L^{\beta}\varphi||$. Our asymptotic results will be valid uniformly on $\mathcal{X}_{\beta}(\Gamma)$. The following theorem gives the asymptotic behavior of \hat{x}_{α} .

Theorem 2. Let us consider the observational model (4) with x_* being the true value of x that generates the data. Under Assumptions A.1-A.3 and B, we have

$$\sup_{(x_*-x_0)\in\mathcal{X}_{\beta}(\Gamma)} \mathbf{E}_{x_*} ||\hat{x}_{\alpha} - x_*||^2 = \mathcal{O}\left(\alpha^{\frac{\tilde{\beta}}{a+s}} + \delta\alpha^{-\frac{a+\gamma(a+s)}{a+s}}\right)$$

with $\tilde{\beta} = \min(\beta, a + 2s)$.

The theorem provides an uniform rate over the ellipsoid $\mathcal{X}_{\beta}(\Gamma)$. It follows from this result that if $(x_* - x_0)$ satisfies Assumption B then $\mathbf{E}_{x_*} ||\hat{x}_{\alpha} - x_*||^2 = \mathcal{O}\left(\alpha^{\frac{\tilde{\beta}}{a+s}} + \delta\alpha^{-\frac{a+\gamma(a+s)}{a+s}}\right)$.

The value a + 2s in the theorem plays the role of the *qualification* in a classical regularization scheme, that is, a + 2s is the maximum degree of regularity of x_* that can be exploited in order to improve the rate. It is equal to the qualification of a Tikhonov regularization in the Hilbert scale $(\mathcal{X}_s)_{s \in \mathbb{R}}$, see e.g. Engl et al. (2000) section 8.5.

The value of α that minimizes the rate given in Theorem 2 is: $\alpha^{\min} := \kappa \delta^{\frac{a+s}{\beta+a+\gamma(a+s)}}$, for some constant $\kappa > 0$. When α is set equal to α^{\min} (i.e. the prior is scaling), then $\sup_{(x_*-x_0)\in\mathcal{X}_{\beta}(\Gamma)} \mathbf{E}_{x_*} || \hat{x}_{\alpha^{\min}} - x_* ||^2 = \mathcal{O}\left(\delta^{\frac{\beta}{\beta+a+\gamma(a+s)}}\right)$ and this rate is equal to the minimax rate $\delta^{\frac{\beta}{\beta+a+\frac{1}{2}}}$ if $\beta \leq a+2s$ and $\gamma = \frac{1}{2(a+s)}$ (which is possible only if $s \geq \frac{1}{2}$)⁴. In this case $(\alpha = \alpha^{\min})$, the prior is: (i) non-scaling if $s = \beta + \frac{1}{2}$ (i.e. the regularity of the prior matches the regularity of the truth); (ii) spreading out if $s > \beta + \frac{1}{2}$; (ii) shrinking if $s < \beta + \frac{1}{2}$. This means that when $s \neq \beta + \frac{1}{2}$ in order to achieve the minimax rate of convergence the prior must either spread out to become rougher (case (iii)) or shrink to become smoother (case (iii)).

When $\alpha \neq \alpha^{\min}$ the rate is slower than $\delta^{\frac{\beta}{\beta+a+\frac{1}{2}}}$ but we still have consistency provided that we set $\alpha \approx \delta^{\epsilon}$ for $0 < \epsilon < \frac{a+s}{a+\gamma(a+s)}$. Remark that since $tr(L^{-2s}) < \infty$ under A.2 then $\gamma \leq \frac{s}{a+s}$ so that $\frac{a+s}{a+\gamma(a+s)} \geq 1$. Thus, 1 is a possible value for ϵ which implies that consistency is always obtained with a non-scaling prior even if the minimax rate is obtained only in particular cases.

The same discussion can be made concerning the rate of contraction of the posterior distribution

³The results of Theorems 2 and 3 hold more generally for the case where K and Σ depend on δ , *i.e.* they are estimated. To not additionally burden the paper we do not provide these results for the general case but only for the IV regression estimation, see Corollary 1 below.

⁴Remark that a γ smaller than $\frac{1}{2(a+s)}$ implies that $tr(L^{-1}) < tr(L^{1/g}) < \infty$ for some g > 1. This means that L^{-1} is very smooth and its spectrum is decreasing fast. Thus, if $(x_* - x_0)$ is not very smooth then Assumption B will be satisfied only with a β very small. A small β will decrease the rate of convergence.

which is given in the next theorem.

Theorem 3. Let the assumptions of Theorem 2 be satisfied. For any sequence $M_{\delta} \to \infty$ the posterior distribution satisfies

$$\mu_{\delta}^{Y}\{x \in \mathcal{X} : ||x - x_{*}|| > \varepsilon_{\delta}M_{\delta}\} \to 0$$

in
$$P^{x_*}$$
-probability as $\delta \to 0$, where $\varepsilon_{\delta} = \left(\alpha^{\frac{\tilde{\beta}}{2(a+s)}} + \delta^{\frac{1}{2}}\alpha^{-\frac{a+\gamma(a+s)}{2(a+s)}}\right)$ and $\tilde{\beta} = \min(\beta, a+2s)$.

We refer to ε_{δ} as the rate of contraction of the posterior distribution. If the prior is fixed, that is, $\alpha \asymp \delta$, then $\varepsilon_{\delta} = \delta^{\frac{\beta \wedge (s - \gamma(a + s))}{2(a + s)}}$. If α is chosen such that the two terms in ε_{δ} are balanced, that is, $\alpha \asymp \alpha^{\min}$, then $\varepsilon_{\delta} = \delta^{\frac{\beta}{2\tilde{\beta} + 2a + 2\gamma(a + s)}}$. The minimax rate $\delta^{\frac{\beta}{2\beta + 2a + 1}}$ is obtained when $\beta \le a + 2s$, $\gamma = \frac{1}{2(a + s)}$ and we set $\alpha = \alpha^{\min}$. In this case, the prior is either fixed or scaling depending whether s equates $\beta + \frac{1}{2}$ or not.

3.2 Example 1: Instrumental variable (IV) regression estimation (continued)

In this section we explicit the rate of Theorem 2 for the IV regression estimation. Remark that $B^*B = \Omega_0^{\frac{1}{2}} \tilde{K}^* [\sigma^2 \hat{f}_W]^{-1} \tilde{K} \Omega_0^{\frac{1}{2}}$ where $\hat{f}_W = \frac{1}{n} \sum_{i=1}^n \frac{K_{W,h}(w_i-w)}{h}$, \tilde{K} has been defined before display (8) and $\tilde{K}^* : L^2[0,1] \to L^2[0,1]$, is the adjoint operator of \tilde{K} that takes the form: $\tilde{K}^*\phi = \frac{1}{n} \sum_{i=1}^n \frac{K_{Z,h}(z_i-z)}{h} < \phi$, $\frac{K_{W,h}(w_i-w)}{h} >$, $\forall \phi \in L^2[0,1]$. The Bayesian estimator of the IV regression is $\hat{x}_\alpha = \Omega_0^{\frac{1}{2}} (\alpha I + B^*B)^{-1} (\Sigma^{-1/2}B)^* (y^\delta - \Lambda \tilde{K} x_0)$. We assume that the true IV regression x_* that generates the data satisfies Assumption B. In order to determine the rate of the MISE associated with \hat{x}_α the proof of Theorem 2 must be slightly modified. This is because the covariance operator of U and K in the inverse problem associated with the IV regression estimation depend on n. Therefore, the rate of the MISE must incorporate the rate of convergence of these operators towards their limits. The crucial issue in order to establish the rate of the MISE associated with \hat{x}_α is the rate of convergence of B^*B towards $\Omega_0^{\frac{1}{2}} \mathcal{R}^* \frac{1}{\sigma^2 f_W} \mathcal{R} \Omega_0^{\frac{1}{2}}$ where $\mathcal{R} = \lim_{n \to \infty} \tilde{K} = f_W \mathbf{E}(\cdot|W)$ and $\mathcal{R}^* = \lim_{n \to \infty} \tilde{K}^* = f_Z \mathbf{E}(\cdot|Z)$. This rate is specified by Assumption (HS) below and we refer to Darolles et al. (2011, Appendix B) for a set of sufficient conditions that justify this assumption.

Assumption HS. There exists $\rho \geq 2$ such that:

(i).
$$\mathbf{E}||B^*B - \Omega_0^{\frac{1}{2}} \mathfrak{K}^* \frac{1}{\sigma^2 f_W} \mathfrak{K} \Omega_0^{\frac{1}{2}}||^2 = \mathcal{O}\left(n^{-1} + h^{2\rho}\right);$$

(ii).
$$\mathbf{E}||\Omega_0^{\frac{1}{2}}(\hat{\mathcal{T}}^* - \mathcal{T}^*)||^2 = \mathcal{O}(\frac{1}{n} + h^{2\rho}), \text{ where } \mathcal{T}^* = \mathbf{E}(\cdot|Z) \text{ and } \hat{\mathcal{T}}^* = \hat{\mathbf{E}}(\cdot|Z).$$

To get rid of σ^2 it is sufficient to specify Ω_0 as $\Omega_0\sigma^2$ so that B^*B does not depend on σ^2 anymore and we do not need to estimate it to get the estimate of x. The next corollary to Theorem 2 gives the rate of the MISE of \hat{x}_{α} under this assumption.

Corollary 1. Let us consider the observational model $y^{\delta} = Kx_* + U$, with y^{δ} and K defined as in (8) and x_* being the true value of x. Under Assumptions A.1-A.3, B and B:

$$\sup_{(x_* - x_0) \in \mathcal{X}_{\beta}(\Gamma)} \mathbf{E}_{x_*} ||\hat{x}_{\alpha} - x_*||^2 = \mathcal{O}\left(\left(\alpha^{\frac{\tilde{\beta}}{a+s}} + n^{-1}\alpha^{-\frac{a+\gamma(a+s)}{a+s}}\right) \left(1 + \alpha^{-2}\left(\frac{1}{n} + h^{2\rho}\right)\right) + \alpha^{-2}\left(\frac{1}{n} + h^{2\rho}\right)\left(\frac{1}{nh} + h^{2\rho}\right)\right)$$

with $\tilde{\beta} = \min(\beta, a + 2s)$.

In this example we have to set two tuning parameters: the bandwidth h and α . We can set h such that $h^{2\rho}$ goes to 0 at least as fast as n^{-1} and $\alpha = \alpha^{\min} \propto n^{-\frac{a+s}{\beta+a+\gamma(a+s)}}$. With this choice, the rate of convergence of B^*B towards $\Omega_0^{\frac{1}{2}}\mathfrak{K}^*\frac{1}{f_W}\mathfrak{K}\Omega_0^{\frac{1}{2}}$ and of $\Omega_0^{\frac{1}{2}}\hat{\mathcal{T}}^*$ towards $\Omega_0^{\frac{1}{2}}\mathcal{T}^*$ will not affect the rate in the MISE (that is, the rate will be the same as the rate given in Theorem 2) if $\tilde{\beta} \geq 2s + a - 2\gamma(a+s)$ and $\rho > \frac{\tilde{\beta}+a+\gamma(a+s)}{\tilde{\beta}+2\gamma(a+s)-2s}$. We remark that when the prior is not scaling, *i.e.* $\alpha \approx n^{-1}$, the condition $[\alpha^2 n]^{-1} = \mathcal{O}(1)$ is not satisfied. The rate of Corollary 1, with $\alpha = \alpha^{\min}$ and $h = \mathcal{O}(n^{-1/(2\rho)})$, is minimax when $a + 2s \geq \beta \geq a + 2s - \frac{1}{2}$ and $\gamma = \frac{1}{2(a+s)}$.

4 Operators with geometric spectra

We analyze now the important case where the inverse problem (4) is mildly ill-posed. We denote with λ_{jK} the singular values of K and with $\lambda_{j\Sigma}$ and λ_{jL} the eigenvalues of Σ and L^{-1} , respectively. Assumption C states that the operators $\Sigma^{-1/2}$ and KK^* (resp. $K^*\Sigma^{-1}K$ and L^{-1}) are diagonalizable in the same eigenbasis and have polynomially decreasing spectra.

Assumption C. The operator $\Sigma^{-1/2}$ (resp. $K^*\Sigma^{-1}K$) has the same eigenfunctions $\{\varphi_j\}$ as KK^* (resp. $\{\psi_j\}$ as L^{-1}). Moreover, the eigenvalues of KK^* , Σ and L^{-1} satisfy

$$\underline{a}j^{-2a_0} \leq \lambda_{jK}^2 \leq \overline{a}j^{-2a_0}, \quad \underline{c}j^{-c_0} \leq \lambda_{j\Sigma} \leq \overline{c}j^{-c_0} \quad and \quad \underline{l}j^{-1} \leq \lambda_{jL} \leq \overline{l}j^{-1}, \quad j = 1, 2, \dots$$

with $a_0 \ge 0, c_0 > 1$ and $\underline{a}, \overline{a}, \underline{c}, \overline{c}, \underline{l}, \overline{l} > 0$.

This assumption implies that KK^* and K^*K are strictly positive definite. In this section we provide the exact rate attained by \hat{x}_{α} in the setting described by Assumption C.

Assumption C is standard in statistical inverse problems literature (see e.g. Assumption A.3 in Hall and Horowitz (2005) or Assumption B3 in Cavalier and Tsybakov (2002)) and, by using the notation defined in section 2.1, it may be rewritten as $\lambda_{jK} \asymp j^{-a_0}$, $\lambda_{j\Sigma} \asymp j^{-c_0}$ and $\lambda_{jL} \asymp j^{-1}$. Under Assumption C we may write $\mathcal{X}_{\beta}(\Gamma)$ as $\mathcal{X}_{\beta}(\Gamma) := \{\varphi \in \mathcal{X} : \sum_{j} j^{2\beta} < \varphi, \psi_{j} >^{2} \le \Gamma \}$ and the parameter γ is equal to $\frac{1}{2(a+s)}$. The following proposition provides necessary and sufficient conditions for the validity of A.1, A.2 and A.3 when C is satisfied.

Proposition 1. Under Assumption C: (i) A.1 is satisfied if and only if $a_0 \ge \frac{c_0}{2}$; (ii) A.2 is satisfied if and only if $a = a_0 - \frac{c_0}{2} > 0$ and $s > \frac{1}{2}$; (iii) A.3 is satisfied if and only if $a_0 \ge c_0 - s$.

The following proposition gives the minimax rate attained by \hat{x}_{α} under Assumption C.

Proposition 2. Let B, C hold, $a = a_0 - \frac{c_0}{2} > 0$, $s > \frac{1}{2}$ and $a_0 \ge c_0 - s$. Then we have

$$\sup_{(x_* - x_0) \in \mathcal{X}_{\beta}(\Gamma)} \mathbf{E}_{x_*} ||\hat{x}_{\alpha} - x_*||^2 \approx \alpha^{\frac{\tilde{\beta}}{a+s}} + \delta \alpha^{-\frac{2a+1}{2(a+s)}}, \tag{12}$$

with $\tilde{\beta} = \min(\beta, 2(a+s))$. Moreover, (i) for $\alpha \approx \delta$ (fixed prior),

$$\sup_{(x_*-x_0)\in\mathcal{X}_{\beta}(\Gamma)} \mathbf{E}_{x_*} ||\hat{x}_{\alpha} - x_*||^2 \asymp \delta^{\frac{\beta \wedge \left(s - \frac{1}{2}\right)}{a + s}};$$

(ii) for $\alpha \simeq \delta^{\frac{a+s}{\beta+a+\frac{1}{2}}}$ (optimal rate of α),

$$\sup_{(x_*-x_0)\in\mathcal{X}_{\beta}(\Gamma)} \mathbf{E}_{x_*} ||\hat{x}_{\alpha} - x_*||^2 \asymp \delta^{\frac{\tilde{\beta}}{\tilde{\beta}+a+\frac{1}{2}}}.$$

The minimax rate of convergence over a Sobolev ellipsoid $\mathcal{X}_{\beta}(\Gamma)$ is of the order $\delta^{\frac{\beta}{\beta+a+\frac{1}{2}}}$. By the results of the proposition, the uniform rate of the MISE associated with \hat{x}_{α} is minimax if: either the parameter s of the prior is chosen such that $s=\beta+\frac{1}{2}$ and the prior is fixed (case (i)) or if $\beta \leq 2(a+s)$ and the prior is scaling at the optimal rate (case (ii), i.e. $\alpha \approx \delta^{(a+s)/(\beta+a+\frac{1}{2})}$). In all the other cases the rate is slower than the minimax rate but consistency is still verified provided that $\alpha \approx \delta^{\epsilon}$ for $0 < \epsilon < \frac{a+s}{a+\frac{1}{2}}$. Remark that since $s > \frac{1}{2}$ then a fixed prior $(\epsilon = 1)$ always guarantees consistency (even if the rate is not always minimax).

This result, similar to that one in Theorem 2, means that when the prior is "correctly specified" ("correct" in the sense that the regularity $s-\frac{1}{2}$ of the trajectories generated by the prior is the same as the regularity of x_*) we obtain a minimax rate without scaling the prior covariance. On the other hand, if $s < \beta + \frac{1}{2}$, i.e. the prior is "undersmoothing", the minimax rate can still be achieved as soon as $\beta \leq 2(a+s)$ and the prior is shrinking at the optimal rate. When the prior is "oversmoothing", i.e. $s > \beta + \frac{1}{2}$, the minimax rate can be achieved if the prior distribution spreads out at the optimal α (the prior has to be more and more dispersed in order to become rougher).

In many cases it is reasonable to assume that the functional parameter x_* has generalized Fourier coefficients (in the basis made of the eigenfunctions $\{\psi_j\}$ of L^{-s}) that are geometrically decreasing, see e.g. Assumption A.3 in Hall and Horowitz (2005) and Theorem 4.1 in Van Rooij and Ruymgaart (1999). Thus, we may consider the following assumption instead of Assumption B.

Assumption B'. For some $b_0 > \frac{1}{2}$ and $\{\psi_j\}$ defined in Assumption $C, <(x_* - x_0), \psi_j > \approx j^{-b_0}$.

Assumption B' is often encountered in statistical inverse problems literature. Assumption B is more general than Assumption B' since it allows to consider the important case of exponentially declining Fourier coefficients $<(x_*-x_0), \psi_j>$. We use Assumption B' to show sharp adaptiveness for our Empirical Bayes procedure. If B' holds then $\exists \Gamma<\infty$ such that Assumption B holds for some $0 \le \beta < b_0 - \frac{1}{2}$. The following result gives the rate of the MISE when Assumption B' holds.

Proposition 3. Let B', C hold with $a = a_0 - \frac{c_0}{2} > 0$, $s > \frac{1}{2}$ and $a_0 \ge c_0 - s$. Then,

$$\alpha^{\frac{2b_0-1}{2(a+s)}}c_1 + \delta\alpha^{-\frac{2a+1}{2(a+s)}}c_2(\bar{t}) \leq \mathbf{E}_{x_*}||\hat{x}_\alpha - x_*||^2 \leq \alpha^2 I(b_0 \geq (2a+2s)) + \alpha^{\frac{2b_0-1}{2(a+s)}}\tilde{c}_1 + \delta\alpha^{-\frac{2a+1}{2(a+s)}}\tilde{c}_2(\bar{t}) \quad (13)$$

where c_1 , \tilde{c}_1 , c_2 , \tilde{c}_2 and \bar{t} are positive constants and $I(\mathcal{A})$ denotes the indicator function of an event \mathcal{A} . Moreover, for $\alpha \approx \delta^{\frac{a+s}{\tilde{b}_0+a}} \equiv \alpha_*$ and $\tilde{b}_0 = \min(b_0, 2a + 2s + 1/2)$,

$$\mathbf{E}_{x_*} ||\hat{x}_{\alpha_*} - x_*||^2 \approx \delta^{\frac{2\tilde{b}_0 - 1}{2(\tilde{b}_0 + a)}}. \tag{14}$$

When $b_0 \leq 2a + 2s + 1/2$ the rate of the lower bound $\alpha^{\frac{2b_0-1}{2(a+s)}} + \delta \alpha^{-\frac{2a+1}{2(a+s)}}$ given in (13) provides, up to a constant, a lower and an upper bound for the rate of the estimator \hat{x}_{α} and so it is optimal. Thus, the minimax-optimal rate $\delta^{\frac{2b_0-1}{2(b_0+a)}}$ is obtained when we set $\alpha \approx \delta^{\frac{a+s}{b_0+a}}$ if: either $s = b_0$ (fixed prior with $\alpha \approx \delta$), or $s < b_0 \leq 2a + 2s + 1/2$ (shrinking prior), or $s > b_0$ (spreading out prior). When $s < b_0$ (resp. $s > b_0$) the trajectories generated by the prior are rougher (resp. smoother) than x_* and so the support of the prior must be shrunk (resp. spread out). When $\delta = n^{-1}$ the rate $n^{-\frac{(2b_0-1)}{2(b_0+a)}}$ is shown to be minimax in Hall and Horowitz (2007).

Moreover, if we set $\beta = \sup \left\{ \check{\beta} \geq 0; \ (x_* - x_0) \text{ satisfies B' and } \sum_j j^{2\check{\beta}} < (x_* - x_0), \psi_j >^2 < \infty \right\}$ then $\beta = b_0 - \frac{1}{2}$ and the rate $\delta^{\frac{2\check{b}_0 - 1}{2(\check{b}_0 + a)}}$ is uniform in x_* over $\mathcal{X}_{\beta}(\Gamma)$ and equal to the optimal rate of proposition 2.

In the following theorem we give the rate of contraction of the posterior distribution under Assumption C.

Theorem 4. Let the assumptions of Proposition 3 be satisfied. For any sequence $M_{\delta} \to \infty$ the posterior probability satisfies

$$\mu_{\delta}^{Y} \{ x \in \mathcal{X} : ||x - x_{*}|| > \varepsilon_{\delta} M_{\delta} \} \to 0$$

in P^{x_*} -probability as $\delta \to 0$, where $\varepsilon_{\delta} = \left(\alpha^{\frac{2\tilde{b}_0 - 1}{4(a+s)}} + \delta^{\frac{1}{2}}\alpha^{-\frac{2a+1}{4(a+s)}}\right)$, $\alpha > 0$, $\alpha \to 0$ and $\tilde{b}_0 = \min(\beta, 2(a+s) + 1/2)$. Moreover, (i) for $\alpha \times \delta$ (fixed prior):

$$\varepsilon_{\delta} = \delta^{\frac{2(b_0 \wedge s) - 1}{4(a+s)}}$$

and (ii) for $\alpha \approx \delta^{\frac{a+s}{\overline{b_0}+a}}$ (optimal rate of α):

$$\varepsilon_{\delta} = \delta^{\frac{2\tilde{b}_0 - 1}{4(\tilde{b}_0 + a)}}.$$

The rate of contraction ε_{δ} is equal to the minimax rate $\delta^{\frac{2b_0-1}{4(b_0+a)}}$ if $b_0 \leq 2(a+s)+1/2$ and $\alpha \approx \delta^{\frac{a+s}{b_0+a}}$. Depending on the relation between s and b_0 the corresponding prior μ is shrinking, spreading out or fixed, see comments after proposition 3.

