# Confidence bands for multivariate and time dependent inverse regression models 

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#### Abstract

Uniform asymptotic confidence bands for a multivariate regression function in an inverse regression model with a convolution-type operator are constructed. The results are derived using strong approximation methods and a limit theorem for the supremum of a stationary Gaussian field over an increasing system of sets. As a particular application, asymptotic confidence bands for a time dependent regression function $f_{t}(x)$ $\left(x \in \mathbb{R}^{d}, t \in \mathbb{R}\right)$ in a convolution-type inverse regression model are obtained. Finally, we demonstrate the practical feasibility of our proposed methods in a simulation study and an application to the estimation of the luminosity profile of the elliptical galaxy NGC5017. To the best knowledge of the authors, the results presented in this paper are the first which provide uniform confidence bands for multivariate nonparametric function estimation in inverse problems.


Keywords: confidence bands; deconvolution; inverse problems; multivariate regression; nonparametric regression; rates of convergence; time dependent regression function; uniform convergence

## 1. Introduction

### 1.1. Inverse regression models

In many applications, it is impossible to observe a certain quantity of interest because only indirect observations are available for statistical inference. Problems of this type are called inverse problems and arise in many fields such as medical imaging, physics and biology. Mathematically the connection between the quantity of interest and the observable one can often be expressed in terms of a linear operator equation. Well-known examples are Positron Emission Tomography, which involves the Radon Transform (Cavalier [9]), the heat equation (Mair and Ruymgaart [22]), the Laplace Transform (Saitoh [29]) and the reconstruction of astronomical and biological images from telescopic and microscopic imaging devices, which is closely connected to convolution-type operators (Adorf [1], Bertero et al. [2]).

Inverse problems have been studied intensively in a deterministic framework and in mathematical physics. See, for example, Engl et al. [13] for an overview of existing methods in numerical analysis of inverse problems or Saitoh [29] for techniques based on reproducing kernel Hilbert spaces. Recently, the investigation of inverse problems has also become of importance from a statistical point of view. Here, a particularly interesting and active field of research is the construction of statistical inference methods such as hypothesis tests or confidence regions.

In this paper, we are interested in the convolution type inverse regression model

$$
\begin{equation*}
Y=(f * \psi)(x)+\varepsilon, \tag{1.1}
\end{equation*}
$$

where $\varepsilon$ is a random error, the operation $*$ denotes convolution, $\psi$ is a given square integrable function and the object of interest is the function $f$ itself. An important and interesting application of the inverse regression model (1.1) is the recovery of images from imaging devices such as astronomical telescopes or fluorescence microscopes in biology. In these cases, the observed, uncorrected image is always at least slightly blurry due to the physical characteristics of the propagation of light at surfaces of mirrors and lenses in the telescope. In this application, the variable $x$ represents the pixel of a CCD and we can only observe a blurred version of the true image modeled by the function $f$. In the corresponding mathematical model, the observed image is (at least approximately) a convolution of the real image with the so-called point spread function $\psi$, that is, an inverse problem with convolution operator.

The inference problem regarding the function $f$ is called inverse problem with stochastic noise. In recent years, the problem of estimating the regression function $f$ has become an important field of research, where the main focus is on a one dimensional predictor. Several authors propose Bayesian methods (Bertero et al. [2], Kaipio and Somersalo [19]) and construct estimators using tools from nonparametric curve estimation (Mair and Ruymgaart [22], Cavalier [10], Bissantz et al. [8]). Further inference methods, in particular the construction of confidence intervals and confidence bands, are much less developed. Birke et al. [5] have constructed uniform confidence bands for the function $f$ with a one-dimensional predictor.

The present work is motivated by the fact that in many applications one has to deal with an at least two-dimensional predictor. A typical example is image reconstruction since a picture is a two-dimensional object. Also in addition to the spatial dimensions, the data often show a dynamical behavior, thus repeated measurements at different times can be used to extend the statistical inference. For example, in astrophysics spectra of different objects like supernovae or variable stars show changes in time on observable timescales. In this case, the function $f$ depends on a further parameter, say $f_{t}$ and the reconstruction problem refers to a multivariate function even if the predictor is univariate.

The purpose of the present paper is the investigation of asymptotic properties of estimators for the function $f$ in model (1.1) with a multivariate predictor. In particular, we present a result on the weak convergence of the sup-norm of an appropriately centered estimate, which can be used to construct asymptotic confidence bands for the regression function $f$. In contrast to other authors (e.g., Cavalier and Tsybakov [11]), we do not assume that the function $\psi$ in model (1.1) is periodic, because in the reconstruction of astronomical or biological images from telescopes or microscopic imaging devices this assumption is often unrealistic.

### 1.2. Confidence bands

In a pioneering work, Bickel and Rosenblatt [4] extended results of Smirnov [30] for a histogram estimate and constructed confidence bands for a density function of independent, identically distributed (i.i.d.) observations. Their method is based on the asymptotic distribution of the supremum of a centered kernel density estimator. Since then, their method has been further developed both in the context of density and regression estimation. For density estimation, Neumann [23] derived bootstrap confidence bands, and Giné and Nickl [16] derived adaptive asymptotic bands
over generic sets. In a regression context, asymptotic confidence bands were constructed by Eubank and Speckman [14] for the Nadaraya-Watson estimator and by Xia [32] for a local polynomial estimator. Bootstrap confidence bands for nonparametric regression were proposed by Hall [18], Neumann and Polzehl [24] and by Claeskens and van Keilegom [12]. For the statistical inverse problem of deconvolution density estimation, Bissantz et al. [7] constructed asymptotic and bootstrap confidence bands, where Lounici and Nickl [21] obtained non-asymptotic confidence bands by using concentration inequalities. Recently, Birke et al. [5] provided uniform asymptotic and bootstrap confidence bands for a spectral cut-off estimator in the one-dimensional indirect regression model with convolution operator.

All these results are limited to the estimation of univariate densities and regression functions, and are not applicable in cases where the quantity of interest depends on a multivariate predictor. In such cases, to the best knowledge of the authors, confidence bands are not available. One reason for this gap is that a well-established way to construct asymptotic uniform confidence bands, which uses a pioneering result of Bickel and Rosenblatt [4] as the standard tool, cannot be extended in a straightforward manner to the multivariate case. There are substantial differences between the multivariate and one-dimensional case, and for multivariate inverse problems the mathematical construction of confidence bands requires different and/or extended methodology.

In the present paper, we will consider the problem of constructing confidence bands for the regression function in an inverse regression model with a convolution-type operator with a multivariate predictor. The estimators and assumptions for our asymptotic theory are presented in Section 2, while Section 3 contains the main results of the paper. In Section 4, we consider the special case of time dependent regression functions with a univariate predictor, which originally motivated our investigations. In Section 5, the finite-sample properties of the proposed asymptotic confidence bands are illustrated by means of a small simulation study and an application to HST data is discussed. The arguments of Sections 6 and 7, which contain all technical details of the proofs, are based on results by Piterbarg [27] who provided a limit theorem for the supremum

$$
\sup _{t \in T_{n}} X(t)
$$

of a stationary Gaussian field $\left\{X(t) \mid t \in \mathbb{R}^{d}\right\}$, where $\left\{T_{n} \subset \mathbb{R}^{d}\right\}_{n \in \mathbb{N}}$ is an increasing system of sets such that $\lambda^{d}\left(T_{n}\right) \rightarrow \infty$ as $n \rightarrow \infty$. This result generalized the multivariate extension in Bickel and Rosenblatt [3], who provided a limit theorem for the supremum $\sup _{t \in[0, T]^{d}} X(t)$, as $T \rightarrow \infty$.

## 2. Notation and assumptions

### 2.1. Model and notations

Suppose that $(2 n+1)^{d}$ observations $\left(x_{\mathbf{k}}, Y_{\mathbf{k}}\right), \mathbf{k}=\left(k_{1}, \ldots, k_{d}\right) \in G_{\mathbf{n}}:=\{-n, \ldots, n\}^{d}$ from the model

$$
\begin{equation*}
Y_{\mathbf{k}}=g\left(x_{\mathbf{k}}\right)+\varepsilon_{\mathbf{k}}:=(f * \psi)\left(x_{\mathbf{k}}\right)+\varepsilon_{\mathbf{k}} \tag{2.1}
\end{equation*}
$$

are available, where the function $f: \mathbb{R}^{d} \rightarrow \mathbb{R}$ is unknown, $\psi: \mathbb{R}^{d} \rightarrow \mathbb{R}$ is a known function and $g:=f * \psi$ denotes the convolution of $f$ and $\psi$, that is

$$
\begin{equation*}
g(x):=(f * \psi)(x):=\int_{\mathbb{R}^{d}} f(s) \psi(x-s) \mathrm{d} s \tag{2.2}
\end{equation*}
$$

The basic assumptions that guarantee the existence of the integral (2.2) and also assure $g \in$ $L^{2}\left(\mathbb{R}^{d}\right)$ is that $f \in L^{2}\left(\mathbb{R}^{d}\right)$ and $\psi \in L^{1}\left(\mathbb{R}^{d}\right) \cap L^{2}\left(\mathbb{R}^{d}\right)$, which will be assumed throughout this paper. In model (2.1), the predictors $x_{\mathbf{k}}:=\mathbf{k} \cdot \frac{1}{n a_{n}}$ are equally spaced fixed design points on a $d$-dimensional grid, with a sequence $\left(a_{n}\right)_{n \in \mathbb{N}}$ satisfying

$$
n a_{n} \rightarrow \infty \quad \text { and } \quad a_{n} \searrow 0 \quad \text { for } n \rightarrow \infty
$$

