

# Model Selection for Scalar-on-Function Regression with Applications to Near-Infrared Spectroscopy

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# Contents

<b>1</b>	<b>Colour Guide</b>	<b>3</b>
<b>2</b>	<b>Introduction</b>	<b>3</b>
<b>3</b>	<b>Theory</b>	<b>4</b>
3.1	Inner Products and Hilbert Spaces . . . . .	4
3.2	Hilbert Space of Square-Integrable FunctionsRandom Function in $L^2[0, 1]$ space . . . .	5
3.3	Functional Data Set . . . . .	5
3.4	Representing a Function in terms of a Basis . . . . .	6
3.5	Approximation and Smoothing via Basis Truncation . . . . .	7
3.6	Functional Data Sets . . . . .	8
3.7	Karhunen-Loève Expansion and Empirical Eigenbases . . . . .	8
3.8	Scalar-on-Function Regression . . . . .	10
3.8.1	Estimation using Basis-Representation . . . . .	11
3.8.2	Estimation using Functional Principal Components . . . . .	13
3.9	Literature . . . . .	15
<b>4</b>	<b>Simulation Study</b>	<b>16</b>
4.1	Motivation . . . . .	16
4.2	Generating Similar Curves . . . . .	16
4.3	Simulation setup . . . . .	17
4.4	Results . . . . .	18
4.4.1	Basis Expansion Regression . . . . .	18
4.4.2	Functional Principal Component Regression . . . . .	19
4.4.3	Interpretation and Relevance for Application . . . . .	19
<b>5</b>	<b>Application</b>	<b>20</b>
5.1	Literature . . . . .	20
<b>6</b>	<b>Outlook</b>	<b>20</b>
6.1	Literature . . . . .	20
<b>7</b>	<b>Appendix</b>	<b>21</b>
7.1	Near-infrared (NIR) Spectroscopy . . . . .	21
7.2	Basis Plots . . . . .	22
7.3	Simulation Study Results . . . . .	23
<b>8</b>	<b>Proofs</b>	<b>27</b>
8.1	Lemma . . . . .	27
8.2	Theorem (Karhunen-Loève expansion) . . . . .	28
<b>9</b>	<b>Bibliography</b>	<b>30</b>
<b>10</b>	<b>Affidavit</b>	<b>32</b>

# 1 Colour Guide

- **RED**: is for general comments for your own text
- **GREEN**: is for Jona's comments
- **ORANGE**: is for Jonghun's comments
- **BLUE**: is for Jakob's comments

# 2 Introduction

Functional Data Analysis (FDA) is (roots in the 1940s Grenander and Karhunen) gaining more attention as researchers from different fields collect data that is functional in nature. Although classical statistical methods can often process this data, but only FDA allows extracting the information given by the smoothness of the underlying process of the functional data (cf. Levitin et al. 2007). As Kokoszka and Reimherr 2017 describe, FDA should be considered when one can view variables or units of a given data set as smooth curves or functions and the interest is in analyzing samples of curves (cf. Kokoszka and Reimherr 2017, S. 17).

In functional data analysis the concept of a data set can include not only realizations of scalar random variables, but also realizations of random functions which could be the absorption of the Near-infrared (NIR) spectrum which will be used for the simulation and application. NIR-spectroscopy uses the near-infrared region of the electromagnetic spectrum (780nm to 2500nm) to measure the absorption of its waves on a sample to analyse the interaction of each other. For more details refer to the appendix SECTION ?!?!?. This dataset can then be used to perform Functional Linear Regression (FLR) on a scalar response variable. To conduct functional linear regression, several new concepts and methods need to be considered and explained that distinguish it multivariate linear regression. The focus of this paper is to introduce FLR in a scalar-on-function setting. We will be using the standard FLR framework, which relates functional predictors to a scalar response as follows:

$$Y(\omega) = \alpha + \int_0^1 X(\omega)(s)\beta(s)ds + \epsilon(\omega), \quad i = 1, \dots, n \quad (1)$$

where the  $X_i$  are realizations of a random function  $\mathbf{X}$ ,  $Y_i$  are the corresponding realizations of the response variable and  $\beta(s)$  is the coefficient function. The distinct feature of this framework is that the regressor is a function, which necessitates a different approach to estimation. As in the well-known framework of scalar linear regression, this is motivated by an interest in  $\beta(s)$  for prediction. For instance, fluctuation in  $X_i(s)$  at a point  $s_0$  will not have any effect on  $Y_i$  if  $\beta(s_0) = 0$ .

Estimation of  $\beta(s)$  is inherently an infinite-dimensional problem. In Section 2, after introducing the necessary theoretical concepts, we describe three methods of estimating a scalar coefficient function using a concept called truncated basis expansion. We report the results of the Monte-Carlo simulation regarding these three different methods in Section 3. Finally, in Section 4, we test the prediction of FLR in a real-world setting. (We may put some simple descriptions of results about each of MC and Application)

### 3 Theory

In multivariate regression, data is often observed in the form of elements from Euclidean space,  $\mathbb{R}^p$ . However, the statistics derived from infinite-dimensional random functions cannot be defined on a finite dimensional space. To understand functional linear regression and the differences between the methods presented in this paper, it is therefore necessary to introduce some concepts and extend known aspects of linear regression theory to include functional objects. One integral concept in inferential statistics are random variables. Paraphrasing a definition by Bauer 2020, a random variable  $X : \Omega \rightarrow \Omega'$  is an  $\mathcal{A}$ - $\mathcal{A}'$ -measurable function, where  $(\Omega, \mathcal{A}, P)$  is a probability space and  $(\Omega', \mathcal{A}')$  is a measure space. A typical case known to every undergraduate student of economics is  $(\Omega', \mathcal{A}') = (\mathbb{R}, \mathcal{B})$ , where  $\mathcal{B}$  is the canonical  $\sigma$ -algebra on the real numbers. As a first intuition, it is possible to imagine a similar concept where a random variable does not realize as an element of the real numbers but as a function in a function space. A formalization of this idea makes some more theoretical considerations necessary. The following theoretical introduction closely follows chapters 2.3 and 2.4 from Hsing and Eubank 2015.

#### 3.1 Inner Products and Hilbert Spaces

Let  $\mathbb{V}$  be a vector space over some field of scalars  $\mathbb{F}$ . In the following we restrict our analysis to the case of  $\mathbb{F} = \mathbb{R}$ . A function  $\langle \cdot, \cdot \rangle : \mathbb{V} \times \mathbb{V} \rightarrow \mathbb{F}$  is called an inner product, if  $\forall v, v_1, v_2 \in \mathbb{V}$  and  $a_1, a_2 \in \mathbb{F}$  the following properties hold.

1.  $\langle v, v \rangle \geq 0$
2.  $\langle v, v \rangle = 0$  if  $v = 0$
3.  $\langle a_1 v_1 + a_2 v_2, v \rangle = a_1 \langle v_1, v \rangle + a_2 \langle v_2, v \rangle$
4.  $\langle v_1, v_2 \rangle = \langle v_2, v_1 \rangle$

A vector space with an associated inner product is called an inner product space. [verbatim quote!] The inner product naturally defines a norm and an associated distance on the vector space as follows.

$$\|v\| = \langle v, v \rangle^{\frac{1}{2}} \quad (2)$$

$$d(v_1, v_2) = \langle v_2 - v_1, v_2 - v_1 \rangle^{\frac{1}{2}} \quad (3)$$

If the inner product space is complete with respect to the induced distance, it is called a Hilbert space, denoted  $\mathbb{H}$  in the following. To extend the known concept of a basis in a finite dimensional space to the potentially infinite Hilbert spaces, it is necessary to define the closed span of a sequence of elements of  $\mathbb{H}$ . Recall that the span of a set of vectors  $S \subseteq \mathbb{R}^P$  is given by

$$\text{span}(S) = \left\{ \sum_{i=1}^k \lambda_i v_i \mid k \in \mathbb{N}, v_i \in S, \lambda_i \in \mathbb{R} \right\} \quad (4)$$

The closed span  $\overline{\text{span}}(S)$  of a sequence  $S$  in  $\mathbb{H}$  is defined as the closure of the span with respect to the distance induced by the norm.  $S$  is called a basis of  $\mathbb{H}$  if  $\overline{\text{span}}(S) = \mathbb{H}$ .

It is called an orthonormal basis, if in addition the following properties hold.

1.  $\langle v_i, v_j \rangle = 0 \quad \forall v_i, v_j \in S \quad i \neq j$
2.  $\|v\| = 1 \quad \forall v \in S$

As in the case of a Banach space, each element of a Hilbert space can be expressed in terms of a corresponding basis. This can be done using a Fourier expansion of an element  $x \in \mathbb{H}$  w.r.t. a basis

$S = \{s_n\}$  as follows.

$$x = \sum_{j=1}^{\infty} \langle x, s_j \rangle s_j \quad (5)$$

As can be seen, differing from the case of a Banach space, these representations can be limits of series as previously hinted at by using the closed span of the basis. As using an infinite number of basis functions is infeasible in applied contexts, an intuitive way to approximate elements of a Hilbert space, is to use a truncated series.

$$x \approx \sum_{j=1}^K \langle x, s_j \rangle s_j \quad (6)$$

### 3.2 Hilbert Space of Square-Integrable Functions **Random Function in $\mathbb{L}^2[0, 1]$ space**

In functional data analysis, one Hilbert space of particular importance is the space of square-integrable functions on  $[0, 1]$  denoted  $\mathbb{L}^2[0, 1]$ . To define it, look first at the measure space given by  $([0, 1], \mathcal{B}, \mu)$  where  $\mathcal{B}$  is the Borel  $\sigma$ -algebra on  $[0, 1]$  and  $\mu$  is the Lebesgue-measure. Then  $\mathbb{L}^2[0, 1]$  is the collection of all measurable functions  $f$  on  $[0, 1]$  that fulfill the following condition.

$$\|f\|_2 = \int_0^1 |f|^2 d\mu < \infty \quad (7)$$

This ensures that a random function has a finite second moment so that the variance and covariance functions can be defined. Its inner product is defined as

$$\langle f_1, f_2 \rangle = \int_0^1 f_1 f_2 d\mu. \quad (8)$$

$\mathbb{L}^2[0, 1]$  is the function space that is most often used for theoretical considerations in functional data analysis but analogous constructions can be made for every closed interval of  $\mathbb{R}$ . A random function defined on  $\mathbb{L}^2[0, 1]$  is a function  $X : \Omega \rightarrow \mathbb{L}^2[0, 1]$  defined on a common probability space  $(\Omega, \mathcal{A}, P)$  where  $\Omega$  is a sample space with  $\sigma$ -algebra  $\mathcal{A}$  and a probability space  $P$ . The realized  $X(\omega)(t)$  for every  $t \in [0, 1]$  is called a sample curve for the process. The collection of such sample curves constitutes a functional data set.

