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Consistency of Functional Learning Methods Based on Derivatives

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Abstract

In some real world applications, such as spectrometry, functional models achieve better predictive performances if they work on the derivatives of order m of their inputs rather than on the original functions. As a consequence, the use of derivatives is a common practice in functional data analysis, despite a lack of theoretical guarantees on the asymptotically achievable performances of a derivative based model. In this paper, we show that a smoothing spline approach can be used to preprocess multivariate observations obtained by sampling functions on a discrete and finite sampling grid in a way that leads to a consistent scheme on the original infinite dimensional functional problem. This work extends Mas and Pumo (2009) to nonparametric approaches and incomplete knowledge. To be more precise, the paper tackles two difficulties in a nonparametric framework: the information loss due to the use of the derivatives instead of the original functions and the information loss due to the fact that the functions are observed through a discrete sampling and are thus also unperfectly known: the use of a smoothing spline based approach solves these two problems. Finally, the proposed approach is tested on two real world datasets and the approach is experimentaly proven to be a good solution in the case of noisy functional predictors.

Keywords: Functional Data Analysis, Consistency, Statistical learning, Derivatives, SVM, Smoothing splines, RKHS, Kernel

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1 1. Introduction

As the measurement techniques are developping, more and more data are high dimensional vectors generated by measuring a continuous process on a discrete sampling grid. Many examples of this type of data can be found in real world applications, in various fields such as spectrometry, voice recognition, time series analysis, etc.

Data of this type should not be handled in the same way as standard 7 multivariate observations but rather analysed as *functional* data: each ob-8 servation is a function coming from an input space with infinite dimension, g sampled on a high resolution sampling grid. This leads to a large number 10 of variables, generally more than the number of observations. Moreover, 11 functional data are frequently smooth and generate highly correlated vari-12 ables as a consequence. Applied to the obtained high dimensional vectors, 13 classical statistical methods (e.g., linear regression, factor analysis) often 14 lead to ill-posed problems, especially when a covariance matrix has to be 15 inverted (this is the case, e.g., in linear regression, in discriminant analysis 16 and also in sliced inverse regression). Indeed, the number of observed values 17 for each function is generally larger than the number of functions itself and 18 these values are often strongly correlated. As a consequence, when these 19 data are considered as multidimensional vectors, the covariance matrix is ill-20 conditioned and leads to unstable and unaccurate solutions in models where 21 its inverse is required. Thus, these methods cannot be directly used. During 22 past years, several methods have been adapted to that particular context 23 and grouped under the generic name of Functional Data Analysis (FDA) 24 methods. Seminal works focused on linear methods such as factorial analysis 25 (Deville (1974); Dauxois and Pousse (1976); Besse and Ramsay (1986); James 26 et al. (2000), among others) and linear models Ramsay and Dalzell (1991); 27 Cardot et al. (1999); James and Hastie (2001); a comprehensive presenta-28 tion of linear FDA methods is given in Ramsay and Silverman (1997, 2002). 29 More recently, nonlinear functional models have been extensively developed 30 and include generalized linear models James (2002); James and Silverman 31 (2005), kernel nonparametric regression Ferraty and Vieu (2006), Functional 32 Inverse Regression Ferré and Yao (2003), neural networks Rossi and Conan-33 Guez (2005); Rossi et al. (2005), k-nearest neighbors Biau et al. (2005); Laloë 34 (2008), Support Vector Machines (SVM), Rossi and Villa (2006), among a 35

³⁶ very large variety of methods.

In previous works, numerous authors have shown that the derivatives 37 of the functions lead sometimes to better predictive performances than the 38 functions themselves in inference tasks, as they provide information about 39 the shape or the regularity of the function. In particular applications such 40 as spectrometry Ferraty and Vieu (2006); Rossi et al. (2005); Rossi and Villa 41 (2006), micro-array data Dejean et al. (2007) and handwriting recognition 42 Williams et al. (2006); Bahlmann and Burkhardt (2004), these characteristics 43 lead to accurate predictive models. But, on a theoretical point of the view, 44 limited results about the effect of the use of the derivatives instead of the 45 original functions are available: Mas and Pumo (2009) studies this problem 46 for a linear model built on the first derivatives of the functions. In the present 47 paper, we also focus on the theoretical relevance of this common practice and 48 extend Mas and Pumo (2009) to nonparametric approaches and incomplete 49 knowledge. 50

More precisely, we address the problem of the estimation of the condi-51 tional expectation $\mathbb{E}(Y|X)$ of a random variable Y given a functional random 52 variable X. Y is assumed to be either real valued (leading to a regression 53 problem) or to take values in $\{-1, 1\}$ (leading to a binary classification prob-54 lem). We target two theoretical difficulties. The first difficulty is the po-55 tential information loss induced by using a derivative instead of the original 56 function: when one replaces X by its order m derivative $X^{(m)}$, consistent 57 estimators (such as kernel models Ferraty and Vieu (2006)) guarantee an 58 asymptotic estimation of $\mathbb{E}(Y|X^{(m)})$ but cannot be used directly to address 59 the original problem, namely estimating $\mathbb{E}(Y|X)$. This is a simple conse-60 quence of the fact that $X \mapsto X^{(m)}$ is not a one to one mapping. The second 61 difficulty is induced by sampling: in practice, functions are never observed 62 exactly but rather, as explained above, sampled on a discrete sampling grid. 63 As a consequence, one relies on approximate derivatives, $\widehat{X}_{\tau}^{(m)}$ (where τ de-64 notes the sampling grid). This approach induces even more information loss 65 with respect to the underlying functional variable X: in general, a consistent 66 estimator of $\mathbb{E}\left(Y|\widehat{X}_{\tau}^{(m)}\right)$ will not provide a consistent estimation of $\mathbb{E}\left(Y|X\right)$ 67 and the optimal predictive performances for Y given $\widehat{X}_{\tau}^{(m)}$ will be lower than 68 the optimal predictive performances for Y given X. 69

We show in this paper that the use of a smoothing spline based approach
solves both problems. Smoothing splines are used to estimate the functions
from their sampled version in a convergent way. In addition, properties of

splines are used to obtain estimates of the derivatives of the functions with no 73 induced information loss. Both aspects are implemented as a preprocessing 74 step applied to the multivariate observations generated via the sampling grid. 75 The preprocessed observations can then be fed into any finite dimensional 76 consistent regression estimator or classifier, leading to a consistent estima-77 tor for the original infinite dimensional problem (in real world applications, 78 we instantiate the general scheme in the particular case of kernel machines 79 Shawe-Taylor and Cristianini (2004)). 80

The remainder of the paper is organized as follows: Section 2 introduces the model, the main smoothness assumption and the notations. Section 3 recalls important properties of spline smoothing. Section 4 presents approximation results used to build a general consistent classifier or a general consistent regression estimator in Section 5. Finally, Section 6 illustrates the behavior of the proposed method for two real world spectrometric problems. The proofs are given at the end of the article.

88 2. Setup and notations

⁸⁹ 2.1. Consistent classifiers and regression functions

We consider a pair of random variables (X, Y) where X takes values in a functional space \mathcal{X} and Y is either a real valued random variable (regression case) or a random variable taking values in $\{-1, 1\}$ (binary classification case). From this, we are given a learning set $S_n = \{(X_i, Y_i)\}_{i=1}^n$ of nindependent copies of (X, Y). Moreover, the functions X_i are not entirely known but sampled according to a non random sampling grid of finite length, $\tau_d = (t_l)_{l=1}^{|\tau_d|}$: we only observe $\mathbf{X}_i^{\tau_d} = (X_i(t_1), \ldots X_i(t_{|\tau_d|}))^T$, a vector of $\mathbb{R}^{|\tau_d|}$ and denote S_{n,τ_d} the corresponding learning set. Our goal is to construct:

1. in the binary classification case: a classifier, ϕ_{n,τ_d} , whose misclassification probability

$$L(\phi_{n,\tau_d}) = \mathbb{P}\left(\phi_{n,\tau_d}(\mathbf{X}^{\tau_d}) \neq Y\right)$$

asymptotically reaches the Bayes risk

$$L^* = \inf_{\phi: \mathcal{X} \to \{-1, 1\}} \mathbb{P}\left(\phi(X) \neq Y\right)$$

101 i.e., $\lim_{|\tau_d| \to +\infty} \lim_{n \to +\infty} \mathbb{E}\left(L(\phi_{n,\tau_d})\right) = L^*$;

2. in the regression case: a regression function, ϕ_{n,τ_d} , whose L^2 error

$$L(\phi_{n,\tau_d}) = \mathbb{E}\left([\phi_{n,\tau_d}(\mathbf{X}^{\tau_d}) - Y]^2 \right)$$

asymptotically reaches the minimal L^2 error

$$L^* = \inf_{\phi: \mathcal{X} \to \mathbb{R}} \mathbb{E} \left([\phi(\mathbf{X}^{\tau_d}) - Y]^2 \right)$$

104 i.e., $\lim_{|\tau_d| \to +\infty} \lim_{n \to +\infty} L(\phi_{n,\tau_d}) = L^*$.

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This definition implicitly requires $\mathbb{E}(Y^2) < \infty$ and as a consequence, corresponds to a L^2 convergence of ϕ_{n,τ_d} to the conditional expectation $\phi^* = \mathbb{E}(Y|X)$, i.e., to $\lim_{|\tau_d| \to +\infty} \lim_{n \to +\infty} \mathbb{E}\left([\phi_{n,\tau_d}(\mathbf{X}^{\tau_d}) - \phi^*(X)]^2\right) = 0.$

Such ϕ_{n,τ_d} are said to be *(weakly) consistent* Devroye et al. (1996); Györfi et al. (2002). We have deliberately used the same notations for the (optimal) predictive performances in both the binary classification and the regression case. We will call L^* the Bayes risk even in the case of regression. Most of the theoretical background of this paper is common to both the regression case and the classification case: the distinction between both cases will be made only when necessary.

As pointed out in the introduction, the main difficulty is to show that the performances of a model built on the $\mathbf{X}_i^{\tau_d}$ asymptotically reach the best performance achievable on the original functions X_i . In addition, we will build the model on derivatives estimated from the $\mathbf{X}_i^{\tau_d}$.