4.1 Example 2: Functional Linear Regression Estimation (continued)

We develop a little further Example 2. Here, the covariance operator Σ is proportional to K (as shown in section 2.2.1): $\delta\Sigma = \frac{\tau^2}{n}K$. The operator $K: L^2([0,1]) \to L^2([0,1])$, defined as $\forall \varphi \in L^2([0,1])$, $K\varphi := \frac{1}{n}\sum_i < Z_i, \varphi > Z_i(t)$, is self-adjoint, i.e. $K = K^*$, and depends on n. It converges to the operator \tilde{K} , defined as $\forall \varphi \in L^2([0,1])$, $\tilde{K}\varphi = \int_0^1 \varphi(s)Cov(Z(s),Z(t))ds$, which is trace-class since $\mathbf{E}||Z_i||^2 < \infty$. Choose $L = \tilde{K}^{-1}$, and suppose that the spectrum of \tilde{K} declines at the rate j^{-a_0} , $a_0 \geq 0$, as in Assumption C (for instance, \tilde{K} could be the covariance operator of a Brownian motion). Then, $j^{-a_0} = j^{-c_0}$, $a_0 > 1$ and $a_0 = c_0$. Moreover, we set $\Omega_0 = K^{2\tilde{s}}$, for some $\tilde{s} > \frac{1}{2a_0}$ (so that $\lim_{n\to\infty} \Omega_0 = \tilde{K}^{2\tilde{s}}$), and assume that x_* satisfies Assumption B'. The posterior mean estimator takes the form: $\hat{x}_\alpha = K^{\tilde{s}}(\alpha I + K^{2\tilde{s}+1})^{-1}K^{\tilde{s}}(y^\delta - Kx_0) + x_0$, for which the following lemma holds.

Lemma 1. Let $\tilde{K}: L^2([0,1]) \to L^2([0,1])$ have eigenvalues $\{\lambda_{jK}\}$ that satisfy $\underline{a}j^{-a_0} \le \lambda_{jK} \le \overline{a}j^{-a_0}$, for $\underline{a}, \overline{a} > 0$ and $a_0 \ge 0$. Assume that $\mathbf{E}||Z_i||^4 < \infty$. Then, under Assumption B' with $b_0 > \max\{a_0, a_0\tilde{s}\}, \tilde{s} > \frac{1}{2a_0}$ and if $\alpha \asymp \alpha_* = n^{-\frac{a_0(2\tilde{s}+1)}{2\tilde{b}_0+a_0}}$, we have

$$\mathbf{E}_{x_*}||\hat{x}_{\alpha_*} - x_*||^2 = \mathcal{O}\left(n^{-\frac{2\tilde{b}_0 - 1}{2\tilde{b}_0 + a_0}}\right), \quad where \quad \tilde{b}_0 = \min\left(b_0, a_0 + 2a_0\tilde{s} + \frac{1}{2}\right).$$

The assumptions and the rate for $\tilde{b}_0 = b_0$ in the lemma are the same as in Hall and Horowitz (2007).

5 Adaptive selection of α : an Empirical Bayes (EB) approach

As shown by Theorem 1 (iii) and Remark 3.1, the parameter α of the prior plays the role of a regularization parameter and $\{\hat{x}_{\alpha}\}_{{\alpha}\geq 0}$ defines a class of estimators for x, which are equal to a Tikhonov-type frequentist estimator. We have shown that for α decreasing at a convenient rate, this estimator converges at the minimax rate. However, this rate and the corresponding value for α are unknown in practice since they depend on the regularity of x_* which is unknown. Thus, it is

very important to have an adaptive data-driven method for selecting α since a suitable value for α is crucial for the implementation of the estimation procedure. We say that a data-driven method selects a "suitable value" $\hat{\alpha}$ for α if the posterior distribution of x computed by using this $\hat{\alpha}$ still satisfies consistency in a frequentist sense at an (almost) minimax rate.

We propose in sections 5.1 and 5.2 a data-driven method based on an EB procedure for selecting $\hat{\alpha}$. This procedure can be easily implemented for general operators K, Σ and L satisfying Assumptions A.1 and A.3.

5.1 Characterization of the likelihood

The marginal distribution of $y^{\delta}|\alpha, s$, which was given in Theorem 1 (i), is

$$y^{\delta}|\alpha, s \sim P_{\alpha}, \quad P_{\alpha} = \mathcal{N}\left(Kx_0, \delta\Sigma + \frac{\delta}{\alpha}K\Omega_0K^*\right)$$
 (15)

and is obtained by marginalizing P^x with respect to the prior distribution of x. The following theorem, which is an application of Theorem 3.3 in Kuo (1975), characterizes a probability measure P_0 which is equivalent to P_{α} for every $\alpha > 0$ and the likelihood of P_{α} with respect to P_0 .

Theorem 5. Let P_0 be a Gaussian measure with mean Kx_0 and covariance operator $\delta\Sigma$, that is, $P_0 = \mathcal{N}(Kx_0, \delta\Sigma)$. Under Assumptions A.1 and A.3, the Gaussian measure P_α defined in (15) is equivalent to P_0 . Moreover, the Radon-Nikodym derivative is given by

$$\frac{dP_{\alpha}}{dP_0}(\{z_j\}) = \prod_{j=1}^{\infty} \sqrt{\frac{\alpha}{\lambda_j^2 + \alpha}} e^{\frac{\lambda_j^2}{2(\lambda_j^2 + \alpha)} z_j^2},$$
(16)

where $\{\varphi_j, \lambda_j^2\}$ are the eigenfunctions and eigenvalues of BB^* , respectively.

In our setting: $z_j = \frac{\langle y^{\delta} - Kx_0, \Sigma^{-1/2}\varphi_j \rangle}{\sqrt{\delta}}$ and $\Sigma^{-1/2}\varphi_j$ is defined under Assumption A.3.

5.2 Adaptive EB procedure

Let ν denote a prior distribution for α such that $\frac{d \log \nu(\alpha)}{d\alpha} = \nu_1 \alpha^{-1} + \nu_2$ for two constants $\nu_1 > 0$, $\nu_2 < 0$. An EB procedure consists in plugging in the prior distribution of $x | \alpha, s$ a value for α selected from the data y^{δ} . We define the marginal maximum a posteriori estimator $\hat{\alpha}$ of α to be the maximizer of the marginal log-posterior of $\alpha | y^{\delta}, s$ which is proportional to $\log \left[\frac{dP_{\alpha}}{dP_0} \nu(\alpha) \right]$:

$$\hat{\alpha} := \arg \max_{\alpha} \bar{S}(\alpha, y^{\delta}) \tag{17}$$

$$\bar{S}(\alpha, y^{\delta}) := \log \left(\frac{dP_{\alpha}}{dP_{0}} (\{z_{j}\}) \nu(\alpha) \right)
= \frac{1}{2} \sum_{j=1}^{\infty} \left[\log \left(\frac{\alpha}{\alpha + \lambda_{j}^{2}} \right) + \frac{\lambda_{j}^{2}}{\alpha + \lambda_{j}^{2}} \frac{\langle y^{\delta} - Kx_{0}, \Sigma^{-1/2} \varphi_{j} \rangle^{2}}{\delta} \right] + \log \nu(\alpha).$$
(18)

In an equivalent way, $\hat{\alpha}$ is defined as the solution of the first order condition $\frac{\partial}{\partial \alpha} \bar{S}(\alpha, y^{\delta}) = 0$. We denote $S_{y^{\delta}}(\alpha) := \frac{\partial}{\partial \alpha} \bar{S}(\alpha, y^{\delta})$ where

$$S_{y^{\delta}}(\alpha) := \frac{1}{2} \left(\sum_{j=1}^{\infty} \frac{\lambda_j^2}{\alpha(\alpha + \lambda_j^2)} - \sum_{j=1}^{\infty} \frac{\lambda_j^2 < y^{\delta} - Kx_0, \Sigma^{-1/2} \varphi_j >^2}{\delta(\alpha + \lambda_j^2)^2} \right) + \frac{d \log \nu(\alpha)}{d\alpha} \quad \text{and} \quad S_{y^{\delta}}(\hat{\alpha}) = 0.$$

$$\tag{19}$$

Our strategy will then be to plug the value $\hat{\alpha}$, found in this way, back into the prior of $x|\alpha, s$ and then compute the posterior distribution $\mu_{\delta,\hat{\alpha}}^Y$ and the posterior mean estimator $\hat{x}_{\hat{\alpha}}$ using this value of α . We refer to $\hat{x}_{\hat{\alpha}}$ as the EB-estimator and to $\mu_{\delta,\hat{\alpha}}^Y$ as the EB-posterior. Examples of suitable priors for α are: either a gamma distribution or a beta distribution on (0,1).

5.3 Posterior consistency for the EB-posterior (mildly ill-posed case)

In this section we study existence of $\hat{\alpha}$ and the rate at which $\hat{\alpha}$ decreases to 0 and show posterior consistency of $\mu_{\delta,\hat{\alpha}}^{Y}$. When the true x_* is not too smooth, with respect to the degree of ill-posedness of the problem and the smoothness of the prior, then $\hat{\alpha}$ is of the same order as α_* ($\equiv \delta^{\frac{a+s}{b_0+a}}$) with probability that approaches 1 as $\delta \to 0$. Moreover, we state that the EB-posterior distribution $\mu_{\delta,\hat{\alpha}}^{Y}$ concentrates around the true x_* in probability. In the next theorem, let $I(\mathcal{A})$ denote the indicator function of an event \mathcal{A}

Theorem 6. Let B' and C hold with $a=a_0-\frac{c_0}{2}>0$, $s>\frac{1}{2}$ and $a_0\geq c_0-s$. Let ν be a prior distribution for α such that $\frac{d\log\nu(\alpha)}{d\alpha}=\nu_1\alpha^{-1}+\nu_2$ for two constants $\nu_1>0$, $\nu_2<0$. Then, with probability approaching 1, a solution $\hat{\alpha}$ to the equation $S_{y^{\delta}}(\alpha)=0$ exists and is of order: $\hat{\alpha}\asymp \delta^{\frac{a+s}{b_0+a+\tilde{\gamma}}}$ for $\tilde{\eta}=\eta I(b_0-a-2s-1/2>0)$ and any $(b_0+a)>\eta>\max\{(b_0-s-1/2),0,(b_0-2s-a-1/2)\frac{1}{4}\}$. Moreover, for any sequence $M_\delta\to\infty$ the EB-posterior distribution satisfies

$$\mu_{\delta,\hat{\alpha}}^{Y} \{ x \in \mathcal{X} : ||x - x_*|| > \varepsilon_{\delta} M_{\delta} \} \to 0$$
 (20)

in P^{x_*} -probability as $\delta \to 0$ where $\varepsilon_{\delta} = \delta^{\frac{\tilde{b}_0 - \frac{1}{2}}{2(b_0 + a + \tilde{\eta})}}$ and $\tilde{b}_0 = \min(b_0, 2(a + s) + 1/2)$.

The consistency of the EB-estimator $\hat{x}_{\hat{\alpha}}$ follows from posterior consistency of $\mu_{\delta,\hat{\alpha}}^Y$. The theorem says that the posterior contraction rate of the EB-posterior distribution is equal to the minimax rate $\delta^{\frac{2b_0-1}{4(b_0+a)}}$ when $b_0 \leq a+2s+1/2$ which is satisfied, for instance, when the prior is very smooth. In all the other cases the rate is slower. In order to have a contraction rate equal to the minimax rate when $b_0 > a+2s+1/2$ we should specify the prior ν on α depending on b_0 and a in some convenient way. However, this prior would be unfeasible in practice since b_0 is never known. For this reason we do not pursue this analysis since it would have an interest only from a theoretical

⁵For simplicity of exposition we limit this analysis to the case where K, Σ and L have geometric spectra (*mildly ill-posed case*, see section 4). It is possible to extend the result of Theorem 6 to the general case at the price of complicate much more the proof and the notation. For this reason we do not show the general result here.

point of view while the main motivation for this section is the practical implementation of our estimator.

Remark 5.1. While this theorem is stated and proved for a Gaussian error term U, this result holds also for the more general case where U is only asymptotically Gaussian. In appendix D we give some hints about how the proof should be modified in this case. Therefore, our EB-approach works also in the case where U is only approximately Gaussian.

Remark 5.2. Alternative ways to choose regularization parameters for inverse problems have been proposed in mathematics, statistics and econometrics. Most of these methods work under assumptions that are not suitable for econometrics applications. Many of these methods focus on series estimators while very little is known about adaptive procedures in Tikhonov estimation. Our methods contributes to the latter. For adaptive procedures for series estimators in the nonparametric IV framework see *e.g.* Loubes and Marteau (2013), Horowitz (2013) and references therein. Assumptions and settings of these adaptive procedures differ significantly from ours.

6 Numerical Implementation

6.1 Instrumental variable regression estimation

This section shows the implementation of our proposed estimation method for the IV regression example 1 and its finite sample properties. We simulate n = 1000 observations from the following model, which involves only one endogenous covariate Z and two instrumental variables $W = (W_1, W_2)$,

$$W_{i} = \begin{pmatrix} W_{1,i} \\ W_{2,i} \end{pmatrix} \sim \mathcal{N} \left(\begin{pmatrix} 0 \\ 0 \end{pmatrix}, \begin{pmatrix} 1 & 0.3 \\ 0.3 & 1 \end{pmatrix} \right)$$

$$v_i \sim \mathcal{N}(0, \sigma_v^2), \qquad Z_i = 0.1 w_{i,1} + 0.1 w_{i,2} + v_i$$
$$\varepsilon_i \sim \mathcal{N}(0, (0.4)^2), \qquad \eta_i = -0.5 v_i + \varepsilon_i$$
$$Y_i = x_*(Z_i) + \eta_i$$

for $i=1,\ldots,n$. Endogeneity is caused by correlation between η_i and the error term v_i affecting the covariates. The true x_* is the parabola $x_*(Z)=Z^2$. In all the simulations we have fixed $\sigma_v=0.27$. We do not transform the data to the interval [0,1] and the spaces of reference are $\mathcal{X}=\mathcal{Y}=L^2(Z)$, where $L^2(Z)$ denotes the space of square integrable functions of Z with respect to its marginal distribution. Moreover, the function a(w,v) is chosen such that $K:L^2(Z)\to L^2(Z)$ is the (estimated) double conditional expectation operator, that is, $\forall \varphi \in L^2(Z), K\varphi = \hat{\mathbf{E}}(\hat{\mathbf{E}}(\varphi|W)|Z)$

and the functional observation y^{δ} takes the form $y^{\delta}(Z) := \hat{\mathbf{E}}(\hat{\mathbf{E}}(Y|W)|Z)$. Here, $\hat{\mathbf{E}}$ denotes the estimated expectation that we compute through a Gaussian kernel smoothing estimator of the joint density of (Y, Z, W). The bandwidth for Z has been set equal to $\sqrt{\widehat{Var}(Z)}n^{-1/5}$ and in a similar way the bandwidths for Y and W.

The sampling covariance operator is $\Sigma \varphi = \sigma^2 \mathbf{E}(\mathbf{E}(\varphi|W)|Z)$, $\forall \varphi \in L^2(Z)$, under the assumption $\mathbf{E}[(Y - x_*(Z))^2|W] = \sigma^2$, and it has been replaced by its empirical counterpart. Following the discussion in section 2.2.1 we specify the prior covariance operator as $\Omega_0 = \omega_0 \sigma^2 K^s$ for $s \geq 1$ and ω_0 a fixed parameter. In this way the conditional variance σ^2 in Σ and Ω_0 simplifies and does not need to be estimated.

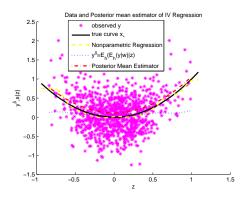
We have performed simulations for two specifications of x_0 , ω_0 and s (where x_0 denotes the prior mean function): Figure 1 refers to $x_0(Z) = 0.95Z^2 + 0.25$, $\omega_0 = 1$ and s = 2 while Figure 2 refers to $x_0(Z) = 0$, $\omega_0 = 2$ and s = 15. We have first performed simulations for a fixed value of α (we have fixed $\alpha = 0.9$ to obtain Figures 1a-1b and 2a-2b) and in a second simulation we have used the α selected through our EB method.

Graphs 1a and 2a represent: the n observed Y's (magenta asterisks), the corresponding y^{δ} (dotted blue line) obtained from the observed sample of (Y, Z, W), the true x_* (black solid line), the nonparametric estimation of the regression function $\mathbf{E}(Y|Z)$ (yellow dashed line), and our posterior mean estimator \hat{x}_{α} (dotted-dashed red line). We show the estimator of $\mathbf{E}(Y|Z)$ with the purpose of making clear the bias due to endogeneity. Graphs 1b and 2b represent: the prior mean function x_0 (magenta dashed line), the observed function y^{δ} (dotted blue line) obtained from the observed sample of (Y, Z, W), the true x_* (black solid line), and our posterior mean estimator \hat{x}_{α} (dotted-dashed red line).

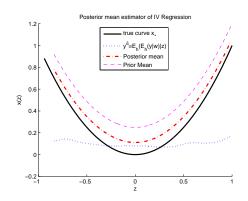
Graphs 1c and 2c draw the log-posterior $\log \left[\frac{dP_{\alpha}}{dP_0}\nu(\alpha)\right]$ against α and show the value of the maximum a posteriori $\hat{\alpha}$. We have specified an exponential prior for α : $\nu(\alpha) = 11e^{-11\alpha}$, $\forall \alpha \geq 0$. Finally, graphs 1d and 2d represent our EB-posterior mean estimator $\hat{x}_{\hat{\alpha}}$ (dotted-dashed red line) – obtained by using the $\hat{\alpha}$ selected with the EB-procedure – together with the prior mean function x_0 (magenta dashed line), the observed function y^{δ} (dotted blue line) and the true x_* (black solid line).

6.2 Geometric Spectrum case

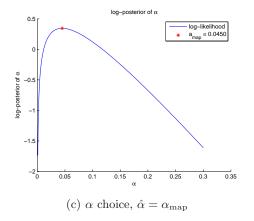
In this simulation we assume $\mathcal{X} = \mathcal{Y} = L^2(\mathbb{R})$ with respect to the measure $e^{-u^2/2}$ so that the operator K is self-adjoint. We use the Hermite polynomials as common eigenbasis for the operators K, Σ and L. The Hermite polynomials $\{H_j\}_{j\geq 0}$ form an orthogonal basis of $L^2(\mathbb{R})$ with respect to the measure $e^{-u^2/2}$. The first few Hermite polynomials are $\{1, u, (u^2 - 1), (u^3 - 3u), \ldots\}$ and an important property of these polynomials is that they are orthogonal with respect to $e^{-u^2/2}$: $\int_{\mathbb{R}} H_l(u)H_j(u)e^{-u^2/2}du = \sqrt{n}(n!)\delta_{lj}$, where δ_{lj} is equal to 1 if l=j and 0 otherwise. Moreover,



(a) Data and posterior mean estimator for $\alpha=0.9$



(b) Posterior mean estimator for $\alpha=0.9$



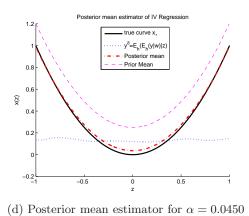
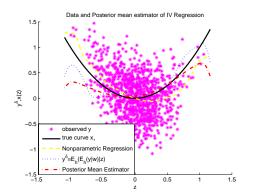
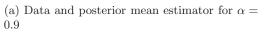
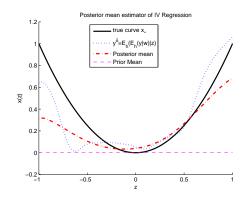


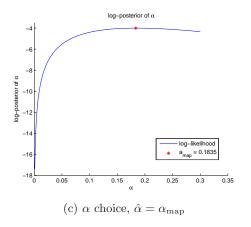
Figure 1: Posterior mean estimator for smooth x_* . Graph for: $x_0(Z) = 0.95Z^2 + 0.25$, $\omega_0 = 1$, s = 2.

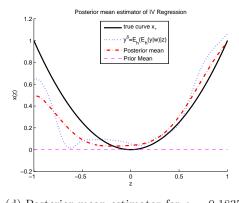






(b) Posterior mean estimator for $\alpha=0.9$





(d) Posterior mean estimator for $\alpha=0.1835$

Figure 2: Posterior mean estimator for smooth x_* . Graph for: $x_0(Z) = 0$, $\omega_0 = 2$, s = 15.

they satisfy the recursion $H_{j+1}(u) = uH_j(u) - jH_{j-1}(u)$ which is used in our simulation. We fix: $\delta = 1/n$, $a_0 = 1$, $c_0 = 1.2$ and s = 1, thus the simulation design is:

$$\Sigma \cdot = \tau \sum_{j=0}^{\infty} \frac{j^{-c_0}}{\sqrt{2\pi}(n!)} < H_j, \cdot > H_j$$

$$K \cdot = \sum_{j=0}^{\infty} \frac{j^{-a_0}}{\sqrt{2\pi}(n!)} < H_j, \cdot > H_j$$

$$\Omega \cdot = \omega_0 \sum_{j=0}^{\infty} \frac{j^{-2s}}{\sqrt{2\pi}(n!)} < H_j, \cdot > H_j$$

$$y^{\delta} = Kx_* + U$$

with $x_*(u) = u^2$ and $U = \frac{1}{\sqrt{n}} \mathcal{N}(0, \Sigma)$. Moreover, we fix: $\tau = 10$ and $\omega_0 = 5$. The inner product is approximated by discretizing the integral $\int_{\mathbb{R}} H_j(u) \cdot e^{-u^2/2} du$ with 1000 discretization points uniformly generated between -3 and 3. The infinite sums are truncated at j = 200.

We have first performed simulations for a fixed value of α (we have fixed $\alpha = 0.9$ to obtain Figure 3a) and in a second simulation we have used the α selected through our EB method.

Graph 3a represents: the prior mean function x_0 (magenta dashed line), the observed function y^{δ} (dotted blue line), the true x_* (black solid line), and our posterior mean estimator \hat{x}_{α} (dotted-dashed red line).

Graph 3b draws the log-posterior $\log \left[\frac{dP_{\alpha}}{dP_0}\nu(\alpha)\right]$ against α and shows the value of the maximum a posteriori $\hat{\alpha}$. We have specified a Gamma prior for α : $\nu(\alpha) \propto \alpha^{11}e^{-10\alpha}$, $\forall \alpha \geq 0$. Finally, graph 3c represents our EB-posterior mean estimator $\hat{x}_{\hat{\alpha}}$ (dotted-dashed red line) – obtained by using the $\hat{\alpha}$ selected with the EB-procedure – together with the prior mean function x_0 (magenta dashed line), the observed function y^{δ} (dotted blue line) and the true x_* (black solid line).

7 Conclusion

This paper develops a Bayesian approach for nonparametric estimation of parameters in econometric models that are characterized as the solution of an inverse problem. We consider a conjugate-Gaussian setting where the "likelihood" is only required to be asymptotically Gaussian. For "likelihood" we mean the sampling distribution of a functional transformation of an n-sample.

We first propose the posterior mean as a point estimator and give sufficient conditions that allow us to show that: (i) it has a closed-form, (ii) it is easy to implement in practice, (iii) it has a pure Bayesian interpretation and (iv) it is consistent in a frequentist sense. Then, we contribute to the literature on inverse problems by proposing an adaptive data-driven method to select the regularization parameter which enters our point estimator. Our data-driven method is based on an Empirical Bayes approach.

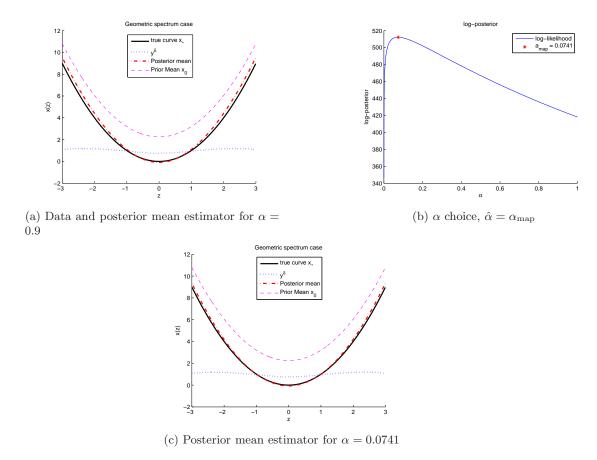


Figure 3: Posterior mean estimator for smooth x_* . Graph for: $x_0(Z) = 0.95Z^2 + 0.25$, $\omega_0 = 5$, s = 1 and $\tau = 10$.

We prove that the posterior mean estimator is, under certain conditions, equal to the Tikhonov estimator which has been largely used in (frequentist) nonparametric estimation based on inverse problems. Because of this equality, our EB data-driven method for selecting the regularization parameter is valid even for the frequentist Tikhonov estimator. Thus, our approach makes an important contribution to the theory of inverse problems in econometrics in general.