### 3.3 Functional Data Set

Consider the case of a data set containing observations  $x_i$  of a random function  $X$ , measured at discrete points  $t_{i,j}$ :

$$x_i(t_{i,j}) \in \mathbb{R}, \quad i = 1, \dots, N, \quad j = 1, \dots, J_i, \quad t_{i,j} \in [T_1, T_2] \quad (9)$$

Each curve  $x_i(t)$  exists  $\forall t \in [T_1, T_2]$ , but is only observed at discrete measurement points. As it can be seen, the measurment points can be differentiated through each sample. However, this is not a primary interest of this paper, and, therefore, we consider the case where every measurement point exists evenly, namely, that  $t_{i,j}$  and  $J_i$  are identical for all  $i$ . In principal, the recorded functional data can be seen as a set of discretized realizations of a random function since it is finitely countable. The several basis systems to make the functional data smooth is introduced from the next sub-chapter. As in the finite dimensional setting, the concept of identically distributed and independent data is

important for the functional data analysis. The random functions share the identical and independent distribution, which makes the statistical inferential procedures possible. An example of such a data set could be the data set of 60 NIR-spectra of gasoline samples in which each spectra was randomly selected and evenly measured.

### 3.4 Representing a Function in terms of a Basis

As previously described, a basis of a Hilbert space can be used to express elements of the space using the corresponding Fourier expansion. Three examples of bases that are often used in practice to approximate elements of  $\mathbb{L}^2[0, 1]$  are explained in the following.

Both of these can be used to express or in the case of the b-spline basis approximate elements of  $\mathbb{L}^2[0, 1]$  as a weighted sum of basis functions. Let therefore  $\{\phi_i(t) \mid i \in \mathcal{I}\}$  be the basis used to express / approximate a realization  $X(\omega_0) = x(t)$  of  $X(\omega)$ .

$$X(\omega_0) = x(t) = \sum_{j \in \mathcal{I}} A_j(\omega_0) \phi_j(t) \quad (10)$$

**Monomial Basis** No shift parameter in implementation in fda. Monomial bases are trivially decompositions of every polynomial function. The class of non-polynomial but enough many differentiable functions can be approximated by Taylor series, a special type of power series. From these concepts, a function can be smooth by a power series. The bases construct

$$x(t) = \sum_{i \in \mathcal{I}} c_i (t - \alpha)^i \quad (11)$$

which forms the basis

$$\phi_i^M(t) = (t - \alpha)^k \quad i \in \mathcal{I} \quad (12)$$

where the shift parameter  $\alpha$  is specified to be the center of the interval subject to approximation. The first order function within the monomial basis system, for  $k = 0$ , is called the constant basis system (cf. Horváth and Kokoszka 2012). There exist collinearity problems of this basis system since the monomial basis functions become more correlated to each other as the degrees increase, which would result in numerically unstable situation. This restricts the number of basis functions. Therefore, regarding the use, this may be useful for relatively simple functions. However, the lower order monomial function makes it impracticable to capture pronounced local peculiarities. Moreover, it leads to an undesirable behavior at the tails such as Taylor expansion. (cf. Ramsay and Silverman 2005)

**Fourier Basis** Fourier series basically decomposes a periodic function with a weighted summation of trigonometric functions. Imagine that an arbitrary function in  $\mathbb{L}^2[0, 1]$  has one or more cycles between 0 and 1, and we can then apply the Fourier series to approximate the function for smoothing. As the series originated from the amplitude-phase form regarding cosine functions, it can be rewritten by sine and cosine form through a simple trigonometry formula, which restricts the use of an odd number of Fourier basis functions. That is why the Fourier basis for  $\mathbb{L}^2[0, 1]$  is given by the following sequence of

functions defined on  $[0, 1]$ .

$$\phi_i^F(x) = \begin{cases} 1 & \text{if } i = 1 \\ \sqrt{2} \cos(\pi i x) & \text{if } i \text{ is even} \\ \sqrt{2} \sin(\pi(i-1)x) & \text{otherwise} \end{cases} \quad (13)$$

Therefore, its basis functions inherits a repeating behavior which is useful to expand functions that represent an periodic or seasonal underlying process over the period  $T$ . Rephrasing from Ramsay and Silverman 2005, the Fourier basis functions are orthonormal when the values  $t_j$  are equally spaced within  $T$ . This basis is suitable to expand functions with a similar curvature order across the domain, resulting generally in uniformly smooth expansions. [Motive from Fourier series](#)

**B-spline Basis** Following chapter 3.5 from Ramsay and Silverman 2005, splines are defined by first dividing the interval of interest  $[\tau_0, \tau_L]$  into  $L$  subintervals of non-negative length divided by a non-decreasing sequence of points  $(\tau_l)_{l=1, \dots, L-1}$  called knots. On each subinterval, a spline is a polynomial of chosen order  $m = n + 1$  where  $n$  is its degree. Additionally, at each  $\tau_l$  the the polynomials on neighbouring subintervals must match derivatives up to order  $m - 2$ . A B-spline is a spline belonging to a basis system developed by Boor 1978. Let  $\phi_{l,m}^{BS}(x)$   $l = 1, \dots, L - 1$  be the B-spline of order  $m$  for an interval  $[\tau_0, \tau_L]$  and knots  $\{\tau_l \mid l = 1, \dots, L - 1\}$ , then it is defined by the Cox-de Boor recursion formula as follows.

$$\phi_{l,0}^{BS}(x) = \begin{cases} 1 & \text{if } x \in [\tau_l, \tau_{l+1}) \\ 0 & \text{otherwise} \end{cases} \quad (14)$$

$$\phi_{l,m}^{BS}(x) = \frac{x - \tau_l}{\tau_{l+m} - \tau_l} \phi_{l,m-1}^{BS}(x) + \frac{\tau_{l+m+1} - x}{\tau_{l+m+1} - \tau_{l+1}} \phi_{l+1,m-1}^{BS}(x)$$

This, however, does not really yield a basis of  $\mathbb{L}^2[0, 1]$  as the closed span of this finite sequence of functions is not equal to  $\mathbb{L}^2[0, 1]$ . To really obtain a basis of  $\mathbb{L}^2[0, 1]$  from B-splines, further theoretical considerations about, for example, infinite series of B-splines and specific knot choices would have to be made. As this is out of the scope of this paper, for the sake of simplicity, we will assume that a B-spline basis representation of a function in  $\mathbb{L}^2[0, 1]$  will serve as a sufficient approximation for an appropriately chosen B-spline basis. Even though, this approach is not theoretically exact, in practice, this is often a reasonable approach and yields satisfactory results in cases where the functional form of B-splines makes them an appropriate approximation tool.

Explain what happens when  $l + m + 1 > L$  !!!

Additionally, modify for multiple knots at the boundaries to get better behavior at the boundaries. Needs some more explanation.

### 3.5 Approximation and Smoothing via Basis Truncation

As above-mentioned, the realized curves can be estimated with basis functions. For the basis expansion, it is technically possible to use all basis functions  $\phi_{j \in \mathcal{I}}$  to fit an completely unbiased estimate. This is not efficient since the expansion with too many basis functions introduces high amounts of variance in

pronounced local variations of the curves or can even led to the approximation of noise in the sample curves, which possibly interrupts the analysis. On the other hand, the important information on the curves could be missed with a too-small number of basis functions. In that sense, this discussion is subject to the Bias-Variance Tradeoff. It challenges the researcher to seek a point at which they truncate the basis function to remove noise and, at the same time, not to introduce too much bias by maintaining significant fluctuations of the curves. The basis expansion with truncation is defined by

$$X(\omega_0) = x(t) = \sum_{j \in \mathcal{I}} A_j(\omega_0) \phi_j(t) = \sum_{j=1}^L A_j(\omega_0) \phi_j(t) + \delta(t) \approx \sum_{j=1}^L A_j(\omega_0) \phi_j(t) \quad (15)$$

where  $\delta(t)$  is the truncation error and  $L \leq \max_{j \in \mathcal{I}}(j)$  for all  $L \in \mathcal{I}$ . The number  $L$  can be chosen subjectively, but also through applying a data-driven method like Cross-Validation, which aims to minimize the Mean Squared Error (MSE) or other criteria. The significance of this choice of truncation becomes more evident in [Link estimation](#) and is the subject of the simulation study.

### 3.6 Functional Data Sets

**Move this part below 3.2?** Consider the case of a data set containing observations  $x_i$  of a random function  $X$ , measured at discrete points  $t_{i,j}$ :

$$x_i(t_{i,j}) \in \mathbb{R}, \quad i = 1, \dots, N, \quad j = 1, \dots, J_i, \quad t_{i,j} \in [T_1, T_2] \quad (16)$$

Each curve  $x_i(t)$  exists  $\forall t \in [T_1, T_2]$ , but is only observed at discrete measurement points. As it can be seen, the measurment points can be differentiated through each sample. However, this is not a primary interest of this paper, and, therefore, we consider the case where every measurement point exists evenly, namely, that  $t_{i,j}$  and  $J_i$  are identical for all  $i$ . In principal, the recorded functional data can be seen as a set of discretized realizations of a random function since it is finitely countable. The several basis systems to make the functional data smooth is introduced from the next sub-chapter.

As in the finite dimensional setting, the concept of identically distributed and independent data is important for the functional data analysis. The random functions share the identical and independent distribution, which makes the statistical inferential procedures possible. An example of such a data set could be the data set of 60 NIR-spectra of gasoline samples in which each spectra was randomly selected and evenly measured.