The noise terms $\left\{\varepsilon_{\mathbf{k}} \mid \mathbf{k} \in G_{\mathbf{n}}\right\}$ are a field of centered i.i.d. random variables with variance $\sigma^{2}:=$ $\mathbb{E}_{\mathbf{k}}^{2}>0$ and existing fourth moments. As a consequence of the convolution theorem and the formula for Fourier inversion, we obtain the representation

$$
\begin{equation*}
f(x)=\frac{1}{(2 \pi)^{d}} \int_{\mathbb{R}^{d}} \frac{\mathcal{F} g(\xi)}{\mathcal{F} \psi(\xi)} \exp \left(\mathrm{i} \xi^{T} x\right) \mathrm{d} \xi \tag{2.3}
\end{equation*}
$$

An estimator for the regression function $f$ can now easily be obtained by replacing the unknown quantity $\mathcal{F} g=\mathcal{F}(f * \psi)$ by an estimator $\mathcal{F} \hat{g}$. The random fluctuations in the estimator $\mathcal{F} \hat{g}$ cause instability of the ratio $\frac{\mathcal{F} \hat{g}(\xi)}{\mathcal{F} \psi(\xi)}$ if at least one of the components of $\xi$ is large. As a consequence, the problem at hand is ill-posed and requires regularization. We address this issue by excluding large values of $\xi_{j}$ for any $j=1, \ldots, d$ from the domain of integration, that is, we multiply the integrand in (2.3) with a sequence of Fourier transforms $\mathcal{F} \eta(h \cdot)$ of smooth functions with compact support $\left[-h^{-1}, h^{-1}\right]^{d}$. Here $h=h_{n}$ is a regularization parameter which corresponds to a bandwidth in nonparametric curve estimation and satisfies $h \rightarrow 0$ if $n \rightarrow \infty$. For the exact properties of the function $\eta$, we refer to Assumption A below.

An estimator $\hat{f}_{n}$ for the function $f$ in model (2.1) is now easily obtained as

$$
\begin{equation*}
\hat{f}_{n}(x)=\frac{1}{(2 \pi)^{d}} \int_{\mathbb{R}^{d}} \frac{\mathcal{F} \hat{g}(\xi)}{\mathcal{F} \psi(\xi)} \exp \left(\mathrm{i} \xi^{T} x\right) \mathcal{F} \eta(h \xi) \mathrm{d} \xi \tag{2.4}
\end{equation*}
$$

where

$$
\mathcal{F} \hat{g}(\xi)=\frac{1}{(2 \pi)^{d / 2} n^{d} a_{n}^{d}} \sum_{\mathbf{k} \in G_{n}} Y_{\mathbf{k}} \exp \left(-\mathrm{i} \xi^{T} x_{\mathbf{k}}\right)
$$

is the empirical analogue of the Fourier transform of $g$. Note that with the definition of the kernel

$$
\begin{equation*}
K_{n}(x)=\frac{1}{(2 \pi)^{d / 2}} \int_{\mathbb{R}^{d}} \frac{\mathcal{F} \eta(\xi)}{\mathcal{F} \psi(\xi / h)} \exp \left(\mathrm{i} \xi^{T} x\right) \mathrm{d} \xi \tag{2.5}
\end{equation*}
$$

the estimator (2.4) has the following representation

$$
\begin{equation*}
\hat{f}_{n}(x)=\frac{1}{(2 \pi)^{d} n^{d} a_{n}^{d} h^{d}} \sum_{\mathbf{k} \in G_{\mathbf{n}}} Y_{\mathbf{k}} K_{n}\left(\left(x-x_{\mathbf{k}}\right) \frac{1}{h}\right) \tag{2.6}
\end{equation*}
$$

Note that the kernel $K_{n}$ can be expressed as a Fourier transform as follows

$$
K_{n}=\overline{\mathcal{F}\left(\frac{\mathcal{F} \eta}{\mathcal{F} \psi(\cdot / h)}\right)}
$$

Also note that the kernel $K_{n}$ is a so-called deconvolution kernel. It is the analogue of a kernel in classical nonparametric kernel estimation with the difference that it depends on $n$ via the bandwidth $h$ in a rather complicated manner. For this reason, we use the notation $K_{n}$ instead of $K_{h}$ which corresponds to $K_{h}(\cdot / h)=\frac{1}{h^{d}} K(\cdot / h)$. Asymptotically, this kernel can be replaced by its limit $K$, see Assumption B, Remark 2 and Example 1 in the following discussion.

The first step of the proof of our main result (see Theorem 1 in Section 3) will consist of a uniform approximation of $\hat{f}_{n}(x)-\mathbb{E} \hat{f}_{n}(x)$ by an appropriate stationary Gaussian field. In the second step, we apply results of Piterbarg [27] and Bickel and Rosenblatt [3] to obtain the desired uniform convergence for the approximation process of the first step. Finally, these results are used to construct uniform confidence regions for $\mathbb{E} \hat{f}_{n}(x)$. Our approach is then based on undersmoothing: the choice of sufficiently small bandwidths assures the same limiting behaviour of $\hat{f}_{n}(x)-\mathbb{E} \hat{f}_{n}(x)$ and $\hat{f}_{n}(x)-f(x)$. This avoids the estimation of higher order derivatives, which often turns out to be difficult in applications. Thus, the limit theorem obtained in the second step will also provide uniform confidence regions for the function $f$ itself. Whereas undersmoothing implies that the rate-optimal bandwidth cannot be used, there has also been some theoretical justification why this choice of the regularization parameter is useful for constructing confidence intervals (see Hall [17]).

### 2.2. Assumptions

We now introduce the necessary assumptions which are required for the proofs of our main results in Section 3. The first assumption refers to the type of (inverse) deconvolution problem describing the shape of the kernel function $\eta$ in the spectral domain.

Assumption A. Let $\mathcal{F} \eta$ denote the Fourier transform of a function $\eta$ such that
A1. $\operatorname{supp}(\mathcal{F} \eta) \subset[-1,1]^{d}$.
A2. $\mathcal{F} \eta \in \mathscr{D}\left(\mathbb{R}^{d}\right)=\left\{f: \mathbb{R}^{d} \rightarrow \mathbb{R} \mid f \in C^{\infty}\left(\mathbb{R}^{d}\right), \operatorname{supp}(f) \subset \mathbb{R}^{d}\right.$ compact $\}$.
A3. There exists a constant $D>0$, such that $\mathcal{F} \eta(\xi)=1$ for all $\xi \in[-D, D]^{d}$ and $|\mathcal{F} \eta(\xi)| \leq$ 1 for all $\xi \in \mathbb{R}^{d}$.

## Remark 1.

1. The decay of the tails of the kernel $K_{n}$ is given in terms of the smoothness of the integrand in (2.5). The choice of a smooth regularizing function $\mathcal{F} \eta$ has the advantage that the smoothness of $1 / \mathcal{F} \psi$ carries over to $\mathcal{F} \eta(h \cdot) / \mathcal{F} \psi$.
2. Functions like $\mathcal{F} \eta$ are called bump functions. Their existence follows from the $C^{\infty}$ Urysohn lemma (see, e.g., Folland [15], Lemma 8.18).
3. Note that $\mathscr{D}\left(\mathbb{R}^{d}\right) \subset \mathscr{S}\left(\mathbb{R}^{d}\right)$, where $\mathscr{S}\left(\mathbb{R}^{d}\right)$ denotes the Schwartz space of smooth and rapidly decreasing functions. Since $\mathcal{F}: \mathscr{S}\left(\mathbb{R}^{d}\right) \rightarrow \mathscr{S}\left(\mathbb{R}^{d}\right)$ is a bijection (see, e.g., Folland [15], Corollary 8.28 ) we know that $\eta \in \mathscr{S}\left(\mathbb{R}^{d}\right)$ as well.
4. For the sake of transparency, we state the conditions and results with the same regularization parameter $h$ for each direction. In practical applications, this might not be the best strategy. The results presented in Sections 3 and 4 also hold for different sequences of bandwidths $h_{1}, \ldots, h_{d}$ as long as the system of rectangles $\left\{\left[0, h_{1}^{-1}\right] \times \cdots \times\left[0, h_{d}^{-1}\right] \mid n \in \mathbb{N}\right\}$ is a blowing up system of sets in the sense of Definition 14.1 in Piterbarg [27]. This is the case if the assumption

$$
\sum_{p=1}^{d}\left(\prod_{j=1, j \neq p}^{d} \frac{1}{h_{j}}\right) \leq L_{1} \cdot\left(\prod_{j=1}^{d} \frac{1}{h_{j}}\right)^{L_{2}}
$$

is satisfied for a constant $L_{1}$ that only depends on $d$ and a constant $L_{2}<1$. This condition is not a restriction in our setting because it holds whenever $h_{j} \cdot n^{\gamma_{j}} \rightarrow C_{j}$ for constants $C_{j}, \gamma_{j}>0, j=1, \ldots, d$.