Appendix

In all the proofs we use the notation $(\lambda_j, \varphi_j, \psi_j)_j$ to denote the singular value decomposition of B (or equivalently of B^*), that is, $B\psi_j = \lambda_j \varphi_j$ and $B^*\varphi_j = \lambda_j \psi_j$ where φ_j and ψ_j are of norm 1, $\forall j$. We also use the notation I(A) to denote the indicator function of an event A and the notation = to mean "equal in distribution". In order to prove several results we make use of Corollary 8.22 in Engl *et al.* (2000). We give here a simplified version of it adapted to our framework and we refer to Engl *et al.* (2000) for the proof of it.

Corollary 2. Let (\mathcal{X}_t) , $t \in \mathbb{R}$ be a Hilbert scale generated by L and let $\Sigma^{-1/2}K : \mathcal{X} \to \mathcal{Y}$ be a bounded operator satisfying Assumption A.2, $\forall x \in \mathcal{X}$ and for some a > 0. Then, for $B = \Sigma^{-1/2}KL^{-s}$, $s \geq 0$ and $|\nu| \leq 1$

$$\underline{c}(\nu)||L^{-\nu(a+s)}x|| \le ||(B^*B)^{\frac{\nu}{2}}x|| \le \overline{c}(\nu)||L^{-\nu(a+s)}x|| \tag{21}$$

holds on $\mathcal{D}((B^*B)^{\frac{\nu}{2}})$ with $\underline{c}(\nu) = \min(\underline{m}^{\nu}, \overline{m}^{\nu})$ and $\overline{c}(\nu) = \max(\underline{m}^{\nu}, \overline{m}^{\nu})$. Moreover, $\mathcal{R}((B^*B)^{\frac{\nu}{2}}) = \mathcal{X}_{\nu(a+s)} \equiv \mathcal{D}(L^{\nu(a+s)})$, where $(B^*B)^{\frac{\nu}{2}}$ has to be replaced by its extension to \mathcal{X} if $\nu < 0$.

A Proofs for Section 3

A.1 Proof of Theorem 1

- (i) See the proof of Theorem 1 (i) and (ii) in Florens and Simoni (2012b).
- (ii) See Theorem 2 and Corollary 2 in Mandelbaum (1984) and their proofs in sections 3.4 and 3.5, page 392.
- (iii) The P_{α} -mlt A is defined as $A := \Omega_0 K^* (\alpha \Sigma + K \Omega_0 K^*)^{-1}$ on $\mathcal{R}((\delta \Sigma + \frac{\delta}{\alpha} K \Omega_0 K^*)^{\frac{1}{2}})$, see Luschgy (1995). Under A.3, the operator $\Omega_0 K^* (\alpha \Sigma + K \Omega_0 K^*)^{-1}$ can equivalently be rewritten as

$$\Omega_0^{\frac{1}{2}}\Omega_0^{\frac{1}{2}}K^*\Sigma^{-1/2}(\alpha I + \Sigma^{-1/2}K\Omega_0K^*\Sigma^{-1/2})^{-1}\Sigma^{-1/2} = \Omega_0^{\frac{1}{2}}B^*(\alpha I + BB^*)^{-1}\Sigma^{-1/2}
= \Omega_0^{\frac{1}{2}}(\alpha I + B^*B)^{-1}(\Sigma^{-1/2}B)^* + \Omega_0^{\frac{1}{2}}[B^*(\alpha I + BB^*)^{-1} - (\alpha I + B^*B)^{-1}B^*]\Sigma^{-1/2}
= \Omega_0^{\frac{1}{2}}(\alpha I + B^*B)^{-1}(\Sigma^{-1/2}B)^*$$
(22)

since $[B^*(\alpha I + BB^*)^{-1} - (\alpha I + B^*B)^{-1}B^*]$ is equal to

$$(\alpha I + B^*B)^{-1}[(\alpha I + B^*B)B^* - B^*(\alpha I + BB^*)](\alpha I + BB^*)^{-1}$$

which is zero. By using expression (22) for $\Omega_0 K^* (\alpha \Sigma + K \Omega_0 K^*)^{-1}$ we show that $\Omega_0 K^* (\alpha \Sigma + K \Omega_0 K^*)^{-1}$ is bounded and continuous on \mathcal{Y} . By Assumption A.3 the operators B and $\Sigma^{-1/2}B$ exist and are bounded. Under A.1 the operator B is compact because it is the product of a compact and a bounded operator. This implies that $B^*B: \mathcal{X} \to \mathcal{X}$ is compact. Because $<(\alpha I + B^*B)\phi, \phi> \geq \alpha ||\phi||^2, \forall \phi \in \mathcal{X}$ we conclude that $(\alpha I + B^*B)$ is injective for $\alpha>0$. Then, from the Riesz Theorem 3.4 in Kress (1999) it follows that the inverse $(\alpha I + B^*B)^{-1}$ is bounded.

Finally, since the product of bounded linear operators is bounded, $\Omega_0 K^* (\alpha \Sigma + K \Omega_0 K^*)^{-1}$ is bounded, i.e. there exists a constant C such that for all $\varphi \in \mathcal{Y}$: $||\Omega_0 K^* (\alpha \Sigma + K \Omega_0 K^*)^{-1} \varphi|| \leq C||\varphi||$. Since $A := \Omega_0 K^* (\alpha \Sigma + K \Omega_0 K^*)^{-1}$ is linear and bounded on \mathcal{Y} , it is continuous on \mathcal{Y} .

A.2 Proof of Theorem 2

The difference $(\hat{x}_{\alpha} - x_*)$ is re-written as

$$\hat{x}_{\alpha} - x_* = -(I - AK)(x_* - x_0) + AU := C_1 + C_2,$$

where the expression of A is given in (9). We consider the MISE associated with \hat{x}_{α} : $\mathbf{E}_{x_*}||\hat{x}_{\alpha}-x_*||^2=||\mathcal{C}_1||^2+\mathbf{E}_{x_*}||\mathcal{C}_2||^2$ and we start by considering the first term. Our proof follows Natterer (1984). Since $0<\frac{s}{a+s}<1$ we can use the left part of inequality (21) in Corollary 2 with $\nu=\frac{s}{a+s}$ and $x=(\alpha I+B^*B)^{-1}L^s(x_*-x_0)$ to obtain (23) below:

$$||\mathcal{C}_{1}||^{2} = ||[I - \Omega_{0}^{\frac{1}{2}}(\alpha I + B^{*}B)^{-1}(\Sigma^{-1/2}B)^{*}K](x_{*} - x_{0})||^{2}$$

$$= ||[I - L^{-s}(\alpha I + B^{*}B)^{-1}(\Sigma^{-1/2}B)^{*}K]L^{-s}L^{s}(x_{*} - x_{0})||^{2}$$

$$= ||L^{-s}[I - (\alpha I + B^{*}B)^{-1}(\Sigma^{-1/2}B)^{*}KL^{-s}]L^{s}(x_{*} - x_{0})||^{2}$$

$$\leq \alpha^{2}\underline{c}^{-2}\left(\frac{s}{a+s}\right)||(B^{*}B)^{\frac{s}{2(a+s)}}(\alpha I + B^{*}B)^{-1}L^{s}(x_{*} - x_{0})||^{2}$$

$$= \alpha^{2}\underline{c}^{-2}\left(\frac{s}{a+s}\right)||(B^{*}B)^{\frac{s}{2(a+s)}}(\alpha I + B^{*}B)^{-1}(B^{*}B)^{\frac{\tilde{\beta}-s}{2(a+s)}}(B^{*}B)^{\frac{s-\tilde{\beta}}{2(a+s)}}L^{s}(x_{*} - x_{0})||^{2}$$

$$= \alpha^{2}\underline{c}^{-2}\left(\frac{s}{a+s}\right)||(B^{*}B)^{\frac{\tilde{\beta}}{2(a+s)}}(\alpha I + B^{*}B)^{-1}w||^{2}$$

$$(24)$$

where $\tilde{\beta} = \min(\beta, a + 2s)$ and $w := (B^*B)^{\frac{s-\beta}{2(a+s)}} L^s(x_* - x_0)$. Now, by using the right part of inequality (21) in Corollary 2 with $\nu = \frac{s-\tilde{\beta}}{a+s}$ (remark that $|\nu| \leq 1$ is verified) and $x = L^s(x_* - x_0)$ we have that

$$\sup_{(x_{*}-x_{0})\in\mathcal{X}_{\beta}(\Gamma)} ||w|| \leq \overline{c} \left(\frac{s-\tilde{\beta}}{a+s}\right) \sup_{(x_{*}-x_{0})\in\mathcal{X}_{\beta}(\Gamma)} ||L^{\tilde{\beta}-s}L^{s}(x_{*}-x_{0})||$$

$$= \overline{c} \left(\frac{s-\tilde{\beta}}{a+s}\right) \sup_{(x_{*}-x_{0})\in\mathcal{X}_{\beta}(\Gamma)} ||(x_{*}-x_{0})||_{\tilde{\beta}}$$

$$\leq \overline{c} \left(\frac{s-\tilde{\beta}}{a+s}\right) \sup_{(x_{*}-x_{0})\in\mathcal{X}_{\beta}(\Gamma)} ||(x_{*}-x_{0})||_{\beta} ||L^{\tilde{\beta}-\beta}||$$

$$= \overline{c} \left(\frac{s-\tilde{\beta}}{a+s}\right) \Gamma^{\frac{1}{2}} ||L^{\tilde{\beta}-\beta}||$$
(25)

where in the penultimate line the equality holds for $\tilde{\beta} = \beta$. Remark that $||L^{\tilde{\beta}-\beta}|| = 1$ if $\tilde{\beta} = \beta$ and is bounded if $\tilde{\beta} < \beta$. Finally,

$$\sup_{(x_{*}-x_{0})\in\mathcal{X}_{\beta}(\Gamma)} ||(B^{*}B)^{\frac{\tilde{\beta}}{2(a+s)}} (\alpha I + B^{*}B)^{-1}w||^{2} = \sup_{(x_{*}-x_{0})\in\mathcal{X}_{\beta}(\Gamma)} \sum_{j=1}^{\infty} \frac{\lambda_{j}^{\frac{2\beta}{a+s}}}{(\alpha + \lambda_{j}^{2})^{2}} \langle w, \psi_{j} \rangle^{2}$$

$$\leq \left(\sup_{\lambda_{j}\geq 0} \frac{\lambda_{j}^{\frac{\tilde{\beta}}{a+s}}}{(\alpha + \lambda_{j}^{2})}\right)^{2} \sup_{(x_{*}-x_{0})\in\mathcal{X}_{\beta}(\Gamma)} ||w||^{2}$$
(26)

and combining (24)-(26) with the fact that $(\sup_{\lambda \geq 0} \lambda^{2b} (\alpha + \lambda^2)^{-1})^2 = \alpha^{2(b-1)} b^{2b} (1-b)^{2(1-b)}$ for $0 \leq b \leq 1$ (and in our case $b = \frac{\tilde{\beta}}{2(a+s)}$), we get the result

$$\sup_{(x_*-x_0)\in\mathcal{X}_{\beta}(\Gamma)} ||\mathcal{C}_1||^2 \le \alpha^{\frac{\tilde{\beta}}{a+s}} b^{2b} (1-b)^{2(1-b)} \underline{c}^{-2} \left(\frac{s}{a+s}\right) \overline{c}^2 \left(\frac{s-\tilde{\beta}}{a+s}\right) ||L^{\tilde{\beta}-\beta}||^2 \Gamma = \mathcal{O}\left(\alpha^{\tilde{\beta}/(a+s)}\right).$$

Next, we address the second term of the MISE. To obtain (27) below, we use the left part of inequality (21) in Corollary 2 with $\nu = \frac{s}{a+s}$ and $x = (\alpha I + B^*B)^{-1}(\Sigma^{-1/2}B)^*U$:

$$\begin{aligned}
\mathbf{E}_{x_{*}} ||\mathcal{C}_{2}||^{2} &= \mathbf{E}_{x_{*}} ||AU||^{2} = \mathbf{E}_{x_{*}} ||\Omega_{0}^{\frac{1}{2}} (\alpha I + B^{*}B)^{-1} (\Sigma^{-1/2}B)^{*}U||^{2} \\
&= \mathbf{E}_{x_{*}} ||L^{-s} (\alpha I + B^{*}B)^{-1} (\Sigma^{-1/2}B)^{*}U||^{2} \\
&\leq \underline{c}^{-2} \left(\frac{s}{a+s}\right) \mathbf{E}_{x_{*}} ||(B^{*}B)^{\frac{s}{2(a+s)}} (\alpha I + B^{*}B)^{-1} (\Sigma^{-1/2}B)^{*}U||^{2} \\
&= \delta \underline{c}^{-2} \left(\frac{s}{a+s}\right) tr \left((B^{*}B)^{\frac{s}{2(a+s)}} (\alpha I + B^{*}B)^{-1} B^{*}B (\alpha I + B^{*}B)^{-1} (B^{*}B)^{\frac{s}{2(a+s)}}\right) \\
&\leq \delta \underline{c}^{-2} \left(\frac{s}{a+s}\right) \left(\sup_{j} \frac{\lambda_{j}^{\frac{s}{(a+s)}+1-\gamma}}{(\alpha + \lambda_{j}^{2})}\right)^{2} \sum_{j=1}^{\infty} \langle (B^{*}B)^{\gamma} \psi_{j}, \psi_{j} \rangle \\
&\leq \delta \underline{c}^{-2} \left(\frac{s}{a+s}\right) \alpha^{2(d-1)} d^{2d} (1-d)^{2(1-d)} \overline{c}(\gamma) tr(L^{-2\gamma(a+s)})
\end{aligned} \tag{28}$$

where $d=\frac{a+2s-\gamma(a+s)}{2(a+s)}$ and which is bounded since $tr(L^{-2\gamma(a+s)})<\infty$ by definition of γ . The last inequality has been obtained by applying the right part of Corollary 2 with $\nu=\gamma$ to $\sum_{j=1}^{\infty}<(B^*B)^{\gamma}\psi_j,\psi_j>$. In fact, $\sum_{j=1}^{\infty}<(B^*B)^{\gamma}\psi_j,\psi_j>=\sum_{j=1}^{\infty}<(B^*B)^{\frac{\gamma}{2}}\psi_j,(B^*B)^{\frac{\gamma}{2}}\psi_j>=\sum_{j=1}^{\infty}||(B^*B)^{\frac{\gamma}{2}}\psi_j||^2\leq \overline{c}(\gamma)\sum_j||L^{-\gamma(a+s)}\psi_j||^2=\overline{c}(\gamma)\sum_j< L^{-\gamma(a+s)}\psi_j,L^{-\gamma(a+s)}\psi_j>=\overline{c}(\gamma)tr(L^{-2\gamma(a+s)}).$ Therefore, $\mathbf{E}_{x_*}||\mathcal{C}_2||^2=\mathcal{O}\left(\delta\alpha^{2(d-1)}\right)$ and $\sup_{(x_*-x_0)\in\mathcal{X}_{\beta}(\Gamma)}\mathbf{E}_{x_*}||\hat{x}_{\alpha}-x_*||^2=\mathcal{O}\left(\alpha^{\frac{\tilde{\beta}}{a+s}}+\delta\alpha^{-\frac{a+\gamma(a+s)}{(a+s)}}\right).$

A.3 Proof of Theorem 3

Let \mathbf{E}_{δ}^{Y} be the expectation taken with respect to the posterior μ_{δ}^{Y} . By the Chebyshev's inequality, for $\varepsilon_{\delta} > 0$ small enough and $M_{\delta} \to \infty$:

$$\mu_{\delta}^{Y}\left\{x \in \mathcal{X}: ||x - x_{*}|| > \varepsilon_{\delta} M_{\delta}\right\} \leq \frac{1}{\varepsilon_{\delta}^{2} M_{\delta}^{2}} \mathbf{E}_{\delta}^{Y} ||x - x_{*}||^{2}$$

$$= \frac{1}{\varepsilon_{\delta}^{2} M_{\delta}^{2}} \left(||\mathbf{E}(x|y^{\delta}, \alpha, s) - x_{*}||^{2} + trVar(x|y^{\delta}, \alpha, s)\right)$$

and we have to determine the rate of $Var(x|y^{\delta},\alpha,s)$. Since $Var(x|y^{\delta},\alpha,s)=\frac{\delta}{\alpha}[\Omega_0-AK\Omega_0]=\delta\Omega_0^{\frac{1}{2}}(\alpha I+B^*B)^{-1}\Omega_0^{\frac{1}{2}}$ we have: $trVar(x|y^{\delta},\alpha,s)=\delta tr(L^{-s}(\alpha I+B^*B)^{-1}L^{-s})=\delta tr(RR^*)$ with $R=L^{-s}(\alpha I+B^*B)^{-\frac{1}{2}}$. Let $||\cdot||_{HS}$ denote the Hilbert-Schmidt norm. By using the left part of Corollary 2 with $\nu=\frac{s}{a+s}$ we get

$$\begin{split} trVar(x|y^{\delta},\alpha,s) &= \delta tr(R^*R) = \delta ||R||_{HS}^2 = \delta \sum_{j=1}^{\infty} ||R\psi_j||^2 \\ &\leq \delta \underline{c}^{-2} \Big(\frac{s}{a+s}\Big) \sum_{j=1}^{\infty} ||(B^*B)^{\frac{s}{2(a+s)}} (\alpha I + B^*B)^{-\frac{1}{2}} \psi_j||^2 \\ &= \delta \underline{c}^{-2} \left(\frac{s}{a+s}\right) \sum_{j=1}^{\infty} <(B^*B)^{\frac{s}{2(a+s)}} (\alpha I + B^*B)^{-\frac{1}{2}} \psi_j, (B^*B)^{\frac{s}{2(a+s)}} (\alpha I + B^*B)^{-\frac{1}{2}} \psi_j > \\ &= \delta \underline{c}^{-2} \Big(\frac{s}{a+s}\Big) \sum_{j=1}^{\infty} \frac{\lambda_j^{\frac{2s}{a+s}-2\gamma}}{\alpha + \lambda_j^2} <(B^*B)^{\gamma} \psi_j, \psi_j > \\ &\leq \delta \underline{c}^{-2} \Big(\frac{s}{a+s}\Big) \sup_j \frac{\lambda_j^{\frac{2s}{a+s}-2\gamma}}{\alpha + \lambda_j^2} \sum_{j=1}^{\infty} <(B^*B)^{\gamma} \psi_j, \psi_j > \\ &= \delta \underline{c}^{-2} \Big(\frac{s}{a+s}\Big) \sup_j \frac{\lambda_j^{\frac{2s}{a+s}-2\gamma}}{\alpha + \lambda_j^2} \sum_{j=1}^{\infty} \Big|\Big|(B^*B)^{\gamma/2} \psi_j\Big|\Big|^2 \\ &\leq \delta \underline{c}^{-2} \Big(\frac{s}{a+s}\Big) \alpha^{v-1} v^v (1-v)^{1-v} \overline{c}(\gamma) tr(L^{-2\gamma(a+s)}), \quad v := \frac{s}{(a+s)} - \gamma \end{split}$$

which is finite because $tr(L^{-2\gamma(a+s)}) < \infty$. The last inequality has been obtained by applying the right part of Corollary 2. Now, by this and the result of Theorem 2:

$$\mu_{\delta}^{Y}\left\{x \in \mathcal{X} : ||x - x_{*}|| > \varepsilon_{\delta} M_{\delta}\right\} = \mathcal{O}_{p}\left(\frac{1}{\varepsilon_{\delta}^{2} M_{\delta}^{2}} \left(\alpha^{\frac{\tilde{\beta}}{a+s}} + \delta \alpha^{-\frac{a+\gamma(a+s)}{(a+s)}}\right)\right), \quad \tilde{\beta} = \min(\beta, a+2s)$$

in P^{x_*} -probability. Hence, for $\varepsilon_{\delta} = (\alpha^{\frac{\tilde{\beta}}{2(a+s)}} + \delta^{\frac{1}{2}} \alpha^{-\frac{a+\gamma(a+s)}{2(a+s)}})$ we have $\mu_{\delta}^{Y} \{x \in \mathcal{X} : ||x - x_*|| > \varepsilon_{\delta} M_{\delta}\} \to 0$.

A.4 Proof of Corollary 1

Let $\widetilde{B^*B} = \Omega_0^{\frac{1}{2}} \mathfrak{K}^* \frac{1}{f_W} \mathfrak{K} \Omega_0^{\frac{1}{2}}$, $R_{\alpha} = (\alpha I + \widetilde{B^*B})^{-1}$, $\hat{R}_{\alpha} = (\alpha I + B^*B)^{-1}$, $\Sigma^{-1/2}B = (\Lambda^*)^{-1} \frac{1}{\hat{f}_W} \tilde{K} \Omega_0^{\frac{1}{2}}$, $\widetilde{\Omega}_0^{\frac{1}{2}}$, $\widehat{R} = f_W \mathbf{E}(\cdot|W) = \lim_{n \to \infty} \tilde{K}$ and $\mathfrak{K}^* = f_Z \mathbf{E}(\cdot|Z) = \lim_{n \to \infty} \tilde{K}^*$. Moreover, we define $\tilde{y}^{\delta} = \Lambda \mathfrak{K} x_* + U$ and $\tilde{x}_{\alpha} = \Omega_0^{\frac{1}{2}} R_{\alpha} (\widetilde{\Sigma^{-1/2}B})^* (\tilde{y}^{\delta} - \Lambda \mathfrak{K} x_0)$. We decompose $(\hat{x}_{\alpha} - x_*)$ as

$$\begin{split} (\hat{x}_{\alpha} - x_{*}) &= (\tilde{x}_{\alpha} - x_{*}) + \Omega_{0}^{\frac{1}{2}} \hat{R}_{\alpha} (\Sigma^{-1/2} B)^{*} (y^{\delta} - \Lambda \tilde{K} x_{0}) - \Omega_{0}^{\frac{1}{2}} R_{\alpha} (\widetilde{\Sigma^{-1/2} B})^{*} (\tilde{y}^{\delta} - \Lambda \mathfrak{K} x_{0}) \\ &= (\tilde{x}_{\alpha} - x_{*}) + \Omega_{0}^{\frac{1}{2}} \left[\hat{R}_{\alpha} (\Sigma^{-1/2} B)^{*} \Lambda \tilde{K} - R_{\alpha} (\widetilde{\Sigma^{-1/2} B})^{*} \Lambda \mathfrak{K} \right] (x_{*} - x_{0}) \\ &+ \Omega_{0}^{\frac{1}{2}} \left[\hat{R}_{\alpha} (\Sigma^{-1/2} B)^{*} - R_{\alpha} (\widetilde{\Sigma^{-1/2} B})^{*} \right] U \\ &:= (\tilde{x}_{\alpha} - x_{*}) + \mathfrak{A}_{1} + \mathfrak{A}_{2} \end{split}$$

I) Convergence of $\mathbf{E}||(\tilde{x}_{\alpha}-x_{*})||^{2}$. This rate is given in Theorem 2 (since this term depends on operators that do not vary with n).