120 2.2. Smoothness assumption

Our goal is to leverage the functional nature of the data by allow-121 ing differentiation operators to be applied to functions prior their submis-122 sion to a more common classifier or regression function. Therefore we as-123 sume that the functional space \mathcal{X} contains only differentiable functions. 124 More precisely, \mathcal{X} is the Sobolev space $\mathcal{H}^m = \left\{h \in L^2([0,1]) \mid \forall j = 1\right\}$ 125 $1, \ldots, m, D^{j}h$ exists in the weak sense, and $D^{m}h \in L^{2}([0, 1])$, where $D^{j}h$ 126 is the *j*-th derivative of h (also denoted by $h^{(j)}$) and for an integer $m \ge 1$. 127 Of course, by a straightforward generalization, any bounded interval can be 128 considered instead of [0, 1]. 129

To estimate the underlying functions X_i and their derivatives from sampled data, we rely on smoothing splines. More precisely, let us consider a deterministic function $x \in \mathcal{H}^m$ sampled on the aforementioned grid. A smoothing spline estimate of x is the solution, \hat{x}_{λ,τ_d} , of

$$\arg\min_{h\in\mathcal{H}^m}\frac{1}{|\tau_d|}\sum_{l=1}^{|\tau_d|}(x(t_l)-h(t_l))^2 + \lambda\int_{[0,1]}(h^{(m)}(t))^2dt,\tag{1}$$

where λ is a regularization parameter that balances interpolation error and 134 smoothness (measured by the L^2 norm of the *m*-th derivative of the esti-135 mate). The goal is to show that a classifier or a regression function built 136 on $\widehat{X}_{\lambda,\tau_d}^{(m)}$ is consistent for the original problem (i.e., the problem defined by 137 the pair (X, Y): this means that using $\widehat{X}_{\lambda, \tau_d}^{(m)}$ instead of X has no dramatic 138 consequences on the accuracy of the classifier or of the regression function. 139 In other words, asymptotically, no information loss occurs when one replaces 140 X by $\widehat{X}_{\lambda,\tau_d}^{(m)}$ 141

The proof is based on the following steps:

1. First, we show that building a classifier or a regression function on 143 $\widehat{X}_{\lambda,\tau_d}^{(m)}$ is approximately equivalent to building a classifier or a regression 144 function on $\mathbf{X}^{\tau_d} = (X(t_l))_{l=1}^{|\tau_d|}$ using a specific metric. This is done by 145 leveraging the Reproducing Kernel Hilbert Space (RKHS) structure of 146 \mathcal{H}^m . This part serves one main purpose: it provides a solution to work 147 with estimation of the derivatives of the original function in a way 148 that preserves all the information available in \mathbf{X}^{τ_d} . In other words, the 149 best predictive performances for Y theoretically available by building a 150 multivariate model on \mathbf{X}^{τ_d} are equal to the best predictive performances 151 obtained by building a functional model on $\widehat{X}_{\lambda,\tau_{\tau}}^{(m)}$. 152

¹⁵³ 2. Then, we link $\mathbb{E}\left(Y|\widehat{X}_{\lambda,\tau_d}\right)$ with $\mathbb{E}\left(Y|X\right)$ by approximation results ¹⁵⁴ available for smoothing splines. This part of the proof handles the ¹⁵⁵ effects of sampling.

¹⁵⁶ 3. Finally, we glue both results via standard $\mathbb{R}^{|\tau_d|}$ consistency results.

157 3. Smoothing splines and differentiation operators

158 3.1. RKHS and smoothing splines

As we want to work on derivatives of functions from \mathcal{H}^m , a natural inner product for two functions of \mathcal{H}^m would be $(u, v) \to \int_0^1 u^{(m)}(t)v^{(m)}(t)dt$. However, we prefer to use an inner product of \mathcal{H}^m $(\int_0^1 u^{(m)}(t)v^{(m)}(t)dt$ only induces a semi-norm on \mathcal{H}^m) because, as will be shown later, such an inner product is related to an inner product between the sampled functions considered as vectors of $\mathbb{R}^{|\tau_d|}$.

This can be done by decomposing \mathcal{H}^m into $\mathcal{H}^m = \mathcal{H}^m_0 \oplus \mathcal{H}^m_1$ Kimeldorf 165 and Wahba (1971), where $\mathcal{H}_0^m = \text{Ker} D^m = \mathbb{P}^{m-1}$ (the space of polynomial 166 functions of degree less or equal to m-1) and \mathcal{H}_1^m is an infinite dimensional 16 subspace of \mathcal{H}^m defined via *m* boundary conditions. The boundary condi-168 tions are given by a full rank linear operator from \mathcal{H}^m to \mathbb{R}^m , denoted B, such 169 that $\operatorname{Ker} B \cap \mathbb{P}^{m-1} = \{0\}$. Classical examples of boundary conditions include 170 the case of "natural splines" (for m = 2, h(0) = h(1) = 0) and constraints 171 that target only the first values of h and its derivatives at a fixed position, for 172 instance the conditions: $h(0) = \ldots = h^{(m-1)}(0) = 0$. Other boundary con-173 ditions can be used Berlinet and Thomas-Agnan (2004); Besse and Ramsay 174 (1986); Craven and Wahba (1978), depending on the application. 175

Once the boundary conditions are fixed, an inner product on both \mathcal{H}_0^m and \mathcal{H}_1^m can be defined:

$$\langle u, v \rangle_1 = \langle D^m u, D^m v \rangle_{L^2} = \int_0^1 u^{(m)}(t) v^{(m)}(t) dt$$

is an inner product on \mathcal{H}_1^m (as $h \in \mathcal{H}_1^m$ and $D^m h \equiv 0$ give $h \equiv 0$). Moreover, if we denote $B = (B^j)_{j=1}^m$, then $\langle u, v \rangle_0 = \sum_{j=1}^m B^j u B^j v$ is an inner product on \mathcal{H}_0^m . We obtain this way an inner product on \mathcal{H}^m given by

$$\langle u, v \rangle_{\mathcal{H}^m} = \int_0^1 u^{(m)}(t) v^{(m)}(t) dt + \sum_{j=1}^m B^j u B^j v$$
$$= \langle \mathcal{P}_1^m(u), \mathcal{P}_1^m(v) \rangle_1 + \langle \mathcal{P}_0^m(u), \mathcal{P}_0^m(v) \rangle_0$$

¹⁸¹ where \mathcal{P}_i^m is the projector on \mathcal{H}_i^m .

Equipped with $\langle ., . \rangle_{\mathcal{H}^m}$, \mathcal{H}^m is a Reproducing Kernel Hilbert Space (RKHS, see e.g. Berlinet and Thomas-Agnan (2004); Heckman and Ramsay (2000); Wahba (1990)). More precisely, it exists a kernel $k : [0, 1]^2 \to \mathbb{R}$ such that, for all $u \in \mathcal{H}^m$ and all $t \in [0, 1]$, $\langle u, k(t, .) \rangle_{\mathcal{H}^m} = u(t)$. The same occurs for \mathcal{H}_0^m and \mathcal{H}_1^m which respectively have reproducing kernels denoted by k_0 and k_1 . We have $k = k_0 + k_1$.

In the most common cases, k_0 and k_1 have already been explicitly calculated (see e.g., Berlinet and Thomas-Agnan (2004), especially chapter 6, sections 1.1 and 1.6.2). For example, for $m \ge 1$ and the boundary conditions $h(0) = h'(0) = \ldots = h^{(m-1)}(0) = 0$, we have:

$$k_0(s,t) = \sum_{k=0}^{m-1} \frac{t^k s^k}{(k!)^2}.$$

192 and

$$k_1(s,t) = \int_0^1 \frac{(t-w)_+^{m-1}(s-w)_+^{m-1}}{(m-1)!^2} \, dw.$$

¹⁹³ 3.2. Computing the splines

We need now to compute to \hat{x}_{λ,τ_d} starting with $\mathbf{x}^{\tau_d} = (x(t))_{t\in\tau_d}^T$. This can be done via a theorem from Kimeldorf and Wahba (1971). We need the following compatibility assumptions between the sampling grid τ_d and the boundary conditions operator B:

198 Assumption 1. The sampling grid $\tau_d = (t_l)_{l=1}^{|\tau_d|}$ is such that

199 1. sampling points are distinct in [0,1] and $|\tau_d| \ge m-1$

200 2. the *m* boundary conditions B^j are linearly independent from the $|\tau_d|$ 201 linear forms $h \mapsto h(t_l)$, for $l = 1, ..., |\tau_d|$ (defined on \mathcal{H}^m)

Then
$$\widehat{x}_{\lambda,\tau_d}$$
 and $\mathbf{x}^{\tau_d} = (x(t))_{t \in \tau_d}^T$ are linked by the following result:

Theorem 1 (Kimeldorf and Wahba (1971)). Under Assumption (A1), the unique solution \hat{x}_{λ,τ_d} to equation (1) is given by:

$$\widehat{x}_{\lambda,\tau_d} = \mathcal{S}_{\lambda,\tau_d} \mathbf{x}^{\tau_d},\tag{2}$$

where S_{λ,τ_d} is a full rank linear operator from $\mathbb{R}^{|\tau_d|}$ to \mathcal{H}^m defined by:

$$S_{\lambda,\tau_d} = \omega^T M_0 + \eta^T M_1 \tag{3}$$

206 with

•
$$M_0 = \left(U(K_1 + \lambda I_d)^{-1} U^T \right)^{-1} U(K_1 + \lambda I_d)^{-1}$$

•
$$M_1 = (K_1 + \lambda I_d)^{-1} (I_d - U^T M_0);$$

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• {
$$\omega_1, \ldots, \omega_m$$
} is a basis of \mathbb{P}^{m-1} , $\omega = (\omega_1, \ldots, \omega_m)^T$ and $U = (\omega_i(t))_{i=1,\ldots,m} \underset{t \in \tau_d}{\ldots}$;

•
$$\eta = (k_1(t, .))_{t \in \tau_d}^T$$
 and $K_1 = (k_1(t, t'))_{t, t' \in \tau_d}$.