In general, two kinds of convolution problems are distinguished in the literature, because the decay of the Fourier transform of the convolution function $\psi$ determines the degree of illposedness. In the case of an exponentially decreasing Fourier transform $\mathcal{F} \psi$ the problem is called severely ill-posed. In the present paper, the class of moderately ill-posed problems is considered, where the Fourier transform of the convolution function decays at a polynomial rate (the precise condition will be specified in Assumption B below). Throughout this paper

$$
\mathcal{W}^{m}\left(\mathbb{R}^{d}\right)=\left\{f \in L^{2}\left(\mathbb{R}^{d}\right) \mid \partial^{(\alpha)} f \in L^{2}\left(\mathbb{R}^{d}\right) \text { exists } \forall \alpha \in \mathbb{N}^{d},|\alpha| \leq m\right\}
$$

denotes the Sobolev space of order $m \in \mathbb{N}$, where $\partial^{(\alpha)} f$ is the weak derivative of $f$ of order $\alpha$. In the subsequent discussion, we will also make use of the Sobolev space for general $m>0$, which is defined by

$$
\mathcal{W}^{m}\left(\mathbb{R}^{d}\right)=\left\{f \in L^{2}\left(\mathbb{R}^{d}\right) \mid\left(1+|\xi|^{2}\right)^{m / 2} \mathcal{F} f \in L^{2}\left(\mathbb{R}^{d}\right)\right\}
$$

Assumption B. We assume the existence of a function $\Psi: \mathbb{R}^{d} \rightarrow \mathbb{R}$ such that the kernel $K=$ $\overline{\mathcal{F}(\Psi \cdot \mathcal{F} \eta)}$ satisfies

B1. $K \neq 0$ and there exist constants $\beta>d / 2, M \in \mathbb{N}$, indices $0<\mu_{1}<\mu_{2}<\cdots<\mu_{M}$ and $L^{2}$-functions $f_{1}, \ldots, f_{M-1}, f_{M}: \mathbb{R}^{d} \rightarrow \mathbb{R}$ with the property

$$
\xi^{\alpha} f_{p} \in \mathcal{W}^{m}\left(\mathbb{R}^{d}\right) \quad(p=1, \ldots, M-1)
$$

for all multi-indices $\alpha \in\{0, \ldots, d\}^{d},|\alpha| \leq d$ and all $m>\frac{d+|\alpha|}{2}$, such that

$$
\begin{equation*}
h^{\beta} K_{n}(x)-K(x)=\sum_{p=1}^{M-1} h^{\mu_{p}} \overline{\mathcal{F} f_{p}}(x)+h^{\mu_{M}} \overline{\mathcal{F} f_{n, M}}(x) \tag{2.7}
\end{equation*}
$$

where $f_{M}$ may depend on $n$, that is, $f_{M}=f_{M, n}$ and $\left\|f_{M, n}\right\|_{L^{1}\left(\mathbb{R}^{d}\right)}=\mathrm{O}(1)$.
B2. $\xi^{\alpha} \Psi \cdot \mathcal{F} \eta, \xi^{\alpha} \frac{h^{\beta}}{\mathcal{F} \psi(\cdot / h)} \cdot \mathcal{F} \eta \in \mathcal{W}^{m}\left(\mathbb{R}^{d}\right)$ for some $m>\frac{d+|\alpha|}{2}$.
B3. $\log (n) \cdot h^{\mu_{M}}\left(a_{n}^{-d / 2} h^{-d / 2}\right) \cdot\left\|f_{M}\right\|_{L^{1}\left(\mathbb{R}^{d}\right)}=\mathrm{o}(1)$ and $h^{\mu_{1}}(\log (n))^{2}=\mathrm{o}(1)$.
Remark 2. Assumption B1 implies $h^{\beta} K_{n} \rightarrow K$ in $L^{2}\left(\mathbb{R}^{d}\right)$ and also specifies the order of this convergence. It can be understood as follows. If the convergence of the difference $h^{\beta} K_{n}-K$ is fast enough, that is,

$$
\begin{equation*}
\log (n) \cdot h^{\mu_{1}}\left(a_{n} h\right)^{-d / 2}=\mathrm{o}(1) \tag{2.8}
\end{equation*}
$$

we have $M=1$. On the other hand, in some relevant situations (see Example 1(ii) below) the rate of convergence $h^{\mu_{1}}$ is given by $h^{2}$ for each $d$ and (2.8) cannot hold for $d \geq 4$. Here, the expansion (2.7) provides a structure, such that our main results remain correct although the rate of convergence is not very fast. We can decompose the difference $h^{\beta} K_{n}-K$ in two parts, where one part depends on $n$ only through the factors $h^{\mu_{p}}$ and the other part converges sufficiently fast (in some cases this term vanishes completely).

Example 1. This example illustrates the construction of the functions in the representation (2.7).
(i) Let $d=2$ and $\psi(x)=\frac{1}{4} \exp \left(-\left|x_{1}\right|\right) \exp \left(-\left|x_{2}\right|\right), x=\left(x_{1}, x_{2}\right)^{T}, \xi=\left(\xi_{1}, \xi_{2}\right)^{T}$. Then we have $\frac{h^{4}}{\mathcal{F} \psi(\xi / h)}=2 \pi\left(h^{4}+h^{2}\left(\xi_{1}^{2}+\xi_{2}^{2}\right)+\xi_{1}^{2} \xi_{2}^{2}\right)$, which implies $\beta=4, M=3$ and

$$
\begin{aligned}
h^{4} \cdot K_{n}(x) & =\int_{\mathbb{R}^{2}}\left(h^{4}+h^{2}\left(\xi_{1}^{2}+\xi_{2}^{2}\right)+\xi_{1}^{2} \xi_{2}^{2}\right) \mathcal{F} \eta(\xi) \exp \left(\mathrm{i} x^{T} \xi\right) \mathrm{d} \xi \\
K(x) & =\int_{\mathbb{R}^{2}} \mathcal{F} \eta(\xi) \xi_{1}^{2} \xi_{2}^{2} \exp \left(\mathrm{i} x^{T} \xi\right) \mathrm{d} \xi
\end{aligned}
$$

With the definitions $f_{1}(\xi)=2 \pi\left(\xi_{1}^{2}+\xi_{2}^{2}\right) \mathcal{F} \eta(\xi), f_{2}(\xi)=2 \pi \mathcal{F} \eta(\xi)$ and $f_{n, 3} \equiv 0$ we obtain

$$
h^{4} \cdot K_{n}(x)-K(x)=h^{2} \cdot \overline{\mathcal{F} f_{1}(\xi)}+h^{4} \cdot \overline{\mathcal{F} f_{2}(\xi)}
$$

In this example, the condition $\log (n) h^{2} / \sqrt{a_{n}^{d} h^{d}}=\mathrm{o}(1)$ is satisfied. However, the following results are valid if the weaker condition of a decomposition of the form (2.7) holds. Furthermore, since the factors of $\mathcal{F} \eta$ in $f_{1}$ and $f_{2}$ are polynomials, we have $\mathcal{F} f_{j}(\xi) \in \mathscr{S}\left(\mathbb{R}^{d}\right)$, which implies $\xi^{\alpha} f_{j} \in \mathcal{W}^{m}\left(\mathbb{R}^{d}\right)$ for all $\alpha$ and all $m \in \mathbb{N}$.
(ii) If $|x|=\sqrt{x_{1}^{2}+\cdots+x_{d}^{2}}$ and $\psi(x)=2^{-(d+1) / 2} \mathrm{e}^{-|x|}$ we have

$$
\mathcal{F} \psi(\xi)=\frac{1}{\sqrt{2 \pi}} \Gamma\left(\frac{d+1}{2}\right) \frac{1}{\left(1+|\xi|^{2}\right)^{(d+1) / 2}}
$$

(see Folland [15], Exercise 13). If $d$ is odd we use the identity

$$
\left(h^{2}+|\xi|^{2}\right)^{(d+1) / 2}=\sum_{j=0}^{(d+1) / 2}\binom{\frac{d+1}{2}}{j} h^{2 j}|\xi|^{(d+1) / 2-2 j}
$$

and an expansion of the form (2.7) is obvious from the definition of $K_{n}$ in (2.5). If the dimension $d$ is even the situation is more complicated. Consider for example the case $d=4$, where

$$
\frac{h^{5}}{\mathcal{F} \psi(\xi / h)} \rightarrow \frac{\sqrt{2 \pi}}{\Gamma(5 / 2)}|\xi|^{5}=\frac{\sqrt{2 \pi}}{\Gamma(5 / 2)} \sqrt{\left(\xi_{1}^{2}+\xi_{2}^{2}+\xi_{3}^{2}+\xi_{4}^{2}\right)^{5}} \quad \text { as } n \rightarrow \infty
$$

It follows that the constant $\beta$ and the functions $\Psi, K_{n}$ and $K$ from Assumption B are given by $\beta=d+1=5, \Psi(\xi)=\frac{\sqrt{2 \pi}}{\Gamma(5 / 2)}|\xi|^{5}$ and

$$
\begin{aligned}
h^{\beta} K_{n}(x) & =\frac{1}{(2 \pi)^{2}} \int_{\mathbb{R}^{d}} \frac{\sqrt{2 \pi}}{\Gamma(5 / 2)}\left(h^{2}+|\xi|^{2}\right)^{(d+1) / 2} \mathcal{F} \eta(\xi) \exp \left(\mathrm{i} \xi^{T} x\right) \mathrm{d} \xi \\
K(x) & =\frac{1}{(2 \pi)^{2}} \int_{\mathbb{R}^{d}} \frac{\sqrt{2 \pi}}{\Gamma(5 / 2)}|\xi|^{d+1} \mathcal{F} \eta(\xi) \exp \left(\mathrm{i} \xi^{T} x\right) \mathrm{d} \xi,
\end{aligned}
$$

respectively. In order to show that Assumption B1 holds in this case we use Taylor's theorem and obtain

$$
\frac{h^{5}}{\mathcal{F} \psi(\xi / h)}-\Psi(\xi)=\frac{\sqrt{2 \pi}}{\Gamma((d+1) / 2)}\left(h^{2} \cdot \frac{5}{2} \cdot|\xi|^{3}+h^{4} \cdot \frac{5}{2} \cdot \frac{3}{2} \cdot\left(|\xi|^{2}+\lambda_{d} h^{2}\right)^{1 / 2}\right)
$$

for some constant $\lambda_{d} \in[0,1)$. Recalling the definition of $K_{n}$ in (2.5) this gives

$$
\left(h^{\beta} K_{n}-K\right)(x)=h^{2} \overline{\mathcal{F} f_{1}(\xi)}+h^{4} \overline{\mathcal{F} f_{2, n}(\xi)},
$$

where the functions $f_{1}$ and $f_{2, n}$ are defined by

$$
\begin{aligned}
f_{1}(\xi) & =\frac{1}{(2 \pi)^{3 / 2} \Gamma(5 / 2)} \cdot|\xi|^{3} \cdot \frac{5}{2} \cdot \mathcal{F} \eta(\xi), \\
f_{2, n}(x) & =\frac{1}{(2 \pi)^{3 / 2} \Gamma(5 / 2)} \frac{5}{2} \cdot \frac{3}{2} h^{4}\left(|\xi|^{2}+\lambda_{d} h^{2}\right)^{1 / 2} \cdot \mathcal{F} \eta(\xi),
\end{aligned}
$$

respectively. It can be shown by a straightforward calculation that $\xi^{\alpha} f_{j} \in \mathcal{W}^{6+|\alpha|}\left(\mathbb{R}^{d}\right)$ for all $\alpha \in\{0, \ldots, d\}^{d}$.