II) Convergence of \mathfrak{A}_1 .

$$\begin{split} \mathbf{E}||\mathfrak{A}_{1}||^{2} &= \mathbf{E}||\Omega_{0}^{\frac{1}{2}}\left(\hat{R}_{\alpha}\left[(\Sigma^{-1/2}B)^{*}\Lambda\tilde{K}L^{-s} - (\widetilde{\Sigma^{-1/2}B})^{*}\Lambda\mathfrak{K}L^{-s}\right] + (\hat{R}_{\alpha} - R_{\alpha})(\widetilde{\Sigma^{-1/2}B})^{*}\Lambda\mathfrak{K}L^{-s}\right)L^{s}(x_{*} - x_{0})||^{2} \\ &= \mathbf{E}||\Omega_{0}^{\frac{1}{2}}\left(\hat{R}_{\alpha}\left[B^{*}B - \widetilde{B^{*}B}\right] + \hat{R}_{\alpha}(\widetilde{B^{*}B} - B^{*}B)R_{\alpha}\widetilde{B^{*}B}\right)L^{s}(x_{*} - x_{0})||^{2} \\ &= \mathbf{E}||\Omega_{0}^{\frac{1}{2}}\hat{R}_{\alpha}\left[B^{*}B - \widetilde{B^{*}B}\right]\alpha R_{\alpha}L^{s}(x_{*} - x_{0})||^{2} \\ &\leq \mathbf{E}||\Omega_{0}^{\frac{1}{2}}\hat{R}_{\alpha}||^{2}\mathbf{E}||\Omega_{0}^{\frac{1}{2}}(\tilde{K}^{*}[\hat{f}_{W}]^{-1}\tilde{K} - \mathfrak{K}^{*}[f_{W}]^{-1}\mathfrak{K})||^{2}||L^{-s}R_{\alpha}L^{s}(x_{*} - x_{0})||^{2}\alpha^{2}. \end{split}$$

The last term $||L^{-s}R_{\alpha}L^{s}(x_{*}-x_{0})||^{2}\alpha^{2}$ is equal to term \mathcal{C}_{1} in the proof of Theorem 2 while $\mathbf{E}||\Omega_{0}^{\frac{1}{2}}\hat{R}_{\alpha}||^{2} = \mathcal{O}(\alpha^{-2})$ and $\mathbf{E}||\Omega_{0}^{\frac{1}{2}}(\tilde{K}^{*}[\hat{f}_{W}]^{-1}\tilde{K}-\mathfrak{K}^{*}[f_{W}]^{-1}\mathfrak{K})||^{2} = \mathcal{O}(n^{-1}+h^{2\rho})$ under Assumption HS since $\mathbf{E}||\Omega_{0}^{\frac{1}{2}}(\tilde{K}^{*}[\hat{f}_{W}]^{-1}\tilde{K}-\mathfrak{K}^{*}[f_{W}]^{-1}\mathfrak{K})||^{2} \times \mathbf{E}||B^{*}B-\widetilde{B^{*}B}||^{2}$. Therefore, $\mathbf{E}||\mathfrak{A}_{1}||^{2} = \mathcal{O}\left(\alpha^{-2}(n^{-1}+h^{2\rho})\alpha^{\frac{\tilde{\beta}}{a+s}}\right)$.

III) Convergence of \mathfrak{A}_2 .

$$\begin{split} \mathbf{E}||\mathfrak{A}_{2}||^{2} &= \mathbf{E}||\Omega_{0}^{\frac{1}{2}}\left(\hat{R}_{\alpha}[(\Sigma^{-1/2}B)^{*} - (\widetilde{\Sigma^{-1/2}B})^{*}] + (\hat{R}_{\alpha} - R_{\alpha})(\widetilde{\Sigma^{-1/2}B})^{*}\right)U||^{2} \\ &\leq 2\mathbf{E}||\Omega_{0}^{\frac{1}{2}}\hat{R}_{\alpha}[(\Sigma^{-1/2}B)^{*} - (\widetilde{\Sigma^{-1/2}B})^{*}]U||^{2} + 2\mathbf{E}||L^{-s}\hat{R}_{\alpha}(\widetilde{B^{*}B} - B^{*}B)R_{\alpha}(\widetilde{\Sigma^{-1/2}B})^{*}U||^{2} \\ &:= 2\mathfrak{A}_{2,1} + 2\mathfrak{A}_{2,2}. \end{split}$$

We start with the analysis of term $\mathfrak{A}_{2,1}$ where we use the notation $\mathcal{T} = \mathbf{E}(\cdot|W)$, $\mathcal{T}^* = \mathbf{E}(\cdot|Z)$, $\hat{\mathcal{T}} = \hat{\mathbf{E}}(\cdot|W)$ and $\hat{\mathcal{T}}^* = \hat{\mathbf{E}}(\cdot|Z)$:

$$\begin{split} \mathfrak{A}_{2,1} & \leq & \mathbf{E}||\Omega_{0}^{\frac{1}{2}}\hat{R}_{\alpha}||^{2}\mathbf{E}||\left((\Sigma^{-1/2}B)^{*} - (\widetilde{\Sigma^{-1/2}B})^{*}\right)U||^{2} \\ & = & \mathbf{E}||\Omega_{0}^{\frac{1}{2}}\hat{R}_{\alpha}||^{2}\mathbf{E}||\left((\Lambda^{*})^{-1}(\hat{\mathcal{T}} - \mathcal{T})\Omega_{0}^{\frac{1}{2}}\right)^{*}U||^{2} \\ & = & \mathbf{E}||\Omega_{0}^{\frac{1}{2}}\hat{R}_{\alpha}||^{2}\mathbf{E}||\Omega_{0}^{\frac{1}{2}}(\hat{\mathcal{T}} - \mathcal{T})^{*}\left[\frac{1}{nh}\sum_{i}(y_{i} - \langle x_{*}, \frac{K_{Z,h}(z_{i} - z)}{h} \rangle)K_{W,h}(w_{i} - w)\right]||^{2} \\ & \leq & \mathbf{E}||\Omega_{0}^{\frac{1}{2}}\hat{R}_{\alpha}||^{2}\mathbf{E}||\Omega_{0}^{\frac{1}{2}}(\hat{\mathcal{T}} - \mathcal{T})^{*}||^{2}\mathbf{E}||\left[\frac{1}{nh}\sum_{i}(y_{i} - \langle x_{*}, \frac{K_{Z,h}(z_{i} - z)}{h} \rangle)K_{W,h}(w_{i} - w)\right]||^{2} \\ & = & \mathcal{O}\left(\alpha^{-2}(n^{-1} + h^{2\rho})((nh)^{-1} + h^{2\rho})\right) \end{split}$$

since $\mathbf{E}||\Omega_0^{\frac{1}{2}}(\hat{\mathcal{T}}-\mathcal{T})^*||^2 = \mathcal{O}(n^{-1}+h^{2\rho})$ under Assumption HS. Finally, term $\mathfrak{A}_{2,2}$ can be developed as follows:

$$\begin{array}{lcl} \mathfrak{A}_{2,2} & = & \mathbf{E}||\Omega_{0}^{\frac{1}{2}}\hat{R}_{\alpha}\Omega_{0}^{\frac{1}{2}}\left(\tilde{K}^{*}\frac{1}{\hat{f}_{W}}\tilde{K}-\mathfrak{K}^{*}\frac{1}{f_{W}}\mathfrak{K}\right)L^{-s}R_{\alpha}(\widetilde{\Sigma^{-1/2}B})^{*}U||^{2} \\ & \leq & \mathbf{E}||\Omega_{0}^{\frac{1}{2}}\hat{R}_{\alpha}||^{2}\mathbf{E}||\Omega_{0}^{\frac{1}{2}}\left(\tilde{K}^{*}\frac{1}{\hat{f}_{W}}\tilde{K}-\mathfrak{K}^{*}\frac{1}{f_{W}}\mathfrak{K}\right)||^{2}\mathbf{E}||L^{-s}R_{\alpha}(\widetilde{\Sigma^{-1/2}B})^{*}U||^{2} \\ & = & \mathcal{O}\left(\alpha^{-2}(n^{-1}+h^{2\rho})n^{-1}\alpha^{-\frac{a+\gamma(a+s)}{a+s}}\right) \end{array}$$

since the last term is equal to term C_2 in the proof of Theorem 2 and $\mathbf{E}||\Omega_0^{\frac{1}{2}}\left(\tilde{K}^*\frac{1}{\hat{f_W}}\tilde{K}-\mathfrak{K}^*\frac{1}{f_W}\mathfrak{K}\right)||^2 = \mathcal{O}(n^{-1}+h^{2\rho})$ under Assumption HS. By writing $\mathbf{E}||\hat{x}_{\alpha}-x_*||^2 \leq 2\mathbf{E}||\tilde{x}_{\alpha}-x_*||^2 + 2\mathbf{E}||\mathfrak{A}_1+\mathfrak{A}_2||^2$ and by putting all these results together we get the result.

B Proofs for Section 4

B.1 Proof of Proposition 2

Let us consider the singular system $\{\lambda_j, \varphi_j, \psi_j\}$ associated with B. Under Assumption C there exist $\underline{\lambda}, \bar{\lambda} > 0$ such that $\underline{\lambda}j^{-(a+s)} \leq \lambda_j \leq \bar{\lambda}j^{-(a+s)}$. Therefore, the risk associated with \hat{x}_{α} can be rewritten as:

$$\mathbf{E}_{x_*}||\hat{x}_{\alpha} - x_*||^2 = \sum_{j} \frac{\alpha^2}{(\alpha + j^{-2(a+s)})^2} < (x_* - x_0), \psi_j >^2 + \sum_{j} \frac{\delta j^{-2s - 2(a+s)}}{(\alpha + j^{-2(a+s)})^2} =: \mathcal{A}1 + \mathcal{A}2.$$

We have that $\sup_{(x_*-x_0)\in\mathcal{X}_{\beta}(\Gamma)}\mathcal{A}1 = \sup_{\sum_j j^{2\beta} < (x_*-x_0), \psi_j > 2 \le \Gamma}\mathcal{A}1$ and

$$\sup_{\sum_{j} j^{2\beta} < (x_{*} - x_{0}), \psi_{j} >^{2} \leq \Gamma} \mathcal{A}1 = \alpha^{2} \sup_{\sum_{j} j^{2\beta} < (x_{*} - x_{0}), \psi_{j} >^{2} \leq \Gamma} \sum_{j} \frac{1}{(\alpha + j^{-2(a+s)})^{2}} < (x_{*} - x_{0}), \psi_{j} >^{2}$$

$$= \alpha^{2} \sup_{\sum_{j} j^{2\beta} < (x_{*} - x_{0}), \psi_{j} >^{2} \leq \Gamma} \sum_{j} \frac{j^{-2\beta}}{(\alpha + j^{-2(a+s)})^{2}} j^{2\beta} < (x_{*} - x_{0}), \psi_{j} >^{2}$$

$$= \alpha^{2} \sup_{j} \frac{j^{-2\beta}}{(\alpha + j^{-2(a+s)})^{2}} \Gamma.$$

The supremum is attained at $j = \alpha^{-\frac{1}{2(a+s)}} \left(\frac{\beta}{2(a+s)-\beta}\right)^{-\frac{1}{2(a+s)}}$ as long as $\beta < 2(a+s)$. If $\beta \ge 2(a+s)$ then $\sup_j \frac{j^{-2\beta}}{(\alpha+j^{-2(a+s)})^2} = 1$. Consequently,

$$\sup_{(x_*-x_0)\in\mathcal{X}_{\beta}(\Gamma)} \mathcal{A}1 = \alpha^{\frac{\tilde{\beta}}{a+s}} b^{2b} (1-b)^{2(1-b)} \Gamma, \quad b = \frac{\tilde{\beta}}{2(a+s)}, \quad \tilde{\beta} = \min(\beta, 2(a+s))$$

by using the convention $(1-b)^{2(1-b)}=1$ if b=1. In order to analyze term $\mathcal{A}2$ we first remark that the summand function $f(j):=\frac{j^{-2s-2(a+s)}}{(\alpha+j^{-2(a+s)})^2}$ defined on \mathbb{R}_+ is increasing for $j<\bar{j}$ and decreasing for $j>\bar{j}$ where $\bar{j}=\left(\frac{a+2s}{a}\alpha\right)^{-\frac{1}{2(a+s)}}$. Thus,

$$\delta \int_{\bar{j}}^{\infty} \frac{j^{-2s-2(a+s)}}{(\alpha+j^{-2(a+s)})^2} dj \le \mathcal{A}_2 \le \delta \sum_{j=1}^{\bar{j}} \frac{j^{-2s-2(a+s)}}{(\alpha+j^{-2(a+s)})^2} + \delta \sum_{j=\bar{j}}^{\infty} \frac{j^{-2s-2(a+s)}}{(\alpha+j^{-2(a+s)})^2}$$

$$\Leftrightarrow \delta \alpha^{-\frac{2a+1}{2(a+s)}} \int_{\bar{t}}^{\infty} \frac{1}{[t^{-a}(t^{2(a+s)}+1)]^2} dt \le \mathcal{A}_2 \le \delta \frac{\bar{j}^{1-2s-2(a+s)}}{(\alpha+\bar{j}^{-2(a+s)})^2} + \delta \alpha^{-\frac{2a+1}{2(a+s)}} \int_{\bar{t}}^{\infty} \frac{1}{[t^{-a}(t^{2(a+s)}+1)]^2} dt$$

where $\bar{t} = \left(\frac{a+2s}{a}\right)^{-\frac{1}{2(a+s)}}$. Denote $c_2(\bar{t}) = \int_{\bar{t}}^{\infty} \frac{1}{[t^{-a}(t^{2(a+s)}+1)]^2} dt$ and replace \bar{j} by its value to obtain

$$\delta\alpha^{-\frac{2a+1}{2(a+s)}}c_2(\bar{t}) \leq \mathcal{A}_2 \leq \delta\alpha^{-\frac{2a+1}{2(a+s)}} \left(c_2(\bar{t}) + \bar{t}^{(1-2s-2(a+s))} \frac{a^2}{4(a+s)^2}\right) =: \delta\alpha^{-\frac{2a+1}{2(a+s)}} \tilde{c}_2(\bar{t}).$$

Therefore, $A_2 \simeq \delta \alpha^{-\frac{2a+1}{2(a+s)}}$. Finally, let $G = b^{2b}(1-b)^{2(1-b)}\Gamma$,

$$\alpha^{\frac{\tilde{\beta}}{(a+s)}}G + \delta\alpha^{-\frac{2a+1}{2(a+s)}}c_2(\bar{t}) \leq \sup_{(x_*-x_0) \in \mathcal{X}_{\beta}(\Gamma)} \mathbf{E}_{x_*}||\hat{x}_{\alpha} - x_*||^2 \quad \leq \quad \alpha^{\frac{\tilde{\beta}}{(a+s)}}G + \delta\alpha^{-\frac{2a+1}{2(a+s)}}\tilde{c}_2(\bar{t}),$$

this proves (12). By replacing $\alpha \approx \delta$ (resp. $\alpha \approx \delta^{\frac{a+s}{\beta+a+1/2}}$) in this expression we get (i) (resp. (ii)).

B.2 Proof of Proposition 3

The proof proceeds similar to the proof of Proposition 2, so we just sketch it. The risk $\mathbf{E}_{x_*}||\hat{x}_{\alpha} - x_*||^2$ associated with \hat{x}_{α} rewrites as the risk in section B.1. The only term that we need to analyze is $\mathcal{A}1$ since the analysis of $\mathcal{A}2$ is the same as in B.1. By using Assumption B'

$$\mathcal{A}1 = \alpha^2 \sum_{j} \frac{1}{(\alpha + j^{-2(a+s)})^2} < (x_* - x_0), \psi_j >^2 = \alpha^2 \sum_{j} \frac{j^{-2b_0}}{(\alpha + j^{-2(a+s)})^2}$$

and the function $f(j) := \frac{j^{-2b_0}}{(\alpha+j^{-2(a+s)})^2}$ defined on \mathbb{R}_+ is decreasing in j if $b_0 \geq 2a+2s$. If $b_0 < 2a+2s$ then f(j) is increasing for $j < \bar{j}$ and decreasing for $j > \bar{j}$ where $\bar{j} = \alpha^{-\frac{1}{2(a+s)}} \left(\frac{b_0}{2a+2s-b_0}\right)^{-\frac{1}{2(a+s)}}$. Therefore, to upper and lower bound \mathcal{A}_1 we have to consider these two cases separately. If $b_0 < 2a+2s$

$$\alpha^{2} \int_{\bar{j}}^{\infty} \frac{j^{-2b_{0}}}{(\alpha + j^{-2(a+s)})^{2}} dj \leq \mathcal{A}_{1} \leq \alpha^{2} \sum_{j=1}^{\bar{j}} \frac{j^{-2b_{0}}}{(\alpha + j^{-2(a+s)})^{2}} + \alpha^{2} \sum_{j=\bar{j}}^{\infty} \frac{j^{-2b_{0}}}{(\alpha + j^{-2(a+s)})^{2}}$$

$$\Leftrightarrow \alpha^{\frac{2b_{0}-1}{2(a+s)}} \int_{\underline{u}}^{\infty} \frac{u^{4(a+s)-2b_{0}}}{(u^{2(a+s)}+1)^{2}} du \leq \mathcal{A}_{1} \leq \alpha^{2} \frac{\bar{j}^{1-2b_{0}}}{(\alpha + \bar{j}^{-2(a+s)})^{2}} + \alpha^{\frac{2b_{0}-1}{2(a+s)}} \int_{\underline{u}}^{\infty} \frac{u^{4(a+s)-2b_{0}}}{(u^{2(a+s)}+1)^{2}} du$$

where $\underline{u} = \left(\frac{b_0}{2a+2s-b_0}\right)^{-\frac{1}{2(a+s)}}$. Denote $c_1(\tilde{u}) = \int_{\tilde{u}}^{\infty} \frac{u^{4(a+s)-2b_0}}{(u^{2(a+s)}+1)^2} du$ for some scalar \tilde{u} and replace \bar{j} by its value. Then,

$$\alpha^{\frac{2b_0-1}{2(a+s)}}c_1(\underline{u}) \le \mathcal{A}_1 \le \alpha^{\frac{2b_0-1}{2(a+s)}} \left(c_1(\underline{u}) + \underline{u}^{1-2b_0} \left(\frac{2a+2s-b_0}{2(a+s)}\right)^2\right) =: \alpha^{\frac{2b_0-1}{2(a+s)}} \tilde{c}_1(\underline{u}).$$

If $b_0 \ge 2a + 2s$:

$$\alpha^{2} \int_{1}^{\infty} \frac{j^{-2b_{0}}}{(\alpha + j^{-2(a+s)})^{2}} dj \leq \mathcal{A}_{1} \leq \alpha^{2} \frac{j^{-2b_{0}}}{(\alpha + j^{-2(a+s)})^{2}} \bigg|_{j=1} + \alpha^{2} \int_{1}^{\infty} \frac{j^{-2b_{0}}}{(\alpha + j^{-2(a+s)})^{2}} dj$$

$$\Leftrightarrow \alpha^{\frac{2b_{0}-1}{2(a+s)}} \int_{\bar{u}}^{\infty} \frac{u^{4(a+s)-2b_{0}}}{(u^{2(a+s)}+1)^{2}} du \leq \mathcal{A}_{1} \leq \alpha^{2} + \alpha^{\frac{2b_{0}-1}{2(a+s)}} \int_{\bar{u}}^{\infty} \frac{u^{4(a+s)-2b_{0}}}{(u^{2(a+s)}+1)^{2}} du$$

where $\bar{u} = \alpha^{\frac{1}{2(a+s)}}$. By using the notation defined above we obtain:

$$\alpha^{\frac{2b_0-1}{2(a+s)}}c_1(\bar{u}) \le \mathcal{A}_1 \le \alpha^2 + \alpha^{\frac{2b_0-1}{2(a+s)}}c_1(\bar{u})$$

Since the integral in $c_1(\bar{u})$ is convergent, the upper bound is of order $\alpha^{\frac{2\bar{b}_0-1}{2(a+s)}}$ where $\tilde{b}_0 = \min(b_0, 2a+2s+1/2)$. Summarizing the two cases and by defining: $c_1 = c_1(\underline{u})I(b_0 < 2(a+s)) + c_1(\bar{u})I(b_0 \geq 2(a+s))$ and $\tilde{c}_1 = \tilde{c}_1(\underline{u})I(b_0 < 2(a+s)) + c_1(\bar{u})I(b_0 \ge 2(a+s))$ we have

$$\alpha^{\frac{2b_0-1}{2(a+s)}}c_1 \le \mathcal{A}_1 \le \alpha^2 I(b_0 \ge (2a+2s)) + \alpha^{\frac{2b_0-1}{2(a+s)}}\tilde{c}_1.$$

By using the upper and lower bounds for A_2 given in the proof of Proposition 2 we get the expression in (13):

$$\alpha^{\frac{2b_0-1}{2(a+s)}}c_1 + \delta\alpha^{-\frac{2a+1}{2(a+s)}}c_2(\bar{t}) \leq \mathbf{E}_{x_*}||\hat{x}_{\alpha} - x_*||^2 \leq \alpha^2 I(b_0 \geq (2a+2s)) + \alpha^{\frac{2b_0-1}{2(a+s)}}\tilde{c}_1 + \delta\alpha^{-\frac{2a+1}{2(a+s)}}\tilde{c}_2(\bar{t}).$$

By replacing $\alpha \asymp \delta^{\frac{a+s}{\overline{b_0}+a}}$ we obtain (14).

B.3 Proof of Theorem 4

We use the same strategy used for the proof of Theorem 3. The trace of the posterior covariance operator is:

$$tr[Var(x|y^{\delta},\alpha,s)] = \delta \sum_{j=1}^{\infty} <\Omega_0^{\frac{1}{2}} (\alpha I + B^*B)^{-1} \Omega_0^{\frac{1}{2}} \psi_j, \psi_j > = \delta \sum_{j=1}^{\infty} \frac{j^{-2s}}{\alpha + j^{-2(a+s)}}.$$

Since $f(j) := \frac{j^{-2s}}{\alpha + j^{-2(a+s)}}$ is increasing in j for $j < \bar{j}$ and decreasing for $j > \bar{j}$ where $\bar{j} = (\alpha s/a)^{-\frac{1}{2(a+s)}}$ then we can upper and lower bound the trace of the posterior covariance operator as:

$$\begin{split} \delta \sum_{j=\bar{j}}^{\infty} f(j) & \leq & tr[Var(x|y^{\delta},\alpha,s)] \leq \delta \sum_{j=1}^{\bar{j}} f(j) + \delta \sum_{j=\bar{j}}^{\infty} f(j) \\ \Leftrightarrow \delta \alpha^{-\frac{2a+1}{2(a+s)}} \int_{\bar{t}}^{\infty} \frac{t^{2a}}{t^{2(a+s)} + 1} dt & \leq & tr[Var(x|y^{\delta},\alpha,s)] \leq \delta \frac{\bar{j}^{1-2s}}{\alpha + \bar{j}^{-2(a+s)}} + \delta \alpha^{-\frac{2a+1}{2(a+s)}} \int_{\bar{t}}^{\infty} \frac{t^{2a}}{t^{2(a+s)} + 1} dt \\ \Leftrightarrow \delta \alpha^{-\frac{2a+1}{2(a+s)}} \kappa_{2}(\bar{t}) & \leq & tr[Var(x|y^{\delta},\alpha,s)] \leq \delta \alpha^{-\frac{2a+1}{2(a+s)}} \bar{t}^{-\frac{(1-2s)}{2(a+s)}} \frac{a}{a+s} + \delta \alpha^{-\frac{2a+1}{2(a+s)}} \kappa_{2}(\bar{t}) \end{split}$$

where $\bar{t} = (s/a)^{-\frac{1}{2(a+s)}}$ and $\kappa_2(\bar{t}) = \int_{\bar{t}}^{\infty} \frac{t^{2a}}{t^{2(a+s)}+1} dt$. Thus, $tr[Var(x|y^{\delta}, \alpha, s)] \approx \delta \alpha^{-\frac{2a+1}{2(a+s)}}$. By the Chebyshev's inequality and Proposition 3, for $\varepsilon_{\delta} > 0$ small enough and $M_{\delta} \to \infty$

$$\mu_{\delta}^{Y}\left\{x \in \mathcal{X}: ||x - x_{*}|| > \varepsilon_{\delta} M_{\delta}\right\} \leq \frac{1}{\varepsilon_{\delta}^{2} M_{\delta}^{2}} \left(||\mathbf{E}(x|y^{\delta}, \alpha, s) - x_{*}||^{2} + tr Var(x|y^{\delta}, \alpha, s)\right)$$

$$= \mathcal{O}_{p}\left(\frac{1}{\varepsilon_{\delta}^{2} M_{\delta}^{2}} \left(\alpha^{\frac{2\tilde{b}_{0} - 1}{2(a + s)}} + \delta \alpha^{-\frac{2a + 1}{2(a + s)}}\right)\right), \quad \tilde{b}_{0} = \min(b_{0}, 2(a + s) + 1/2)$$

in P^{x_*} -probability. Hence, for $\varepsilon_{\delta} = (\alpha^{\frac{2\bar{b}_0 - 1}{4(a+s)}} + \delta^{\frac{1}{2}}\alpha^{-\frac{2a+1}{4(a+s)}})$ we conclude that $\mu_{\delta}^Y \{x \in \mathcal{X} : ||x - x_*|| > \varepsilon_{\delta} M_{\delta}\} \to 0$. For $\alpha \asymp \delta$ we obtain (i). To obtain (ii): $\arg\inf_{\alpha} \varepsilon_{\delta} = \delta^{\frac{a+s}{b_0 + a}}$ and $\inf_{\alpha} \varepsilon_{\delta} = \delta^{\frac{2\bar{b}_0 - 1}{4(b_0 + a)}}$.