212 3.3. No information loss

The first important consequence of Theorem 1 is that building a model on $\widehat{X}_{\lambda,\tau_d}$ or on \mathbf{X}^{τ_d} leads to the same optimal predictive performances (to the same Bayes risk). This is formalized by the following corollary:

²¹⁶ Corollary 1. Under Assumption (A1), we have

• in the binary classification case:

$$\inf_{\substack{\phi:\mathcal{H}^m \to \{-1,1\}}} \mathbb{P}\left(\phi(\widehat{X}_{\lambda,\tau_d}) \neq Y\right) = \\\inf_{\substack{\phi:\mathbb{R}^{|\tau_d|} \to \{-1,1\}}} \mathbb{P}\left(\phi(\mathbf{X}^{\tau_d}) \neq Y\right)$$
(4)

• in the regression case:

$$\inf_{\phi:\mathcal{H}^m \to \mathbb{R}} \mathbb{E}\left(\left[\phi\left(\widehat{X}_{\lambda,\tau_d}\right) - Y\right]^2\right) = \inf_{\substack{\phi:\mathbb{R}^{|\tau_d|} \to \mathbb{R}}} \mathbb{E}\left(\left[\phi\left(\mathbf{X}^{\tau_d}\right) - Y\right]^2\right)$$
(5)

219 3.4. Differentiation operator

The second important consequence of Theorem 1 is that the inner product $\langle ., . \rangle_{\mathcal{H}^m}$ is equivalent to a specific inner product on $\mathbb{R}^{|\tau_d|}$ given in the following corollary:

²²³ Corollary 2. Under Assumption (A1) and for any $\mathbf{u}^{\tau_d} = (u(t))_{t \in \tau_d}^T$ and ²²⁴ $\mathbf{v}^{\tau_d} = (v(t))_{t \in \tau_d}^T$ in $\mathbb{R}^{|\tau_d|}$,

$$\langle \widehat{u}_{\lambda,\tau_d}, \widehat{v}_{\lambda,\tau_d} \rangle_{\mathcal{H}^m} = (\mathbf{u}^{\tau_d})^T \mathbf{M}_{\lambda,\tau_d} \mathbf{v}^{\tau_d}$$
 (6)

where $\mathbf{M}_{\lambda,\tau_d} = M_0^T W M_0 + M_1^T K_1 M_1$ with $W = (\langle w_i, w_j \rangle_0)_{i,j=1,...,m}$. The matrix $\mathbf{M}_{\lambda,\tau_d}$ is symmetric and positive definite and defines an inner product on $\mathbb{R}^{|\tau_d|}$. The corollary is a direct consequence of equations (2) and (3).

In practice, the corollary means that the euclidean space $\left(\mathbb{R}^{|\tau_d|}, \langle ., . \rangle_{\mathbf{M}_{\lambda,\tau_d}}\right)$ is isomorphic to $(\mathcal{I}_{\lambda,\tau_d}, \langle ., . \rangle_{\mathcal{H}^m})$, where $\mathcal{I}_{\lambda,\tau_d}$ is the image of $\mathbb{R}^{|\tau_d|}$ by $\mathcal{S}_{\lambda,\tau_d}$. As a consequence, one can use the Hilbert structure of \mathcal{H}^m directly in $\mathbb{R}^{|\tau_d|}$ via $\mathbf{M}_{\lambda,\tau_d}$: as the inner product of \mathcal{H}^m is defined on the order m derivatives of the functions, this corresponds to using those derivatives instead of the original functions.

²³⁵ More precisely, let $\mathbf{Q}_{\lambda,\tau_d}$ be the transpose of the Cholesky triangle of $\mathbf{M}_{\lambda,\tau_d}$ ²³⁶ (given by the Cholesky decomposition $\mathbf{Q}_{\lambda,\tau_d}^T \mathbf{Q}_{\lambda,\tau_d} = \mathbf{M}_{\lambda,\tau_d}$). Corollary 2 ²³⁷ shows that $\mathbf{Q}_{\lambda,\tau_d}$ acts as an approximate differentiation operation on sampled ²³⁸ functions.

Let us indeed consider an estimation method for multivariate inputs based only on inner products or norms (that are directly derived from the inner products), such as, e.g., Kernel Ridge Regression Saunders et al. (1998); Shawe-Taylor and Cristianini (2004). In this latter case, if a Gaussian kernel is used, the regression function has the following form:

$$u \mapsto \sum_{i=1}^{n} T_i \alpha_i e^{-\gamma \|U_i - u\|_{\mathbb{R}^p}^2} \tag{7}$$

where $(U_i, T_i)_{1 \le i \le n}$ are learning examples in $\mathbb{R}^p \times \{-1, 1\}$ and the α_i are non negative real values obtained by solving a quadratic programming problem and γ is a parameter of the method. Then, if we use Kernel Ridge Regression on the training set $\{(\mathbf{Q}_{\lambda,\tau_d}\mathbf{X}_i^{\tau_d}, Y_i)\}_{i=1}^n$ (rather than the original training set $\{(\mathbf{X}_i^{\tau_d}, Y_i)\}_{i=1}^n$), it will work on the norm in L^2 of the derivatives of order m of the spline estimates of the X_i (up to the boundary conditions). More precisely, the regression function will have the following form:

$$\mathbf{x}^{\tau_d} \mapsto \sum_{i=1}^n Y_i \alpha_i e^{-\gamma \left\| \mathbf{Q}_{\lambda, \tau_d} \mathbf{x}_i^{\tau_d} - \mathbf{Q}_{\lambda, \tau_d} \mathbf{x}^{\tau_d} \right\|_{\mathbb{R}^{|\tau_d|}}^2}$$
$$\mapsto \sum_{i=1}^n Y_i \alpha_i e^{-\gamma \left\| D^m \widehat{X}_{i\lambda, \tau_d} - D^m \widehat{x}_{\lambda, \tau_d} \right\|_{L^2}^2}$$
$$\times e^{-\gamma \sum_{j=1}^m \left(B^j \widehat{X}_{i\lambda, \tau_d} - B^j \widehat{x}_{\lambda, \tau_d} \right)^2}$$

In other words, up to the boundary conditions, an estimation method based solely on inner products, or on norms derived from these inner products, can be given modified inputs that will make it work on an estimation of thederivatives of the observed functions.

Remark 1. As shown in Corollary 1 in the previous section, building a model on \mathbf{X}^{τ_d} or on $\widehat{X}_{\lambda,\tau_d}$ leads to the same optimal predictive performances. In addition, it is obvious that given any one-to-one mapping f from $\mathbb{R}^{|\tau_d|}$ to itself, building a model on $f(\mathbf{X}^{\tau_d})$ gives also the same optimal performances than building a model on \mathbf{X}^{τ_d} . Then as $\mathbf{Q}_{\lambda,\tau_d}$ is invertible, the optimal predictive performances achievable with $\mathbf{Q}_{\lambda,\tau_d}\mathbf{X}^{\tau_d}$ are equal to the optimal performances achievable with $\widehat{\mathbf{X}}_{\lambda,\tau_d}$.

In practice however, the actual preprocessing of the data can have a strong influence on the obtained performances, as will be illustrated in Section 6. The goal of the theoretical analysis of the present section is to guarantee that no systematic loss can be observed as a consequence of the proposed functional preprocessing scheme.

²⁶⁷ 4. Approximation results

The previous section showed that working on \mathbf{X}^{τ_d} , $\mathbf{Q}_{\lambda,\tau_d}\mathbf{X}^{\tau_d}$ or $\widehat{X}_{\lambda,\tau_d}$ makes no difference in terms of optimal predictive performances. The present section addresses the effects of sampling: asymptotically, the optimal predictive performances obtained on $\widehat{X}_{\lambda,\tau_d}$ converge to the optimal performances achievable on the original and unobserved functional variable X.

273 4.1. Spline approximation

From the sampled random function $\mathbf{X}^{\tau_d} = (X(t_1), \dots, X(t_{|\tau_d|}))$, we can build an estimate, \hat{X}_{λ,τ_d} , of X. To ensure consistency, we must guarantee that \hat{X}_{λ,τ_d} converges to X. In the case of a deterministic function x, this problem has been studied in numerous papers, such as Craven and Wahba (1978); Ragozin (1983); Cox (1984); Utreras (1988); Wahba (1990) (among others). Here we recall one of the results which is particularly well adapted to our context.

Obviously, the sampling grid must behave correctly, whereas the information contained in \mathbf{X}^{τ_d} will not be sufficient to recover X. We need also the regularization parameter λ to depend on τ_d . Following Ragozin (1983), a sampling grid τ_d is characterized by two quantities:

$$\Delta_{\tau_d} = \max\{t_1, t_2 - t_1, \dots, 1 - t_{|\tau_d|}\}$$

$$\underline{\Delta}_{\tau_d} = \min_{1 \le i \le |\tau_d|} \{t_{i+1} - t_i\}.$$
(8)

One way to control the distance between X and $\widehat{X}_{\lambda,\tau_d}$ is to bound the ratio $\overline{\Delta}_{\tau_d}/\underline{\Delta}_{\tau_d}$ so as to ensure quasi-uniformity of the sampling grid. More precisely, we will use the following assumption:

Assumption 2. There is R such that $\overline{\Delta}_{\tau_d} / \underline{\Delta}_{\tau_d} \leq R$ for all d.

²⁸⁹ Then we have:

Theorem 2 (Ragozin (1983)). Under Assumptions (A1) and (A2), there are two constants $A_{R,m}$ and $B_{R,m}$ depending only on R and m, such that for any $x \in \mathcal{H}^m$ and any positive λ :

$$\|\widehat{x}_{\lambda,\tau_d} - x\|_{L^2}^2 \le \left(A_{R,m}\lambda + B_{R,m}\frac{1}{|\tau_d|^{2m}}\right) \|D^m x\|_{L^2}^2.$$

This result is a rephrasing of Corollary 4.16 from Ragozin (1983) which is itself a direct consequence of Theorem 4.10 from the same paper.

²⁹⁵ Convergence of \hat{x}_{λ,τ_d} to x is then obtained by the following simple as-²⁹⁶ sumptions:

Assumption 3. The series of sampling points τ_d and the series of regularization parameters, λ , depending on τ_d and denoted by $(\lambda_d)_{d\geq 1}$, are such that $\lim_{d\to+\infty} |\tau_d| = +\infty$ and $\lim_{d\to+\infty} \lambda_d = 0$.