Remark 3. In the one-dimensional regression model (2.1), Birke et al. [5] assume that the kernel $K$ has exponentially decreasing tails in order to obtain asymptotic confidence bands, which, in combination with the other assumptions only allows for kernels that are Fourier transforms of $C^{\infty}$-functions with square integrable derivatives. Our Assumption B is already satisfied if $K$ is the Fourier transform of a once weakly differentiable function with square integrable weak derivative, such that all indices of ill-posedness $\beta$ that satisfy $\beta>\frac{1}{2}$ are included if $d=1$. Moreover, the assumptions regarding the bandwidths are less restrictive compared to Birke et al. [5].

Our final assumptions refer to the smoothness of the function $f$ and to the decay of the convolution $f * \psi$.

Assumption C. We assume that
C1. There exist constants $\gamma>2, m>\gamma+\frac{d}{2}$ such that $f \in \mathcal{W}^{m}\left(\mathbb{R}^{d}\right)$.
C2. There exists a constant $v>0$ such that

$$
\int_{\mathbb{R}}|(f * \psi)(z)|^{2}\left(1+|z|^{2}\right)^{v} \mathrm{~d} z<\infty
$$

## 3. Asymptotic confidence regions

In this section, we construct asymptotic confidence regions for the function $f$ on the unit cube $[0,1]^{d}$. These results can easily be generalized to arbitrary rectangles $X_{j=1}^{d}\left[a_{j}, b_{j}\right]$ for fixed constants $a_{j}<b_{j}(j=1, \ldots, d)$ and the details are omitted for the sake of brevity. We investigate the limiting distribution of the supremum of the process $\left\{\tilde{Y}_{n}(x) \mid x \in[0,1]^{d}\right\}$, where

$$
\begin{align*}
\tilde{Y}_{n}(x) & =\frac{(2 \pi)^{d} h^{\beta} \sqrt{h^{d} n^{d} a_{n}^{d}}}{\sigma\|K\|_{L^{2}\left(\mathbb{R}^{d}\right)}}\left[\hat{f}_{n}(x)-\mathbb{E} \hat{f}_{n}(x)\right]  \tag{3.1}\\
& =\frac{(2 \pi)^{d} h^{\beta}}{\sigma\|K\|_{L^{2}\left(\mathbb{R}^{d}\right)} \sqrt{h^{d} n^{d} a_{n}^{d}}} \sum_{\mathbf{k} \in G_{n}} K_{n}\left(\left(x-x_{\mathbf{k}}\right) \frac{1}{h}\right) \varepsilon_{\mathbf{k}}
\end{align*}
$$

and the kernel $K_{n}$ is defined in (2.5). Note that

$$
\sup _{x \in[0,1]^{d}}\left|\tilde{Y}_{n}(x)\right|=\sup _{x \in\left[0, h^{-1}\right]^{d}}\left|Y_{n}(x)\right|,
$$

where the process

$$
\begin{equation*}
Y_{n}(x):=\frac{(2 \pi)^{d} h^{\beta}}{\sigma\|K\|_{L^{2}\left(\mathbb{R}^{d}\right)} \sqrt{h^{d} n^{d} a_{n}^{d}}} \sum_{\mathbf{k} \in G_{n}} K_{n}\left(x-x_{\mathbf{k}} \frac{1}{h}\right) \varepsilon_{\mathbf{k}} \tag{3.2}
\end{equation*}
$$

can be approximated by a stationary Gaussian field uniformly with respect to $\left[0, h^{-1}\right]^{d}$. Thus the desired limiting distribution corresponds to the limiting distribution of the supremum of a stationary Gaussian process over a system of increasing smooth sets with sufficient similarity of their speed of increase, and is therefore of Gumbel-type. The precise result is given in the following theorem.

Theorem 1. Assume that for some fixed constant $\delta \in(0,1], \delta<d$ and a constant $r>\frac{2 d}{d-\delta}$ the $r$ th moment of the errors exists, that is, $\mathbb{E}\left|\varepsilon_{\mathbf{k}}\right|^{r}<\infty$. If additionally Assumptions A and B are satisfied and $\frac{\log (n)}{n^{\delta} a_{h}^{\delta} h^{d}}=\mathrm{o}(1)$, then we have

$$
\lim _{n \rightarrow \infty} \mathbb{P}\left(\sup _{x \in[0,1]^{d}}\left(\left|\tilde{Y}_{n}(x)\right|-C_{n, 3}\right) \cdot C_{n, 3}<\kappa\right)=\mathrm{e}^{-2 \mathrm{e}^{-\kappa}}
$$

where

$$
\begin{aligned}
C_{1} & =\operatorname{det}\left(\left[\frac{(2 \pi)^{2 d}}{\|K\|_{2}^{2}} \int_{\mathbb{R}^{d}}|\Psi(v) \mathcal{F} \eta(v)|^{2} v_{i} v_{j} \mathrm{~d} v\right], i, j=1, \ldots, d\right), \\
C_{n, 2} & =\sqrt{\frac{C_{1}}{(2 \pi)^{d+1}}} \frac{1}{h^{d}}, \\
C_{n, 3} & =\sqrt{2 \ln \left(C_{n, 2}\right)}+\frac{(d-1) \ln \left(2 \ln \left(C_{n, 2}\right)\right)}{2 \sqrt{2 \ln \left(C_{n, 2}\right)}} .
\end{aligned}
$$

The proof of this result is long and complicated and therefore deferred to Sections 6 and 7. In the following, we apply Theorem 1 to construct uniform confidence regions for the function $f$ by choosing the bandwidth such that the bias decays to zero sufficiently fast. More precisely, if the condition

$$
\log (n) \sup _{x \in[0,1]^{d}}\left|f(x)-\mathbb{E} \hat{f}_{n}(x)\right|=\mathrm{o}\left(\left(h^{\beta} \sqrt{h^{d} n^{d} a_{n}^{d}}\right)^{-1}\right)
$$

is satisfied, it follows directly that the random quantities $\sup _{x \in[0,1]^{d}}\left|\tilde{Y}_{n}(x)\right|$ and

$$
\frac{(2 \pi)^{d} h^{\beta} \sqrt{h^{d} n^{d} a_{n}^{d}}}{\|K\|_{L^{2}\left(\mathbb{R}^{d}\right)}^{\sigma}} \sup _{x \in[0,1]^{d}}\left|f(x)-\hat{f}_{n}(x)\right|
$$

have the same limiting behavior.
Corollary 1. Assume that the conditions of Theorem 1, Assumption C and the condition

$$
\sqrt{h^{d} n^{d} a_{n}^{d}} \sqrt{\log (n)}\left(\frac{1}{n^{3} a_{n}^{3} h^{2}}+\frac{a_{n}^{\nu}}{n}+a_{n}^{\nu+d / 2}+h^{\gamma+\beta}\right)=\mathrm{o}(1) \quad \text { for } n \rightarrow \infty
$$

are satisfied. Then we have for any $\kappa \in \mathbb{R}$

$$
\lim _{n \rightarrow \infty} \mathbb{P}\left(\hat{f}_{n}(x)-\Phi_{n, \kappa} \leq f(x) \leq \hat{f}_{n}(x)+\Phi_{n, \kappa} \text { for all } x \in[0,1]^{d}\right)=\mathrm{e}^{-2 \mathrm{e}^{-\kappa}}
$$

where the sequence $\Phi_{n, \kappa}$ is defined by

$$
\Phi_{n, \kappa}=\frac{\left(\kappa / C_{n, 3}+C_{n, 3}\right) \sigma\|K\|_{L^{2}\left(\mathbb{R}^{d}\right)}}{(2 \pi)^{d} h^{\beta} \sqrt{h^{d} n^{d} a_{n}^{d}}} .
$$

As a consequence of Corollary 1 an asymptotic uniform confidence region for the function $f$ with confidence level $1-\alpha$ is given by

$$
\begin{equation*}
\left\{\left[\hat{f}_{n}(x)-\Phi_{n,-\ln (-0.5 \ln (1-\alpha))}, \hat{f}_{n}(x)+\Phi_{n,-\ln (-0.5 \ln (1-\alpha))}\right] \mid x \in[0,1]^{d}\right\} . \tag{3.3}
\end{equation*}
$$

 the ill-posedness of the inverse problem (see Assumption B). It does not appear in corresponding results for the direct regression case. On the other hand, the factor $a_{n}^{-d / 2}$ arises from the design on the growing system of sets $\left\{\left[-a_{n}^{-1}, a_{n}^{-1}\right]^{d} \mid n \in \mathbb{N}\right\}$. In the case of a regression on a fixed interval it does not appear as well. The width of the asymptotic point-wise confidence intervals in the multivariate indirect regression case as obtained in Bissantz and Birke [6] is of order $\frac{1}{h^{\beta} \sqrt{N h^{d} a_{n}^{d}}}$, where $N$ is the total number of observations. Their point-wise confidence intervals are smaller than the uniform ones obtained in Corollary 1. The price for uniformity is an additional factor of logarithmic order, which is typical for results of this kind.