B.4 Proof of Lemma 1

Let us denote: $s = a_0 \tilde{s}$, $\tilde{R}_{\alpha} = (\alpha I + \tilde{K}^{1+2\tilde{s}})^{-1}$, $R_{\alpha} = (\alpha I + K^{1+2\tilde{s}})^{-1}$, $\hat{x}_{\alpha} = K^{\tilde{s}} R_{\alpha} K^{\tilde{s}} (y^{\delta} - K x_0) + x_0$. Then,

$$\begin{split} \mathbf{E}_{x_*} || \hat{x}_{\alpha} - x_* ||^2 &= \mathbf{E}_{x_*} || \alpha R_{\alpha}(x_* - x_0) + K^{\tilde{s}} R_{\alpha} K^{\tilde{s}} U ||^2 \\ &= \mathbf{E}_{x_*} || \alpha \tilde{R}_{\alpha}(x_* - x_0) + \alpha (R_{\alpha} - \tilde{R}_{\alpha})(x_* - x_0) + \tilde{R}_{\alpha} \tilde{K}^{2\tilde{s}} U + (R_{\alpha} K^{2\tilde{s}} - \tilde{R}_{\alpha} \tilde{K}^{2\tilde{s}}) U ||^2 \\ &\leq 2 \mathbf{E}_{x_*} || \alpha \tilde{R}_{\alpha}(x_* - x_0) + \tilde{R}_{\alpha} \tilde{K}^{2\tilde{s}} U ||^2 + 4 \mathbf{E}_{x_*} || \alpha (R_{\alpha} - \tilde{R}_{\alpha})(x_* - x_0) ||^2 \\ &+ 4 \mathbf{E}_{x_*} || (R_{\alpha} K^{2\tilde{s}} - \tilde{R}_{\alpha} \tilde{K}^{2\tilde{s}}) U ||^2 =: 2 \mathcal{A}_1 + 4 \mathcal{A}_2 + 4 \mathcal{A}_3. \end{split}$$

The rate of term \mathcal{A}_1 is given by Proposition 3 with $\delta = n^{-1}$, $a = \frac{a_0}{2}$ and $s = a_0\tilde{s}$: $\mathcal{A}_1 \asymp \left(\alpha^{\frac{2\tilde{b}_0 - 1}{2s + a_0}} + n^{-1}\alpha^{-\frac{a_0 + 1}{2s + a_0}}\right)$. By using the result of Lemma 2 and a Taylor expansion of the second order of the function $\tilde{K}^{2\tilde{s}+1}$ around K (that is, $\tilde{K}^{2\tilde{s}+1} = K^{2\tilde{s}+1} + (2\tilde{s}+1)K^{2\tilde{s}}(|\tilde{K}-K|) + \mathcal{O}(|\tilde{K}-K|^2)$) we obtain

$$\begin{split} \mathcal{A}_{2} &= \mathbf{E}_{x_{*}} ||\alpha(R_{\alpha} - \tilde{R}_{\alpha})(x_{*} - x_{0})||^{2} = \mathbf{E}_{x_{*}} ||R_{\alpha}(\tilde{K}^{2\tilde{s}+1} - K^{2\tilde{s}+1})\alpha\tilde{R}_{\alpha})(x_{*} - x_{0})||^{2} \\ &= \mathbf{E}_{x_{*}} ||R_{\alpha} \left((2\tilde{s}+1)K^{2s}(|\tilde{K} - K|) + \mathcal{O}(|\tilde{K} - K|^{2}) \right) \alpha\tilde{R}_{\alpha})(x_{*} - x_{0})||^{2} \\ &\leq \mathcal{O} \left(\mathbf{E}_{x_{*}} ||R_{\alpha}K^{2s}||^{2} \mathbf{E}_{x_{*}} ||\tilde{K} - K||^{2} \mathbf{E}_{x_{*}} ||\alpha\tilde{R}_{\alpha}(x_{*} - x_{0})||^{2} \right) \\ &+ \mathcal{O} \left(\mathbf{E}_{x_{*}} ||R_{\alpha}||^{2} \mathbf{E}_{x_{*}} ||(\tilde{K} - K)^{2}||^{2} \mathbf{E}_{x_{*}} ||\alpha\tilde{R}_{\alpha}(x_{*} - x_{0})||^{2} \right) \\ &= \mathcal{O} \left(n^{-1} \alpha^{-\frac{2a_{0}}{2s+a_{0}}} \alpha^{\frac{2\tilde{b}_{0}-1}{2s+a_{0}}} \right) + \mathcal{O} \left(n^{-2} \alpha^{-2} \alpha^{\frac{2\tilde{b}_{0}-1}{2s+a_{0}}} \right) \end{split}$$

where we have used the facts that: $\mathbf{E}_{x_*}||R_{\alpha}K^{2s}||^2 = \mathcal{O}\left(\alpha^{-\frac{2a_0}{2s+a_0}}\right)$ and $\mathbf{E}_{x_*}||(\tilde{K}-K)^2||^2 = \mathcal{O}\left(n^{-2}\right)$. The latter rate is obtained as follows: for some basis $\{\varphi_i\}$:

$$\begin{split} ||(\tilde{K}-K)^{2}||^{2} &:= \sup_{||\phi|| \leq 1} \sum_{j=1}^{\infty} <(\tilde{K}-K)^{2}\phi, \varphi_{j}>^{2} = \sup_{||\phi|| \leq 1} \sum_{j=1}^{\infty} <(\tilde{K}-K)\phi, (\tilde{K}-K)\varphi_{j}>^{2} \\ &\leq \sup_{||\phi|| \leq 1} \sum_{j=1}^{\infty} ||(\tilde{K}-K)\phi||^{2} ||(\tilde{K}-K)\varphi_{j}||^{2} \\ &= ||(\tilde{K}-K)||^{2} \sum_{j=1}^{\infty} ||(\tilde{K}-K)\varphi_{j}||^{2} = ||(\tilde{K}-K)||^{2} ||(\tilde{K}-$$

Thus, the proof of Lemma 2 implies: $\mathbf{E}_{x_*}||(\tilde{K}-K)^2||^2 \le \left(\mathbf{E}_{x_*}||(\tilde{K}-K)^2||_{HS}^2\right) = \mathcal{O}\left(n^{-2}\right)$. Next,

$$\begin{split} \mathcal{A}_{3} &= \mathbf{E}_{x_{*}} || (R_{\alpha}K^{2\tilde{s}} - \tilde{R}_{\alpha}\tilde{K}^{2\tilde{s}})U ||^{2} = \mathbf{E}_{x_{*}} || (R_{\alpha}K^{2\tilde{s}} - \tilde{K}^{2\tilde{s}}\tilde{R}_{\alpha})U ||^{2} \\ &= \mathbf{E}_{x_{*}} || R_{\alpha} (\alpha K^{2\tilde{s}} + K^{2\tilde{s}}\tilde{K}^{2\tilde{s}+1} - \alpha \tilde{K}^{2\tilde{s}} - K^{2\tilde{s}+1}\tilde{K}^{2\tilde{s}})\tilde{R}_{\alpha}U ||^{2} \\ &= \mathbf{E}_{x_{*}} || R_{\alpha} \left[\alpha (K^{2\tilde{s}} - \tilde{K}^{2\tilde{s}}) + K^{2\tilde{s}}(\tilde{K} - K)\tilde{K}^{2\tilde{s}} \right] \tilde{R}_{\alpha}U ||^{2} \\ &\leq 2\mathbf{E}_{x_{*}} || R_{\alpha} \alpha (K^{2\tilde{s}} - \tilde{K}^{2\tilde{s}})\tilde{R}_{\alpha}U ||^{2} + 2\mathbf{E}_{x_{*}} || R_{\alpha}K^{2\tilde{s}}(\tilde{K} - K)\tilde{K}^{2\tilde{s}}\tilde{R}_{\alpha}U ||^{2} \end{split}$$

and $2\mathbf{E}_{x_*}||R_{\alpha}K^{2\tilde{s}}(\tilde{K}-K)\tilde{K}^{2\tilde{s}}\tilde{R}_{\alpha}U||^2 = \mathcal{O}\left(\alpha^{-\frac{2a_0}{2s+a_0}}n^{-2}\alpha^{-\frac{a_0+1}{2s+a_0}}\right)$ by Lemma 2 and since $\mathbf{E}_{x_*}||R_{\alpha}K^{2\tilde{s}}||^2 = 2\mathbf{E}_{x_*}||R_{\alpha}K^{2\tilde{s}}||^2$

 $\mathcal{O}\left(\alpha^{-\frac{2a_0}{2s+a_0}}\right)$. To obtain the rate for the first term we use a Taylor expansion of the second order of the function $K^{2\tilde{s}}$ around \tilde{K} (that is, $K^{2\tilde{s}} = \tilde{K}^{2\tilde{s}} + 2\tilde{s}\tilde{K}^{2\tilde{s}-1}(|K-\tilde{K}|) + \mathcal{O}(|K-\tilde{K}|^2)$) and the \mathcal{O} term is negligible as shown above):

$$\mathbf{E}_{x_*} || R_{\alpha} \alpha (K^{2\tilde{s}} - \tilde{K}^{2\tilde{s}}) \tilde{R}_{\alpha} U ||^2 \quad = \quad \mathcal{O} \left(\mathbf{E}_{x_*} || R_{\alpha} \alpha || \mathbf{E}_{x_*} || K - \tilde{K} ||^2 \mathbf{E}_{x_*} || \tilde{K}^{2\tilde{s} - 1} \tilde{R}_{\alpha} U ||^2 \right) = \mathcal{O} \left(n^{-1} \alpha^{-\frac{4a_0 + 1}{a_0(2\tilde{s} + 1)}} n^{-1} \right)$$

by the result of Lemma 2 and since $\mathbf{E}_{x_*}||\alpha R_{\alpha}||^2 = \mathcal{O}(1)$ and $\mathbf{E}_{x_*}||\tilde{K}^{2\tilde{s}-1}\tilde{R}_{\alpha}U||^2 = \mathcal{O}\left(n^{-1}\alpha^{-\frac{4a_0+1}{2s+a_0}}\right)$. Remark that to recover the rate for $\mathbf{E}_{x_*}||\tilde{K}^{2\tilde{s}-1}\tilde{R}_{\alpha}U||^2$ a procedure similar to the one for recovering the rate for term \mathcal{A}_2 in proposition 2 has been used. By putting all these results together we obtain

$$\begin{aligned} \mathbf{E}_{x_*} || \hat{x}_{\alpha} - x_* ||^2 &= \mathcal{O} \Big(\alpha^{\frac{2\tilde{b}_0 - 1}{a_0 + 2s}} + n^{-1} \alpha^{-\frac{a_0 + 1}{a_0 + 2s}} + n^{-1} \alpha^{-\frac{2a_0}{a_0 + 2s}} \alpha^{\frac{2\tilde{b}_0 - 1}{2s + a_0}} + n^{-2} \alpha^{-2} \alpha^{\frac{2\tilde{b}_0 - 1}{2s + a_0}} \\ &+ n^{-1} \alpha^{-\frac{4a_0 + 1}{a_0(2\tilde{s} + 1)}} n^{-1} + \alpha^{-\frac{2a_0}{2s + a_0}} n^{-2} \alpha^{-\frac{a_0 + 1}{a_0 + 2s}} \Big). \end{aligned}$$

Finally, by replacing α with $\alpha_* \approx n^{-\frac{a_0+2s}{2b_0+a_0}}$, the third to sixth terms are negligible with respect to the first and second terms if $b_0 > s$ and $b_0 > a_0$. This concludes the proof.

C Proofs for Section 5

C.1 Proof of Theorem 5

To prove Theorem 5 we use Theorem 3.3 p.123 in Kuo (1975). We first rewrite this theorem and then show that the conditions of this theorem are verified in our case. The proof of Theorem 7 is given in Kuo (1975).

Theorem 7. Let P_2 be a Gaussian measure on \mathcal{Y} with mean m and covariance operator S_2 and P_1 be another Gaussian measure on the same space with mean m and covariance operator S_1 . If there exists a positive definite, bounded, continuously invertible operator \mathcal{H} such that $S_2 = S_1^{\frac{1}{2}} \mathcal{H} S_1^{\frac{1}{2}}$ and $\mathcal{H} - I$ is Hilbert-Schmidt, then P_2 is equivalent to P_1 . Moreover, the Radon-Nikodym derivative is given by

$$\frac{dP_2}{dP_1}(\{z_j\}) = \prod_{j=1}^{\infty} \sqrt{\frac{\alpha}{\lambda_j^2 + \alpha}} e^{\frac{\lambda_j^2}{2(\lambda_j^2 + \alpha)} z_j^2},$$
(29)

with $\frac{\lambda_j^2}{\alpha}$ the eigenvalues of $\mathcal{H}-I$ and z_j a sequence of real numbers.

In our case: $P_2 = P_{\alpha}$, $m = Kx_0$, $S_2 = \delta \Sigma + \frac{\delta}{\alpha} K \Omega_0 K^*$, $P_1 = P_0$ and $S_1 = \delta \Sigma$. We rewrite S_2 as

$$S_2 \equiv \left(\delta \Sigma + \frac{\delta}{\alpha} K \Omega_0 K^*\right) = \sqrt{\delta} \Sigma^{\frac{1}{2}} \left[I + \frac{1}{\alpha} \Sigma^{-1/2} K \Omega_0 K^* \Sigma^{-1/2} \right] \Sigma^{\frac{1}{2}} \sqrt{\delta} = S_1^{\frac{1}{2}} \mathcal{H} S_1^{\frac{1}{2}}$$

with $\mathcal{H} = [I + \frac{1}{\alpha} \Sigma^{-1/2} K \Omega_0 K^* \Sigma^{-1/2}] = (I + \frac{1}{\alpha} B B^*)$. In the following four points we show that \mathcal{H} satisfies all the properties required in Theorem 7.

1) \mathcal{H} is positive definite. In fact, $(I + \frac{1}{\alpha}BB^*)$ is self-adjoint, i.e. $(I + \frac{1}{\alpha}BB^*)^* = (I + \frac{1}{\alpha}BB^*)$ and $\forall \varphi \in \mathcal{Y}$,

$$\varphi \neq 0$$

$$<(I + \frac{1}{\alpha}BB^*)\varphi, \varphi > = <\varphi, \varphi > + \frac{1}{\alpha} < B^*\varphi, B^*\varphi > = ||\varphi||^2 + \frac{1}{\alpha}||B^*\varphi||^2 > 0.$$

- 2) \mathcal{H} is bounded. The operators B and B^* are bounded if Assumption A.3 holds, the operator I is bounded by definition and a linear combination of bounded operators is bounded, see Remark 2.7 in Kress (1999).
- 3) \mathcal{H} is continuously invertible. To show this, we first recall that $(I + \frac{1}{\alpha}BB^*)$ is continuously invertible if its inverse is bounded, i.e. there exists a positive number C such that $||(I + \frac{1}{\alpha}BB^*)^{-1}\varphi|| \leq C||\varphi||$, $\forall \varphi \in \mathcal{Y}$. We have $||(I + \frac{1}{\alpha}BB^*)^{-1}\varphi|| \leq (\sup_j \frac{\alpha}{\alpha + \lambda_j^2})||\varphi|| = ||\varphi||$, $\forall \varphi \in \mathcal{Y}$.
- 4) $(\mathcal{H} I)$ is Hilbert-Schmidt. To show this we have to consider the Hilbert-Schmidt norm $||\frac{1}{\alpha}BB^*||_{HS} = \frac{1}{\alpha}\sqrt{tr((BB^*)^2)}$. Now, $tr((BB^*)^2) = tr(\Omega_0T^*T\Omega_0T^*T) \leq tr(\Omega_0)||T^*T\Omega_0T^*T|| < \infty$ since $T := \Sigma^{-1/2}K$ has a bounded norm under Assumption A.1.

C.2 Proof of Theorem 6

We start by showing the first statement of the theorem and then we proceed with the proof of (20). Let \tilde{c}_3 , c_4 and \tilde{c}_4 be the constants defined in Lemma 3. Fix $\epsilon_\delta = \delta^{\frac{1-r}{4(b_0+a)}}$ for 0 < r < 1, $\alpha_1 = \left(\frac{\tilde{c}_3}{c_4}(1+\epsilon_\delta)\right)^{\frac{a+s}{b_0+a}}\delta^{\tilde{p}}$, $\alpha_2 = \left(\frac{\tilde{c}_3}{\tilde{c}_4}(1-\epsilon_\delta)\right)^{\frac{a+s}{b_0+a}}\delta^p$ where: $\tilde{p} = \frac{a+s}{b_0+a+\tilde{\eta}}$, $p = \frac{a+s}{b_0+a-\eta}$, $\tilde{\eta} = \eta I\{b_0-a-2s-1/2>0\}$ and $(b_0+a) > \eta > \max\{(b_0-s-1/2),0,(b_0-2s-a-1/2)\frac{1}{4}\}$. Remark that, since $b_0-s-1/2>(b_0-2s-a-1)$ and $0 < \epsilon_\delta < 1$ (because $0 < \delta < 1$) then the assumptions of Lemmas 5 and 6 are satisfied. Moreover, $\tilde{c}_4 \geq c_4$, see the proof of Lemma 3. Let $G := \{\exists \hat{\alpha}; S_{y^\delta}(\hat{\alpha}) = 0 \cap \hat{\alpha} \in (\alpha_2, \alpha_1)\}$ and G^c be its complement. In order to prove the first statement of the theorem we make the following decomposition:

$$P\{\exists \hat{\alpha}; S_{y^{\delta}}(\hat{\alpha}) = 0\} = P\{\exists \hat{\alpha}; S_{y^{\delta}}(\hat{\alpha}) = 0 \cap \hat{\alpha} \in (\alpha_2, \alpha_1)\} + P\{\exists \hat{\alpha}; S_{y^{\delta}}(\hat{\alpha}) = 0 \cap \hat{\alpha} \notin (\alpha_2, \alpha_1)\}$$
$$= P\{G\} + P\{G^c\}.$$

From Lemma 7 we conclude that as $\delta \to 0$: $P(G) \to 1$ and so $\hat{\alpha} \in (\alpha_2, \alpha_1)$ with probability approaching 1. This implies that we can write $\hat{\alpha}$ as a (random) convex combination of α_2 and α_1 : $\hat{\alpha} = \rho \alpha_2 + (1 - \rho)\alpha_1$ for ρ a random variable with values in (0, 1). Since $\alpha_2 \to 0$ faster than $\alpha_1 \to 0$ then:

$$\hat{\alpha} = o_p \left(\delta^{\frac{a+s}{b_0 + a + \tilde{\eta}}} \right) + \delta^{\frac{a+s}{b_0 + a + \tilde{\eta}}} (1 - \rho) \left[\frac{\tilde{c}_3}{c_4} (1 + \epsilon_\delta) \right]^{\frac{a+s}{b_0 + a}} =: o_p \left(\delta^{\frac{a+s}{b_0 + a + \tilde{\eta}}} \right) + \delta^{\frac{a+s}{b_0 + a + \tilde{\eta}}} (1 - \rho) \kappa_1. \tag{30}$$

Now, we proceed with the proof of (20). Denote by I_G and I_{G^c} the indicator functions of the events G and G^c , respectively. By the Markov's inequality, to show that $\mu_{\delta,\hat{\alpha}}^Y\{||x-x_*|| > \varepsilon_\delta M_\delta\} \to 0$ in probability we can show that its expectation, with respect to P^{x_*} , converges to 0. Then,

$$\mathbf{E}_{x_*} \left(\mu_{\delta,\hat{\alpha}}^Y \{ ||x - x_*|| > \varepsilon_{\delta} M_{\delta} \} \right) = \mathbf{E}_{x_*} \left(\mu_{\delta,\hat{\alpha}}^Y \{ ||x - x_*|| > \varepsilon_{\delta} M_{\delta} \} (1 - I_G + I_G) \right)$$

$$= \mathbf{E}_{x_*} \left(\mu_{\delta,\hat{\alpha}}^Y \{ ||x - x_*|| > \varepsilon_{\delta} M_{\delta} \} I_{G^c} \right) + \mathbf{E}_{x_*} \left(\mu_{\delta,\hat{\alpha}}^Y \{ ||x - x_*|| > \varepsilon_{\delta} M_{\delta} \} I_G \right)$$

$$\leq \mathbf{E}_{x_*} \left(I_{G^c} \right) + \mathbf{E}_{x_*} \left(\frac{||\hat{x}_{\hat{\alpha}} - x_*||^2}{\varepsilon_{\delta}^2 M_{\delta}^2} I_G + \frac{tr Var(x|y^{\delta}, \hat{\alpha}, s)}{\varepsilon_{\delta}^2 M_{\delta}^2} I_G \right)$$

$$(31)$$

where the last inequality follows from the Chebyshev's inequality. By Lemma 7: $\mathbf{E}(I_{G^c}) = P(G^c) \to 0$. We then analyze the second term of (31). Let α_* denote the optimum value of α given in proposition 3, that is, $\alpha_* = \delta^{\frac{a+s}{b_0+a}}$ with $\tilde{b}_0 = \min\{b_0, 2(a+s) + 1/2\}$. By using the notation $B(\alpha) := (\alpha I + B^*B)^{-1}$ and $\tilde{B}(\alpha) := (\alpha I + \Omega_0 K^* \Sigma^{-1} K)^{-1}$ and by adding and subtracting $\Omega_0^{\frac{1}{2}} B(\alpha_*) (\Sigma^{-1/2} B)^* (y^{\delta} - Kx_0)$ we obtain:

$$\begin{split} \mathbf{E}_{x_*} \Big(|| \hat{x}_{\hat{\alpha}} - x_* ||^2 I_G \Big) &= \mathbf{E}_{x_*} || \Omega_0^{\frac{1}{2}} B(\hat{\alpha}) (\Sigma^{-1/2} B)^* (y^{\delta} - K x_0) + x_0 - x_* ||^2 I_G \\ &\leq 2 \mathbf{E}_{x_*} || \Omega_0^{\frac{1}{2}} B(\hat{\alpha}) (\alpha_* - \hat{\alpha}) B(\alpha_*) (\Sigma^{-1/2} B)^* (y^{\delta} - K x_0) ||^2 I_G \\ &\quad + 2 \mathbf{E}_{x_*} || \Omega_0^{\frac{1}{2}} B(\alpha_*) (\Sigma^{-1/2} B)^* (y^{\delta} - K x_0) + x_0 - x_* ||^2 I_G \\ &\leq 4 \mathbf{E}_{x_*} || \tilde{B}(\hat{\alpha}) (\alpha_* - \hat{\alpha}) \Big[\Omega_0^{\frac{1}{2}} B(\alpha_*) (\Sigma^{-1/2} B)^* (y^{\delta} - K x_0) + x_0 - x_* \Big] ||^2 I_G \\ &\quad + 4 \mathbf{E}_{x_*} || \tilde{B}(\hat{\alpha}) (\alpha_* - \hat{\alpha}) (x_0 - x_*) ||^2 I_G + 2 \mathbf{E}_{x_*} || \hat{x}_{\alpha_*} - x_* ||^2 \\ &\leq 4 \mathbf{E}_{x_*} || \tilde{B}(\hat{\alpha}) (\alpha_* - \hat{\alpha}) ||^2 I_G \mathbf{E}_{x_*} || \hat{x}_{\alpha_*} - x_* ||^2 + 4 \mathbf{E}_{x_*} (\alpha_* - \hat{\alpha})^2 || \tilde{B}(\hat{\alpha}) (x_0 - x_*) ||^2 I_G + 2 \mathbf{E}_{x_*} || \hat{x}_{\alpha_*} - x_* ||^2, \end{split}$$

where the last inequality is due to the Cauchy-Schwartz inequality. By the first part of this proof, $\hat{\alpha}$ can be written as in (30) so that there exists a positive constant κ_1 such that:

$$\mathbf{E}_{x_{*}} || \tilde{B}(\hat{\alpha})(\alpha_{*} - \hat{\alpha}) ||^{2} I_{G} = \mathbf{E}_{x_{*}} (\alpha_{*} - \hat{\alpha})^{2} \sup_{||\phi||=1} \sum_{j=1}^{\infty} \frac{1}{(\hat{\alpha} + j^{-2(a+s)})^{2}} < \phi, \psi_{j} >^{2} I_{G} \leq \mathbf{E}_{x_{*}} \left(\frac{\alpha_{*} - \hat{\alpha}}{\hat{\alpha}} I_{G} \right)^{2}$$

$$= \mathbf{E}_{x_{*}} \left(\frac{\left(\delta^{\tilde{p} \frac{(b_{0} - \tilde{b}_{0} + \tilde{\eta})}{(a+\tilde{b}_{0})}} - o_{p}(1) - (1-\rho)\kappa_{1} \right)^{2}}{(o_{p}(1) + (1-\rho)\kappa_{1})^{2}} \right) = \mathcal{O}(1). \tag{32}$$