300 4.2. Conditional expectation approximation

The next step consists in relating the optimal predictive performances for the regression and the classification problem (X, Y) to the performances associated to $(\widehat{X}_{\lambda_d,\tau_d}, Y)$ when d goes to infinity, i.e., relating L^* to

1. binary classification case:

$$L_d^* = \inf_{\phi:\mathcal{H}^m \to \{-1,1\}} \mathbb{P}\left(\phi(\widehat{X}_{\lambda_d,\tau_d}) \neq Y\right),$$

305 2. regression case:

$$L_d^* = \inf_{\phi: \mathcal{H}^m \to \mathbb{R}} \mathbb{E}\left(\left[\phi(\widehat{X}_{\lambda_d, \tau_d}) - Y \right]^2 \right)$$

Two sets of assumptions will be investigated to provide the convergence 307 of the Bayes risk L_d^* to L^* :

308 Assumption 4. Either

³⁰⁹
$$(A4a) \mathbb{E}(\|D^m X\|_{L^2}^2)$$
 is finite and $Y \in \{-1, 1\}$,
³¹⁰ **or**

311 (A4b) $\tau_d \subset \tau_{d+1}$ and $\mathbb{E}(Y^2)$ is finite.

The first assumption (A4a) requires an additional smoothing property for the predictor functional variable X and is only valid for a binary classification problem whereas the second assumption (A4a) requires an additional property for the sampling point series: they have to be growing sets. Theorem 2 then leads to the following corollary:

317 Corollary 3. Under Assumptions (A1)-(A4), we have:

$$\lim_{d\to+\infty}L_d^*=L^*$$

³¹⁸ 5. General consistent functional classifiers and regression functions

³¹⁹ 5.1. Definition of classifiers and regression functions on derivatives

Let us now consider any consistent classification or regression scheme for 320 standard multivariate data based either on the inner product or on the Eu-32 clidean distance between observations. Examples of such classifiers are Sup-322 port Vector Machine Steinwart (2002), the kernel classification rule Devroye 323 and Krzyżak (1989) and k-nearest neighbors Devroye and Györfi (1985); 324 Zhao (1987) to name a few. In the same way, multilayer perceptrons Lu-325 gosi and Zeger (1990), kernel estimates Devroye and Krzyżak (1989) and 326 k-nearest neighbors regression Devroye et al. (1994) are consistent regression 327 estimators. Additional examples of consistent estimators in classification and 328 regression can be found in Devroye et al. (1996); Györfi et al. (2002). 329

We denote $\psi_{\mathcal{D}}$ the estimator constructed by the chosen scheme using a dataset $\mathcal{D} = \{(U_i, T_i)_{1 \leq i \leq n}\}$, where the $(U_i, T_i)_{1 \leq i \leq n}$ are *n* independent copies of a pair of random variables (U, T) with values in $\mathbb{R}^p \times \{-1, 1\}$ (classification) or $\mathbb{R}^p \times \mathbb{R}$ (regression).

The proposed functional scheme consists in choosing the estimator ϕ_{n,τ_d} as $\psi_{\mathcal{E}_{n,\tau_d}}$ with the dataset \mathcal{E}_{n,τ_d} defined by:

$$\mathcal{E}_{n,\tau_d} = \{ (\mathbf{Q}_{\lambda_d,\tau_d} \mathbf{X}_i^{\tau_d}, Y_i)_{1 \le i \le n} \}$$

As pointed out in Section 3.4, the linear transformation $\mathbf{Q}_{\lambda_d,\tau_d}$ is an approximate multivariate differentiation operator: up to the boundary conditions, an estimator based on $\mathbf{Q}_{\lambda_d,\tau_d} \mathbf{X}^{\tau_d}$ is working on the *m*-th derivative of $\widehat{X}_{\lambda_d,\tau_d}$. In more algorithmic terms, the estimator is obtained as follows:

 $_{340}$ 1. choose an appropriate value for λ_d

2. compute $\mathbf{M}_{\lambda_d,\tau_d}$ using Theorem 1 and Corollary 2;

342 3. compute the Cholesky decomposition of $\mathbf{M}_{\lambda_d,\tau_d}$ and the transpose of 343 the Cholesky triangle, $\mathbf{Q}_{\lambda_d,\tau_d}$ (such that $\mathbf{Q}_{\lambda_d,\tau_d}^T \mathbf{Q}_{\lambda_d,\tau_d} = \mathbf{M}_{\lambda_d,\tau_d}$);

4. compute $\mathbf{Q}_{\lambda_d,\tau_d} \mathbf{X}_i^{\tau_d}$ to obtain the transformed dataset \mathcal{E}_{n,τ_d} ;

5. build a classifier/regression function $\psi_{\mathcal{E}_{n,\tau_d}}$ with a multivariate method in $\mathbb{R}^{|\tau_d|}$ applied to the dataset \mathcal{E}_{n,τ_d} ;

6. associate to a new sampled function $\mathbf{X}_{n+1}^{\tau_d}$ the prediction $\psi_{\mathcal{E}_{n,\tau_d}}(\mathbf{Q}_{\lambda,\tau_d}\mathbf{X}_{n+1}^{\tau_d}).$

Figure 5.1 illustrates the way the method performs: instead of relying on an approximation of the function and then on the derivation preprocessing of this estimates, it directly uses an equivalent metric by applying the $\mathbf{Q}_{\lambda_d,\tau_d}$ matrix to the sampled function. The consistency result proved in Theorem 3 shows that, combined with any consistent multidimensional learning algorithm, this method is (asymptotically) equivalent to using the original function drawn at the top left side of Figure 5.1.

On a practical point of view, Wahba (1990) demonstrates that cross validated estimates of λ achieve suitable convergence rates. Hence, steps 1 and 2 can be computed simultaneously by minimizing the total cross validated error for all the observations, given by

$$\sum_{i=1}^{n} \frac{1}{|\tau_d|} \sum_{t \in \tau_d} \frac{\left(x_i(t) - \widehat{x}_{i\lambda,\tau_d}(t)\right)^2}{\left(1 - A_{tt}(\lambda)\right)^2},$$

where A is a $|\tau_d| \times |\tau_d|$ matrix called the *influence matrix* (see Wahba (1990)), over a finite number of λ values.

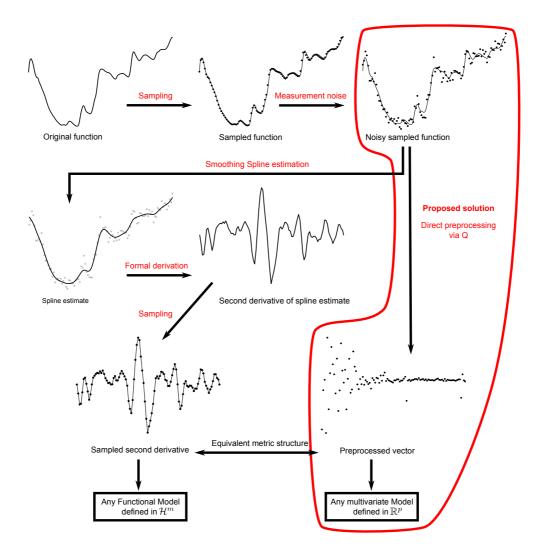


Figure 1: Method scheme and its equivalence to the usual approach for using derivatives in learning algorithms.

362 5.2. Consistency result

Corollary 1 and Corollary 3 guarantee that the estimator proposed in the
 previous section is consistent:

Theorem 3. Under assumptions (A1)-(A4), the series of classifiers/regression functions $(\phi_{n,\tau_d})_{n,d}$ is consistent:

$$\lim_{d \to +\infty} \lim_{n \to +\infty} \mathbb{E} \left(L \phi_{n, \tau_d} \right) = L^*$$

367 5.3. Discussion

While Theorem 3 is very general, it could be easily extended to cover special cases such as additional hypothesis needed by the estimation scheme or to provide data based parameter selections. We discuss briefly those issues in the present section.

It should first be noted that most estimation schemes, $\psi_{\mathcal{D}}$, depend on parameters that should fulfill some assumptions for the scheme to be consistent. For instance, in the Kernel Ridge Regression method in \mathbb{R}^p , with Gaussian kernel, $\psi_{\mathcal{D}}$ has the form given in Equation (7) where the (α_i) are the solutions of

$$\arg\min_{\alpha\in\mathbb{R}^n}\sum_{i=1}^n \left(T_i - \sum_{j=1}^n T_j\alpha_j e^{-\gamma \|U_i - U_j\|_{\mathbb{R}^p}^2}\right)^2 + \delta_n \sum_{i,j=1}^n T_i T_j \alpha_i \alpha_j e^{-\gamma \|U_i - U_j\|_{\mathbb{R}^p}^2}.$$

The method thus depends on the parameter of the Gaussian kernel, γ and of the regularization parameter δ_n . This method is known to be consistent if (see Theorem 9.1 of Steinwart and Christmann (2008)):

$$\delta_n \xrightarrow{n \to +\infty} 0$$
 and $n\delta_n^4 \xrightarrow{n \to +\infty} +\infty$.

Additional conditions of this form can obviously be directly integrated in Theorem 3 to obtain consistency results specific to the corresponding algorithms.

Moreover, practitioners generally rely on data based selection of the parameters of the estimation scheme $\psi_{\mathcal{D}}$ via a validation method: for instance, rather than setting δ_n to e.g., n^{-5} for *n* observations (a choice which is compatible with theoretical constraints on δ_n), one chooses the value of δ_n that ³⁸⁷ optimizes an estimation of the performances of the regression function ob-³⁸⁸ tained on an independent data set (or via a re-sampling approach).