In applications the standard deviation is unknown but can be estimated easily from the data, because this does not require the estimation of the function $f$. In particular, (3.3) remains an asymptotic $(1-\alpha)$-confidence band, if $\sigma$ is replaced by an estimator satisfying $\hat{\sigma}-\sigma=\mathrm{o}_{P}(1 / \log (n))$.

## 4. Time dependent regression functions

In this section, we extend model (2.1) to include a time dependent regression function, that is

$$
\begin{equation*}
Y_{j, \mathbf{k}, n}=\left(T_{\psi} f_{t_{j}}\right)\left(x_{\mathbf{k}}\right)+\varepsilon_{\mathbf{k}}, \quad \mathbf{k} \in G_{\mathbf{n}}, j=-m, \ldots, m, \tag{4.1}
\end{equation*}
$$

where $x_{\mathbf{k}}=\frac{\mathbf{k}}{n a_{n}}$ and $t_{j}=\frac{j}{m b_{m}}, m=m(n)$, such that $m(n) \rightarrow \infty$ and $b_{m(n)} \searrow 0$ as $n \rightarrow \infty$.
We assume that $\psi$ does not depend on the time and the operator $T_{\psi}$ is defined by

$$
\left(T_{\psi} f_{t}\right)=\int_{\mathbb{R}^{d}} f_{t}(y) \psi(\cdot-y) \mathrm{d} y
$$

This assumption is reasonable in the context of imaging where the function $\psi$ corresponds to the point spread function (Bertero et al. [2]). If it is not satisfied, that is, the convolution operator effects all coordinates, the problem can be modeled as in Section 2.

For a precise statement of the results, we will add an index to the Fourier operator $\mathcal{F}$ which gives the dimension of the space under consideration. We will write $\mathcal{F}_{d+1}$ if the Fourier transform is taken over the whole space $\mathbb{R}^{d+1}$ and $\mathcal{F}_{d}$ to denote Fourier transformation with respect to the spatial dimensions. By the same considerations as given in Section 2, we obtain an estimator $\check{f}$ for the function $f_{t}$

$$
\begin{aligned}
\check{f}_{n}(x ; t) & =\frac{1}{(2 \pi)^{(d+1) / 2}} \int_{\mathbb{R}^{d+1}} \frac{\mathcal{F}_{d+1}(\widehat{f * \psi})(\xi, \tau)}{(2 \pi)^{d / 2} \mathcal{F}_{d} \psi(\xi)} \mathcal{F}_{d} \check{\eta}\left(\xi h, \tau h_{t}\right) \exp \left(\mathrm{i} t \tau+\mathrm{i} x^{T} \xi\right) \mathrm{d}(\xi, \tau) \\
& =\frac{1}{(2 \pi)^{d+1 / 2} n^{d} m a_{n}^{d} b_{m}} \sum_{(\mathbf{k}, j) \in G_{(\mathbf{n}, m)}^{d+1}} Y_{\mathbf{k}, j} \check{K}_{n}\left(\frac{x-x_{\mathbf{k}}}{h}, \frac{t-t_{j}}{h_{t}}\right),
\end{aligned}
$$

where $G_{(n, m)}^{d+1}$ denotes the grid $\{-n, \ldots, n\}^{d} \times\{-m, \ldots, m\}$ and the kernel $\check{K}_{n}$ is given by

$$
\begin{equation*}
\check{K}_{n}(x ; t)=\frac{1}{(2 \pi)^{(d+1) / 2}} \int_{\mathbb{R}^{d+1}} \frac{\exp \left(\mathrm{i} \tau t+\mathrm{i} \xi^{T} x\right)}{\mathcal{F}_{d} \psi(\xi / h)} \mathcal{F}_{d+1} \check{\eta}(\xi, \tau) \mathrm{d}(\xi, \tau) . \tag{4.2}
\end{equation*}
$$

Here the function $\check{\eta}: \mathbb{R}^{d+1} \rightarrow \mathbb{R}$ satisfies condition A and $h_{t}=h_{t}(n)$ is an additional sequence of bandwidths referring to the time domain. For the asymptotic analysis, we require a modified version of Assumption B.

Assumption B̌. Let Assumptions B1 (with corresponding kernel Ǩ) and B2 hold and additionally assume that
¿̌3. $\log (n+m(n)) \cdot h^{\mu_{M}}\left(a_{n}^{-d / 2} h^{-d / 2} b_{m(n)}^{1 / 2} m(n)^{1 / 2}\right)=\mathrm{o}(1)$ and for $p=1, \ldots, M-1$

$$
h^{\mu_{p}}(\log (n+m))^{2}=\mathrm{o}(1)
$$

Theorem 2. Define

$$
\check{Y}_{n}(x ; t):=\frac{(2 \pi)^{d+1} h^{\beta} \sqrt{h^{d} h_{t} n^{d} m b_{m} a_{n}^{d}}}{\sigma\|\check{K}\|_{L^{2}\left(\mathbb{R}^{d+1}\right)} \|}\left[\check{f}_{n}(x ; t)-\mathbb{E} \check{f}_{n}(x ; t)\right]
$$

and let the moment condition of Theorem 1 and Assumptions A and $\overline{\mathrm{B}}$ be satisfied. We further assume that the bandwidths $h_{t}$ and $h$, and the sequences $\left(a_{n}\right)_{n \in \mathbb{N}}$ and $\left(b_{m(n)}\right)_{n \in \mathbb{N}}$ satisfy

$$
\begin{gathered}
\log (n+m)\left(\sqrt{\frac{n a_{n}}{m b_{m}}} \frac{1}{\sqrt{n^{\delta} h_{t} a_{n}^{\delta} h^{d}}}+\left(\frac{m b_{m}}{n a_{n}}\right)^{d / 2} \frac{1}{\sqrt{m^{\delta} h_{t} h^{d}}}\right)=\mathrm{o}(1) \quad \text { for } n \rightarrow \infty \\
h_{t}+h \leq L_{1} \cdot h^{d\left(1-L_{2}\right)} h_{t}^{\left(1-L_{2}\right)}
\end{gathered}
$$

for some constants $L_{1}<\infty$ and $L_{2} \in(0,1)$. Then we have for each $\kappa \in \mathbb{R}$,

$$
\lim _{n \rightarrow \infty} \mathbb{P}\left(\sup _{x \in[0,1]^{d}}\left(\left|\check{Y}_{n}(x ; t)\right|-D_{n, 3}\right) \cdot D_{n, 3}<\kappa\right)=\mathrm{e}^{-2 \mathrm{e}^{-\kappa}}
$$

where

$$
\begin{aligned}
D_{1} & =\operatorname{det}\left(\left[\frac{(2 \pi)^{2(d+1)}}{\|\check{K}\|_{L^{2}\left(\mathbb{R}^{d+1}\right)}^{2}} \int_{\mathbb{R}^{d+1}}\left|\Psi\left(v_{1}, \ldots, v_{d}\right) \mathcal{F}_{d+1} \check{\eta}(v)\right|^{2} v_{i} v_{j} \mathrm{~d} v\right], i, j=1, \ldots, d+1\right) \\
D_{n, 2} & =\sqrt{\frac{D_{1}}{(2 \pi)^{d+2}}} \frac{1}{h^{d} h_{t}} \quad \text { and } \\
D_{n, 3} & =\sqrt{2 \ln \left(D_{n, 2}\right)}+\frac{(d-1) \ln \left(2 \ln \left(D_{n, 2}\right)\right)}{2 \sqrt{2 \ln \left(D_{n, 2}\right)}} .
\end{aligned}
$$

Corollary 2. If the assumptions of Theorem 2 are satisfied, the limit kernel $\check{K}$ is defined by

$$
\begin{equation*}
\check{K}(x, t)=\frac{1}{(2 \pi)^{(d+1) / 2}} \int_{\mathbb{R}^{d+1}} \Psi(\xi) \mathcal{F}_{d+1} \check{\eta}(\xi, \tau) \exp \left(\mathrm{i} \xi^{T} x+\mathrm{i} \tau t\right) \mathrm{d}(\xi, \tau) \tag{4.3}
\end{equation*}
$$

and the function $f_{(\cdot)}(\cdot): \mathbb{R}^{d+1} \rightarrow \mathbb{R}^{1}$ satisfies Assumption C , then it follows that

$$
\lim _{n \rightarrow \infty} \mathbb{P}\left(\check{f}_{n}(x ; t)-\check{\Phi}_{n, \kappa} \leq f(x ; t) \leq \check{f}_{n}(x ; t)+\check{\Phi}_{n, \kappa} \text { for all }(x, t) \in[0,1]^{d+1}\right)=\mathrm{e}^{-2 \mathrm{e}^{-\kappa}}
$$

where the constant $\check{\Phi}_{n, \kappa}$ is defined by

$$
\check{\Phi}_{n, \kappa}=\frac{\left(\kappa / D_{n, 3}+D_{n, 3}\right) \sigma\|\check{K}\|_{L^{2}\left(\mathbb{R}^{d+1}\right)}}{h^{\beta} \sqrt{h^{d} n^{d} a_{n}^{d} m b_{m} h_{t}}(2 \pi)^{d+1}} .
$$

Asymptotic confidence bands for the function $f_{t}(x)$ at level $1-\alpha$ are hence given by

$$
\left\{\left[\check{f}_{n}(x ; t)-\check{\Phi}_{n,-\ln (-0.5 \ln (1-\alpha))}, \check{f}_{n}(x ; t)+\check{\Phi}_{n,-\ln (-0.5 \ln (1-\alpha))}\right] \mid(x, t) \in[0,1]^{d+1}\right\} .
$$

## 5. Finite sample properties

In this section, we investigate the finite sample properties of the proposed asymptotic confidence bands by means of a small simulation study and illustrate the procedure in an example analyzing the luminosity profile of the elliptical galaxy NGC5017.