### 3.7 Karhunen-Loève Expansion and Empirical Eigenbases

Given a realization of a random function realizing in  $\mathbb{L}^2[0, 1]$ , it is possible to represent this realization in terms of its generating stochastic process. To obtain the analogous concept for a random function, it is necessary to define the covariance function of a random function realizing in  $\mathbb{L}^2[0, 1]$ . Therefore, let  $X : \Omega \mapsto \mathbb{L}^2[0, 1]$  be such a random function. Then the mean and covariance functions of  $X$  are defined as follows.

$$\mu(t) = \mathbb{E}[X(\omega)(t)] \quad (17)$$

$$c(t, s) = \mathbb{E}[(X(\omega)(t) - \mu(t))(X(\omega)(s) - \mu(s))] \quad (18)$$



where the  $c(t, s)$  are Hilbert-Schmidt Kernels defined through  $c : [0, 1] \times [0, 1] \rightarrow \mathbb{R}$ . Let  $K$  be the integral operator on  $\mathbb{L}^2[0, 1]$  such that  $K : \nu \rightarrow K\nu$  for  $\nu \in \mathbb{L}^2[0, 1]$ , by

$$[K\nu](t) = \int_0^1 c(t, s)\nu(s)ds = \lambda\nu(t) \quad (19)$$

Then, the operator  $K$  has orthonormal basis functions  $\nu^m \in \mathbb{L}^2[0, 1]$  corresponding to eigenvalues  $\lambda^m$  for all  $m$  since it is a positive compact self-adjoint operator (cf. Alexanderian 2015). Moreover,  $K$  holds that the eigenvalues can be ordered in nonincreasing order as follows  $\lambda^1 \geq \lambda^2 \geq \dots \geq 0$  where the superscript is not the power but index. Therefore, the functions  $X$  are approximated enough well by first few principal components since the order of them are sorted in descending order of eigenvalues corresponding to the eigenfunctions (e.g.  $Var(\xi^m) \geq Var(\xi^n)$  for all  $m < n$ ). Theoretical considerations lead to the result that  $X$  can be represented in the following form, called its Karhunen-Lo  ve expansion. The proofs are provided at 8.1 and 8.2.

$$X(\omega)(t) = \mu(t) + \sum_{m \in \mathbb{N}} \xi^m(\omega) \nu^m(t), \quad \xi^m(\omega) = \int_0^1 (X(\omega)(s) - \mu(s)) \nu^m(s) ds \quad (20)$$

where the  $\nu^m$  are defined by the countable set of solutions  $\{(\lambda^m, \nu^m) \mid m \in \mathbb{N}\}$  of (19). The random variables  $\xi^m(\omega)$  satisfy following properties.

1.  $\mathbb{E}[\xi^m(\omega)] = 0$
2.  $Cov(\xi^m(\omega), \xi^n(\omega)) = \delta^{m,n} \lambda^m$
3.  $Var(\xi^m(\omega)) = \lambda^m$

where  $\delta^{m,n} = 0$  if  $m \neq n$ , otherwise 1. In the typical scalar setting, a similar consideration leads to the concept of principal components. This is also possible in a functional setting. Let  $\{x_1(t), \dots, x_n(t)\}$  be a set of i.i.d. realizations generated by a random function  $X(\omega) \mapsto \mathbb{L}^2[0, 1]$ . Define the following sample analogues for the mean and covariance functions.

$$\hat{\mu}(t) = \frac{1}{n} \sum_{i=1}^n x_i(t) \quad (21)$$

$$\hat{c}(t, s) = \frac{1}{n} \sum_{i=1}^n (x_i(t) - \hat{\mu}(t)) (x_i(s) - \hat{\mu}(s)) \quad (22)$$

With these it is possible to derive a set of sample analogs  $\{(\hat{\lambda}^m, \hat{\nu}^m) \mid m \in \mathcal{M}\}$  for  $\{(\lambda^m, \nu^m) \mid m \in \mathbb{N}\}$  as the solutions of the following equation. **I think the number of principal components has to be smaller than the number of observations. So I can give more information about  $\mathcal{M}$ .**

$$\int_0^1 \hat{c}(t, s) \hat{\nu}(s) ds = \hat{\lambda} \hat{\nu}(t) \quad (23)$$

This naturally leads to the following representation.

$$x_i(t) = \hat{\mu}(t) + \sum_{j=1}^{!!!} \hat{\xi}_i^j \hat{\nu}^j(t) \quad (24)$$

where the  $\hat{\xi}_i^m$  are derived as

$$\hat{\xi}_i^m(\omega) = \langle x_i - \hat{\mu}, \hat{\nu}^m \rangle = \int_0^1 (x_i(s) - \hat{\mu}(s)) \hat{\nu}^m(s) ds \quad (25)$$

In reality these calculations are often done using basis representations of both the functional principal components  $\hat{\nu}^m$  and the observations  $x_i(t)$  leading to the following representation. For the sake of clarity the following equation assumes that the bases used for the expansion of both the observations and the coefficient function are true bases of  $\mathbb{L}^2[0, 1]$  and can therefore be used to express the corresponding objects exactly.

$$\begin{aligned} \hat{\xi}_i^m &= \int_0^1 (x_i(s) - \hat{\mu}(s)) \hat{\nu}^m(s) ds = \int_0^1 \left( \sum_{j \in \mathcal{I}} a_{i,j} \phi_j(s) \right) \left( \sum_{k \in \mathcal{L}} b_k^m \psi_k(s) \right) ds \\ &= \int_0^1 \left( \sum_{j=1}^J a_{i,j} \phi_j(s) + \delta_i^J(s) \right) \left( \sum_{k=1}^K b_k^m \psi_k(s) + \delta_\beta^K(s) \right) ds \\ &= \sum_{j=1}^J \left[ a_{i,j} \sum_{k=1}^K b_k^m \int_0^1 \phi_j(s) \psi_k(s) ds \right] + \sum_{k=1}^K b_k^m \int_0^1 \delta_i^J(s) \psi_k(s) ds + \sum_{j=1}^J a_{i,j} \int_0^1 \phi_j(s) \delta_\beta^K(s) ds \end{aligned} \quad (26)$$

In practice, a typical choice is to use the same basis  $(\phi_j(t))_{j \in \mathcal{I}}$  and the same truncation parameter  $L$  for the basis expansion of both the demeaned observations  $(x_i(t) - \hat{\mu}(t))$  and the functional principal components  $\hat{\nu}^m$ . This leads to the following simplification of Equation 26.

$$\hat{\xi}_i^m = \sum_{j=1}^L \left[ a_{i,j} \sum_{k=1}^L b_k^m \int_0^1 \phi_j(s) \psi_k(s) ds \right] + \sum_{k=1}^L b_k^m \int_0^1 \delta_i^L(s) \psi_k(s) ds + \sum_{j=1}^L a_{i,j} \int_0^1 \phi_j(s) \delta_\beta^L(s) ds \quad (27)$$

And we can define the following objects:

$$\tilde{\xi}_i^{m,L} := \sum_{j=1}^L \left[ a_{i,j} \sum_{k=1}^L b_k^m \int_0^1 \phi_j(s) \psi_k(s) ds \right] \quad \delta_{\xi,i}^L := \hat{\xi}_i^m - \tilde{\xi}_i^{m,L} \quad (28)$$

$$\tilde{\nu}^{m,L}(t) := \sum_{k=1}^L b_k^m \phi_k(t) \quad \delta_{\nu,m}^L(t) := \hat{\nu}^m(t) - \tilde{\nu}^{m,L}(t) \quad (29)$$

### 3.8 Scalar-on-Function Regression

In the simple scalar setting one of the most important tools in econometrics is the linear regression. Its goal is to predict the value of a dependent variable given a set of associated variables. For reference assume a data generating process as follows. **Structure is bad... This will be changed.**

$$Y = X\beta + \epsilon \quad (30)$$

Where  $Y$  is the vector of response variables,  $X$  is the matrix containing the corresponding regressors in its columns and  $\beta = (\beta_0, \beta_1, \dots, \beta_p)'$  is the vector containing the unknown coefficients. In this finite dimensional setting one important question is how to estimate the unknown coefficients  $\beta$ . The most well known estimator in all of econometrics, the Ordinary Least Squares (OLS) estimator, fulfills this purpose under a set of assumptions. **List Assumptions? Then we need to list the assumptions for**

functional linear regression as well I think.

$$\hat{\beta}_{OLS} = (X'X)^{-1}X'Y \quad (31)$$

The concept of linear regression can be extended to a setting of functional data, where a scalar response variable is supposed to be predicted from a functional variable. A general data generating process in this functional scenario could look like the following equation.

$$Y(\omega) = \alpha + \Psi(X(\omega)) + \epsilon(\omega) \quad (32)$$

Here  $\Psi$  is a functional that maps a realization of a random function in  $\mathbb{L}^2[0, 1]$  into  $\mathbb{R}$ . One simple example to illustrate the principle is the maximum  $\Psi(f) = \max_{x \in [0, 1]} f(x)$ . However, the typical setup is often as follows mimicking the structure of the multivariate linear model extended from summation to integration. This structure is crucial for the extension of linear regression to the case of functional regressors. Therefore, in this paper we always implicitly assume that a data generating process has the following structure. **maybe a better motivation is the one Jona gave: from very dense observations that run into problems of colinearity when using the original OLS estimator**

$$Y(\omega) = \alpha + \int_0^1 \beta(s)X(\omega)(s)ds + \epsilon(\omega) \quad (33)$$

Where  $x(t)$  is the realization of a random function in  $\mathbb{L}^2[0, 1]$  and  $\beta(t)$  is an unknown coefficient function. Similar to the finite dimensional setting, an interesting question is how to estimate the unknown function  $\beta(t)$  given a data set containing realizations of a random function and associated scalar response variables. However, a simple extension of the OLS estimator to allow for infinite dimensional objects is not possible. Therefore, other options have to be considered.