Moreover, $||\tilde{B}(\hat{\alpha})(x_0 - x_*)||^2 \approx \sum_j \frac{j^{-2b_0}}{(\hat{\alpha} + j^{-2(a+s)})^2} = \hat{\alpha}^{-2} \mathcal{A}_1$ where \mathcal{A}_1 is the term defined in the proof of proposition 3 with the only difference that α must be replaced by $\hat{\alpha}$. With reference to the notation of that proof we have:

$$\hat{\alpha}^{\frac{2b_0-1}{2(a+s)}-2}c_1 \le ||\tilde{B}(\hat{\alpha})(x_0-x_*)||^2 \le I(b_0 \ge 2(a+s)) + \hat{\alpha}^{\frac{2b_0-1}{2(a+s)}-2}\tilde{c}_1.$$

Therefore,

$$\begin{split} \delta^{\frac{2b_0-1}{2(b_0+a+\bar{\eta})}} \mathbf{E}_{x_*} [o_p(1) + (1-\rho)\kappa_1]^{\frac{2b_0-1}{2(a+s)}-2} [\delta^{\tilde{p}\frac{b_0-\bar{b}_0+\bar{\eta}}{(a+b_0)}} - o_p(1) - (1-\rho)\kappa_1]^2 c_1 &\leq \mathbf{E}_{x_*} (\alpha_* - \hat{\alpha})^2 \, ||\tilde{B}(\hat{\alpha})(x_0 - x_*)||^2 I_G \\ &\leq \left(\delta^{\frac{2(a+s)}{b_0+a+\bar{\eta}}} I(b_0 \geq 2(a+s)) + \delta^{\frac{2b_0-1}{2(b_0+a+\bar{\eta})}} \mathbf{E}_{x_*} [o_p(1) + (1-\rho)\kappa_1]^{\frac{2b_0-1}{2(a+s)}-2} \tilde{c}_1\right) \times \\ & [\delta^{\tilde{p}\frac{b_0-\bar{b}_0+\bar{\eta}}{(a+\bar{b}_0)}} - o_p(1) - (1-\rho)\kappa_1]^2, \end{split}$$

that is, $\mathbf{E}_{x_*}(\alpha_* - \hat{\alpha})^2 ||\tilde{B}(\hat{\alpha})(x_0 - x_*)||^2 I_G \approx \delta^{\frac{2b_0 - 1}{2(b_0 + a + \bar{\eta})}}$ then, we conclude that

$$\mathbf{E}_{x_*} \Big(||\hat{x}_{\hat{\alpha}} - x_*||^2 I_G \Big) \le \mathbf{E}_{x_*} ||\hat{x}_{\alpha_*} - x_*||^2 (2 + 4\mathcal{O}(1)) + 4\mathcal{O}\left(\delta^{\frac{2\tilde{b}_0 - 1}{2(b_0 + a + \tilde{\eta})}}\right). \tag{33}$$

Let us turn to the analysis of the variance term:

$$trVar(x|y^{\delta}, \hat{\alpha}, s) = \delta tr \left[\Omega_{0}(\hat{\alpha}I + B^{*}B)^{-1}\right]$$

$$= \delta tr \left[\Omega_{0}(\alpha_{*}I + B^{*}B)^{-1} + \Omega_{0}[(\hat{\alpha}I + B^{*}B)^{-1} - (\alpha_{*}I + B^{*}B)^{-1}]\right]$$

$$= \delta tr \left[\Omega_{0}(\alpha_{*}I + B^{*}B)^{-1}\right] + \delta tr \left[(\hat{\alpha}I + B^{*}B)^{-1}(\alpha_{*} - \hat{\alpha})(\alpha_{*}I + B^{*}B)^{-1}\Omega_{0}\right]$$

$$= \delta tr \left[\Omega_{0}(\alpha_{*}I + B^{*}B)^{-1}\right] + \delta \sum_{j=1}^{\infty} \frac{(\alpha_{*} - \hat{\alpha})}{(\hat{\alpha} + j^{-2(a+s)})} \frac{j^{-2s}}{(\alpha_{*} + j^{-2(a+s)})}$$

$$\leq \delta \left(1 + \frac{|\alpha_{*} - \hat{\alpha}|}{\hat{\alpha}}\right) \sum_{j=1}^{\infty} \frac{j^{-2s}}{(\alpha_{*} + j^{-2(a+s)})}.$$

by using the expression for the trace of the variance in the proof of Theorem 4. From (32) it results that $\mathbf{E}_{x_*}\left(\frac{|\alpha_*-\hat{\alpha}|}{\hat{\alpha}}I_G\right) = \mathcal{O}(1)$. Therefore,

$$\mathbf{E}_{x_*} trVar(x|y^{\delta}, \hat{\alpha}, s)I_G = (1 + \mathcal{O}(1)) tr[Var(x|y^{\delta}, \alpha_*, s)]$$

and by using (30) and the upper and lower bounds for $tr[Var(x|y^{\delta}, \alpha_*, s)]$ derived in the proof of Theorem 4 we conclude that $\mathbf{E}_{x_*}trVar(x|y^{\delta}, \hat{\alpha}, s)I_G \simeq (1 + \mathcal{O}(1)) \delta \alpha_*^{-\frac{2a+1}{2(a+s)}}$. By replacing these rates in (31) we have:

$$\begin{split} \mathbf{E}_{x_*} \left(\mu_{\delta,\hat{\alpha}}^Y \{ ||x - x_*|| > \varepsilon_{\delta} M_{\delta} \} \right) & \leq & o(1) + \mathbf{E}_{x_*} ||\hat{x}_{\alpha_*} - x_*||^2 \frac{(2 + 4\mathcal{O}(1))}{\varepsilon_{\delta}^2 M_{\delta}^2} + \frac{4\mathcal{O}\left(\delta^{\frac{2\bar{b}_0 - 1}{2(b_0 + a + \bar{\eta})}}\right)}{\varepsilon_{\delta}^2 M_{\delta}^2} + (1 + \mathcal{O}(1)) \frac{\delta \alpha_*^{-\frac{2a + 1}{2(a + s)}}}{\varepsilon_{\delta}^2 M_{\delta}^2} \\ & = & \mathcal{O}\left(\frac{1}{\varepsilon_{\delta}^2 M_{\delta}^2} \left(\delta^{\frac{2\bar{b}_0 - 1}{2(b_0 + a)}} + \delta^{\frac{2\bar{b}_0 - 1}{2(b_0 + a + \bar{\eta})}}\right)\right) \end{split}$$

which converges to 0 for $\varepsilon_{\delta} \simeq \delta^{\frac{2\tilde{b}_0 - 1}{4(b_0 + a + \tilde{\eta})}}$.

D Technical Lemmas

Lemma 2. Let $\mathcal{X} = \mathcal{Y} = L^2([0,1])$ and K and \tilde{K} be operators from \mathcal{X} to \mathcal{Y} such that $K\varphi = \frac{1}{n}\sum_{i=1}^n \langle Z_i, \varphi \rangle Z_i(t)$ and $\tilde{K}\varphi = \int_0^1 \varphi(s)Cov(Z(s), Z(t))ds$, $\forall \varphi \in \mathcal{X}$, where $Z \in L^2([0,1])$ is a centered random function. Then, if $\mathbf{E}||Z||^4 < \infty$ we have

$$\mathbf{E}_{x_*}||K - \tilde{K}||^2 = \mathcal{O}\left(\frac{1}{n}\right).$$

Proof. Let $||\cdot||_{HS}$ denote the Hilbert-Schmidt norm. Since $Cov(Z(s), Z(t)) = \mathbf{E}(Z(s)Z(t))$ we have

$$||K - \tilde{K}||^2 \le ||K - \tilde{K}||_{HS}^2 := \int_0^1 \int_0^1 \left[\frac{1}{n} \sum_{i=1}^n Z_i(s) Z_i(t) - \mathbf{E}(Z(s)Z(t))\right]^2 ds dt$$

so that

$$\begin{split} \mathbf{E}_{x_*} || K - \tilde{K} ||^2 & \leq \int_0^1 \int_0^1 \mathbf{E}_{x_*} \Big[\frac{1}{n} \sum_{i=1}^n Z_i(s) Z_i(t) - \mathbf{E}(Z(s) Z(t)) \Big]^2 ds dt \\ & = \int_0^1 \int_0^1 Var \Big[\frac{1}{n} \sum_{i=1}^n Z_i(s) Z_i(t) \Big] ds dt \\ & = \int_0^1 \int_0^1 \Big[\frac{1}{n^2} \sum_{i=1}^n Var(Z_i(s) Z_i(t)) + \frac{2n(n-1)}{n^2} \sum_{i>j} Cov(Z_i(s) Z_i(t), Z_j(s) Z_j(t)) \Big] ds dt \\ & = \int_0^1 \int_0^1 \frac{1}{n} \mathbf{E}(Z^2(s) Z^2(t)) ds dt = \frac{1}{n} \mathbf{E} \int_0^1 \int_0^1 (Z^2(s) Z^2(t)) ds dt = \frac{1}{n} \mathbf{E} ||Z||^4 \end{split}$$

since $Cov(Z_i(s)Z_i(t), Z_j(s)Z_j(t)) = 0.$

Lemma 3. Let the assumptions of Theorem 6 be satisfied. Then,

$$S_{y^{\delta}}(\alpha) = S_1 - S_2 - S_3 - (S_{4a} + S_{4b}) + \nu_1 \frac{1}{\alpha} + \nu_2$$

where

$$\frac{1}{2\delta} \alpha^{\frac{b_0 - s - \frac{1}{2}}{a + s}} c_4 \leq \mathcal{S}_2 \leq \frac{1}{2} \delta^{-1} I(b_0 \geq s) + \frac{1}{2\delta} \alpha^{\frac{b_0 - s - \frac{1}{2}}{a + s}} \tilde{c}_4,
\frac{1}{2} \alpha^{-\frac{2(a+s)+1}{2(a+s)}} \tilde{c}_3 \leq (\mathcal{S}_1 - \mathcal{S}_{4b}) \leq \frac{1}{2} \alpha^{-\frac{2(a+s)+1}{2(a+s)}} \tilde{c}_3 + \frac{1}{2\alpha},$$

$$S_3 \sim \xi_1 \frac{1}{\sqrt{\delta}} \left(\sum_{j=1}^{\infty} \frac{\lambda_j^4}{(\alpha + \lambda_j^2)^4} j^{-2(a+b_0)} \right)^{\frac{1}{2}}, \qquad S_{4a} \sim \xi_2 \frac{1}{\sqrt{2}} \left(\sum_{j=1}^{\infty} \frac{\lambda_j^4}{(\alpha + \lambda_j^2)^4} \right)^{\frac{1}{2}}$$

where $\xi_i \sim \mathcal{N}(0,1)$ for i = 1, 2, $Cov(S_3, S_{4a}) = 0$ and \tilde{c}_3 , c_4 , \tilde{c}_4 are positive constants.

Proof. We develop $S_{y^{\delta}}(\alpha)$ by using the fact that under Assumption C there exist $\underline{\lambda}, \overline{\lambda} > 0$ such that $\underline{\lambda}j^{-(a+s)} \leq \lambda_j \leq \overline{\lambda}j^{-(a+s)}$ for $j = 1, 2, \ldots$ Then,

$$\begin{split} S_{y^{\delta}}(\alpha) &= \frac{1}{2} \sum_{j=1}^{\infty} \frac{\lambda_j^2}{\alpha(\alpha + \lambda_j^2)} - \frac{1}{2} \sum_{j=1}^{\infty} \frac{\lambda_j^2}{\delta(\alpha + \lambda_j^2)^2 \lambda_{j\Sigma}} < K(x_* - x_0), \varphi_j >^2 \\ &- \frac{1}{2} \sum_{j=1}^{\infty} \frac{2\lambda_j^2}{\delta(\alpha + \lambda_j^2)^2 \lambda_{j\Sigma}} < K(x_* - x_0), \varphi_j > < U, \varphi_j > \\ &- \frac{1}{2} \sum_{j=1}^{\infty} \frac{\lambda_j^2}{\delta(\alpha + \lambda_j^2)^2 \lambda_{j\Sigma}} < U, \varphi_j >^2 := \mathcal{S}_1 - \mathcal{S}_2 - \mathcal{S}_3 - \mathcal{S}_4. \end{split}$$

Let us start by computing S_2 :

$$S_{2} = \frac{1}{2} \sum_{j=1}^{\infty} \frac{\lambda_{j}^{2}}{\delta(\alpha + \lambda_{j}^{2})^{2} \lambda_{j\Sigma}} < K(x_{*} - x_{0}), \varphi_{j} >^{2}$$

$$\approx \frac{1}{2\delta} \sum_{j=1}^{\infty} \frac{j^{-2(a+s)-2a-2b_{0}}}{(\alpha + j^{-2(a+s)})^{2}} = \frac{1}{2\delta} \sum_{j=1}^{\infty} \frac{j^{-4a-2s-2b_{0}}}{(\alpha + j^{-2(a+s)})^{2}}$$

The function $f(j) = \frac{j^{-4a-2s-2b_0}}{(\alpha+j^{-2(a+s)})^2}$ defined on \mathbb{R}_+ is always decreasing in j when $b_0 \geq s$. When $b_0 < s$, f(j) is increasing for $j < \bar{j}$ and decreasing for $j > \bar{j}$ where $\bar{j} = \left(\frac{2a+s+b_0}{s-b_0}\alpha\right)^{-\frac{1}{2(a+s)}}$. Therefore, to upper and lower bound \mathcal{S}_2 we have to consider these two cases separately. First, if $b_0 < s$:

$$\frac{1}{2\delta} \sum_{j=1}^{\infty} \frac{j^{-4a-2s-2b_0}}{(\alpha+j^{-2(a+s)})^2} \geq \frac{1}{2\delta} \sum_{j=\bar{j}}^{\infty} \frac{j^{-4a-2s-2b_0}}{(\alpha+j^{-2(a+s)})^2} \geq \frac{1}{2\delta} \alpha^{\frac{b_0-s-\frac{1}{2}}{a+s}} \int_{\underline{u}}^{\infty} \frac{u^{2(s-b_0)}}{(u^{2(a+s)}+1)^2} du := \frac{1}{2\delta} \alpha^{\frac{b_0-s-\frac{1}{2}}{a+s}} c_4(\underline{u});$$

$$\frac{1}{2\delta} \sum_{j=1}^{\infty} \frac{j^{-4a-2s-2b_0}}{(\alpha+j^{-2(a+s)})^2} \leq \frac{1}{2\delta} \left(\sum_{j=1}^{\bar{j}} \frac{j^{-4a-2s-2b_0}}{(\alpha+j^{-2(a+s)})^2} + \alpha^{\frac{b_0-s-\frac{1}{2}}{a+s}} \int_{\underline{u}}^{\infty} \frac{u^{2(s-b_0)}}{(u^{2(a+s)}+1)^2} du \right)$$

$$\leq \frac{1}{2\delta} \left(\bar{j} \frac{\bar{j}^{-4a-2s-2b_0}}{(\alpha+\bar{j}^{-2(a+s)})^2} + \alpha^{\frac{b_0-s-\frac{1}{2}}{a+s}} \int_{\underline{u}}^{\infty} \frac{u^{2(s-b_0)}}{(u^{2(a+s)}+1)^2} du \right)$$

$$= \frac{1}{2\delta} \left(\frac{\left(\alpha^{\frac{2a+s+b_0}{(s-b_0)}}\right)^{-\frac{(1-4a-2s-2b_0)}{2(a+s)}}}{(\alpha+\alpha^{\frac{2a+s+b_0}{2(a+s)}})^2} + \alpha^{\frac{b_0-s-\frac{1}{2}}{a+s}} c_4(\underline{u}) \right) = \frac{1}{2\delta} \alpha^{\frac{b_0-s-1/2}{(a+s)}} \tilde{c}_4(\underline{u}),$$

where $\underline{u} = \left(\frac{2a+s+b_0}{s-b_0}\right)^{-\frac{1}{2(a+s)}}$, $c_4(\underline{u}) = \int_{\underline{u}}^{\infty} \frac{u^{2(s-b_0)}}{(u^{2(a+s)}+1)^2} du$ and $\tilde{c}_4(\underline{u}) := (2a+s+b_0)^{-\frac{1-4a-2s-2b_0}{2(a+s)}} \frac{(s-b_0)^{\frac{1+2s-2b_0}{2(a+s)}}}{4(a+s)^2} + c_4(\underline{u})$. Second, if $b_0 \ge s$:

$$\frac{1}{2\delta} \sum_{j=1}^{\infty} \frac{j^{-4a-2s-2b_0}}{(\alpha+j^{-2(a+s)})^2} \geq \frac{1}{2\delta} \alpha^{\frac{b_0-s-\frac{1}{2}}{a+s}} \int_{\bar{u}}^{\infty} \frac{u^{2(s-b_0)}}{(u^{2(a+s)}+1)^2} du =: \frac{1}{2\delta} \alpha^{\frac{b_0-s-\frac{1}{2}}{a+s}} c_4(\bar{u});$$

$$\frac{1}{2\delta} \sum_{j=1}^{\infty} \frac{j^{-4a-2s-2b_0}}{(\alpha+j^{-2(a+s)})^2} \leq \frac{1}{2\delta} \left(\frac{j^{-4a-2s-2b_0}}{(\alpha+j^{-2(a+s)})^2} \Big|_{j=1} + \alpha^{\frac{b_0-s-\frac{1}{2}}{a+s}} \int_{\bar{u}}^{\infty} \frac{u^{2(s-b_0)}}{(u^{2(a+s)}+1)^2} du + \right)$$

$$\leq \frac{1}{2\delta} \alpha^{\frac{b_0-s-1/2}{a+s}} c_4(\bar{u}) + \frac{1}{2\delta}.$$

where $\bar{u} = \alpha^{\frac{1}{2(a+s)}}$. By defining $c_4 = c_4(\underline{u})I(b_0 < s) + c_4(\bar{u})I(b_0 \ge s)$ and $\tilde{c}_4 = \tilde{c}_4(\underline{u})I(b_0 < s) + c_4(\bar{u})I(b_0 \ge s)$, the first line of inequalities of the lemma is proved.

We analyze now S_3 . Let $\{\xi_j\}$ denote a sequence of independent $\mathcal{N}(0,1)$ random variables; we rewrite S_3 as

$$S_{3} = \sum_{j=1}^{\infty} \frac{\lambda_{j}^{2}}{\sqrt{\delta}(\alpha + \lambda_{j}^{2})^{2}\sqrt{\lambda_{j}\Sigma}} < K(x_{*} - x_{0}), \varphi_{j} > \frac{< U, \varphi_{j} >}{\sqrt{\delta}\lambda_{j}\Sigma}$$
$$= \frac{1}{\sqrt{\delta}} \sum_{j=1}^{\infty} \frac{\lambda_{j}^{2}}{(\alpha + \lambda_{j}^{2})^{2}\sqrt{\lambda_{j}\Sigma}} < K(x_{*} - x_{0}), \varphi_{j} > \xi_{j}.$$

This series is convergent if, for a fixed α , $\mathbf{E}||\mathcal{S}_3||^2 < \infty$. This is always verified because:

$$\mathbf{E}||\mathcal{S}_{3}||^{2} \approx \frac{1}{\delta} \sum_{j=1}^{\infty} \frac{\lambda_{j}^{4}}{(\alpha + \lambda_{j}^{2})^{4}} j^{-2(a+b_{0})}$$

$$\leq \frac{1}{\delta} \left(\sup_{j} \frac{\lambda_{j}^{2}}{(\alpha + \lambda_{j}^{2})^{2}} \right)^{2} \sum_{j=1}^{\infty} j^{-2(a+b_{0})} = \frac{1}{\delta \alpha} \sum_{j=1}^{\infty} j^{-2(a+b_{0})}$$

and it is finite if and only if $b_0 > \frac{1}{2} - a$ which is verified by assumption. Therefore, S_3 is equal in distribution to a Gaussian random variable with 0 mean and variance $\frac{1}{\delta} \sum_{j=1}^{\infty} \frac{\lambda_j^4}{(\alpha + \lambda_j^2)^4} j^{-2(a+b_0)}$. 6

Consider now term S_4 . Since $\frac{\langle U, \varphi_j \rangle^2}{\delta \lambda_{j\Sigma}} \sim i.i.d. \chi_1^2$, we center term S_4 around its mean and apply the Lyapunov Central Limit Theorem to term S_{4a} below:

$$S_4 = \frac{1}{2} \sum_{j=1}^{\infty} \frac{\lambda_j^2}{(\alpha + \lambda_j^2)^2} \left(\frac{\langle U, \varphi_j \rangle^2}{\delta \lambda_{j\Sigma}} - 1 \right) + \frac{1}{2} \sum_{j=1}^{\infty} \frac{\lambda_j^2}{(\alpha + \lambda_j^2)^2} := S_{4a} + S_{4b}.$$

For that, we verify the Lyapunov condition. Denote $\vartheta_j := \left(\frac{\langle U, \varphi_j \rangle^2}{\delta \lambda_{j\Sigma}} - 1\right)$ and it is easy to show that $\mathbf{E}|\vartheta_j|^3 = 8.6916$, thus $\sum_{j=1}^{\infty} \mathbf{E} \left|\frac{\lambda_j^2 \vartheta_j}{(\alpha + \lambda_j^2)^2}\right|^3 = 8.6916 \sum_{j=1}^{\infty} \frac{j^{-6(a+s)}}{(\alpha + j^{-2(a+s)})^6}$ which is upper bounded by

$$\sum_{j=1}^{\infty} \mathbf{E} \left| \frac{\lambda_{j}^{2} \vartheta_{j}}{(\alpha + \lambda_{j}^{2})^{2}} \right|^{3} \leq 8.6916 \left[\sum_{j=1}^{\bar{j}} \frac{j^{-6(a+s)}}{(\alpha + j^{-2(a+s)})^{6}} + \sum_{j=\bar{j}}^{\infty} \frac{j^{-6(a+s)}}{(\alpha + j^{-2(a+s)})^{6}} \right] \\
= 8.6916 \left[\frac{\bar{j}^{-6(a+s)+1}}{(\alpha + \bar{j}^{-2(a+s)})^{6}} + \alpha^{-3 - \frac{1}{2(a+s)}} \int_{1}^{-\infty} \frac{u^{6(a+s)}}{(u^{2(a+s)} + 1)^{6}} du \right] \\
= 8.6916 \alpha^{-3 - \frac{1}{2(a+s)}} \left[2^{-6} + \int_{1}^{-\infty} \frac{u^{6(a+s)}}{(u^{2(a+s)} + 1)^{6}} du \right]$$

since the function $f(j) = \frac{j^{-6(a+s)}}{(\alpha+j^{-2(a+s)})^6}$ is increasing in j for $0 < j < \bar{j} := \alpha^{-\frac{1}{2(a+s)}}$ and decreasing for $j > \bar{j}$. By using the lower bound of $Var(S_{4a})$ given in Lemma 4 below we obtain:

$$(Var(\mathcal{S}_{4a}))^{-3/2} \sum_{j=1}^{\infty} \mathbf{E} \left| \frac{\lambda_j^2 \vartheta}{(\alpha + \lambda_j^2)^2} \right|^3 \leq 8.6916 \left(\frac{1}{2} \alpha^{-\frac{4(a+s)+1}{2(a+s)}} c_6 \right)^{-3/2} \alpha^{-3-\frac{1}{2(a+s)}} \left[2^{-6} + \int_1^{-\infty} \frac{u^{6(a+s)}}{(u^{2(a+s)}+1)^6} du \right]$$

which converges to 0 so that the Lyapunov condition is satisfied and $S_{4a} = \frac{1}{2} \left(2 \sum_{j=1}^{\infty} \frac{\lambda_j^4}{(\alpha + \lambda_j^2)^4} \right)^{\frac{1}{2}} \xi_2$ where

$$(Var(S_3))^{-(2+\varrho)/2} \sum_{i=1}^{\infty} \left(\frac{\lambda_j^2}{\delta(\alpha + \lambda_j^2)^2 \lambda_{j\Sigma}} < K(x_* - x_0), \varphi_j > \right)^{2+\varrho} \mathbf{E} \left| < U, \varphi_j > -\mathbf{E}(< U, \varphi_j >) \right|^{2+\varrho} = 0.$$

If this condition is satisfied then S_3 is equal in distribution to a Gaussian random variable with mean $\sum_{j=1}^{\infty} \frac{j^{-2(2a+s-b_0)}}{\delta(\alpha+j^{-2(a+s)})^2} \mathbf{E} < U, \varphi_j > \text{and variance } \frac{1}{\delta} \sum_{j=1}^{\infty} \frac{\lambda_j^4}{(\alpha+\lambda_j^2)^4} j^{-2(a+b_0)+c_0} Var\left(\frac{< U, \varphi_j>}{\sqrt{\delta}}\right)$. Asymptotically, this mean is 0 and the variance is equal to the expression for $\mathbf{E}||S_3||^2$ given above.