### 5.1. Simulation study

All results are based on 5000 simulation runs. We simulate data from the bivariate regression model (1.1) where the errors $\varepsilon_{\left(k_{1}, k_{2}\right)}$ are independent and $\mathcal{N}\left(0, \sigma^{2}\right)$-distributed, where $\sigma=$ $0.1,\left(k_{1}, k_{2}\right) \in\{-n, \ldots, n\}^{2}, n \in\{150,250,350,500,650\}$ and our dataset provides $n=150$, which corresponds to a grid of $301 \times 301$ data points. For the unknown regression function we consider both, a unimodal function

$$
f_{1}(x, y)=4 \exp \left(-(3.5 x-1.5)^{2}-(3.5 y-1.5)^{2}\right)
$$

and a bimodal function

$$
\begin{equation*}
f_{2}(x, y)=4 \exp \left(-(3(x-0.1))^{2}-3.5(y-0.5)^{2}\right)+3 \exp \left(-(2(x-1))^{2}-3.5(y-0.5)^{2}\right) \tag{5.1}
\end{equation*}
$$

As convolution function $\psi$ we consider $\psi(x, y)=\frac{1}{4} \exp (-(|x|+|y|))$ and the values for the sequence $\left(a_{n}\right)_{n \in \mathbb{N}}$ are chosen such that most of the signal considerably different from 0 is included in the observations, that is $\left(a_{150}, a_{250}, a_{350}, a_{500}, a_{650}\right)=(0.29,0.28,0.27,0.26,0.25)$. A difference-based variance estimator is used to estimate $\sigma$. Figure 1 shows exemplary one simulated dataset and the reconstruction of the bimodal regression function $f_{2}$ from this dataset in comparison to the function $f_{2}$ itself and the convolution $f_{2} * \psi$.

For computational feasibility, we determine at first for each scenario a bandwidth by a small preliminary simulation study. For this purpose, we applied the $L_{\infty}$-motivated bandwidth selection method introduced in Bissantz et al. [7] and the estimated bandwidth is used in all 5000 runs.


Figure 1. Upper panel: bimodal regression function $f_{2}$ (left) and convolution $f_{2} * \psi$ (right) defined in (5.1). Lower panel: simulated data from $f_{2}$ observations according to the model (1.1) (left) and the corresponding reconstruction of $f_{2}$ (right). (See Section 5.1 for details of the choice of the (slightly undermoothing) bandwidth.)

Table 1 shows the simulated coverage probabilities and the average half-widths of the bands normalized with respect to the maximum of the respective function. Figure 2 illustrates the decrease of the normalized average half-widths of the confidence bands plotted against the sample sizes

Table 1. Simulated coverage probabilities and mean half-lengths of the corresponding confidence bands for the bivariate Gaussian function $f_{1}$ and the bimodal function $f_{2}$

| $n$ | $a_{n}$ | $f$ | 90\% nominal coverage |  | 95\% nominal coverage |  | 99\% nominal coverage |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  |  | 100-Cov. (\%) | Length | 100-Cov. (\%) | Length | 100-Cov. (\%) | Length |
| 150 | 0.29 | $f_{1}$ | 11.5 | 0.235 | 3.6 | 0.240 | 0.1 | 0.285 |
|  |  | $f_{2}$ | 9.0 | 1.219 | 4.4 | 0.324 | 0.8 | 0.368 |
| 250 | 0.28 | $f_{1}$ | 10.6 | 0.159 | 3.9 | 0.169 | 0.2 | 0.192 |
|  |  | $f_{2}$ | 9.7 | 0.265 | 4.8 | 0.290 | 0.6 | 0.319 |
| 350 | 0.27 | $f_{1}$ | 10.0 | 0.109 | 5.0 | 0.146 | 0.5 | 0.166 |
|  |  | $f_{2}$ | 10.4 | 0.240 | 5.9 | 0.254 | 0.6 | 0.287 |
| 500 | 0.26 | $f_{1}$ | 8.6 | 0.108 | 4.3 | 0.115 | 0.6 | 0.130 |
|  |  | $f_{2}$ | 9.1 | 0.191 | 5.4 | 0.203 | 0.8 | 0.229 |
| 650 | 0.25 | $f_{1}$ | 10.3 | 0.412 | 5.3 | 0.437 | 0.6 | 0.495 |
|  |  | $f_{2}$ | 9.7 | 0.648 | 4.8 | 0.687 | 0.6 | 0.775 |



Figure 2. Confidence band average half-lengths for both regression functions, $f_{1}$ ( $\alpha=0.01$ : $\boldsymbol{\square}, \alpha=0.05$ : $\Delta, \alpha=0.1: \bullet)$ and $f_{2}(\alpha=0.01: \square, \alpha=0.05: \Delta, \alpha=0.1: \bigcirc)$.
for both, the unimodal and the bimodal function. We conclude that for larger sample sizes (note that $n=150$ corresponds to $301 \times 301$ observations) and relatively small variances the simulated coverage probabilities are close to the nominal values. For $n=150$ the bands are rather wide, especially for the regression function $f_{2}$, that requires a smaller bandwidth for the estimation than the function $f_{1}$, which results in considerably wider bands for the function $f_{2}$ than for $f_{1}$. For increasing sample sizes, the widths of the bands decrease significantly. For illustrational purposes, Figure 3 shows a cross-section of the bivariate function $f_{2}$, estimators and a set of $90 \%$ confidence bands for $n=150$ and $n=650$ respectively, where $y=0.5$ and the corresponding confidence bands have been obtained in the bivariate setting. This figure clearly demonstrates the increase of precision of the bands for increasing sample size. Note that the sample size $301 \times 301$ in this example is rather small compared to the sample sizes which are usually available for astronomical images. Moreover, in these applications the signal-to-noise ratio is often much smaller and the point spread function is usually more sharply peaked than the one used in the simulations.

### 5.2. The luminosity profile of the galaxy NGC5017

In this section, we use the methodology derived in this paper to analyze the shape of luminosity profiles of elliptical galaxies. Figure 4 shows a $301 \times 301$ pixel section of an HST/WFPC2 Rband image of the elliptical galaxy NGC5017.

It is well known that images, taken with telescopes, are usually at least slightly blurry which can be modelled as convolution of the sharp image with a so-called point spread function (PSF), $\psi$, of the optical instrument. The detector used is a digital imaging device (CCD, charge-coupled device). We use a dataset of the size $301 \times 301$, where each data point corresponds to a pixel on an equally spaced grid. Hence, the two dimensional model (1.1) is suitable to describe the data.


Figure 3. Cross-section of bivariate, bimodal regression function $f_{2}$ (solid line), cross-section of reconstructions of $f_{2}$ (thin dashed lines) and corresponding $90 \%$-confidence bands (thick dashed lines) for $n=650$ (left) and $n=150$ (right).

In the analysis of elliptical galaxies the luminosity profile, that is, the decrease of the brightness with increasing distance from the galactic centre, is of particular interest. We use the confidence bands based on the nonparametric estimator (2.4) to narrow down the region for the parameter $\kappa$ in the Sérsic (1968) model for the luminosity profile of NGC5017. This model is defined as

$$
\begin{equation*}
f^{\kappa}(r)=I_{0} \cdot \exp \left(-b_{\kappa}\left(r / r_{e}\right)^{1 / \kappa}\right), \tag{5.2}
\end{equation*}
$$

where $r$ is the distance from the centre of the galaxy, $I_{0}$ is the brightness in the centre $(r=0), r_{e}$ is the scale radius (i.e., the half-light radius) and $b_{\kappa}$ is a normalization constant that is uniquely


Figure 4. Hubble Space Telescope/Wide Field Planetary Camera 2 [HST/WFPC2] R-band image of the elliptical galaxy NGC5017. Image source: NASA/STScI, dataset U3CM1A01R.