### 3.8.1 Estimation using Basis-Representation

The most common way to make this problem tractable is via a basis representation of  $\beta(t)$ . Therefore, let  $\{b_i(t) \mid i \in \mathcal{I}\}$  be a basis of  $\mathbb{L}^2[0, 1]$  and represent  $\beta(t)$  in terms of this basis.

$$\beta(t) = \sum_{j \in \mathcal{I}} b_j \phi_j(t) \quad (34)$$

This enables us to write equation 33 with  $\beta(t)$  represented in this way to obtain a formulation as a sum of scalar random variables  $Z_j(\omega)$ .

$$\begin{aligned} Y(\omega) &= \alpha + \int_0^1 \beta(s)X(\omega)(s)ds + \epsilon(\omega) = \alpha + \int_0^1 \left[ \left( \sum_{j \in \mathcal{I}} b_j \phi_j(s) \right) X(\omega)(s) \right] ds + \epsilon(\omega) \\ &= \alpha + \sum_{j \in \mathcal{I}} \left[ b_j \int_0^1 X(\omega)(s) \phi_j(s) ds \right] + \epsilon(\omega) = \alpha + \sum_{j \in \mathcal{I}} b_j Z_j(\omega) + \epsilon(\omega) \end{aligned} \quad (35)$$

This representation translates the original problem of regressing a scalar on a continuously observed function to a problem where a scalar is regressed on what is possibly a countably infinite sequence of regressors. Using a truncation of the basis at some parameter  $L$  can be used to make this problem tractable with typical theory from multivariate regression while staying reasonably accurate.

$$\begin{aligned}
Y(\omega) &= \alpha + \int_0^1 \left[ \left( \sum_{j=1}^J b_j \phi_j(s) + \delta_\beta^J(s) \right) X(\omega)(s) \right] ds + \epsilon(\omega) \\
&= \alpha + \sum_{j=1}^J b_j \int_0^1 \phi_j(s) X(\omega)(s) ds + \int_0^1 \delta_\beta^J(s) X(\omega)(s) ds + \epsilon(\omega)
\end{aligned} \tag{36}$$

In practice it is common to not only express the coefficient function in terms of a basis but also the observations. Therefore two bases  $((\phi_j(t))_{j \in \mathcal{I}})$  and  $(\psi_k(t))_{k \in \mathcal{L}}$  and two corresponding truncation parameters ( $J$  and  $K$ ) can be chosen. This leads to the following representation.

$$\begin{aligned}
Y(\omega) &= \alpha + \int_0^1 \beta(s) X(\omega)(s) ds + \epsilon(\omega) = \alpha + \int_0^1 \left[ \left( \sum_{j \in \mathcal{I}} b_j \phi_j(s) \right) \left( \sum_{k \in \mathcal{L}} a_k(\omega) \psi_k(s) \right) \right] ds + \epsilon(\omega) \\
&= \alpha + \int_0^1 \left[ \left( \sum_{j=1}^J b_j \phi_j(s) + \delta_\beta^J(s) \right) \left( \sum_{k=1}^K a_k(\omega) \psi_k(s) + \delta_X^K(\omega)(s) \right) \right] ds + \epsilon(\omega) \\
&= \alpha + \sum_{j=1}^J b_j \left[ \sum_{k=1}^K a_k(\omega) \int_0^1 \phi_j(s) \psi_k(s) ds \right] + \sum_{j=1}^J b_j \int_0^1 \phi_j(s) \delta_X^K(\omega)(s) ds \\
&\quad + \sum_{k=1}^K a_k(\omega) \int_0^1 \delta_\beta^J(s) \phi_j(s) ds + \epsilon(\omega)
\end{aligned} \tag{37}$$

A typical choice in this scenario is to use the same functional basis  $(\phi_j(t))_{j \in \mathcal{I}}$  and the same truncation parameter  $L$  for both the coefficient function and the approximation of the observations. Defining the following notation

$$\tilde{Z}_j(\omega) = \sum_{k=1}^L \left[ a_k(\omega) \int_0^1 \phi_j(s) \phi_k(s) ds \right] \quad j = 1, \dots, L \tag{38}$$

This leads to a considerable simplification of Equation 37.

$$\begin{aligned}
Y(\omega) &= \alpha + \sum_{j=1}^J b_j \tilde{Z}_j(\omega) + \sum_{j=1}^J b_j \int_0^1 \phi_j(s) \delta_X^K(\omega)(s) ds + \sum_{k=1}^K a_k \int_0^1 \delta_\beta^J(s) \phi_j(s) ds + \epsilon(\omega) \\
&\approx \alpha + \sum_{j=1}^J b_j \tilde{Z}_j(\omega) + \epsilon(\omega)
\end{aligned} \tag{39}$$

A model in the form of Equation 39 lends itself to be estimated using theory from multivariate linear regression. Define therefore the following objects

$$Y = \begin{pmatrix} y_1 \\ \vdots \\ y_n \end{pmatrix}, \quad Z = \begin{pmatrix} 1 & \tilde{Z}_{1,1} & \dots & \tilde{Z}_{1,J} \\ \vdots & \vdots & \ddots & \vdots \\ 1 & \tilde{Z}_{n,1} & \dots & \tilde{Z}_{n,J} \end{pmatrix} \tag{40}$$

Then an OLS estimator can be calculated in the usual way to obtain an estimate for the values of  $\alpha$  and  $b_j$  and an estimate of the coefficient function can be derived accordingly.

$$b^L = (Z'Z)^{-1} Z'Y \in \mathbb{R}^{L+1} \quad \hat{\alpha} = b_1^L \quad \hat{\beta}^L(t) = \sum_{j=1}^J b_{j+1}^L \phi_j(t) \quad (41)$$

### 3.8.2 Estimation using Functional Principal Components

Using the Karhunen-Loève Expansion to represent  $X(\omega)$ , it is also possible to express the data generating process in a slightly different way.

$$\begin{aligned} Y(\omega) &= \alpha + \int_0^1 \textcolor{red}{X}(\omega)(s) \beta(s) ds + \epsilon(\omega) = \alpha + \int_0^1 \left( \mu(s) + \sum_{m=1}^{\infty} \xi^m(\omega) \nu^m(s) \right) \beta(s) ds + \epsilon(\omega) \\ &= \alpha + \int_0^1 \mu(s) \beta(s) ds + \sum_{m=1}^{\infty} \xi^m(\omega) \int_0^1 \nu^m(s) \beta(s) ds + \epsilon(\omega) = \bar{\alpha} + \sum_{m=1}^{\infty} \xi^m(\omega) \beta^m + \epsilon(\omega) \end{aligned} \quad (42)$$

As these theoretical Eigenfunctions and Eigenvalues are typically unknown, the corresponding equation in sample analogues is more interesting as a representation of an observation.

$$\begin{aligned} y_i &= \alpha + \int_0^1 \textcolor{red}{x}_i(s) \beta(s) ds + \epsilon_i = \alpha + \int_0^1 \left( \hat{\mu}(s) + \sum_{m \in \mathcal{M}} \hat{\xi}_i^m \hat{\nu}^m(s) \right) \beta(s) ds + \epsilon_i \\ &= \alpha + \int_0^1 \hat{\mu}(s) \beta(s) ds + \sum_{m \in \mathcal{M}} \hat{\xi}_i^m \int_0^1 \hat{\nu}^m(s) \beta(s) ds + \epsilon_i = \bar{\alpha} + \sum_{m \in \mathcal{M}} \hat{\xi}_i^m \hat{\beta}^m + \epsilon_i \end{aligned} \quad (43)$$

This, however, is a simplification for the purposes of real-world estimation as in most implementations the coefficient function and the principal components are also expressed or derived in terms of a basis that can be chosen freely. Introducing both concepts one step at a time leads to the following complication if we first introduce an expansion of the coefficient function.

$$\begin{aligned} y_i &= \alpha + \int_0^1 \textcolor{red}{x}_i(s) \textcolor{blue}{\beta}(s) ds + \epsilon_i = \alpha + \int_0^1 \left( \hat{\mu}(s) + \sum_{m \in \mathcal{M}} \hat{\xi}_i^m \hat{\nu}^m(s) \right) \left( \sum_{j \in \mathcal{I}} b_j \phi_j(s) \right) ds + \epsilon_i \\ &= \alpha + \int_0^1 \left[ \sum_{j \in \mathcal{I}} b_j \phi_j(s) \hat{\mu}(s) + \sum_{m \in \mathcal{M}} \left[ \hat{\xi}_i^m \sum_{j \in \mathcal{I}} b_j \hat{\nu}^m(s) \phi_j(s) \right] \right] ds + \epsilon_i \\ &= \alpha + \sum_{j \in \mathcal{I}} b_j \int_0^1 \phi_j(s) \hat{\mu}(s) ds + \sum_{m \in \mathcal{M}} \left[ \hat{\xi}_i^m \sum_{j \in \mathcal{I}} b_j \int_0^1 \hat{\nu}^m(s) \phi_j(s) ds \right] + \epsilon_i \end{aligned} \quad (44)$$

Truncating the basis used for expansion of the coefficient function already introduces an approximation error.

$$\begin{aligned}
y_i &= \alpha + \int_0^1 \left( \hat{\mu}(s) + \sum_{m \in \mathcal{M}} \hat{\xi}_i^m \hat{\nu}^m(s) \right) \left( \sum_{j=1}^J b_j \phi_j(s) + \delta_\beta^J(s) \right) ds + \epsilon_i \\
&= \alpha + \sum_{j=1}^J b_j \int_0^1 \phi_j(s) \hat{\mu}(s) ds + \int_0^1 \delta_\beta^J(s) \hat{\mu}(s) ds + \sum_{m \in \mathcal{M}} \left[ \hat{\xi}_i^m \sum_{j=1}^J b_j \int_0^1 \hat{\nu}^m(s) \phi_j(s) ds \right] \\
&\quad + \sum_{m \in \mathcal{M}} \left[ \hat{\xi}_i^m \int_0^1 \hat{\nu}^m(s) \delta_\beta^J(s) ds \right] + \epsilon_i
\end{aligned} \tag{45}$$

If we additionally derive and approximate the principal components and corresponding scores using a truncated basis representation as in Equation 27 we obtain the following. To not complicate things more than necessary, the following equation assumes that the same basis  $(\phi_j(t))_{j \in \mathcal{I}}$  was used in the derivation of the principal components and the expansion of the coefficient function. Additionally, the following approximation also truncates the basis for the expansion of the coefficient function at the same parameter  $L$  that was used for the approximation of the principal components and scores.

For convenience, define the following notation:

$$\tilde{\alpha}^L = \alpha + \sum_{j=1}^L b_j \int_0^1 \phi_j(s) \hat{\mu}(s) ds + \int_0^1 \delta_\beta^L(s) \hat{\mu}(s) ds \tag{46}$$

Then Equation 45 can be expressed as follows.

$$\begin{aligned}
y_i &= \tilde{\alpha}^L + \sum_{m \in \mathcal{M}} \left[ \left( \tilde{\xi}_i^{m,L} + \delta_{\xi,i}^{m,L} \right) \sum_{j=1}^L b_j \int_0^1 \left( \tilde{\nu}^{m,L}(s) + \delta_{\nu,m}^L(s) \right) \phi_j(s) ds \right] + \epsilon_i \\
&= \tilde{\alpha}^L + \sum_{m \in \mathcal{M}} \left[ \tilde{\xi}_i^{m,L} \sum_{j=1}^L b_j \int_0^1 \tilde{\nu}^{m,L}(s) \phi_j(s) ds \right] + \sum_{m \in \mathcal{M}} \left[ \tilde{\xi}_i^{m,L} \sum_{j=1}^L b_j \int_0^1 \delta_{\nu,m}^L(s) \phi_j(s) ds \right] \\
&\quad + \sum_{m \in \mathcal{M}} \left[ \delta_{\xi,i}^{m,L} \sum_{j=1}^L b_j \int_0^1 \tilde{\nu}^{m,L}(s) \phi_j(s) ds \right] + \sum_{m \in \mathcal{M}} \left[ \delta_{\xi,i}^{m,L} \sum_{j=1}^L b_j \int_0^1 \delta_{\nu,m}^L(s) \phi_j(s) ds \right] + \epsilon_i \\
&\approx \tilde{\alpha}^L + \sum_{m \in \mathcal{M}} \left[ \tilde{\xi}_i^{m,L} \sum_{j=1}^L b_j \int_0^1 \tilde{\nu}^{m,L}(s) \phi_j(s) ds \right] + \epsilon_i
\end{aligned} \tag{47}$$

The parameter  $M$  corresponds to the chosen number of principal components and therefore constitutes another choice in the approximation.