In addition to the parameters of the estimation scheme, functional data 389 raise the question of the convenient order of the derivative, m, and of the 390 sampling grid optimality. In practical applications, the number of available 39: sampling points can be unnecessarily large (see Biau et al. (2005) for an ex-392 ample with more than 8 000 sampling points). The preprocessing performed 393 by $\mathbf{Q}_{\lambda_d, \tau_d}$ do not change the dimensionality of the data which means that 394 overfitting can be observed in practice when the number of sampling points 395 is large compared to the number of functions. Moreover, processing very 396 high dimensional vectors is time consuming. It is there quite interesting in 39 practice to use a down-sampled version of the original grid. 398

To select the parameters of ψ_D , the order of the derivative and/or the 399 down-sampled grid, a validation strategy, based on splitting the dataset into 400 training and validation sets, could be used. A simple adaptation of the idea 401 of Berlinet et al. (2008); Biau et al. (2005); Laloë (2008); Rossi and Villa 402 (2006) shows that a penalized validation method can be used to choose any 403 combination of those parameters consistently. According to those papers, 404 the condition for the consistency of the validation strategy would simply 405 relate the shatter coefficients of the set of classifiers in \mathbb{R}^d to the penalization 406 parameter of the validation. Once again, this type of results is a rather direct 407 extension of Theorem 3. 408

409 6. Applications

In this section, we show that the proposed approach works as expected on 410 real world spectrometric examples: for some applications, the use of deriva-411 tives leads to more accurate models than the direct processing of the spectra 412 (see e.g. Rossi et al. (2005); Rossi and Villa (2006) for other examples of such 413 a behavior based on ad hoc estimators of the spectra derivatives). It should 414 be noted that the purpose of this section is only to illustrate the behavior 415 of the proposed method on finite datasets. The theoretical results of the 416 present paper show that all consistent schemes have asymptotically identical 417 performances, and therefore that using derivatives is asymptotically useless. 418 On a finite dataset however, preprocessing can have strong influence on the 419 predictive performances, as will be illustrated in the present section. In ad-420 dition, schemes that are not universally consistent, e.g., linear models, can 42 lead to excellent predictive performances on finite datasets; such models are 422

423 therefore included in the present section despite the fact the theory does not 424 apply to them.

425 6.1. Methodology

The methodology followed for the two illustrative datasets is roughly the same:

1. the dataset is randomly split into a training set on which the model is 428 estimated and a test set on which performances are computed. The split 429 is repeated several times. The Tecator dataset (Section 6.2) is rather 430 small (240 spectra) and exhibits a rather large variability in predic-431 tive performances between different random splits. We have therefore 432 used 250 random splits. For the Yellow-berry dataset (Section 6.3), we 433 used only 50 splits as the relative variability in performances is far less 434 important. 435

2. λ is chosen by a global leave-one-out strategy on the spectra contained 436 in training set (as suggested in Section 5.1). More precisely, a leave-one-437 out estimate of the reconstruction error of the spline approximation of 438 each training spectrum is computed for a finite set of candidate values 439 for λ . Then a common λ is chosen by minimizing the average over 440 the training spectra of the leave-one-out reconstruction errors. This 441 choice is relevant as cross validation estimates of λ are known to have 442 favorable theoretical properties (see Craven and Wahba (1978); Utreras 443 (1981) among others). 444

3. for regression problems, a Kernel Ridge Regression (KRR) Saunders 445 et al. (1998); Shawe-Taylor and Cristianini (2004) is then performed to 446 estimate the regression function; this method is consistent when used 447 with a Gaussian kernel under additional conditions on the parameters 448 (see Theorem 9.1 of Steinwart and Christmann (2008)); as already ex-449 plained, in the applications, Kernel Ridge Regression is performed both 450 with a Gaussian kernel and with a linear kernel (in that last case, the 451 model is essentially a ridge regression model). Parameters of the models 452 (a regularization parameter, δ_n , in all cases and a kernel parameter, γ 453 for Gaussian kernels) are chosen by a grid search that minimizes a vali-454 dation based estimate of the performances of the model (on the training 455 set). A leave-one-out solution has been chosen: in Kernel Ridge Re-456 gression, the leave-one-out estimate of the performances of the model is 457

⁴⁵⁸ obtained as a by-product of the estimation process, without additional ⁴⁵⁹ computation cost, see e.g. Cawley and Talbot (2004).

Additionally, for a sake of comparison with a more traditional approach 460 in FDA, Kernel Ridge Regression is compared with a nonparametric 461 kernel estimate for the Tecator dataset (Section 6.2.1). Nonparametric 462 kernel estimate is the first nonparametric approach introduced in Func-463 tional Data Analysis Ferraty and Vieu (2006) and can thus be seen as 464 a basis for comparison in the context of regression with functional pre-465 dictors. For this method, the same methodology as with Kernel Ridge 466 Regression was used: the parameter of the model (i.e., the bandwidth) 467 was selected on a grid search minimizing a cross-validation estimate of 468 the performances of the model. In this case, a 4-fold cross validation 469 estimate was used instead of a leave-one-out estimate to avoid a large 470 computational cost. 471

4. for the classification problem, a Support Vector Machine (SVM) is used 472 Shawe-Taylor and Cristianini (2004). As KRR, SVM are consistent 473 when used with a Gaussian kernel Steinwart (2002). We also use a 474 SVM with a linear kernel as this is quite adapted for classification in 475 high dimensional spaces associated to sampled function data. We also 476 use a K-nearest neighbor model (KNN) for reference. Parameters of the 477 models (a regularization parameter for both SVM, a kernel parameter, 478 γ for Gaussian kernels and number of neighbors K for KNN) are chosen 479 by a grid search that minimizes a validation based estimate of the 480 classification error: we use a 4-fold cross-validation to get this estimate. 481

- 5. We evaluate the models obtained for each random split on the test set.
 We report the mean and the standard deviation of the performance
 index (classification error and mean squared error, respectively) and
 assess the significance of differences between the reported figures via
 paired Student tests (with level 1%).
- 6. Finally, we compare models estimated on the raw spectra and on spectra transformed via the $\mathbf{Q}_{\lambda_d,\tau_d}$ matrix for m = 1 (first derivative) and m = 2 (second derivative). For both values of m, we used the most classical boundary conditions (x(0) = 0 and Dx(0) = 0). Depending of the problem, other boundary conditions could be investigated but this is outside the scope of the present paper (see Besse and Ramsay (1986); Heckman and Ramsay (2000) for discussion on this subject). For the

⁴⁹⁴ Tecator problem, we also compare these approaches with models es-⁴⁹⁵ timated on first and second derivatives based on interpolating splines ⁴⁹⁶ (i.e. with $\lambda = 0$) and on first and second derivatives estimated by finite ⁴⁹⁷ differences.

Note that the kind of preprocessing used has almost no impact on 498 the computation time. In general, selecting the parameters of the 499 model with leave-one-out or cross-validation will use significantly more 500 computing power than constructing the splines and calculating their 501 derivatives. For instance, computing the optimal λ with the approach 502 described above takes less than 0.1 second for the Tecator dataset on a 503 standard PC using our R implementation which is negligible compared 504 to the several minutes used to select the optimal parameters of the 505 models used on the prepocessed data. 506

507 6.2. Tecator dataset

The first studied dataset is the standard Tecator dataset Thodberg (1996) 508 ¹. It consists in spectrometric data from the food industry. Each of the 509 240 observations is the near infrared absorbance spectrum of a meat sample 510 recorded on a Tecator Infratec Food and Feed Analyzer. Each spectrum is 511 sampled at 100 wavelengths uniformly spaced in the range 850–1050 nm. 512 The composition of each meat sample is determined by analytic chemistry 513 and percentages of moisture, fat and protein are associated this way to each 514 spectrum. 515

The Tecator dataset is a widely used benchmark in Functional Data Anal-516 ysis, hence the motivation for its use for illustrative purposes. More precisely, 517 in Section 6.2.1, we address the original regression problem by predicting the 518 percentage of fat content from the spectra with various regression method 519 and various estimates of the derivative preprocessing: this analysis shows 520 that both the method and the use of derivative have a strong effect on the 521 performances whereas the way the derivatives are estimated has almost no 522 effect. Additionally, in Section 6.2.2, we apply a noise (with various vari-523 ances) to the original spectra in order to study the influence of smoothing 524 in the case of noisy predictors: this section shows the relevance of the use of 525 a smoothing spline approach when the data are noisy. Finally, Section 6.2.3 526 deals with a classification problem derived from the original Tecator problem 527

¹Data are available on statlib at http://lib.stat.cmu.edu/datasets/tecator

(in the same way as what was done in Ferraty and Vieu (2003)): conclusions of this section are similar to the ones of the regression study.

530 6.2.1. Fat content prediction

As explained above, we first address the regression problem that consists 531 in predicting the fat content of peaces of meat from the Tecator dataset. The 532 parameters of the model are optimized with a grid search using the leave-one-533 out estimate of the predictive performances (both models use a regularization 534 parameter, with an additional width parameter in the Gaussian kernel case). 535 The original data set is split randomly into 160 spectra for learning and 80 536 spectra for testing. As shown in the result Table 1, the data exhibit a rather 537 large variability; we use therefore 250 random split to assess the differences 538 between the different approaches. 539

The performance indexes are the mean squared error (M.S.E.) and the $R^{2.2}$ As a reference, the target variable (fat) has a variance equal to 14.36. Results are summarized in Table 1.

The first conclusion is that the method itself has a strong effect on the 543 performances of the prediction: for this application, a linear method is not 544 appropriate (mean squared errors are much greater for linear methods than 545 for the kernel ridge regression used with a Gaussian kernel) and the non-546 parametric kernel estimate gives worse performances than the kernel ridge 547 regression (indeed, they are about 10 times worse). Nevertheless, for non-548 parametric approaches (Gaussian KKR and NKE), the use of derivatives 549 has also a strong impact on the performances: for kernel ridge regression, 550 e.g., preprocessing by estimating the first order derivative leads to a strong 551 decrease of the mean squared error. 552

⁵⁵³ Differences between the average MSEs are not always significant, but ⁵⁵⁴ we can nevertheless rank the methods in increasing order of modeling error ⁵⁵⁵ (using notations explained in Table 1) for Gaussian kernel ridge regression:

$$FD1 \le IS1 \le S1 < DF2 \le SS2 < IS2 < O$$

where < corresponds to a significant difference (for a paired Student test with level 1%) and \leq to a non significant one. In this case, the data are very smooth and thus the use of smoothing splines instead of a finite differences

 $²R^2 = 1 - \frac{\text{M.S.E}}{\text{Var}(y)}$ where Var(y) is the (empirical) variance of the target variable on the test set.