Figure 5. Upper and lower $90 \%$-confidence bands for the luminosity profile of the galaxy (solid lines) and corresponding data points $(*)$. The left plot shows the central region and the right plot the region in which the profile is fitted with Sérsic-curves for $\kappa=4$ (dotted line), $\kappa=5.11$ (dashed-dotted line) and $\kappa=7$ (dashed line).
determined by the choice of $\kappa$ which is the shape parameter, controlling the curvature of the profile. Note that $I_{0}$ and $r_{e}$ are model independent quantities that can be found in the literature. For $\kappa=4$, model (5.2) coincides with the famous de Vaucouleur-profile (1959). For details see, for example, Trujillo et al. [31] who already classified the galaxy under consideration as Sérsictype galaxy with $\kappa=5.11$. To analyse the data, we first fit the PSF, given by

$$
\psi(x, y)=\frac{\lambda}{2} \exp (-\lambda(|x|+|y|))
$$

to the (probably stellar) point source on the left middle of the image since the PSF is not completely known in advance. The bandwidth was chosen according to the procedure described in Section 5.1. We compute the estimator and $90 \%$-confidence bands for the unknown luminosity profile which, in combination, suggest that the actual brightness gradient is clearly steeper in the central region than the (convolved), unprocessed data tell us, see Figure 5, left panel, for a cross-section. The oscillations are artifacts due to our use of a Fourier estimator whose severity is partly caused by the use of one constant bandwidth for the reconstrucion of the whole image. In order to reconstruct the central region of the galaxy with sufficient precision, the bandwidth is chosen too small for the shallower regions of the image. For the analysis, we restrict ourselves to the square $[-0.55,0.45]^{2}$ and check for which parameters $\kappa \in[4,8]$ the resulting profile $f^{\kappa}$ according to formula (5.2) is completely contained in the $90 \%$-confidence band. We find that this is satisfied for the parameters $\kappa \in[4.81,5.93]$, which suggests that the profile for $\kappa=4$, corresponding to the de Vaucouleur law is not appropriate to describe the data well. On the other hand, the data do not provide evidence against $\kappa=5.11$, as proposed by Trujillo et al. [31].

## 6. Proofs of Theorem 1 and Corollary 1

### 6.1. Notation, preliminaries and remarks

First, we introduce some notation which is used extensively in the following proofs. Define for $a=\left(a_{1}, \ldots, a_{d}\right), b=\left(b_{1}, \ldots, b_{d}\right) \in \mathbb{R}^{d}$ the $d$-dimensional cube $[a, b]:=X_{j=1}^{d}\left[a_{j}, b_{j}\right]$. Let $\mathbf{k}=\left(k_{1}, \ldots, k_{d}\right) \in \mathbb{Z}^{d}, \alpha=\left(\alpha_{1}, \ldots, \alpha_{d}\right) \in\{0,1\}^{d}$ be multi-indices, $\mathbf{0}:=(0, \ldots, 0)^{T} \in \mathbb{R}^{d}$ and $\mathbf{1}:=(1, \ldots, 1)^{T} \in \mathbb{R}^{d}$ and define $G_{\mathbf{k}}:=\mathbb{Z}^{d} \cap[-\mathbf{k}, \mathbf{k}]$. For $j \in\{1, \ldots, d\}$ we denote by $G_{\mathbf{k}}^{j}$ the canonical projection of $G_{\mathbf{k}}$ onto $\mathbb{Z}^{j}$, i.e., $G_{\mathbf{k}}^{j}$ is a $j$-dimensional grid of integers with possibly different length in each direction. For $j \in \mathbb{N}$ let $G_{\mathbf{k}}^{j,+}:=G_{\mathbf{k}}^{j} \cap \mathbb{N}^{j}$ denote the part of the grid $G_{\mathbf{k}}^{j}$ whose vectors contain only positive components and write $G_{\mathbf{k}}^{+}$for $G_{\mathbf{k}}^{d,+}$. We further introduce the bijective map

$$
E_{d}:\left\{\begin{array}{l}
\{0,1\}^{d} \rightarrow \mathcal{P}(\{1, \ldots, d\}), \\
\left(\alpha_{1}, \ldots, \alpha_{d}\right) \mapsto v=\left\{v_{1}, \ldots, v_{|\alpha|}\right\} ; \quad \alpha_{v_{j}}=1, j=1, \ldots,|\alpha|=\sum_{i=1}^{d} \alpha_{i},
\end{array}\right.
$$

that maps each $\alpha$ to the set $v \subset\{1, \ldots, d\}$ that contains the positions of its ones. For $\alpha \in\{0,1\}^{d}$ and $\left\{v_{1}, \ldots, v_{|\alpha|}\right\}=E_{d}(\alpha)$ let $(x)_{\alpha}:=\left(x_{v_{1}}, x_{v_{2}}, \ldots, x_{v_{|\alpha|}}\right)$ denote the projection of $x \in \mathbb{R}^{d}$ onto the space spanned by the coordinate axes given by the positions of ones of the multi-index $\alpha$. For $a, b \in \mathbb{R}^{d}$ let $(a)_{\alpha}:(b)_{1-\alpha}=(a: b)_{\alpha}:=\left(a_{1}^{\alpha_{1}} \cdot b_{1}^{1-\alpha_{1}}, \ldots, a_{d}^{\alpha_{d}} \cdot b_{d}^{1-\alpha_{d}}\right)$ denote the vector of the components of $a$ and $b$ specified by the index $\alpha$. The following example illustrates these notations.

Example 2. For $d=2$ we have $\{0,1\}^{2}=\{(1,1),(1,0),(0,1),(0,0)\}$ and the mapping $E_{2}$ is defined by

$$
E_{2}((1,1))=\{1,2\}, \quad E_{2}((1,0))=\{1\}, \quad E_{2}((0,1))=\{2\} \quad \text { and } \quad E_{2}((0,0))=\emptyset .
$$

For any $x=\left(x_{1}, x_{2}\right) \in \mathbb{R}^{2}$ we have

$$
(x)_{(1,1)}=x, \quad(x)_{(1,0)}=x_{1} \quad \text { and } \quad(x)_{(0,1)}=x_{2} .
$$

For $a=\left(a_{1}, a_{2}\right), b=\left(b_{1}, b_{2}\right) \in \mathbb{R}^{2}$ we have

$$
\begin{aligned}
& (a: b)_{(1,1)}=\left(a_{1}, a_{2}\right)=a, \quad(a: b)_{(1,0)}=\left(a_{1}, b_{2}\right), \\
& (a: b)_{(0,1)}=\left(b_{1}, a_{2}\right), \quad(a: b)_{(0,0)}=\left(b_{1}, b_{2}\right)=b .
\end{aligned}
$$

For the approximation of the integrals by Riemann sums we define for multi-indices $\tilde{\alpha}, \alpha \in$ $\{0,1\}^{d} \backslash\{\mathbf{0}\}$

$$
\begin{align*}
\Delta_{\alpha}(f ; a, b) & :=\sum_{\tilde{\alpha} \in\{0,1\}^{d}, \tilde{\alpha} \leq \alpha}(-1)^{|\tilde{\alpha}|} f\left((a: b)_{\tilde{\alpha}}\right) \\
& =\sum_{\tilde{\alpha} \in\{0,1\}^{d}, \tilde{\alpha} \leq \alpha}(-1)^{d-|\tilde{\alpha}|} f\left((a)_{1-\tilde{\alpha}}:(b)_{\tilde{\alpha}}\right), \tag{6.1}
\end{align*}
$$

where the symbol $\tilde{\alpha} \leq \alpha$ means $\tilde{\alpha}_{j} \leq \alpha_{j}$ for $j=1, \ldots, d$. Note that for $\alpha=\mathbf{1} \in \mathbb{R}^{d}$ we obtain the special case of the $d$-fold alternating sum, that is,

$$
\Delta(f ; a, b):=\Delta_{\mathbf{1}}(f ; a, b)=\sum_{\alpha \in\{0,1\}^{d}}(-1)^{|\alpha|} f\left((a: b)_{\alpha}\right)=\sum_{\alpha \in\{0,1\}^{d}}(-1)^{d-|\alpha|} f\left((a: b)_{1-\alpha}\right)
$$

Note that $\Delta_{\alpha}(f ; a, b)$ can be regarded as the increment of the function $f_{\alpha}\left((x)_{\alpha}\right):=f\left((x: b)_{\alpha}\right)$ over the interval $\left[(a)_{\alpha},(b)_{\alpha}\right]$ which also gives rise to the alternative notation

$$
\begin{equation*}
\Delta_{\alpha}(f ; a, b)=\Delta\left(f_{\alpha},(a)_{\alpha},(b)_{\alpha}\right) \tag{6.2}
\end{equation*}
$$

### 6.2. Proof of Theorem 1

To prove the assertion of Theorem 1 we decompose the index set $G_{\mathbf{n}}=\{-n, \ldots, n\}^{d}$ of the sum in (3.1) into $2^{d}+1$ parts: the respective intersections with the $2^{d}$ orthants of the origin in $\mathbb{R}^{d}$ and the marginal intersections with the coordinate axes. Our first auxiliary result shows that the contribution of the term representing the marginals is negligible (here and throughout the paper we use the convention $0^{0}=1$ ).