As in the previous section, this equation again lends itself for estimation with OLS. Define the following objects:

$$Y = \begin{pmatrix} y_1 \\ \vdots \\ y_n \end{pmatrix}, \quad Z = \begin{pmatrix} 1 & \tilde{\xi}_1^{1,L} & \dots & \tilde{\xi}_1^{M,L} \\ \vdots & \vdots & \ddots & \vdots \\ 1 & \tilde{\xi}_n^{1,L} & \dots & \tilde{\xi}_n^{M,L} \end{pmatrix} \tag{48}$$

We can then derive the following estimators.

$$\tilde{b}^{L,M} = (Z'Z)^{-1} Z'Y \in \mathbb{R}^{M+1} \quad \hat{\alpha} = \tilde{b}_1^{L,M} \quad \hat{\beta}^{m,L} = \tilde{b}_{m+1}^{L,M} \tag{49}$$

### 3.9 Literature

- Alexanderian 2015
- Kokoszka and Reimherr 2017
- Hsing and Eubank 2015
- Ramsay and Silverman 2005
- Horváth and Kokoszka 2012
- Cai and Hall 2006
- Levitin et al. 2007

## 4 Simulation Study

### 4.1 Motivation

In the simulation study, we deviate from the standard simulation setting. Instead of generating data by ourselves, we use the gasoline data which consists of 60 samples of Near-infrared (NIR) spectra measured by 2-nm from 900 to 1,700 nm, and a response variable, the octane rating. We chose this setup to improve the approach towards the application in which we predict the octane ratings from the gasoline dataset.

To exploit the regularity of the curves of the spectra, we introduced different basis functions in [Link](#) and demonstrated the importance of the truncation parameter  $K$  for the estimation in [Link](#). For the simulation study, we rely on the introduced estimation strategies with the introduced basis functions and focus on selecting the truncation parameter  $K$  as well as the number of FPC, which is as well affected by  $K$ , by ten fold cross-validation using the prediction mean-squared error. While cross-validation is common for selecting  $K$ , the number of FPC in practice is often truncated after a pre specified amount of explained variability Kokoszka and Reimherr 2017, which might not be optimal since FPC with smaller eigenvalues may have greater influence on the prediction (c.f Jolliffe 1982). This might apply for this simulation as well, since certain eigenfunctions could be resulting from certain chemical combinations and overtones in the absorption bands of the spectra, that could have high predictive power, but explain only little variability of the spectra [Link to NIR](#).

This setup is opposing to the often used penalized functional regression as described by Goldsmith et al. 2011 in which an explicit smoothness constraint  $\lambda$  is used to tune the smoothness of the estimator  $\hat{\beta}(t)$  while setting the  $K$  sufficiently high. This would avoid the heavy computing of validating the best value for  $K$  which we will conduct in the simulation. To provide intuition in this approach, let

$$PSSE_{\lambda}(\alpha, \beta) = \sum_{i=1}^N (Y_i - \alpha - \int_0^1 \beta(t) X_i(t) dt)^2 + \lambda * \int (D^m \beta(t))^2 dt \quad (50)$$

denote the penalized residual sum of squares as notated by Ramsay and Silverman 2005 for the derivative of order  $m$ . A typical choice is the second derivative as highly variable functions are expected to yield large second derivatives and therefore a larger penalty. The smoothing parameter  $\lambda$  is set to minimize the  $PSSE_{\lambda}(\alpha, \beta)$ , which can be archived by different criteria as shown in Thomas C.M. Lee 2003.

### 4.2 Generating Similar Curves

To avoid small sample problems, we generated 200 similar curves from the gasoline dataset, motivated by Karhunen-Loève Expansion. First, the initial curves are expressed in terms of a generated bspline basis which is created using 50 knots. These smooth curves are then centered, before applying the Karhunen-Loève Expansion. It is assumed that the scores follow a normal distribution, thus, the new realizations for the scores are drawn from a multivariate normal  $\tilde{\xi} = (\tilde{\xi}_1, \dots, \tilde{\xi}_J)' \sim \mathcal{N}(0_J, \text{diag}(\hat{\lambda}_1, \dots, \hat{\lambda}_J))$ . Finally, we obtain the generated curves  $NIR_{sim}$

$$\tilde{X}(\omega)(t) = \hat{\mu}(t) + \sum_{j=1}^J \tilde{\xi}_j(\omega) \hat{\nu}_j(t)$$

where  $\tilde{X}(\omega)(t)$ ,  $\hat{\mu}(t)$  and  $\hat{\nu}_j(t)$  are approximated as vectors in  $\mathbb{R}^{401}$ .



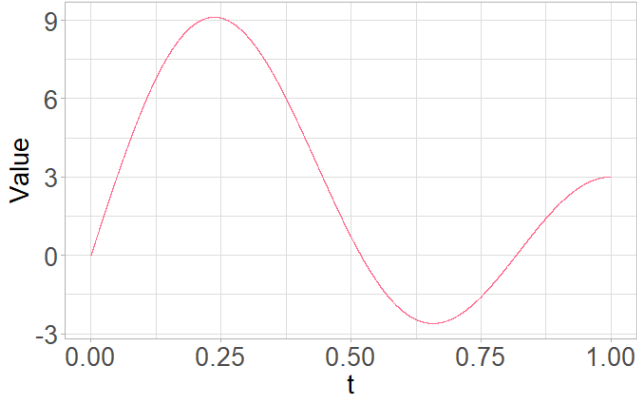


Figure 1:  $f_1(t)$ , smooth function

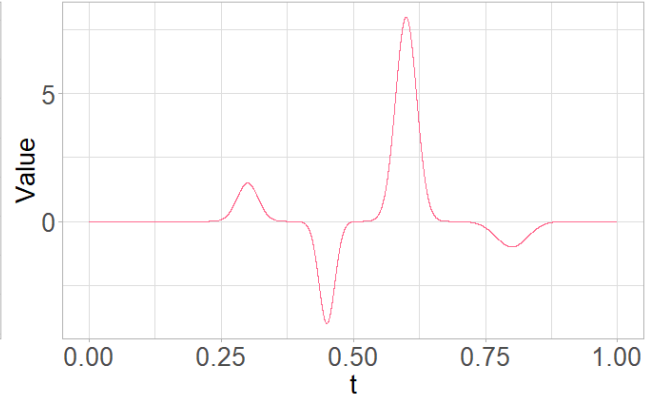


Figure 2:  $f_2(t)$ , bumpy function

### 4.3 Simulation setup

The simulation study follows Reiss and Ogden 2007 as a guideline. Two different true coefficient functions,  $f_1(t)$  and  $f_2(t)$ , are created that differ in their smoothness, to compare the introduced methods with differing true coefficient functions:

$$f_1(t) = 2 \sin(0.5\pi t) + 4 \sin(1.5\pi t) + 5 \sin(2.5\pi t) \quad (51)$$

$$\begin{aligned} f_2(t) = & 1.5 \exp\left(\frac{-0.5(t-0.3)^2}{0.02^2}\right) - 4 \exp\left(\frac{-0.5(t-0.45)^2}{0.015^2}\right) \\ & + 8 \exp\left(\frac{-0.5(t-0.6)^2}{0.02^2}\right) - \exp\left(\frac{-0.5(t-0.8)^2}{0.03^2}\right) \end{aligned} \quad (52)$$

The bumpy function,  $f_2(t)$ , was generated by referring to Cardot 2002. The smooth function  $f_1(t)$  follows Reiss and Ogden 2007 and its inner product  $\langle NIR_{sim}, f_1 \rangle$  creates responses that are similar to the original octane numbers of the gasoline dataset.

Two different error-terms  $\epsilon$  were created by first generating an *i.i.d.* standard normal error term and then multiplying it by two error variations  $\sigma_e$ . The error variations represent different signal-to-noise ratios to test the methods both with a low and a high amount of noise. They are created such that the squared multiple correlation coefficient  $R^2 = \text{var}(Xf) / (\text{var}(Xf) + \sigma_e^2)$  is equal to 0.9 and 0.6. The two error-terms are then used to generate two sets of responses for  $f \in \{f_1(t), f_2(t)\}$

$$\begin{aligned} Y_{1,f} &= \langle NIR, f \rangle + Z \left[ \frac{\text{var}(\langle NIR, f \rangle)}{0.9} - \text{var}(\langle NIR, f \rangle) \right] \\ Y_{2,f} &= \langle NIR, f \rangle + Z \left[ \frac{\text{var}(\langle NIR, f \rangle)}{0.6} - \text{var}(\langle NIR, f \rangle) \right] \end{aligned} \quad (53)$$

where  $Z \sim \mathcal{N}(0, 1)$ . In total, we created four combinations for the simulations, using the two true coefficient functions and the two sets of responses. These four combinations are then used with a different number of monomial basis functions  $\in \{1, 2, \dots, 6\}$ , cubic bspline basis-function  $\{5, 6, \dots, 18\}$  and fourier functions  $\{1, 3, \dots, 25\}$  to predict the generated responses using the basis expansion approach

and the FPCR approach. For the evaluation, we used the prediction MSE calculated by 10 fold cross-validation. To obtain valid out of sample properties for the FPCR, within each of the ten 10 fold cross-validation splits, we first calculate the first  $nharm$  FPC's for  $nharm \in \{2, 3, 4\}$  of the training-set  $\mathcal{T}$  for each curve  $i \in \{1, 2, \dots, 200\}$ , which was smoothed with the respective basis function specification. The obtained eigenfunctions  $\nu^{m,\mathcal{T}}$  are then used to estimate the scores of the holdout set  $\mathcal{H}$ ,  $\hat{\xi}_i^{m,\mathcal{H}}$  by the equation

I guess we don't use the scores of the training set to estimate scores of test set but eigenfunction of traningset...Adressed We should mention how we derived the original principal components. (bspline basis and number of basis functions and harmonics.) As we chose a high number for both to capture all features of the data set. As we explain the approximations for the estimation, this question could come up, so we have to make it specific. Adressed

$$\begin{aligned}\hat{\xi}_i^{m,\mathcal{H}} &= \int_0^1 (X_i^{\mathcal{H}}(s) - \hat{\mu}^{\mathcal{T}}(s)) \hat{\nu}^{m,\mathcal{T}}(s) ds = \int_0^1 \left( \sum_{j \in \mathcal{I}} a_{i,j}^{\mathcal{H}} \phi_j(s) \right) \left( \sum_{k \in \mathcal{L}} b_k^{m,\mathcal{T}} \psi_k(s) \right) ds \\ &= \sum_{j \in \mathcal{I}} \left[ a_{i,j}^{\mathcal{H}} \sum_{k \in \mathcal{L}} b_k^{m,\mathcal{T}} \int_0^1 \phi_j(s) \psi_k(s) ds \right]\end{aligned}\tag{54}$$

with truncation parameter  $m$ .