Method	Data	Average M.S.E. and SD	Average \mathbb{R}^2
KRR Linear	0	8.69 (4.47)	95.7%
	S1	8.09(3.85)	96.1%
	IS1	8.09(3.85)	96.1%
	FD1	8.27(4.17)	96.0%
	S2	9.64(4.98)	95.3%
	IS2	9.87(5.84)	95.2%
	FD2	8.45(4.18)	95.9%
KRR Gaussian	0	5.02(11.47)	97.6%
	S1	$0.485\ (0.385)$	99.8%
	IS1	$0.485\ (0.385)$	99.8%
	FD1	0.484 (0.387)	99.8%
	S2	$0.584\ (0.303)$	99.7%
	IS2	$0.586\ (0.303)$	99.7%
	FD2	$0.569\ (0.281)$	99.7%
NKE	0	73.1(16.5)	64.2%
	S1	4.59(1.09)	97.7%
	IS1	4.59(1.09)	97.7%
	FD1	4.59(1.09)	97.7%
	S2	3.75(1.22)	98.2%
	IS2	3.75(1.22)	98.2%
	FD2	3.67(1.18)	98.2%

Table 1: Summary of the performances of the chosen models on the test set (fat Tecator regression problem) when using either a kernel ridge regression (KRR) with linear kernel or with Gaussian kernel or when using a nonparametric kernel estimate (NKE) with various inputs: O (original data), S1 (smoothing splines with order 1 derivatives), IS1 (interpolating splines with order 1 derivatives), FD1 (order 1 derivatives estimated by finite differences) and S2, IS2 and FD2 (the same as previously with order 2 derivatives).

⁵⁵⁹ approximation does not have a significant impact on the predictions. How-⁵⁶⁰ ever, in this case, the roughest approach, consisting in the estimation of the ⁵⁶¹ derivatives by finite differences, gives the best performances.

562 6.2.2. Noisy spectra

This section studies the situation in which functional data observations are corrupted by noise. This is done by adding a noise to each spectrum of the Tecator dataset. More precisely, each spectrum has been corrupted by

$$X_i^b(t) = X_i(t) + \epsilon_{it} \tag{9}$$

where (ϵ_{it}) are i.i.d. Gaussian variables with standard deviation equal to either 0.01 (small noise) or to 0.2 (large noise). 10 observations of the data generated this way are given in Figure 2.

The same methodology as for the non noisy data has been applied to (X_i^b) to predict the fat content. The experiments have been restricted to the use of kernel ridge regression with a Gaussian kernel (according to the nonlinearity of the problem shown in the previous section). Results are summarized in Table 2 and Figure 3.

In addition, the results can be ranked this way:

Noise with sd equal to 0.01

$$S2 < S1 < IS1 \le O < FD1 < IS2 \le FD2$$

Noise with sd equal to 0.2

$$S1 < O < S2 < FD1 < IS1 < IS2 \le FD2$$

where < corresponds to a significant difference (for a paired Student test with level 1%).

The first conclusion of these experiments is that, even though the deriva-577 tives are the relevant predictors, their performances are strongly affected by 578 the noise (compared to the ones of the original data: note that the average 579 M.S.E. reported in Table 1 are more 10 times lower that the best ones from 580 Table 2 and that, in the best cases, R^2 is slightly greater than 50% for the 581 most noisy dataset). In particular, using interpolating splines or finite differ-582 ence derivatives leads to highly deteriorated performances. In this situation, 583 the approach proposed in the paper is particularly useful and helps to keep 584

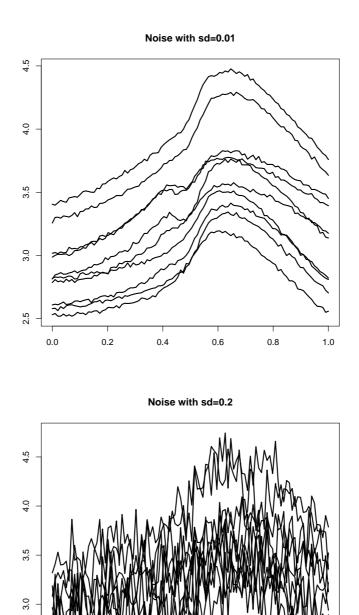


Figure 2: 10 observations of the noisy data generated from the Tecator spectra as in Equation 9 \$24\$

0.6

0.8

1.0

0.4

0.2

2.5

0.0

Noise	Data	Average M.S.E.	Average R^2
		and SD	
sd = 0.01	Ο	13.3(13.5)	93.5%
	S1	7.45(1.5)	96.4%
	IS1	12.72(2.2)	93.8%
	FD1	20.03(2.8)	90.3%
	S2	6.83 (1.4)	$\mathbf{96.7\%}$
	IS2	31.23(5.9)	84.9%
	FD2	31.10(5.9)	84.9%
sd = 0.2	0	87.9 (13.9)	57.4%
	S1	85.0 (12.5)	58.8 %
	IS1	210.1(36.1)	-1.9%
	FD1	209.1(33.0)	-1.4%
	S2	95.9(12.8)	53.5%
	IS2	213.7(33.1)	-3.6%
	FD2	235.1(222.7)	-14.0%

Table 2: Summary of the performances of the chosen models on the test set (fat Tecator regression problem) with noisy spectra.

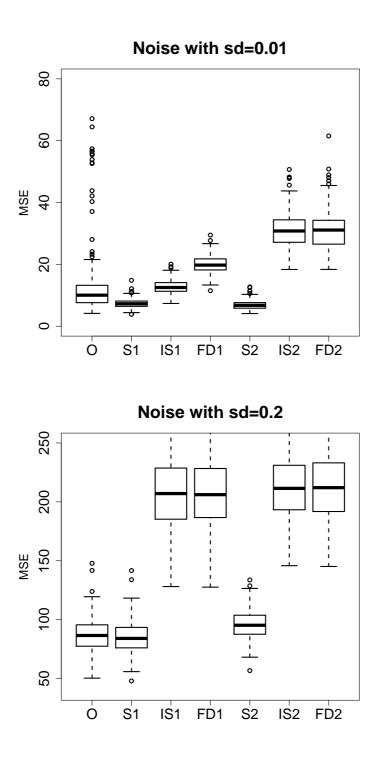


Figure 3: Mean squared errors boxplot for the noisy fat Tecator regression problem with Gaussian kernel (the worst test samples for S and FD have been removed for a sake of clarity)

better performances than with the original data. Indeed, the differences of 585 the smoothing splines approach with the original data is still significant (for 586 both derivatives in the "small noise" case and for the first order derivative 587 in the "high noise" case), even though, the most noisy the data are, the 588 most difficult it is to estimate the derivatives in an accurate way. That is, 589 except for smoothing spline derivatives, the estimation of the derivatives for 590 the most noisy dataset is so bad that it leads to negative R^2 when used in 591 the regression task. 592

⁵⁹³ 6.2.3. Fat content classification

In this section, the fat content regression problem is transformed into a classification problem. To avoid imbalance in class sizes, the median value of the fat in the dataset is used as the splitting criterion: the first class consists in 119 samples with strictly less than 13.5 % of fat, while the second class contains the other 121 samples with a fat content equal or higher than 13.5 %.

As in previous sections, the analysis is conducted on 250 random splits of the dataset into 160 learning spectra and 80 test spectra. We used stratified sampling: the test set contains 40 examples from each class. The 4 fold cross-validation used to select the parameters of the models on the learning set is also stratified with roughly 20 examples of each class in each fold.

The performance index is the mis-classification rate (MCR) on the test 605 set, reported in percentage and averaged over the 250 random splits. Results 606 are summarized in Table 3. As in the previous sections, both the model 607 and the preprocessing have some influence on the results. In particular, 608 using derivatives always improves the classification accuracy while the actual 609 method used to compute those derivatives has no particular influence on the 610 results. Additionally, using interpolation splines leads, in this particular 611 problem, to results that are exactly identical to the ones obtained with the 612 smoothing splines: they are not reported in Table 3. 613

More precisely, for the three models (linear SVM, Gaussian SVM and 614 KNN), differences in mis-classification rates between the smoothing spline 615 preprocessing and the finite differences calculation is never significant, ac-616 cording to a Student test with level 1 %. Additionally while the actual aver-617 age mis-classification rates might seem quite different, the large variability of 618 the results (shown by the standard deviations) leads to significant differences 619 only for the most obvious cases. In particular, SVM models using derivatives 620 (of order one or two) are indistinguishable one from another using a Student 621

Method	Data	Average MCR	SD of MCR
Linear SVM	0	1.41	1.55
	S1	0.73	1.15
	FD1	0.74	1.15
	S2	0.94	1.27
	FD2	0.92	1.23
Gaussian SVM	0	3.39	2.57
	S1	0.97	1.41
	FD1	0.98	1.42
	S2	0.99	2.00
	FD2	0.97	1.27
KNN	0	22.0	5.02
	S1	6.67	2.55
	FD1	6.57	2.55
	S2	1.93	1.65
	FD2	1.93	1.63

Table 3: Summary of the performances of the chosen models on the test set (Tecator fat classification problem). See Table 1 for notations. MCR stands for mis-classification rate, SD for standard deviation.

test with level 1 %: all methods with less than 1 % of mean mis-classification rate perform essentially identically. Other differences are significant: for instance the linear SVM used on raw data performs significantly worse than any SVM model used on derivatives.

It should be noted that the classification task studied in the present section is obviously simpler than the regression task from which it is derived. This explains the very good predictive performances obtained by simple models such as a linear SVM, especially with the proper preprocessing.

630 6.3. Yellow-berry dataset

The goal of the last experiment is to predict the presence of yellow-berry in durum wheat (*Triticum durum*) kernels via a near infrared spectral analysis (see Figure 4). Yellow-berry is a defect of the durum wheat seeds that reduces the quality of the flour produced from affected wheat. The traditional way to assess the occurrence of yellow-berry is by visual analysis of a sample of the seed stock. In the current application, a quality measure related to the occurrence of yellow-berry is predicted from the spectrum of the seed.

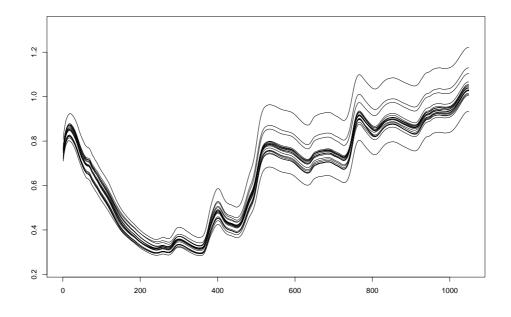


Figure 4: 20 observations of NIR spectra of durum wheat

The dataset consists in 953 spectra sampled at 1049 wavelengths uniformly spaced in the range 400–2498 nm. The dataset is split randomly into 640 600 learning spectra and 353 test spectra. Comparatively to the Tecator 641 dataset, the variability of the results is smaller in the present case. We used 642 therefore 50 random splits rather than 250 in the previous section.