⁶In the case where U is only asymptotically $\mathcal{N}(0, \delta\Sigma)$ then the analysis of \mathcal{S}_3 requires to use the Lyapunov Central Limit Theorem. In particular, we have to check the Lyapunov condition: for some $\varrho > 0$

 $\xi_2 \sim \mathcal{N}(0,1)$. To show that the $Cov(\mathcal{S}_3, \mathcal{S}_{4a})$ is zero, rewrite $\mathcal{S}_3 = \sum_{j=1}^{\infty} a_j \xi_j$ and $\mathcal{S}_{4a} = \sum_{j=1}^{\infty} b_j (\xi_j^2 - 1)$ with $\xi_j \sim \mathcal{N}(0,1)$. Thus, $Cov(\mathcal{S}_3, \mathcal{S}_{4a}) = \sum_{j=1}^{\infty} \sum_{l=1}^{\infty} a_j b_j \mathbf{E}(\xi_j(\xi_j^2 - 1)) = 0$.

Term S_{4b} is non random and we subtract it from S_1 to obtain:

$$S_{1} - S_{4b} = \frac{1}{2} \left[\sum_{j=1}^{\infty} \frac{\lambda_{j}^{2}}{\alpha(\alpha + \lambda_{j}^{2})} - \sum_{j=1}^{\infty} \frac{\lambda_{j}^{2}}{(\alpha + \lambda_{j}^{2})^{2}} \right] = \frac{1}{2\alpha} \sum_{j=1}^{\infty} \frac{\lambda_{j}^{4}}{(\alpha + \lambda_{j}^{2})^{2}}$$
$$\approx \frac{1}{2\alpha} \sum_{j=1}^{\infty} \frac{j^{-4(a+s)}}{(\alpha + j^{-2(a+s)})^{2}} = \frac{1}{2\alpha} \sum_{j=1}^{\infty} \frac{1}{(\alpha j^{2(a+s)} + 1)^{2}}.$$

Since $\frac{j^{-4(a+s)}}{(\alpha+j^{-2(a+s)})^2}$ is a decreasing function of j we have

$$\frac{1}{2\alpha} \int_{1}^{\infty} \frac{1}{(\alpha t^{2(a+s)}+1)^2} dt \leq \frac{1}{2\alpha} \sum_{j=1}^{\infty} \frac{1}{(\alpha j^{2(a+s)}+1)^2} \leq \frac{1}{2\alpha} \int_{1}^{\infty} \frac{1}{(\alpha t^{2(a+s)}+1)^2} dt + \frac{1}{2\alpha(\alpha+1)^2} dt$$

By denoting $c_3(\bar{u}) = \int_{\bar{u}}^{\infty} \frac{1}{(u^{2(a+s)}+1)^2} du$, with \bar{u} defined above, and since $\frac{1}{(\alpha+1)^2} < 1$ for a fixed α , after some algebra we conclude that $\frac{1}{2}\alpha^{-\frac{2(a+s)+1}{2(a+s)}}c_3(\bar{u}) \leq (\mathcal{S}_1 - \mathcal{S}_{4b}) \leq \frac{1}{2}\alpha^{-\frac{2(a+s)+1}{2(a+s)}}c_3(\bar{u}) + \frac{1}{2\alpha}$. By defining $\tilde{c}_3 = c_3(\bar{u})$, the second line of inequalities of the lemma is proved.

Lemma 4. Let the assumptions of Theorem 6 be satisfied and S_3 , S_{4a} be as defined in Lemma 3. Then,

$$\frac{1}{\delta} \alpha^{-\frac{a+2s-b_0+\frac{1}{2}}{(a+s)}} c_5 \le Var(\mathcal{S}_3) \le \frac{1}{\delta} I(b_0 \ge a+2s) + \frac{1}{\delta} \alpha^{-\frac{a+2s-b_0+\frac{1}{2}}{(a+s)}} \tilde{c}_5, \tag{34}$$

$$\frac{1}{2}\alpha^{-\frac{4(a+s)+1}{2(a+s)}}c_6 \le Var(\mathcal{S}_{4a}) \le \frac{17\tilde{c}_6}{32}\alpha^{-\frac{4(a+s)+1}{2(a+s)}}$$
(35)

where c_5 , \tilde{c}_5 , c_6 and \tilde{c}_6 are positive constants.

Proof. We start by considering $Var(S_3) = \frac{1}{\delta} \sum_{j=1}^{\infty} \frac{\lambda_j^4}{(\alpha + \lambda_j^2)^4} j^{-2(a+b_0)}$. Under Assumption C, we can rewrite $Var(S_3) \approx \frac{1}{\delta} \sum_{j=1}^{\infty} \frac{j^{-4(a+s)-2a-2b_0}}{(\alpha + j^{-2(a+s)})^4}$. The function $f(j) = \frac{j^{-4(a+s)-2a-2b_0}}{(\alpha + j^{-2(a+s)})^4}$ defined on \mathbb{R}_+ is always decreasing in j when $b_0 \geq a + 2s$. When $b_0 < a + 2s$, f(j) is increasing for $j < \bar{j}$ and decreasing for $j > \bar{j}$ where $\bar{j} = \left(\frac{3a+2s+b_0}{a+2s+b_0}\alpha\right)^{-\frac{1}{2(a+s)}}$. Therefore, in order to find an upper and a lower bound for $Var(S_3)$ we consider these two cases separately. Let us start with the case $b_0 < a + 2s$:

$$Var(\mathcal{S}_{3}) \geq \frac{1}{\delta} \sum_{j=\bar{j}}^{\infty} \frac{j^{-4(a+s)-2a-2b_{0}}}{(\alpha+j^{-2(a+s)})^{4}} \geq \frac{1}{\delta} \alpha^{-\frac{a+2s-b_{0}+\frac{1}{2}}{(a+s)}} \int_{\underline{u}}^{\infty} \frac{u^{2a+4s-2b_{0}}}{(u^{2(a+s)}+1)^{4}} du =: \frac{1}{\delta} \alpha^{-\frac{a+2s-b_{0}+\frac{1}{2}}{(a+s)}} c_{5}(\underline{u});$$

$$Var(\mathcal{S}_{3}) \leq \frac{1}{\delta} \sum_{j=1}^{\bar{j}} \frac{j^{-4(a+s)-2a-2b_{0}}}{(\alpha+j^{-2(a+s)})^{4}} + \frac{1}{\delta} \alpha^{-\frac{a+2s-b_{0}+\frac{1}{2}}{(a+s)}} c_{5}(\underline{u}) \leq \frac{1}{\delta} \bar{j} \frac{\bar{j}^{-4(a+s)-2a-2b_{0}}}{(\alpha+\bar{j}^{-2(a+s)})^{4}} + \frac{1}{\delta} \alpha^{-\frac{a+2s-b_{0}+\frac{1}{2}}{(a+s)}} c_{5}(\underline{u})$$

$$= \frac{1}{\delta} \alpha^{-\frac{(1+2a+4s-2b_{0})}{2(a+s)}} \left(\underline{u}^{(1+2a+4s-2b_{0})} \left(\frac{a+2s+b_{0}}{4(a+s)} \right) + c_{5}(\underline{u}) \right) =: \frac{1}{\delta} \alpha^{-\frac{(1+2a+4s-2b_{0})}{2(a+s)}} \tilde{c}_{5}(\underline{u})$$

⁷In the case where U is only asymptotically $\mathcal{N}(0,\delta\Sigma)$ then a remark similar to the one for term \mathcal{S}_3 applies.

where $\underline{u} = \left(\frac{3a+2s+b_0}{a+2s+b_0}\right)^{-\frac{1}{2(a+s)}}$. Next, we consider the case $b_0 \ge a+2s$:

$$Var(S_{3}) \geq \frac{1}{\delta}\alpha^{-\frac{a+2s-b_{0}+\frac{1}{2}}{(a+s)}} \int_{\bar{u}}^{\infty} \frac{u^{2a+4s-2b_{0}}}{(u^{2(a+s)}+1)^{4}} du =: \frac{1}{\delta}\alpha^{-\frac{a+2s-b_{0}+\frac{1}{2}}{(a+s)}} c_{5}(\bar{u});$$

$$Var(S_{3}) \leq \frac{1}{\delta} \left(\frac{j^{-4(a+s)-2a-2b_{0}}}{(\alpha+j^{-2(a+s)})^{4}} \Big|_{j=1} + \alpha^{-\frac{a+2s-b_{0}+\frac{1}{2}}{(a+s)}} \int_{\bar{u}}^{\infty} \frac{u^{2a+4s-2b_{0}}}{(u^{2(a+s)}+1)^{4}} du \right)$$

$$\leq \frac{1}{\delta} + \frac{1}{\delta}\alpha^{-\frac{(1+2a+4s-2b_{0})}{2(a+s)}} c_{5}(\bar{u})$$

where $\bar{u} = \alpha^{\frac{1}{2(a+s)}}$. By defining $c_5 = c_5(\underline{u})I(b_0 < a+2s) + c_5(\bar{u})I(b_0 \ge a+2s)$ and $\tilde{c}_5 = \tilde{c}_5(\underline{u})I(b_0 < a+2s) + c_5(\bar{u})I(b_0 \ge a+2s)$, the first inequalities in (34) is proved.

Under Assumption C, the variance of S_{4a} rewrites as $Var(S_{4a}) \approx \frac{1}{2} \sum_{j=1}^{\infty} \frac{j^{-4(a+s)}}{(\alpha+j^{-2(a+s)})^4}$. The function $f(j) = \frac{j^{-4(a+s)}}{(\alpha+j^{-2(a+s)})^4}$ defined on \mathbb{R}_+ is increasing in j for $j < \bar{j}$ and decreasing for $j > \bar{j}$ where $\bar{j} = \alpha^{-\frac{1}{2(a+s)}}$. We can lower and upper bound $Var(S_{4a})$ as follows:

$$Var(S_{4a}) \geq \frac{1}{2} \sum_{j=\bar{j}}^{\infty} \frac{j^{-4(a+s)}}{(\alpha+j^{-2(a+s)})^4} \geq \frac{1}{2} \alpha^{-2-\frac{1}{2(a+s)}} \int_{1}^{\infty} \frac{u^{4(a+s)}}{(u^{2(a+s)}+1)^4} du =: \frac{1}{2} \alpha^{-\frac{4(a+s)+1}{2(a+s)}} c_{6s}$$

$$Var(S_{4a}) \leq \frac{1}{2} \sum_{j=1}^{\bar{j}} \frac{j^{-4(a+s)}}{(\alpha+j^{-2(a+s)})^4} + \frac{1}{2} \alpha^{-\frac{4(a+s)+1}{2(a+s)}} c_{6s} \leq \frac{1}{2} \bar{j} \frac{\bar{j}^{-4(a+s)}}{(\alpha+\bar{j}^{-2(a+s)})^4} + \frac{1}{2} \alpha^{-\frac{4(a+s)+1}{2(a+s)}} c_{6s}$$

$$= \frac{1}{32} \alpha^{-\frac{1+4a+4s}{2(a+s)}} + \frac{c_{6}}{2} \alpha^{-\frac{1+4a+4s}{2(a+s)}} = \frac{17c_{6}+1}{32} \alpha^{-\frac{1+4(a+s)}{2(a+s)}} =: \frac{17\tilde{c}_{6}}{32} \alpha^{-\frac{1+4(a+s)}{2(a+s)}}.$$

Lemma 5. Let the assumptions of Theorem 6 be satisfied and $\alpha_1 = \left(\frac{\tilde{c}_3}{c_4}(1+\epsilon_\delta)\right)^{\frac{a+s}{b_0+a}} \delta^{\tilde{p}}$ where: $\tilde{p} = \frac{a+s}{b_0+a+\tilde{\eta}}$ for $\tilde{\eta} = \eta I\{b_0 - a - 2s - 1/2 > 0\}$, $\eta > \max\{(b_0 - 2s - a - 1), (b_0 - 2s - a - 1/2)\frac{1}{4}\}$, $\epsilon_{\delta} = \delta^{\frac{1-r}{4(b_0+a)}}$, for every $0 < r < \frac{4(a+s)(b_0+a)}{b_0+a+\tilde{\eta}}$ and \tilde{c}_3 , c_4 be as in Lemma 3. Then,

$$P(S_{y^{\delta}}(\alpha_1) < 0) \to 1 \quad as \quad \delta \to 0.$$

Proof. By using the notation of Lemma 3 we write

$$S_{y^{\delta}}(\alpha) := -S_2 - S_3 - S_{4a} + (S_1 - S_{4b}) + \frac{\nu_1}{\alpha} + \nu_2$$

where S_3 and S_{4a} are independent zero-mean Gaussian random variables with variances equal to

$$Var(S_3) = \frac{1}{\delta} \sum_{j=1}^{\infty} \frac{\lambda_j^4}{(\alpha + \lambda_j^2)^4} j^{-2(a+b_0)}$$
 and $Var(S_{4a}) = \frac{1}{2} \sum_{j=1}^{\infty} \frac{\lambda_j^4}{(\alpha + \lambda_j^2)^4}$

respectively. By using the lower bound of S_2 and the upper bound of $(S_1 - S_{4b})$ provided in Lemma 3 and

by denoting $D(\alpha) = Var(\mathcal{S}_3) + Var(\mathcal{S}_4)$ we obtain:

$$P(S_{y^{\delta}}(\alpha_{1}) < 0) = P(S_{3} + S_{4a} > -S_{2} + (S_{1} - S_{4b}) + \frac{\nu_{1}}{\alpha} + \nu_{2})$$

$$\geq P\left(\xi > \frac{-\frac{1}{2}\delta^{-1}\alpha_{1}^{\frac{b_{0}-s-1/2}{a+s}}c_{4} + \nu_{2} + \alpha_{1}^{-\frac{2(a+s)+1}{2(a+s)}}\frac{1}{2}(\tilde{c}_{3} + \alpha_{1}^{\frac{1}{2(a+s)}}(1+2\nu_{1}))}\right) \quad (36)$$

where ξ denotes a $\mathcal{N}(0,1)$ random variable. Moreover, let $D^u(\alpha)$ be an upper bound for $D(\alpha)$ and $D_l(\alpha)$ be a lower bound for $D(\alpha)$ for every α . By the result in Lemma 4 we can take

$$D^{u}(\alpha) = \delta^{-1}I(b_0 \ge a + 2s) + \tilde{c}_5 \delta^{-1} \alpha^{-\frac{a+2s-b_0+1/2}{a+s}} + \frac{17}{32} \alpha^{-\frac{4(a+s)+1}{2(a+s)}} \tilde{c}_6$$
 (37)

$$D_l(\alpha) = \delta^{-1} \alpha^{-\frac{a+2s-b_0+1/2}{a+s}} c_5 + \frac{1}{2} \alpha^{-\frac{4(a+s)+1}{2(a+s)}} c_6$$
(38)

and by replacing the value of α_1 and after some algebras we get:

$$D^{u}(\alpha_{1}) = \delta^{-1}I(b_{0} \geq a + 2s) + \tilde{c}_{5}\delta^{-1}\delta^{-\frac{a+2s-b_{0}+1/2}{b_{0}+a+\tilde{\eta}}} \left[\frac{\tilde{c}_{3}}{c_{4}}(1+\epsilon_{\delta})\right]^{-\frac{a+2s-b_{0}+1/2}{b_{0}+a}} + \frac{17\tilde{c}_{6}}{32}\delta^{-\frac{4(a+s)+1}{2(b_{0}+a+\tilde{\eta})}} \left[\frac{\tilde{c}_{3}}{c_{4}}(1+\epsilon_{\delta})\right]^{-\frac{4(a+s)+1}{2(b_{0}+a)}}.$$

$$D_{l}(\alpha_{1}) = \delta^{-1}\delta^{-\frac{a+2s-b_{0}+1/2}{b_{0}+a+\tilde{\eta}}} \left(\left[\frac{\tilde{c}_{3}}{c_{4}}(1+\epsilon_{\delta})\right]^{-\frac{a+2s-b_{0}+1/2}{b_{0}+a}} c_{5} + \frac{1}{2}\delta^{\frac{\tilde{\eta}}{(b_{0}+a+\tilde{\eta})}} \left[\frac{\tilde{c}_{3}}{c_{4}}(1+\epsilon_{\delta})\right]^{-\frac{2(a+s)+1/2}{b_{0}+a}} c_{6}\right).$$

Remark that: (1) $D^u(\alpha_1) = \mathcal{O}(\delta^{-\frac{4(a+s)+1}{2(b_0+a)}})$ when $b_0 \leq \frac{1}{2} + a + 2s$ and (2) $D^u(\alpha_1) = \mathcal{O}(\delta^{-1})$ when $b_0 > \frac{1}{2} + a + 2s$. Therefore, we analyze these two cases separately.

CASE I: $b_0 \leq \frac{1}{2} + a + 2s$. By substituting the value of α_1 in the numerator of (36), factorizing the term $\delta^{-\frac{2(a+s)+1}{2(b_0+a)}}$ (in the first and third term of the numerator) and after some algebra we obtain

$$\begin{split} P(S_{y^{\delta}}(\alpha_{1}) < 0) & \geq & P\left(\xi > \frac{-\frac{1}{2}\delta^{-\frac{2(a+s)+1}{2(b_{0}+a)}}\left(\tilde{c}_{3}^{\frac{b_{0}-s-1/2}{b_{0}+a}}c_{4}^{\frac{2(a+s)+1}{2(b_{0}+a)}}\right)\epsilon_{\delta}(1+\epsilon_{\delta})^{-\frac{s+a+1/2}{b_{0}+a}} + \nu_{2}}{\left[D(\alpha_{1})\right]^{1/2}} + \\ & \frac{\frac{1}{2}\delta^{-\frac{a+s}{b_{0}+a}}\tilde{c}_{3}^{-\frac{a+s}{(b_{0}+a)}}c_{4}^{\frac{a+s}{b_{0}+a}}(1+\epsilon_{\delta})^{-\frac{a+s}{b_{0}+a}}(1+2\nu_{1})}{\left[D(\alpha_{1})\right]^{1/2}}\right) =: P\left(\xi > \frac{N_{1}}{\left[D(\alpha_{1})\right]^{\frac{1}{2}}} + \frac{N_{2}}{\left[D(\alpha_{1})\right]^{\frac{1}{2}}}\right) \\ \geq & P\left(\xi > \frac{N_{1}}{\left[D^{u}(\alpha_{1})\right]^{\frac{1}{2}}} + \frac{N_{2}}{\left[D_{l}(\alpha_{1})\right]^{\frac{1}{2}}}\right) = 1 - \Phi\left(\frac{N_{1}}{\left[D^{u}(\alpha_{1})\right]^{\frac{1}{2}}} + \frac{N_{2}}{\left[D_{l}(\alpha_{1})\right]^{\frac{1}{2}}}\right) \end{split}$$

since $N_1 < 0$ and $N_2 > 0$, where $\Phi(\cdot)$ denotes the cumulative distribution function of a $\mathcal{N}(0,1)$ distribution. Finally, $\frac{N_1}{[D^u(\alpha_1)]^{\frac{1}{2}}} + \frac{N_2}{[D_l(\alpha_1)]^{\frac{1}{2}}} \approx -\delta^{-\frac{1}{4(b_0+a)}}\epsilon_{\delta} + \delta^{\frac{1}{4(b_0+a)}}$, which converges to $-\infty$ as $\delta \to 0$ if we choose $\epsilon_{\delta} = \delta^{\frac{1-r}{4(b_0+a)}}$, for every r > 0. This proves that $P(S_{y^{\delta}}(\alpha_1) < 0) \to 1$.