The simulation was done with R (version...). In total, 5000 repetitions were done for each set of simulations.

## 4.4 Results

### 4.4.1 Basis Expansion Regression

The following results origin from the **Estimation using Basis-Representation**, in which we transform the observed functions to perform regression of a scalar on a countable sequence of regressors, which is then tractable with typical multivariate regression theory.

**Monomial Basis** Due to the in the number of basis functions increasing collinearity of these basis(How about stressing the higher order? Are they already highly correlated at lower orders?)Adressed , simulations were conducted up until the sixth monomial basis, excluding the first one since this just represents a constant. For reasons outlined in Link chapter, they seem suited for  $f_1$  and show a better performance than bsplines for this coefficient function, but show the weakest performance for the coefficient function  $f_2$ . For  $f_1$ , the simulation selects 5(3) and for  $f_2$  5(5) monomial basis functions for the high(low) signal-to-noise ratios. This weakness is especially pronounced in the setup  $f_2, Y_1$ ; for the high noise setup,  $f_2, Y_2$  this weakness is still visible, but less strongly pronounced.

**Bspline Basis** Simulations with bspline basis functions were possible from 5 to 18, since from 18 onwards the simulations were running into problems concerning what happend here exactly? The smooth function  $f_1$  requires only 5(4) Bspline basis functions, for the high(low) signal-to-noise ratio to obtain the best fit for the bspline basis, which performs the worst for  $f_1$ . For  $f_2$ , 11(6) bspline basis functions are needed for the high(low) signal-to-noise ratios. For  $f_2$ , the bspline basis functions outperform the monomial basis, but come second to the fourier basis.

**Fourier Basis** For  $f_1$ , the simulation chooses a smaller number of fourier basis functions, 5(3) and a higher number for the bumpy function  $f_2$ , 9(7) for the setup with the high(low) signal-to-noise ratio. With the low signal-to-noise ratio, the simulation chooses a smaller number to account for the higher noise in the response variable. The fourier basis functions performs the best in each setup for the basis expansion regression. When recalling  $f_1$  and  $f_2$ , both, but especially  $f_1$  shows similar curvature order across the domain which contributes to the strong performance of this basis.

For all specifications, the effect of the Bias-Variance tradeoff can be observed:: The function  $f_1$  is rather smooth and not much bias is introduced here when choosing a small number of basis functions since  $f_1$  lacks pronounced peculiarities. For  $f_2$ , higher number of basis functions are needed, resulting from the inherent peculiarities of  $f_2$ : We can choose a higher number of basis functions since the amount of bias is decreasing much faster in the number of basis functions than the variance is increasing, compared to  $f_1$ .

#### 4.4.2 Functional Principal Component Regression

The model which is used for the FPCR is described in **Estimation using Functional Principal Components**. Additionally to the selection of the truncation parameter  $K$ , the choice of the number of FPC's adds to the complexity of the model since the eigenvalues and -functions from the decomposition are influenced by the choice of  $K$  that was used to approximate the function.

**Monomial Basis** TBD

**Bspline Basis** For  $f_1, Y_1$ , the CV suggests a better model the more FPC's are used. While the number of basis functions stays at 5 for two and three used FPC, it increases to 6 basis functions for  $nharm = 4$ . In  $f_1, Y_1$ , the cross-validated MSE is only slightly affected by the number of FPC's, but lowest for  $nharm = 3$ , indicating the possibility of over-fitting for four FPC's.  $f_2, Y_1$  shows the strongest decrease of MSE in increasing FPC's of all FPCR models fitted with BSplines (0.69614, 0.42504, 0.09497 for 2,3,4 FPC's). For  $f_2, Y_2$  the same holds true as in the setup with  $f_2$  and the high signal-to-noise ration, but less pronounced. Similar to the two setups with  $f_1$ , the cross-validated number of basis functions is increasing in the number of FPC's (4,6,8 for 2,3,4 FPC's)

The results overall might indicate that a higher number of FPC's has a strong effect in the setup with high signal-to-noise ratio, while not reducing as much error in the setting with the higher noise.

**Fourier Basis** We observe similar behaviour as for the FPCR using BSplines. In  $f_1, Y_1$ , the MSE is decreasing in the number of FPC's while the results in  $f_1, Y_2$  might again indicate over fitting when using 4 FPC's. Since it was not possible to calculate the cross-validated MSE for 3 basis functions for 4 FPC's, the increasing number of basis functions as seen in the FPCR with BSplines can not be observed here.  $f_2, Y_1$  displays again the greatest decline of MSE in the increasing number of FPC's. The same behaviour can be seen in  $f_2, Y_2$ , but less pronounced. Both configurations of  $f_2$  are using 5,15,7 basis functions for 2,3,4 FPC's: This behaviour is opposing to the previous reported increasing nature of basis functions in the number of FPC's.

#### 4.4.3 Interpretation and Relevance for Application

## 5 Application

The application uses the methods from the previous sections to predict the octane ratings of the introduced gasoline dataset. The simulation study will contribute towards the reading of the results.

- incorporate results of simulation study: number of basis, components, etc...
- point out difficulties estimating sderror
- describe setup and results

### 5.1 Literature

- Carey et al. 2002

## 6 Outlook

### 6.1 Literature

- James, Wang, and Zhu 2009 (shape-restrictions)

## 7 Appendix

### 7.1 Near-infrared (NIR) Spectroscopy

NIR- spectroscopy is a spectroscopic method that uses the near-infrared region of the electromagnetic spectrum (From 780 nm to 2500 nm). It therefore measures the absorption and interaction of this spectrum of radiation with the sample. NIR-spectroscopy is not only faster and cheaper than the standard test procedure – another big advantage is that it does not need a reagent and thus does not destroy the sample. It is used for analysis in different sectors and fields, like the agrochemical industry but also in healthcare. Its non-invasive nature makes it also to an asset for medical applications like the monitoring of diabetes in which NIR-spectroscopy is able to detect the worsening of the blood glucose metabolic dysfunction (cf . Li et al. 2020). In the context of this paper, the gasoline dataset which is used for the simulation and the application is constructed using NIR spectroscopy. According to Gy. Bohács, Z. Ovádi, A. Salgó 1998 NIR-spectroscopy is a feasible method for the analysis of gasoline since most of the absorption that is observed within the described interval of wavelengths is due to overtones and interactions of the radiation with chemical combinations (eg.: carbon–hydrogen, carbon–carbon, carbon–oxygen, carbonyl associated groups, aromatic stretching, and deformation vibration of the hydrocarbon molecules). While this paper focuses on the prediction of the octane number of gasoline, other research focuses on different properties of gasoline such as the olefin, naphtaenic and aromatic content (Parisi et al. 1990, as cited in Gy. Bohács, Z. Ovádi, A. Salgó 1998) or the distillation characktersitics (Pauls 1985, as cited in Gy. Bohács, Z. Ovádi, A. Salgó 1998)