The regression models were build via a Kernel Ridge Regression approach using a linear kernel and a Gaussian kernel. In both cases, the regularization parameter of the model is optimized by a leave-one-out approach. In addition, the width parameter of the Gaussian kernel is optimized via the same procedure at the same time.

The performance index is the mean squared error (M.S.E.). As a reference, the target variable has a variance of 0.508. Results are summarized in Table 4 and Figure 5.

Kernel and Data	Average M.S.E.	Standard deviation	Average R^2
Linear-O	0.122	8.7710^{-3}	76.1%
Linear-S1	0.138	9.5310^{-3}	73.0%
Linear-S2	0.122	8.4110^{-3}	76.1%
Gaussian-O	0.110	20.210^{-3}	78.5%
Gaussian-S1	0.0978	7.9210^{-3}	80.9%
Gaussian-S2	0.0944	8.3510^{-3}	81.5%

Table 4: Summary of the performances of the chosen models on the test set (durum wheat regression problem)

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As in the previous section, we can rank the methods in increasing order of modelling error, we obtain the following result:

$$G-S2 < G-S1 < G-O < L-O \le L-S2 < L-S1$$
,

where G stands for Gaussian kernel and L for linear kernel (hence G-S2 stands 653 for kernel ridge regression with gaussian kernel and smoothing splines with 654 order 2 derivatives); < corresponds to a significant difference (for a paired 655 Student test with level 1%) and < to a non significant one. For this appli-656 cation, there is a significant gain in using a non linear model (the Gaussian 657 kernel). In addition, the use of derivatives leads to less contrasted perfor-658 mances that the ones obtained in the previous section but it still improves 659 the quality of the non linear model in a significant way. In term of normal-660 ized mean squared error (mean squared error divided by the variance of the 661

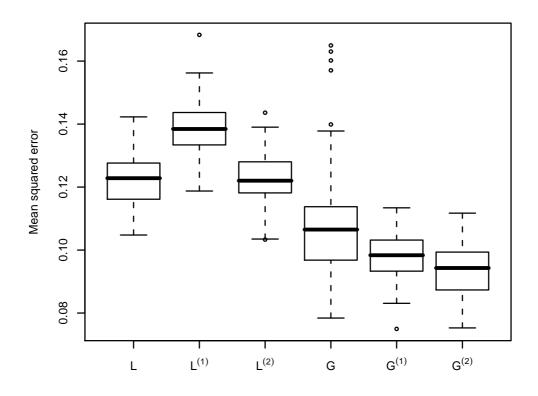


Figure 5: Mean squared error boxplots for the "durum wheat" regression problem (see Table 4 for the full names of the regression models)

target variable), using a non linear model with the second derivatives of the spectra corresponds to an average gain of more than 5% (i.e., a reduction of the normalised mean squared error from 24% for the standard linear model to 18.6%).

666 7. Conclusion

In this paper we proposed a theoretical analysis of a common practice that 667 consists in using derivatives in classification or regression problems when the 668 predictors are curves. Our method relies on smoothing splines reconstruction 669 of the functions which are known only via a discrete deterministic sampling. 670 The method is proved to be consistent for very general classifiers or regres-671 sion schemes: it reaches asymptotically the best risk that could have been 672 obtained by constructing a regression/classification model on the true ran-673 dom functions. 674

We have validated the approach by combining it with nonparametric regression and classification algorithms to study two real-world spectrometric datasets. The results obtained in these applications confirm once again that relying on derivatives can improve the quality of predictive models compared to a direct use of the sampled functions. The way the derivatives are estimated does not have a strong impact on the performances except when the data are noisy. In this case, the use of smoothing splines is quite relevant.

In the future, several issues could be addressed. An important practical 682 problem is the choice of the best order of the derivative, m. We consider 683 that a model selection approach relying on a penalized error loss could be 684 used, as is done, in e.g., Rossi and Villa (2006), to select the dimension of 685 truncated basis representation for functional data. Note that in practice, 686 such parameter selection method could lead to select m = 0 and therefore to 687 automatically exclude derivative calculation when it is not needed. This will 688 extend the application range of the proposed model. 689

A second important point to study it the convergence rate for the method. It would be very convenient for instance, to be able to relate the size of the sampling grid to the number of functions. But, this latter issue would require the use of additional assumptions on the smoothness of the regression function whereas the result presented in this paper, even if more limited, only needs mild conditions.

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823 9. Proofs

824 9.1. Theorem 1

In the original theorem (Lemma 3.1) in Kimeldorf and Wahba (1971), one has to verify that $(k_0(t_l,.))_l$ spans \mathcal{H}_0^m and that $(k_1(t_l,.))_l$ are linearly independent. These are consequences of Assumption (A1).

independent. These are consequences of Assumption (A1). First, $k_0(s,t) = \sum_{i,j=0}^{m-1} b_{ij}^{(-1)} s^i t^j$ where $\widetilde{B} = (b_{i,j}^{(-1)})_{i,j}$ is the inverse of $(\sum_{l=1}^m B^l s^i B^l t^j)_{i,j}$ (see Heckman and Ramsay (2000)). Then $(k_0(t_1,s),\ldots,k_0(t_{|\tau_d|},s)) = (1,s,\ldots,s^{m-1})\widetilde{B}[V_{m-1}(t_1,\ldots,t_{|\tau_d|})]^T$ where $V_{m-1}(t_1,\ldots,t_{|\tau_d|})$ is the Vandermonde matrix with m-1 columns and $|\tau_d|$ rows associated to values $t_1,\ldots,t_{|\tau_d|}$. If the $(t_l)_l$ are distinct, this matrix is of full rank.

Moreover the reproducing property shows that $\sum_{l=1}^{|\tau_d|} a_l k_1(t_l, .) \equiv 0$ implies $\sum_{l=1}^{|\tau_d|} a_l f(t_l) \equiv 0$ for all $f \in \mathcal{H}_1^m$. Hence, $\mathcal{H}_1^m = \text{Ker} \left(B^T, \sum_{l=1}^{\tau_d} a_l \zeta_l\right)^T$ where ζ_l denotes the linear form $h \in \mathcal{H}^m \to h(t_l)$. As the co-dimension of \mathcal{H}_1^m is $\dim \mathcal{H}_0^m = m$ and as, by Assumption (A1), B is linearly independent of $\sum_{l=1}^{\tau_d} a_l \zeta_l$, we thus have $\sum_{l=1}^{\tau_d} a_l \zeta_l \equiv 0$ (or codim Ker $\left(B^T, \sum_{l=1}^{\tau_d} a_l \zeta_l\right)^T =$ dim Im $\left(B^T, \sum_{l=1}^{\tau_d} a_l \zeta_l\right)$ would be m+1). Thus, we obtain that $\sum_{l=1}^{|\tau_d|} a_l f(t_l) \equiv$ 0 for all f in \mathcal{H}^m and, as (t_l) are distinct, that $a_l = 0$ for all l, leading to the independence conclusion for the $(k_1(t_l, .))_l$.

Finally, we prove that S_{λ,τ_d} is of full rank. Indeed, if $S_{\lambda,\tau_d} \mathbf{x}^{\tau_d} = 0$, $\omega^T M_0 \mathbf{x}^{\tau_d} = 0$ and $\eta^T M_1 \mathbf{x}^{\tau_d} = 0$. As $(\omega_k)_k$ is a basis of \mathcal{H}_0^m , $\omega^T M_0 \mathbf{x}^{\tau_d} = 0$ implies $M_0 \mathbf{x}^{\tau_d} = 0$ and therefore $M_1 = (K_1 + \lambda I_d)^{-1}$. As shown above, the $(k_1(t_l, .))_l$ are linearly independent and therefore $\eta M_1 \mathbf{x}^{\tau_d} = 0$ implies $M_1 \mathbf{x}^{\tau_d} = 0$, which in turns leads to $\mathbf{x}^{\tau_d} = 0$ via the simplified formula for M_1 .

847 9.2. Corollary 1

We give only the proof for the classification case, the regression case is identical.

According to Theorem 1, there is a full rank linear mapping from $\mathbb{R}^{|\tau_d|}$ to \mathcal{H}^m , $\mathcal{S}_{\lambda,\tau_d}$, such that for any function $x \in \mathcal{H}^m$, $\hat{x}_{\lambda,\tau_d} = \mathcal{S}_{\lambda,\tau_d} \mathbf{x}^{\tau_d}$. Let us denote $\mathcal{I}_{\lambda,\tau_d}$ the image of $\mathbb{R}^{|\tau_d|}$ by $\mathcal{S}_{\lambda,\tau_d}$, $\mathbf{P}_{\lambda,\tau_d}$ the orthogonal projection from \mathcal{H}^m to $\mathcal{I}_{\lambda,\tau_d}$ and $\mathcal{S}_{\lambda,\tau_d}^{-1}$ the inverse of $\mathcal{S}_{\lambda,\tau_d}$ on $\mathcal{I}_{\lambda,\tau_d}$. Obviously, we have $\mathcal{S}_{\lambda,\tau_d}^{-1} \circ \mathbf{P}_{\lambda,\tau_d}(\hat{x}_{\lambda,\tau_d}) = \mathbf{x}^{\tau_d}$.