CASE II: $b_0 > \frac{1}{2} + a + 2s$. So, $\tilde{\eta} = \eta$. By substituting the value of α_1 in the numerator of (36) we obtain

$$P(S_{y^{\delta}}(\alpha_{1}) < 0) \geq P\left(\xi > \frac{-\frac{1}{2}\delta^{-\frac{(\alpha+s+1/2+\eta)}{b_{0}+a+\eta}} \left[\tilde{c}_{3}(1+\epsilon_{\delta})\right]^{\frac{b_{0}-s-1/2}{b_{0}+a}} c_{4}^{\frac{a+s+1/2}{b_{0}+a}} + \nu_{2}}{\left[D(\alpha_{1})\right]^{1/2}} + \frac{\frac{1}{2}\delta^{-\frac{a+s+1/2}{b_{0}+a+\eta}} \left[\frac{\tilde{c}_{3}}{c_{4}}(1+\epsilon_{\delta})\right]^{-\frac{a+s+1/2}{b_{0}+a+1/2}} \left(\tilde{c}_{3}+\delta^{\frac{1}{2(b_{0}+a+\eta)}} \left[\frac{\tilde{c}_{3}}{c_{4}}(1+\epsilon_{\delta})\right]^{\frac{1}{2(b_{0}+a)}}\right)(1+2\nu_{1})}{\left[D(\alpha_{1})\right]^{\frac{1}{2}}}$$

$$=: P\left(\xi > \frac{\tilde{N}_{1}}{\left[D(\alpha_{1})\right]^{\frac{1}{2}}} + \frac{\tilde{N}_{2}}{\left[D(\alpha_{1})\right]^{\frac{1}{2}}}\right)$$

$$\geq P\left(\xi > \frac{\tilde{N}_{1}}{\left[D^{u}(\alpha_{1})\right]^{\frac{1}{2}}} + \frac{\tilde{N}_{2}}{\left[D_{l}(\alpha_{1})\right]^{\frac{1}{2}}}\right) = 1 - \Phi\left(\frac{\tilde{N}_{1}}{\left[D^{u}(\alpha_{1})\right]^{\frac{1}{2}}} + \frac{\tilde{N}_{2}}{\left[D_{l}(\alpha_{1})\right]^{\frac{1}{2}}}\right)$$

since $\tilde{N}_1 < 0$ and $\tilde{N}_2 > 0$, where $\Phi(\cdot)$ denotes the cumulative distribution function of a $\mathcal{N}(0,1)$ distribution. Remark that in this case we can rewrite $D_l(\alpha_1)$ as

$$D_l(\alpha_1) = \delta^{-1} \delta^{-\frac{a+2s-b_0+1/2}{2(b_0+a+\eta)}} \left(\left[\frac{\tilde{c}_3}{c_4} (1+\epsilon_\delta) \right]^{-\frac{a+2s-b_0+1/2}{b_0+a}} c_5 + \frac{1}{2} \delta^{1-\frac{b_0+a}{b_0+a+\eta}} \left[\frac{\tilde{c}_3}{c_4} (1+\epsilon_\delta) \right]^{-\frac{2(a+s)+1/2}{b_0+a}} c_6 \right)$$

and therefore, $\frac{\tilde{N}_1}{[D^u(\alpha_1)]^{\frac{1}{2}}} + \frac{\tilde{N}_2}{[D_l(\alpha_1)]^{\frac{1}{2}}} \approx -\delta^{-\frac{2s+a-b_0+1+\eta}{2(b_0+a+\eta)}} + \delta^{\frac{b_0-a-2s-3/2+2\eta}{4(b_0+a+\eta)}} = -\delta^{-\frac{2s+a-b_0+1+\eta}{2(b_0+a+\eta)}} (1-o(1))$ if $\tilde{\eta}$ is sufficiently big so that $\eta > \frac{1}{4}(b_0-2s-a-\frac{1}{2})$. This quantity converges to $-\infty$ if $\eta > b_0-2s-a-1$. Therefore, the condition which guarantees convergence is: $\eta > \max\{(b_0-2s-a-1), (b_0-2s-a-1/2)^{\frac{1}{4}}\}$.

Lemma 6. Let the assumptions of Theorem 6 be satisfied and $\alpha_2 = \left(\frac{\tilde{c}_3}{\tilde{c}_4}(1-\epsilon_\delta)\right)^{\frac{a+s}{b_0+a}}\delta^p$ where: $p = \frac{a+s}{b_0+a-\eta}$ for $(b_0+a) > \eta > \max\{b_0-s-1/2,0\}$, $0 < \epsilon_\delta < 1$, \tilde{c}_3 , \tilde{c}_4 be as defined in Lemma 3. Then,

$$P(S_{y^\delta}(\alpha_2)>0)\to 1 \quad \text{ as } \quad \delta\to 0.$$

Proof. This proof follows the line of the proof of Lemma 5, so some details are omitted. By using the upper bound of S_2 and the lower bound of $(S_1 - S_{4b})$ provided in Lemma 3 we obtain:

$$P(S_{y^{\delta}}(\alpha_{2}) > 0) = P(S_{3} + S_{4a} < -S_{2} + (S_{1} - S_{4b}) + \frac{\nu_{1}}{\alpha} + \nu_{2})$$

$$\geq P\left(\xi < \frac{-\frac{1}{2}\delta^{-1}\alpha_{2}^{\frac{b_{0}-s-1/2}{a+s}}\tilde{c}_{4} + \nu_{2} + \frac{\nu_{1}}{\alpha_{2}} + \alpha_{2}^{-\frac{2(a+s)+1}{2(a+s)}}\frac{1}{2}\tilde{c}_{3} - \frac{1}{2}\delta^{-1}I(b_{0} \geq s)}{[D(\alpha_{2})]^{1/2}}\right)$$
(39)

where ξ denotes a $\mathcal{N}(0,1)$ random variable. Moreover, $\forall \alpha$ let $D^u(\alpha)$ (resp. $D_l(\alpha)$) denote the upper (resp.

the lower) bound for $D(\alpha)$ defined in (37)-(38). By replacing the value of α_2 and after some algebras we get:

$$D^{u}(\alpha_{2}) = \delta^{-1}I(b_{0} \geq a + 2s) + \tilde{c}_{5}\delta^{-1 - \frac{a + 2s - b_{0} + 1/2}{b_{0} + a - \eta}} \left[\frac{\tilde{c}_{3}}{\tilde{c}_{4}} (1 - \epsilon_{\delta}) \right]^{-\frac{a + 2s - b_{0} + 1/2}{b_{0} + a}} + \frac{17\tilde{c}_{6}}{32} \delta^{-\frac{4(a + s) + 1}{2(b_{0} + a - \eta)}} \left[\frac{\tilde{c}_{3}}{\tilde{c}_{4}} (1 - \epsilon_{\delta}) \right]^{-\frac{4(a + s) + 1}{2(b_{0} + a)}}$$

$$D_{l}(\alpha_{2}) = \delta^{-\frac{2a + 2s + 1/2}{b_{0} + a - \eta}} \left[\frac{1}{2} c_{6} \left[\frac{\tilde{c}_{3}}{\tilde{c}_{4}} (1 - \epsilon_{\delta}) \right]^{-\frac{2(a + s) + 1/2}{b_{0} + a}} + \delta^{-1 + \frac{b_{0} + a}{b_{0} + a - \eta}} \left[\frac{\tilde{c}_{3}}{\tilde{c}_{4}} (1 - \epsilon_{\delta}) \right]^{-\frac{a + 2s - b_{0} + 1/2}{b_{0} + a}} c_{5} \right].$$

Remark that $D^u(\alpha_2) = \mathcal{O}(\delta^{-\frac{2(a+s)+1/2}{(b_0+a-\eta)}})$ if $b_0 \leq \frac{1}{2} + a + 2s + \eta$. This condition is satisfied by assumption because $\eta > b_0 - s - \frac{1}{2} > b_0 - a - 2s - \frac{1}{2}$. By substituting the value of α_2 in the numerator and after some algebra we obtain

$$P(S_{y^{\delta}}(\alpha_{2}) > 0) \geq P\left(\xi < \frac{\frac{1}{2}\delta^{-\frac{a+s+1/2}{b_{0}+a-\eta}} \left[\frac{\tilde{c}_{3}}{\tilde{c}_{4}}(1-\epsilon_{\delta})\right]^{-\frac{a+s+1/2}{b_{0}+a}} \tilde{c}_{3} + \nu_{1} \left[\frac{\tilde{c}_{3}}{\tilde{c}_{4}}(1-\epsilon_{\delta})\right]^{-\frac{a+s}{b_{0}+a}} \delta^{-\frac{a+s}{b_{0}+a-\eta}}}{[D(\alpha_{2})]^{1/2}} + \frac{-\frac{1}{2}\delta^{-1} \left[I(b_{0} \geq s) + \tilde{c}_{4}\delta^{\frac{b_{0}-s-1/2}{b_{0}+a-\eta}} \left[\frac{\tilde{c}_{3}}{\tilde{c}_{4}}(1-\epsilon_{\delta})\right]^{\frac{b_{0}-s-1/2}{b_{0}+a}}\right] + \nu_{2}}{[D(\alpha_{2})]^{1/2}}\right)}$$

$$=: P\left(\xi < \frac{N_{3}}{[D(\alpha_{2})]^{\frac{1}{2}}} + \frac{N_{4}}{[D(\alpha_{2})]^{\frac{1}{2}}}\right)$$

$$\geq P\left(\xi < \frac{N_{3}}{[D^{u}(\alpha_{2})]^{\frac{1}{2}}} + \frac{N_{4}}{[D_{l}(\alpha_{2})]^{\frac{1}{2}}}\right) = \Phi\left(\frac{N_{3}}{[D^{u}(\alpha_{2})]^{\frac{1}{2}}} + \frac{N_{4}}{[D_{l}(\alpha_{2})]^{\frac{1}{2}}}\right)$$

since $N_3 > 0$ and $N_4 < 0$. Term $N_3 / [D^u(\alpha_2)]^{\frac{1}{2}}$ converges to $+\infty$, as $\delta \to 0$, since

$$\frac{N_3}{[D^u(\alpha_2)]^{\frac{1}{2}}} \ = \ \frac{\frac{1}{2} \delta^{-\frac{a+s+1/2}{b_0+a-\eta} + \frac{2a+2s+1/2}{2(b_0+a-\eta)}} \left[\frac{\tilde{c}_3}{\tilde{c}_4}(1-\epsilon_\delta)\right]^{-\frac{a+s+1/2}{b_0+a}} \left(\tilde{c}_3 + 2\left[\frac{\tilde{c}_3}{\tilde{c}_4}(1-\epsilon_\delta)\right]^{\frac{a+s+1/2}{b_0+a}} \delta^{\frac{1}{2(b_0+a-\eta)}}\right)}{\left[o(1)I(b_0 \geq a+2s) + \left[\frac{\tilde{c}_3}{\tilde{c}_4}(1-\epsilon_\delta)\right]^{-\frac{2(a+s)+1/2}{b_0+a}} \left(\frac{17\tilde{c}_6}{32} + o(1)\right)\right]^{1/2}} \asymp \delta^{-\frac{1}{4(b_0+a-\eta)}}.$$

The asymptotic behavior of $N_4/[D_l(\alpha_2)]^{\frac{1}{2}}$ is different depending on the sign of $(b_0 - s - 1/2)$. So, we treat the two cases separately.

CASE I: $b_0 < s + 1/2$.

$$\frac{N_4}{[D_l(\alpha_2)]^{\frac{1}{2}}} = \frac{-\frac{1}{2} \delta^{\frac{2(a+s)+1/2}{2(b_0+a-\eta)}} \left(\delta^{-1+\frac{b_0-s-1/2}{b_0+a-\eta}} \left[\left[\frac{\tilde{c}_3}{\tilde{c}_4} (1-\epsilon_{\delta}) \right]^{\frac{b_0-s-1/2}{b_0+a}} \tilde{c}_4 + \delta^{-\frac{(b_0-s-1/2)}{b_0+a-\eta}} I(b_0 \ge s) \right] - \nu_2 \right)}{\left[o(1) + \left[\frac{\tilde{c}_3}{\tilde{c}_4} (1-\epsilon_{\delta}) \right]^{-\frac{2(a+s)+1/2}{b_0+a}} \frac{c_6}{2} \right]^{1/2}} \right] \\
\approx -\delta^{-1+\frac{2a+2b_0-1/2}{2(b_0+a-\eta)}}$$

which is bounded if $\eta \geq 1/4$. If $\eta < 1/4$ then $\frac{N_4}{[D_l(\alpha_2)]^{\frac{1}{2}}} \to -\infty$ but slower than $\frac{N_3}{[D^u(\alpha_2)]^{\frac{1}{2}}} \to \infty$.

CASE II: $b_0 \ge s + 1/2$.

$$\frac{N_4}{[D_l(\alpha_2)]^{\frac{1}{2}}} \ = \ \frac{-\frac{1}{2} \left(\delta^{-1 + \frac{2(a+s)+1/2}{2(b_0+a-\eta)}} \left[I(b_0 \geq s) + \delta^{\frac{b_0-s-1/2}{b_0+a-\eta}} \left[\frac{\tilde{c}_3}{\tilde{c}_4} (1-\epsilon_\delta) \right]^{\frac{b_0-s-1/2}{b_0+a}} \tilde{c}_4 - \nu_2 \delta \right] \right)}{\left[o(1) + \left[\frac{\tilde{c}_3}{\tilde{c}_4} (1-\epsilon_\delta) \right]^{-\frac{2(a+s)+1/2}{b_0+a}} \frac{c_6}{2} \right]^{1/2}} \lesssim -\delta^{-1 + \frac{2(a+s)+1/2}{2(b_0+a-\eta)}}$$

which converges to $-\infty$ if $b_0 > s + \eta + 1/4$. However, in this case, $\frac{N_4}{[D_t(\alpha_2)]^{\frac{1}{2}}} \to -\infty$ slower than $\frac{N_3}{[D^u(\alpha_3)]^{\frac{1}{2}}} \to \infty$ if $b_0 < s + 1/2 + \eta$.

Lemma 7. Let the assumptions of Lemmas 5 and 6 be satisfied and α_1, α_2 be as defined in the proof of Theorem 6. Define the event $G := \{\exists \hat{\alpha}; S_{y^{\delta}}(\hat{\alpha}) = 0 \cap \hat{\alpha} \in (\alpha_2, \alpha_1)\}$. Then $P(G) \to 1$ as $\delta \to 0$.

Proof. Because $S_{y^{\delta}}(\cdot)$ is continuous on $[\alpha_2, \alpha_1]$, in order to prove the existence of an $\hat{\alpha}$ such that $S_{y^{\delta}}(\hat{\alpha}) = 0$ it is sufficient to prove that $P\{S_{y^{\delta}}(\alpha_2) > 0 \text{ and } S_{y^{\delta}}(\alpha_1) < 0\} \to 1 \text{ when } \delta \to 0$. Thus, we can upper and lower bound P(G) by

$$1 \ge P(G) \ge P\{S_{y^{\delta}}(\alpha_2) > 0 \cap S_{y^{\delta}}(\alpha_1) < 0\}.$$

We now analyze $P\{S_{\eta^{\delta}}(\alpha_2) > 0 \cap S_{\eta^{\delta}}(\alpha_1) < 0\}$. By using the notation used in Lemmas 3, 5 and 6 define

$$\mathcal{H}(\alpha_i) = \frac{\mathcal{S}_1 - \mathcal{S}_2 - \mathcal{S}_{4b} + \nu_1/\alpha + \nu_2}{Var(\mathcal{S}_3) + Var(\mathcal{S}_4)} \bigg|_{\alpha = \alpha_i}, \quad \text{for } i = 1, 2.$$

Hence, by denoting with $\Phi(\cdot)$ the cumulative distribution function of a $\mathcal{N}(0,1)$,

$$\begin{split} P\{S_{y^{\delta}}(\alpha_{2}) > 0 \cap S_{y^{\delta}}(\alpha_{1}) < 0\} &= P\left(\xi_{1} + \xi_{2} < \mathcal{H}(\alpha_{2}) \text{ and } \xi_{1} + \xi_{2} > \mathcal{H}(\alpha_{1})\right) \\ &= P\left(\mathcal{H}(\alpha_{1}) < \xi_{1} + \xi_{2} < \mathcal{H}(\alpha_{2})\right) = \Phi(\mathcal{H}(\alpha_{2})) - \Phi(\mathcal{H}(\alpha_{1})) \\ &= P\{S_{y^{\delta}}(\alpha_{2}) > 0\} - (1 - P\{S_{y^{\delta}}(\alpha_{1}) < 0\}) \to 1 \text{ as } \delta \to 0 \end{split}$$

where the convergence to 1 follows from Lemmas 5 and 6.

References

Ai, C. and X., Chen (2003), Efficient estimation of models with conditional moment restrictions containing unknown functions, *Econometrica*, **71**, 1795-1843.

Agapiou, S., Larsson, S. and A.M., Stuart (2013), Posterior consistency of the Bayesian approach to linear ill-posed inverse problems, *Stochastic Processes and Applications*, **123**, 3828-3860.

Berger, J. O. and X., Wang (2011), Estimating shape constrained functions using Gaussian processes, JSM proceedings.

- Bissantz, N., Hohage, T., Munk, A. and F., Ruymgaart (2007), Convergence rates of general regularization methods for statistical inverse problems and applications, *SIAM J. Numer. Anal.*, **45**, 2610-2636.
- Blundell, R., Chen, X. and D., Kristensen (2007), Semi-nonparamtric IV estimation of shape-invariant engel curves, *Econometrica*, **75**, 1613-1669.
- Carrasco, M. and J.P., Florens (2000), Generalization of GMM to a continuum of moment conditions, *Econometric Theory*, **16**, 797-834.
- Carrasco, M., Florens, J.P., and E., Renault (2005), Estimation based on spectral decomposition and regularization, *Handbook of Econometrics*, J.J. Heckman and E. Leamer, eds., **6**, Elsevier, North Holland.
- Cavalier, L. and A.B., Tsybakov (2002), Sharp adaptation for inverse problems with random noise, *Proba. Theory and Related Fields*, **123**, 323–354.
- Chen, X. and D., Pouzo (2012), Estimation of nonparametric conditional moment models with possibly nonsmooth generalized residuals, *Econometrica*, **80**, 277-321.
- Chen, X. and M., Reiss, (2011), On rate optimality for ill-posed inverse problems in econometrics, *Econometric Theory*, **27**, Special Issue 03, 497-521.
- Daouia, A., Florens, J-P. and L. Simar (2009), Regularization of nonparametric frontier estimators, Journal of Econometrics, 168, 285-299.
- Darolles, S., Fan, Y., Florens, J.P., and E., Renault (2011), Nonparametric instrumental regression, *Econometrica*, **79**, 1541-1565.
- Diaconis, P.W., and D., Freedman (1986), On the consistency of Bayes estimates, *Annals of Statistics*, **14**, 1-26.
- Diaconis, P.W. and D., Freedman (1998), Consistency of Bayes estimates for nonparametric regression: normal theory, *Bernoulli*, 4, 411-444.
- Engl, H.W., Hanke, M. and A., Neubauer (2000), Regularization of inverse problems, *Kluwer Academic, Dordrecht*.
- Florens, J-P. (2003), Inverse problems and structural econometrics: the example of instrumental variables. Invited Lectures to the World Congress of the Econometric Society, Seattle 2000. In: M., Dewatripont, L.-P., Hansen, and S.J., Turnovsky, (Eds.), Advances in Economics end econometrics: theory and applications, Vol.II, 284-311. Cambridge University Press.

- Florens, J-P., Johannes, J. and S., Van Bellegem (2011), Identification and estimation by penalization in nonparametric instrumental regression, *Econometric Theory*, **27**, 472-496.
- Florens, J-P. and E. Sbaï (2010), Local identification in empirical games of incomplete information, *Econometric Theory*, **26**, 1638-1662.
- Florens, J.P. and A., Simoni (2010), Regularizing priors for linear inverse problems, IDEI Working paper, n. 621.
- Florens, J.P. and A. Simoni (2012a), Nonparametric estimation of an instrumental regression: a quasi-Bayesian approach based on regularized posterior, *Journal of Econometric*, **170**, 458-475.
- Florens, J.P. and A., Simoni (2012b), Regularized posteriors in linear ill-posed inverse problems, Scandinavian Journal of Statistics, 39, 214-235.
- Gautier, E. and Y. Kitamura (2012), Nonparametric estimation in random coefficients binary choice models, *Econometrica*, **81**, 581-607.
- Ghosh, J.K. and R.V., Ramamoorthi (2003), Bayesian nonparametrics, Springer Series in Statistics.
- Goldenshluger, A. and S.V., Pereverzev (2003), On adaptive inverse estimation of linear functionals in Hilbert scales, *Bernoulli*, **9**, 783-807.
- Hall, P. and J., Horowitz (2005), Nonparametric methods for inference in the presence of instrumental variables, *Annals of Statistics*, **33**, 2904-2929.
- Hall, P. and J.L., Horowitz (2007), Methodology and convergence rates for functional linear regression, *Annals of Statistics*, **35**, 70 91.
- Helin, T. (2009), On infinite-dimensional hierarchical probability models in statistical inverse problems, *Inverse Probl. Imaging*, **3**, pp. 567-597.
- Hoderlein, S., Nesheim, L. and A. Simoni (2013), Semiparametric estimation of random coefficients in structural economic models, CEMMAP working paper, CWP09/12.
- Hofinger, A. and H.K., Pikkarainen (2007), Convergence rate for the Bayesian approach to linear inverse problems, *Inverse Problems*, 23, 2469-2484.
- Hofinger, A. and H.K., Pikkarainen (2009), Convergence rate for linear inverse problems in the presence of an additive normal noise, *Stochastic Analysis and Applications*, **27**, 240-257.
- Horowitz, J. (2013), Adaptive nonparametric instrumental variables estimation: empirical choice of the regularisation parameter, CEMMAP working paper, CWP30/13.

- Ito, K. (1970), The topological support of Gauss measure on Hilbert space, Nagoya Math. J., 38, 181 183.
- Johannes, J. (2008), Nonparametric estimation in functional linear model, in: Functional and Operatorial Statistics, S. Dabo-Niang and F. Ferraty, editors, 215-221, Physica-Verlag.
- Johannes, J., Van Bellegem, S. and A., Vanhems (2011), Convergence rates for ill-posed inverse problems with an unknown operator, *Econometric Theory*, **27**, 522-545.
- Kaipio, J., and E., Somersalo (2004), Statistical and computational inverse problems, *Applied Mathematical Series*, **160**, Springer, Berlin.
- Kato, T. (1995), Perturbation theory for linear operators, Springer.
- Knapik, B.T., Van der Vaart, A.W. and J.H., Van Zanten (2011), Bayesian inverse problems, *Annals of Statistics*, **39**, 2626-2657.
- Krein, S.G. and J.I., Petunin (1966), Scales of Banach spaces, Russian Math. Surveys, 21, 85 160.
- Kress, R. (1999), Linear integral equation, Springer.
- Kuo, H.H. (1975), Gaussian measures in Banach spaces, Springer.
- Lassas, M., Saksman, E. and S., Siltanen (2009), Discretization invariant Bayesian inversion and Besov space priors, *Inverse Probl. Imaging*, **3**, 87-122.
- Lehtinen, M.S., Päivärinta, L. and E., Somersalo (1989), Linear inverse problems for generalised random variables, *Inverse Problems*, 5, 599-612.
- Liao, Y. and W., Jiang (2011), Posterior consistency of nonparametric conditional moment restricted models, *Annals of Statistics*, **39**, 3003-3031
- Loubes, J-M. and C. Marteau (2013), Oracle Inequality for Instrumental Variable Regression, working paper, arXiv:0901.4321v1.
- Luschgy, H. (1995), Linear estimators and radonifying operators, *Theory Probab. Appl.*, **40**, 167-175.
- Mandelbaum, A. (1984), Linear estimators and measurable linear transformations on a Hilbert space, Z. Wahrcheinlichkeitstheorie, 3, 385-98.
- Nair, M.T., Pereverzev, S.V. and U., Tautenhahn (2005), Regularization in Hilbert scales under general smoothing conditions, *Inverse Problems*, **21**, 1851-1969.

- Natterer, F. (1984), Error bounds for Tikhonov regularization in Hilbert scale, *Applicable Analysis*, 18, 29-37.
- Newey, W.K. and J.L., Powell (2003), Instrumental variable estimation of nonparametric models. Econometrica, Vol.71, 5, 1565-1578.
- Neubauer, A. and H.K., Pikkarainen (2008), Convergence results for the Bayesian inversion theory, J. Inverse Ill-posed Probl., Vol. 16, 601-613.
- Nussbaum, M. (1996), Asymptotic equivalence of density estimation and gaussian white noise, *Annals of Statistics*, **24**, No.6, 2399–2430.
- Ruymgaart, F.H. (1998), A note on weak convergence of density in hilbert spaces, *Statistics*, **30**, 331-343.
- Van der Vaart, A.W. and J.A., Wellner (1996), Weak convergence and empirical processes, Springer, New York.
- Van Rooij, A.C.M. and F.H., Ruymgaart (1999), On inverse estimation, in: Ghosh, S. (Ed) Asymptotic, Nonparametric and Time Series, 579-613 (Dekker).
- Zellner, A. (1986), On Assessing Prior distributions and Bayesian regression analysis with g-prior distribution, in: Goel, P.K. and Zellner, A. (Eds) *Bayesian Inference and Decision Techniques:* essays in honour of Bruno de Finetti, pp. 233-243 (Amsterdam, North Holland).