## 7.2 Basis Plots

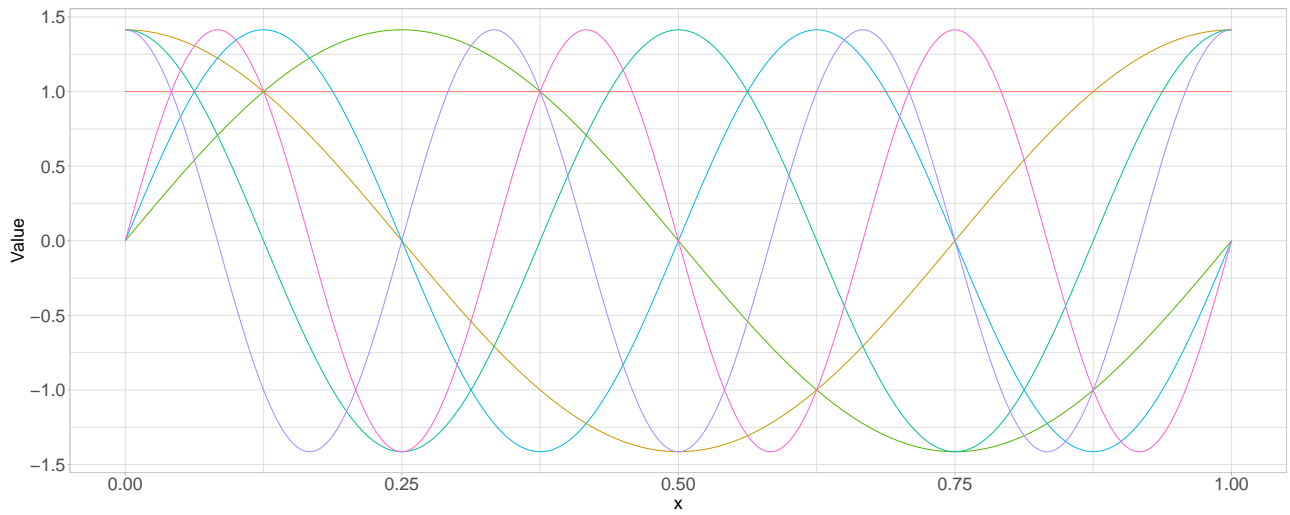


Figure 3: Fourier basis functions for  $i = 1, \dots, 7$

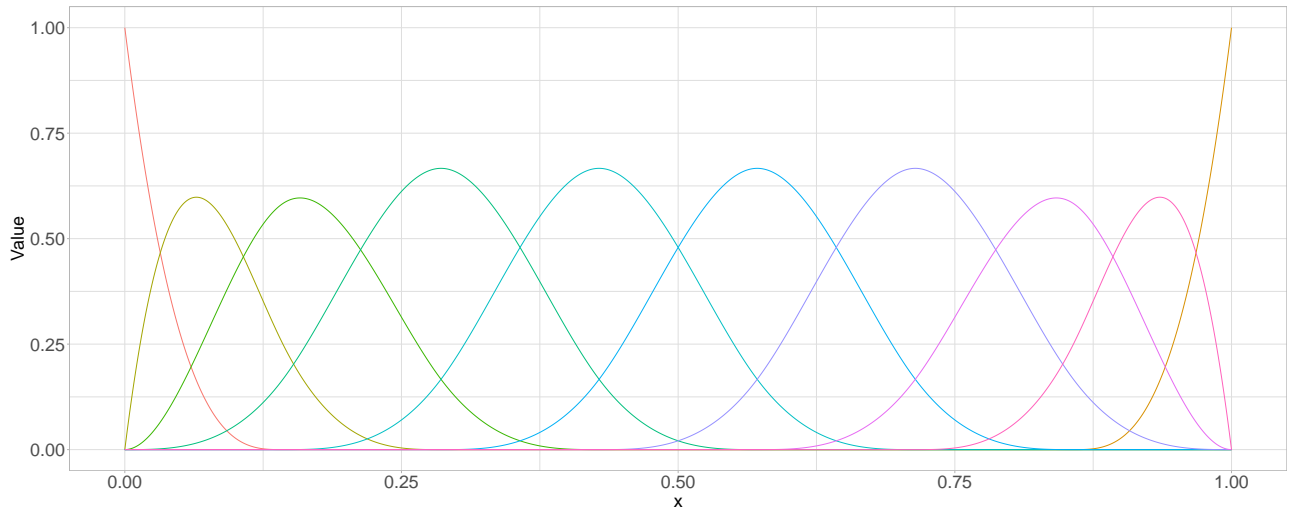


Figure 4: B-spline basis functions of order 4 for 8 equidistant knots on  $[0, 1]$

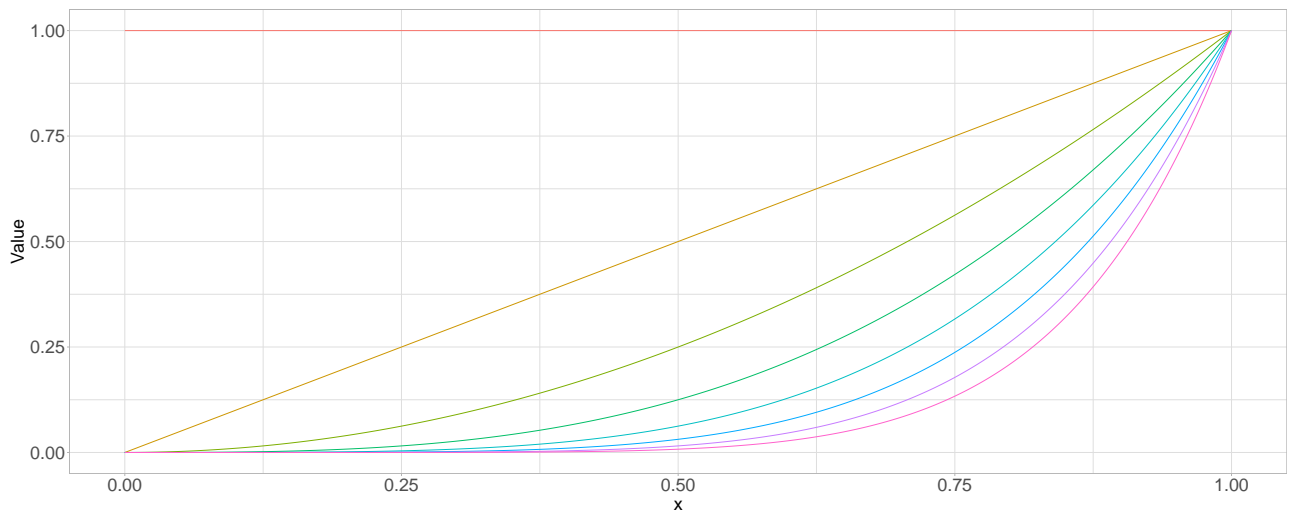


Figure 5: Monomial basis functions of degree 0 to 7

### 7.3 Simulation Study Results

Table 1: Monomial Basis Expansion Regression

$f_1, Y_1$	$f_1, Y_2$	$f_2, Y_1$	$f_2, Y_2$	n_basis
7.17723	131.92717	0.89367	2.58395	2
4.17039	129.64259	0.81224	2.51177	3
3.91275	130.11708	0.38079	2.09047	4
3.6385	130.59817	0.09217	1.81343	5
6.01644	201.15535	0.76208	3.39373	6

Table 2: BSpline Basis Expansion Regression

$f_1, Y_1$	$f_1, Y_2$	$f_2, Y_1$	$f_2, Y_2$	n_basis
3.91275	130.11708	0.38079	2.09047	4
3.64305	130.61095	0.09512	1.81643	5
3.65426	131.35355	0.07775	1.80913	6
3.67705	132.14205	0.07518	1.8168	7
3.71074	133.36537	0.0581	1.8165	8
3.71921	133.68149	0.0564	1.81931	9
3.74201	134.51217	0.05218	1.82576	10
3.7644	135.29727	0.0519	1.8361	11
3.80109	136.57581	0.05204	1.85315	12
3.83436	137.78407	0.05271	1.86997	13
3.86217	138.73431	0.05307	1.88269	14
3.86856	138.97427	0.0526	1.88482	15
3.88506	139.57419	0.05283	1.89335	16
3.91267	140.5149	0.05322	1.90619	17
3.94813	141.80885	0.05358	1.92312	18

Table 3: Fourier Basis Expansion Regression

$f_1, Y_1$	$f_1, Y_2$	$f_2, Y_1$	$f_2, Y_2$	n_basis
3.69752	129.19134	0.69524	2.3944	3
3.6347	130.59282	0.07418	1.79582	5
3.67623	132.08248	0.05147	1.79343	7
3.71885	133.67575	0.05105	1.81291	9
3.76451	135.26219	0.05146	1.83463	11
3.81095	136.90282	0.05197	1.85724	13
3.85674	138.57258	0.05252	1.88021	15
3.90619	140.29178	0.05304	1.90283	17
3.95517	142.05727	0.05365	1.92718	19

Table 4: BSpline FPCR,  $nharm = 2$ 

$f_1, Y_1$	$f_1, Y_2$	$f_2, Y_1$	$f_2, Y_2$	n_basis
10.93556	135.66335	0.69614	2.38597	4
6.01366	130.75986	0.77796	2.46801	5
6.36762	131.10562	0.79035	2.48029	6
6.76471	131.49507	0.7684	2.45834	7
7.21393	131.93622	0.7958	2.48568	8
7.77885	132.48984	0.75447	2.44437	9
8.48517	133.18046	0.71664	2.40654	10
8.95142	133.63714	0.70481	2.39475	11
9.21314	133.89421	0.71456	2.4045	12
9.29854	133.97794	0.73961	2.42952	13
9.31792	133.99732	0.74046	2.43037	14
9.32858	134.00809	0.74354	2.43344	15
9.27722	133.95811	0.77399	2.46383	16
9.23813	133.91976	0.77175	2.46159	17
9.33983	134.01705	0.76653	2.45632	18
9.24067	133.92097	0.78389	2.47367	19
9.34246	134.01837	0.75304	2.44285	20
9.38735	134.06288	0.76148	2.45127	21
9.35218	134.02839	0.75842	2.44823	22
9.50715	134.17875	0.74067	2.43049	23
9.51344	134.18666	0.7552	2.44502	24
9.52277	134.19463	0.75125	2.44107	25

Table 5: BSpline FPCR,  $nharm = 3$ 

$f_1, Y_1$	$f_1, Y_2$	$f_2, Y_1$	$f_2, Y_2$	n_basis
4.26915	129.73814	0.65618	2.35517	4
3.99153	129.47951	0.75902	2.45877	5
4.13265	129.64227	0.42504	2.126	6
4.15027	129.66158	0.52603	2.22684	7
4.28266	129.7862	0.71107	2.41134	8
4.30464	129.8084	0.65968	2.36019	9
4.32579	129.82882	0.60062	2.3013	10
4.29797	129.80089	0.65658	2.35703	11
4.32782	129.82922	0.68072	2.38108	12
4.34042	129.84041	0.71583	2.41605	13
4.3526	129.8512	0.72262	2.42279	14
4.3577	129.85642	0.72621	2.42637	15
4.37694	129.87493	0.75746	2.45749	16
4.35973	129.85893	0.75291	2.45296	17
4.37223	129.87141	0.69875	2.39878	18
4.34433	129.84395	0.74346	2.44353	19
4.36545	129.86595	0.67405	2.37412	20
4.36942	129.86886	0.67941	2.37944	21
4.34272	129.843	0.6919	2.39205	22
4.38437	129.88403	0.64702	2.34708	23
4.35012	129.8495	0.68435	2.38446	24
4.36866	129.86795	0.67512	2.37514	25



Table 6: BSpline FPCR,  $nharm = 4$ 

$f_1, Y_1$	$f_1, Y_2$	$f_2, Y_1$	$f_2, Y_2$	n_basis
3.91274	130.1171	0.38077	2.09045	4
3.97061	130.18108	0.75664	2.46627	5
3.83646	130.05682	0.19244	1.90433	6
3.8814	130.0934	0.3326	2.04419	7
4.2679	130.49905	0.09497	1.8071	8
4.31978	130.54252	0.15828	1.87039	9
4.32413	130.5527	0.10538	1.81745	10
4.31263	130.53896	0.13116	1.84322	11
4.33997	130.56608	0.10762	1.81971	12
4.34107	130.56642	0.11868	1.83076	13
4.35627	130.57893	0.12832	1.84037	14
4.35369	130.57747	0.11173	1.82379	15
4.35313	130.57695	0.11652	1.82861	16
4.3476	130.57119	0.11283	1.82497	17
4.31541	130.54208	0.08832	1.80044	18
4.32121	130.54666	0.10087	1.81304	19
4.2968	130.52361	0.08341	1.79552	20
4.30044	130.52787	0.08785	1.79994	21
4.29859	130.52514	0.08833	1.80045	22
4.28391	130.51071	0.08189	1.79394	23
4.29418	130.52054	0.08928	1.80134	24
4.28911	130.51512	0.0856	1.79766	25

Table 7: Fourier FPCR,  $nharm = 2$ 

$f_1, Y_1$	$f_1, Y_2$	$f_2, Y_1$	$f_2, Y_2$	n_basis
5.04859	129.81003	0.78889	2.4789	3
5.08647	129.84459	0.69756	2.38778	5
5.29235	130.04946	0.80393	2.49395	7
5.32414	130.07859	0.80074	2.49072	9
5.34403	130.09777	0.79714	2.48713	11
5.43601	130.18816	0.8141	2.50404	13
5.51333	130.26036	0.78908	2.47902	15
5.70153	130.44275	0.79659	2.48646	17
5.87259	130.6101	0.77783	2.46771	19