Let ψ be a measurable function from $\mathbb{R}^{|\tau_d|}$ to $\{-1,1\}$. Then ζ_{ψ} defined on \mathcal{H}^m by $\zeta_{\psi}(u) = \psi \left(\mathcal{S}_{\lambda,\tau_d}^{-1} \circ \mathbf{P}_{\lambda,\tau_d}(u) \right)$ is a measurable function from \mathcal{H}^m to $\{-1,1\}$ (because $\mathcal{S}_{\lambda,\tau_d}^{-1}$ and $\mathbf{P}_{\lambda,\tau_d}$ are both continuous). Then for any measurable ψ , $\inf_{\phi:\mathcal{H}^m \to \{-1,1\}} \mathbb{P} \left(\phi(\widehat{X}_{\lambda,\tau_d}) \neq Y \right) \leq \mathbb{P} \left(\zeta_{\psi}(\widehat{X}_{\lambda,\tau_d}) \neq Y \right) =$ $\mathbb{P} \left(\psi(\mathbf{X}^{\tau_d}) \neq Y \right)$, and therefore

$$\inf_{\substack{\phi:\mathcal{H}^m\to\{-1,1\}}} \mathbb{P}\left(\phi(\widehat{X}_{\lambda,\tau_d})\neq Y\right) \leq \inf_{\substack{\phi:\mathbb{R}^{|\tau_d|}\to\{-1,1\}}} \mathbb{P}\left(\phi(\mathbf{X}^{\tau_d})\neq Y\right).$$
(10)

Conversely, let ψ be a measurable function from \mathcal{H}^m to $\{-1, 1\}$. Then ζ_{ψ} defined on $\mathbb{R}^{|\tau_d|}$ by $\zeta_{\psi}(\mathbf{u}) = \psi(\mathcal{S}_{\lambda,\tau_d}(\mathbf{u}))$, is measurable. Then for any measurable ψ , $\inf_{\phi:\mathbb{R}^{|\tau_d|}\to\{-1,1\}} \mathbb{P}(\phi(\mathbf{X}^{\tau_d})\neq Y) \leq \mathbb{P}(\zeta_{\psi}(\mathbf{X}^{\tau_d})\neq Y) = \mathbb{P}(\psi(\widehat{X}_{\lambda,\tau_d})\neq Y)$, and therefore

$$nf_{\phi:\mathbb{R}^{|\tau_d|} \to \{-1,1\}} \mathbb{P}\left(\phi(\mathbf{X}^{\tau_d}) \neq Y\right) \leq \inf_{\phi:\mathcal{H}^m \to \{-1,1\}} \mathbb{P}\left(\phi(\widehat{X}_{\lambda,\tau_d}) \neq Y\right).$$
(11)

The combination of equations (10) and (11) gives equality (4).

865 9.3. Corollary 3

1. Suppose assumption (A4a) is fullfilled

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- The proof is based on Theorem 1 in Faragó and Györfi (1975). This theorem relates the Bayes risk of a classification problem based on (X, Y) with the Bayes risk of the problem $(T_d(X), Y)$ where (T_d) is a series of transformations on X.
- More formally, for a pair of random variables (X, Y), where X takes values in \mathcal{X} , an arbitrary metric space, and Y in $\{-1, 1\}$, let us

denote for any series of functions T_d from \mathcal{X} to itself, $L^*(T_d) =$ 873 $\inf_{\phi: \mathcal{X} \to \{-1,1\}} \mathbb{P}(\phi(T_d(X)) \neq Y)$. Theorem 1 from Faragó and Györfi 874 (1975) states that $\mathbb{E}\left(\delta(T_d(X), X)\right) \xrightarrow{d \to +\infty} 0$ implies $L^*(T_d) \xrightarrow{d \to +\infty} L^*$, 875 where δ denotes the metric on \mathcal{X} . 876

This can be applied to $\mathcal{X} = (\mathcal{H}^m, \langle ., . \rangle_{L^2})$ with $T_d(X) =$ 877 $\widehat{X}_{\lambda_d,\tau_d} = S_{\lambda_d,\tau_d} \mathbf{X}^{\tau_d}$: under Assumptions (A1) and (A2), Theo-878 rem 2 gives: $||T_d(X) - X||_{L^2}^2 \leq \left(A_{R,m}\lambda_d + B_{R,m}\frac{1}{|\tau_d|^{2m}}\right) ||D^m X||_{L^2}^2$. Taking the expectation of both sides gives $\mathbb{E}\left(||T_d(X) - X||_{L^2}\right) \leq$ 879 880 $\left(A_{R,m}\lambda_d + B_{R,m}\frac{1}{|\tau_d|^{2m}}\right) \mathbb{E}\left(\|D^m X\|_{L^2}^2\right)$, using the fact that the constants are independent of the function under analysis. Then under Assump-882 tions (A4a) and (A3), $\mathbb{E}(||T_d(X) - X||_{L^2}) \xrightarrow{d \to +\infty} 0.$ According to 883 Faragó and Györfi (1975), this implies $\lim_{d\to\infty} L_d^* = L^*$. 884

2. Suppose assumption (A4b) is fullfilled

The conclusion will follow both for classification case and for regression case. The proof follows the general ideas of Biau et al. (2005); Rossi and Conan-Guez (2006); Rossi and Villa (2006); Laloë (2008). Under assumption (A1), by Theorem 1 and with an argument similar to those developed in the proof of Corollary 1, $\sigma(\widehat{X}_{\lambda_d,\tau_d}) = \sigma(\{X(t)\}_{t \in \tau_d})$. From assumption (A4b), $\sigma({X(t)}_{t \in \tau_d})$ is clearly a filtration. Moreover, as $\mathbb{E}(Y)$ and thus $\mathbb{E}(Y^2)$ are finite, $\mathbb{E}(Y|\widehat{X}_{\lambda_d,\tau_d})$ is a uniformly bounded martingal for this filtration (see Lemma 35 of Pollard (2002)). This martingale converges in L^1 -norm to $\mathbb{E}\left(Y|\sigma\left(\bigcup_d \sigma(\widehat{X}_{\lambda_d,\tau_d})\right)\right)$; we have

- $\sigma\left(\cup_d \sigma(\widehat{X}_{\lambda_d,\tau_d})\right) \subset \sigma(X)$ as $\widehat{X}_{\lambda_d,\tau_d}$ is a function of X (via Theorem 1);
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- by Theorem 2, $\widehat{X}_{\lambda_d,\tau_d} \xrightarrow{d \to +\infty, \text{ surely}} X$ in L^2 which proves that X is $\sigma\left(\bigcup_d \sigma(\widehat{X}_{\lambda_d,\tau_d})\right)$ -measurable.
- $\mathbb{E}\left(Y|\sigma\left(\bigcup_{d}\sigma(\widehat{X}_{\lambda_{d},\tau_{d}})\right)\right)$ Finally, $\mathbb{E}(Y|X)$ and 899 $\mathbb{E}\left(Y|\widehat{X}_{\lambda_d,\tau_d}\right) \xrightarrow{d \to +\infty, \ L^1} \mathbb{E}\left(Y|X\right).$ 900
- The conclusion follows from the fact that: 901

(a) binary classification case: the bound $L_d^* - L^* \leq 2\mathbb{E}\left(\left|\mathbb{E}\left(Y|\widehat{X}_{\lambda_d,\tau_d}\right) - \mathbb{E}\left(Y|X\right)\right|\right)$ (see Theorem 2.2 of Devroye et al. (1996)) concludes the proof;

(b) regression case: as $\mathbb{E}(Y^2)$ is finite, $\mathbb{E}\left(\mathbb{E}\left(Y|\widehat{X}_{\lambda_d,\tau_d}\right)^2\right)$ is also finite and the convergence also happens for the quadratic norm (see Corollary 6.22 in Kallenberg (1997)), i.e.,

$$\lim_{d \to +\infty} \mathbb{E}\left(\left(\mathbb{E}\left(Y|X\right) - \mathbb{E}\left(Y|\widehat{X}_{\lambda_d, \tau_d}\right)\right)^2\right) = 0$$

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Hence, as $L_d^* - L^* = \mathbb{E}\left(\left(\mathbb{E}\left(Y|X\right) - \mathbb{E}\left(Y|\widehat{X}_{\lambda_d,\tau_d}\right)\right)^2\right)$, the conclusion follows.

910 9.4. Theorem 3

911 We have

$$L(\phi_{n,d}) - L^* = L\phi_{n,\tau_d} - L^*_d + L^*_d - L^*.$$
 (12)

⁹¹² Let ϵ be a positive real. By Corollary 3, it exists $d_0 \in \mathbb{N}^*$ such that, for all ⁹¹³ $d \geq d_0$,

$$L_d^* - L^* \le \epsilon. \tag{13}$$

Moreover, as shown in Corollary 1 and as $\mathbf{Q}_{\lambda_d,\tau_d}$ is invertible, we have in the binary classification case: $L_d^* = \inf_{\phi:\mathbb{R}^{|\tau_d|} \to \{-1,1\}} \mathbb{P}(\phi(\mathbf{X}^{\tau_d}) \neq Y) =$ $\inf_{\phi:\mathbb{R}^{|\tau_d|} \to \{-1,1\}} \mathbb{P}(\phi(\mathbf{Q}_{\lambda_d,\tau_d}\mathbf{X}^{\tau_d}) \neq Y)$, and in the regression case: $L_d^* =$ $\inf_{\phi:\mathbb{R}^{|\tau_d|} \to \mathbb{R}} \mathbb{E}\left([\phi(\mathbf{X}^{\tau_d}) - Y]^2\right) = \inf_{\phi:\mathbb{R}^{|\tau_d|} \to \mathbb{R}} \mathbb{E}\left([\phi(\mathbf{Q}_{\lambda_d,\tau_d}\mathbf{X}^{\tau_d}) - Y]^2\right)$. By hypothesis, for any fixed d, ϕ_{n,τ_d} is consistent, that is

$$\lim_{n \to +\infty} \mathbb{E} \left(L(\phi_{n,\tau_d}) \right) = \inf_{\phi: \mathbb{R}^{|\tau_d|} \to \{-1,1\}} \mathbb{P} \left(\phi \left(\mathbf{Q}_{\lambda_d,\tau_d} \mathbf{X}^{\tau_d} \right) \neq Y \right),$$

⁹¹⁹ in the classification case and

$$\lim_{n \to +\infty} \mathbb{E}\left(L(\phi_{n,\tau_d})\right) = \inf_{\phi: \mathbb{R}^{|\tau_d|} \to \mathbb{R}} \mathbb{E}\left(\left[\phi\left(\mathbf{Q}_{\lambda_d,\tau_d}\mathbf{X}^{\tau_d}\right) - Y\right]^2\right),$$

in the regression case. and therefore for any fixed d_0 , 920 $\lim_{n \to +\infty} \mathbb{E}\left(L(\phi_{n,\tau_{d_0}})\right)$ $= L_{d_0}^*.$ Combined with equations (12) and 921 (13), this concludes the proof. 922