Table 8: Fourier FPCR,  $nharm = 3$ 

$f_1, Y_1$	$f_1, Y_2$	$f_2, Y_1$	$f_2, Y_2$	n_basis
3.69752	129.19134	0.69524	2.3944	3
5.06895	130.54881	0.1376	1.83947	5
5.21758	130.69839	0.13994	1.84156	7
5.22636	130.70697	0.12715	1.82896	9
5.27531	130.74943	0.14282	1.84453	11
5.29036	130.76091	0.13007	1.83175	13
5.16722	130.64051	0.10365	1.80524	15
4.98841	130.46747	0.11563	1.81733	17
4.92246	130.40322	0.14766	1.84937	19

Table 9: Fourier FPCR,  $nharm = 4$

$f_1, Y_1$	$f_1, Y_2$	$f_2, Y_1$	$f_2, Y_2$	n_basis
3.61833	129.83284	0.07736	1.78876	5
3.66209	129.88835	0.06311	1.77442	7
3.68812	129.91351	0.0717	1.78321	9
4.10715	130.34635	0.08317	1.79518	11
4.17569	130.4147	0.08837	1.80046	13
4.21112	130.44645	0.09965	1.81164	15
4.2129	130.45091	0.08775	1.79975	17
4.21884	130.45501	0.09032	1.80229	19

## 8 Proofs

### 8.1 Lemma

The curves  $X(t) \in \mathbb{L}^2[0, 1]$  is expanded by the eigenfunctions  $\{\nu^m\}$  as (20). The coefficients  $\xi^m$  of basis functions satisfy the following properties:

1.  $\mathbb{E}[\xi^m(\omega)] = 0$
2.  $Cov(\xi^m(\omega), \xi^n(\omega)) = \delta^{m,n} \lambda^m$
3.  $Var(\xi^m(\omega)) = \lambda^m$

Remind that  $\delta^{m,n} = 0$  if  $m \neq n$ , otherwise 1.

*Proof.* Assume that  $F(t)$  is the centered process of  $X(t)$ , namely,  $F(t) = X(t) - \int_{\Omega} X(t) dP(\omega)$ . To obtain the first result, we can show that

$$\begin{aligned}
 \mathbb{E}[\xi^m] &= \mathbb{E} \left[ \int_0^1 F(t) \nu_j(t) dt \right] \\
 &= \int_{\Omega} \int_0^1 F(t) \nu^m(t) dt dP(\omega) \\
 &= \int_0^1 \int_{\Omega} F(t) \nu^m(t) dP(\omega) dt \quad (\text{Fubini}) \\
 &= \int_0^1 \int_{\Omega} F(t) dP(\omega) \nu^m(t) dt \\
 &= \int_0^1 \mathbb{E}[F(t)] \nu^m(t) dt = 0
 \end{aligned} \tag{55}$$

where  $\mathbb{E}[F(t)]$  is 0 since  $F(t)$  is a centered process. The second claim is proved as:

$$\begin{aligned}
 \mathbb{E}[\xi^m \xi^n] &= \mathbb{E} \left[ \int_0^1 F(s) \nu^m(s) ds \int_0^1 F(t) \nu^n(t) dt \right] \\
 &= \mathbb{E} \left[ \int_0^1 \int_0^1 F(s) \nu^m(s) F(t) \nu^n(t) ds dt \right] \quad (\text{Fubini}) \\
 &= \int_0^1 \int_0^1 \mathbb{E}[F(s) F(t)] \nu^m(s) \nu^n(t) ds dt \\
 &= \int_0^1 \left( \int_0^1 c(s, t) \nu^m(s) ds \right) \nu^n(t) dt \\
 &= \int_0^1 [K \nu^m](t) \nu^n(t) dt \\
 &= \langle K \nu^m, \nu^n \rangle \\
 &= \langle \lambda^m \nu^m, \nu^n \rangle = \delta^{m,n} \lambda^m
 \end{aligned} \tag{56}$$

where  $\delta^{m,n} = 1$  if  $m = n$ , otherwise 0. The result is produced from orthonormality of the eigenfunctions.

$$Cov(\xi^m, \xi^n) = \mathbb{E}[\xi^m \xi^n] - \mathbb{E}[\xi^m] \mathbb{E}[\xi^n] = \delta^{m,n} \lambda^m \tag{57}$$

where  $\mathbb{E}[\xi^m] = \mathbb{E}[\xi^n] = 0$  as the first property. The last assertion is confirmed from the above two properties.

$$Var[\xi^m] = \mathbb{E}[(\xi^m - \mathbb{E}[\xi^m])^2] = \mathbb{E}[(\xi^m)^2] = \lambda^m \tag{58}$$

The original process also has the same properties as the centered one since

$$X(t) = F(t) + \mathbb{E}[X(t)] = \mu(t) + \sum_{m=1}^{\infty} \xi^m \nu^m(t) \quad (59)$$

□

## 8.2 Theorem (Karhunen-Loève expansion)

Let  $X : [0, 1] \rightarrow \mathbb{R}$  be a mean-square continuous stochastic process with  $X \in \mathbb{L}^2[0, 1]$ . Then there exists a basis  $\xi^m$  of  $\mathbb{L}^2[0, 1]$  such that for all  $t \in [0, 1]$ ,

$$X(t) = \mu(t) + \sum_{m=1}^{\infty} \xi^m \nu^m(t), \quad (60)$$

where  $\mu(t)$  is the mean function of  $X(t)$  and coefficients  $\xi^m$  are given by  $\int_0^1 (X(t) - \mu(t)) \nu^m(t) dt$ . These coefficients satisfy the following conditions.

1.  $\mathbb{E}[\xi^m(\omega)] = 0$
2.  $Cov(\xi^m(\omega), \xi^n(\omega)) = \delta^{m,n} \lambda^m$
3.  $Var(\xi^m(\omega)) = \lambda^m$

*Proof.* We know that  $K$  has a complete set of eigenfunctions  $\nu^m$  in  $\mathbb{L}^2[0, 1]$  and non-negative eigenvalues  $\lambda^m$ . With the reminder that  $\xi^m$  satisfy the three conclusions by Lemma 8.1, we prove this expansion by considering

$$\epsilon_N(t) := \mathbb{E} \left[ \left( X(t) - \mu(t) - \sum_{m=1}^N \xi^m \nu^m(t) \right)^2 \right] = \mathbb{E} \left[ \left( F(t) - \sum_{m=1}^N \xi^m \nu^m(t) \right)^2 \right] \quad (61)$$

where  $F(t)$  is the centered process of  $X(t)$ . Once it is shown that  $\lim_{N \rightarrow \infty} \epsilon_N(t) = 0$  uniformly in  $[0, 1]$ , the proof is completed.

$$\begin{aligned} \epsilon_N(t) &= \mathbb{E} \left[ \left( F(t) - \sum_{m=1}^N \xi^m \nu^m(t) \right)^2 \right] \\ &= \mathbb{E}[F(t)^2] - 2\mathbb{E} \left[ F(t) \sum_{m=1}^N \xi^m \nu^m(t) \right] + \mathbb{E} \left[ \sum_{m=1}^N \sum_{n=1}^N \xi^m \xi^n \nu^m(t) \nu^n(t) \right] \end{aligned} \quad (62)$$

Here,  $\mathbb{E}[F(t)^2] = c(t, t)$  as in (18) since  $F(t)$  is the centered process. Now, take the second term

$$\begin{aligned}
\mathbb{E}\left[F(t) \sum_{m=1}^N \xi^m \nu^m(t)\right] &= \mathbb{E}\left[F(t) \sum_{m=1}^N \left(\int_0^1 F(s) \nu^m(s) ds\right) \nu^m(t)\right] \\
&= \mathbb{E}\left[\sum_{m=1}^N \left(\int_0^1 F(t) F(s) \nu^m(s) ds\right) \nu^m(t)\right] \\
&= \sum_{m=1}^N \left(\int_0^1 \mathbb{E}[F(t) F(s)] \nu^m(s) ds\right) \nu^m(t) \\
&= \sum_{m=1}^N \left(\int_0^1 c(t, s) \nu^m(s) ds\right) \nu^m(t) \\
&= \sum_{m=1}^N [K \nu^m](t) \nu^m(t) \\
&= \sum_{m=1}^N \lambda^m \nu^m(t) \nu^m(t) = \sum_{m=1}^N \lambda^m \nu^m(t)^2
\end{aligned} \tag{63}$$

where the covariance function  $c(t, s)$  has the Hilbert-Schmidt operator as in (19). It turns out the product of the eigenfunction and the corresponding eigenvalue. For the last term, we derive from (56) that

$$\begin{aligned}
\mathbb{E}\left[\sum_{m=1}^N \sum_{n=1}^N \xi^m \xi^n \nu^m(t) \nu^n(t)\right] &= \sum_{m=1}^N \sum_{n=1}^N \mathbb{E}[\xi^m \xi^n] \nu^m(t) \nu^n(t) \\
&= \sum_{m=1}^N \sum_{n=1}^N \delta_{m,n} \lambda^m \nu^m(t) \nu^n(t) = \sum_{m=1}^N \lambda^m \nu^m(t)^2
\end{aligned} \tag{64}$$

where  $\delta_{m,n} = 1$  if  $m = n$ , otherwise 0. Therefore, by (62), (63), and (64) we obtain

$$\epsilon_N(t) = c(t, t) - \sum_{m=1}^N \lambda^m \nu^m(t) \nu^m(t) \tag{65}$$

implementing Mercer's Theorem this proof is concluded by

$$\lim_{N \rightarrow \infty} \epsilon_N(t) = \lim_{n \rightarrow \infty} \mathbb{E}\left[\left(F(t) - \sum_{m=1}^N \xi^m \nu^m(t)\right)^2\right] = 0 \tag{66}$$

□

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## 10 Affidavit

"I hereby confirm that the work presented has been performed and interpreted solely by myself except for where I explicitly identified the contrary. I assure that this work has not been presented in any other form for the fulfillment of any other degree or qualification. Ideas taken from other works in letter and in spirit are identified in every single case."

Bonn, 11.02.2021 \_\_\_\_\_  
Jonghun Baek

Bonn, 11.02.2021 \_\_\_\_\_  
Jakob R. Juergens

Bonn, 11.02.2021 \_\_\_\_\_  
Jonathan Willnow