## Lemma 1.

$$
\sup _{x \in\left[0, h^{-1}\right]^{d}}\left|\frac{h^{\beta}}{\sqrt{h^{d} n^{d} a_{n}^{d}}} \sum_{\alpha \in\{0,1\}^{d} \backslash\{\mathbf{1}\}} \sum_{(\mathbf{k}: \mathbf{0})_{\alpha}, \mathbf{k} \in G_{n}^{+}} K_{n}\left(x-\frac{1}{h} x_{\mathbf{k}}\right) \varepsilon_{\mathbf{k}}\right|=\mathrm{o}_{P}\left(\frac{1}{\log (n)}\right) .
$$

We obtain from its proof in Section 7 that Lemma 1 holds under the weaker condition $\frac{\log (n)}{\sqrt{n a_{n} h}}=$ o(1), which follows from the assumptions of Theorem 1. Next we consider the "positive" orthant $G_{n}^{+}$and show in three steps that

$$
\begin{equation*}
\sup _{x \in\left[0, h^{-1}\right]^{d}}\left|Y_{n}^{(+)}(x)-Y^{(+)}(x)\right|=\mathrm{o}_{p}(1) \tag{6.3}
\end{equation*}
$$

where the processes $Y_{n}^{(+)}$and $Y^{(+)}$are defined by

$$
\begin{align*}
Y_{n}^{(+)}(x) & :=\frac{(2 \pi)^{d} h^{\beta}}{\sigma\|K\|_{L^{2}} \sqrt{h^{d} n^{d} a_{n}^{d}}} \sum_{\mathbf{k} \in G_{n}^{+}} K_{n}\left(x-\frac{1}{h} x_{\mathbf{k}}\right) \varepsilon_{\mathbf{k}},  \tag{6.4}\\
Y^{(+)}(x) & :=\frac{(2 \pi)^{d}}{\|K\|_{L^{2}}} \int_{\mathbb{R}_{+}^{d}} K(x-u) \mathrm{d} B(u) \tag{6.5}
\end{align*}
$$

respectively, $B$ is a standard Brownian sheet on $\mathbb{R}^{d}$ (see the proof of Lemma 2 for details) and $K$ denotes the kernel defined in Assumption B. The final result is then derived using Theorem 14.1 in Piterbarg [27]. To be precise note that it can easily be shown that

$$
\lim _{n \rightarrow \infty} n^{d} a_{n}^{d} h^{d} h^{2 \beta} \cdot \operatorname{Var}\left(\hat{f_{n}}(x)\right)=\frac{\sigma^{2}}{(2 \pi)^{2 d}} \int_{\mathbb{R}^{d}}\left|K\left(\frac{x}{h}-u\right)\right|^{2} \mathrm{~d} u=\frac{\sigma^{2}\|K\|_{L_{2}}^{2}}{(2 \pi)^{2 d}}
$$

(in particular the limit is independent of the variable $x$, which is typical for kernel estimates in homoscedastic regression models with equidistant design). We further obtain for the function $r(t)=(2 \pi)^{2 d}\|K\|_{L^{2}}^{-2} \int_{\mathbb{R}^{d}} K(v+t) K(v) \mathrm{d} v$ that

$$
\|r\|_{L^{1}}=\frac{(2 \pi)^{2 d}}{\|K\|_{L^{2}}^{2}} \int_{\mathbb{R}^{d}}\left|\int_{\mathbb{R}^{d}} K(v+t) K(v) \mathrm{d} v\right| \mathrm{d} t \leq \frac{(2 \pi)^{2 d}\|K\|_{L^{1}}^{2}}{\|K\|_{L^{2}}^{2}}<\infty
$$

Therefore the conditions of Theorem 14.1 in Piterbarg [27] are satisfied and the assertion of Theorem 1 follows.

The remaining proof of the uniform approximation (6.3) will be accomplished showing the following auxiliary results. For this purpose we introduce the process

$$
Y_{n, 1}^{(+)}(x):=\frac{(2 \pi)^{d} h^{\beta}}{\sqrt{n^{d} a_{n}^{d} h^{d}}\|K\|_{L^{2}\left(\mathbb{R}^{d}\right)}} \sum_{\alpha \in\{0,1\}^{d}}(-1)^{|\alpha|} \sum_{\mathbf{j} \in G_{n}^{|\alpha|,+}} \Delta_{\alpha}\left(K_{n} \circ \tau_{x}, I_{\mathbf{j}}\right) B\left(\mathbf{j}:(\mathbf{n})_{\mathbf{1}-\alpha}\right),
$$

where the function $\tau_{x}$ is defined by $\tau_{x}(u):=x-\frac{u+1}{n a_{n} h}$,

$$
\begin{equation*}
I_{\mathbf{j}}:=\left[(\mathbf{j}-\mathbf{1}):(\mathbf{n})_{1-\alpha}, \mathbf{j}:(\mathbf{n})_{1-\alpha}\right] \subset \mathbb{R}_{+}^{d} \tag{6.6}
\end{equation*}
$$

and we use the notation (6.2).
Lemma 2. There exists a Brownian sheet B on $\mathbb{R}^{d}$ such that

$$
\sup _{x \in\left[0, h^{-1}\right]^{d}}\left|Y_{n}^{(+)}(x)-Y_{n, 1}^{(+)}(x)\right|=\mathrm{o}\left(\frac{1}{\sqrt{\log (n)}}\right) \quad \text { a.s. }
$$

We obtain from the proof in Section 7.1 that Lemma 2 holds under the condition $\frac{\log (n)}{n^{\delta / 2} a_{n}^{\delta / 2} h^{d / 2}}=$ $\mathrm{o}(1)$, which follows from the assumptions of Theorem 1. The next step consists of the replacement of the kernel $K_{n}$ in the process $Y_{n, 1}$ by its limit.

## Lemma 3.

$$
\sup _{x \in\left[0, h^{-1}\right]}\left|Y_{n, 1}(x)-Y_{n, 2}(x)\right|=\mathrm{o}_{P}\left(\frac{1}{\log (n)}\right),
$$

where the process $Y_{n, 2}$ is given by

$$
Y_{n, 2}(x):=\frac{(2 \pi)^{d}}{\sqrt{n^{d} a_{n}^{d} h^{d}}\|K\|_{L^{2}}} \sum_{\alpha, \gamma \in\{0,1\}^{d}}(-1)^{|\alpha|} \sum_{\mathbf{j} \in G_{n}^{|\alpha|,+}} \Delta_{\alpha}\left(K \circ \tau_{x},(-1)^{\gamma} I_{\mathbf{j}}\right) B\left((-1)^{\gamma} \mathbf{j}:(\mathbf{n})_{\mathbf{1}-\alpha}\right) .
$$

As described in Section 6.1 for fixed $\alpha \in\{0,1\}^{d}, j \in G_{n}^{|\alpha|,+}$ the quantity $\Delta_{\alpha}\left(K \circ \tau_{x} ; I_{\mathbf{j}}\right)$ can be regarded as the increment of the function $\left(K_{n} \circ \tau_{x}\right)_{\alpha}\left((u)_{\alpha}\right)=K_{n} \circ \tau_{x}\left((u: n)_{\alpha}\right)$ on the cube $[(\mathbf{j}-\mathbf{1}), \mathbf{j}]$. This point of view is the basic step in the approximation by the Riemann-Stieltjes integral of $B\left(((\cdot): \mathbf{n})_{1-\alpha}\right)$ with respect to the function $\left(K_{n} \circ \tau_{x}\right)_{\alpha}$ for each $\alpha \in\{0,1\}^{d}$.

## Lemma 4.

$$
\sup _{x \in\left[0, h^{-1}\right]^{d}}\left|Y_{n, 2}^{(+)}(x)-Y_{n, 3}^{(+)}(x)\right|=\mathrm{o}_{P}\left(\frac{1}{\log (n)}\right),
$$

where the process $Y_{n, 3}^{(+)}$is defined by

$$
\begin{equation*}
Y_{n, 3}^{(+)}(x) \stackrel{\mathcal{D}}{=} \frac{(2 \pi)^{d}}{\|K\|_{L^{2}}} \int_{\left[0,\left(a_{n} h\right)^{-1}\right]^{d}} K(x-u) \mathrm{d} B(u) \tag{6.7}
\end{equation*}
$$

We obtain from its proof in Section 7.2 that Lemma 4 holds under the condition $\frac{\log (n)}{n h^{d}}=\mathrm{o}(1)$, which follows from the assumptions of Theorem 1. In the final step we show that the difference

$$
Y^{(+)}(x)-Y_{n, 3}^{(+)}(x)=\frac{(2 \pi)^{d}}{\|K\|_{L^{2}}} \int_{\mathbb{R}_{+}^{d}} I_{\mathbb{R}_{+}^{d} \backslash\left[0,\left(a_{n} h\right)^{-1}\right]^{d}}(u) K(x-u) \mathrm{d} B(u)
$$

is asymptotically negligible.

## Lemma 5.

$$
\sup _{x \in\left[0, h^{-1}\right]^{d}}\left|Y_{n, 3}(x)-Y(x)\right|=\mathrm{o}_{P}\left(\left(\log (n)^{-1}\right)\right.
$$

### 6.3. Proof of Corollary 1

The assertion follows from the estimate

$$
\begin{equation*}
\sup _{[0,1]^{d}}\left|f(x)-\mathbb{E} \hat{f}_{n}(x)\right|=\mathrm{o}\left(h^{-\beta}\left(h^{d} n^{d} a_{n}^{d}\right)^{-1 / 2}\right) \tag{6.8}
\end{equation*}
$$

To prove (6.8) we use the representation (2.6) and obtain by a straightforward calculation

$$
\begin{aligned}
\mathbb{E} \hat{f}_{n}(x) & =\frac{1}{(2 \pi)^{d} n^{d} a_{n}^{d} h^{d}} \sum_{\mathbf{k} \in G_{\mathbf{n}}}(f * \psi)\left(x_{\mathbf{k}}\right) \cdot K_{n}\left(\left(x-x_{\mathbf{k}}\right) \frac{1}{h}\right) \\
& =\frac{1}{(2 \pi)^{d} h^{d}} \int_{\left[-1 / a_{n}, 1 / a_{n}\right]^{d}}(f * \psi)(z) \cdot K_{n}\left((x-z) \frac{1}{h}\right) \mathrm{d} z+R_{n, 1}(x) \\
& =\frac{1}{(2 \pi)^{d}} \int_{\mathbb{R}^{d}}(f * \psi)(z \cdot h) \cdot K_{n}\left(\frac{x}{h}-z\right) \mathrm{d} z+R_{n, 1}(x)+R_{n, 2}(x)
\end{aligned}
$$

where the term

$$
\begin{align*}
R_{n, 1}(x)=\frac{1}{(2 \pi)^{d} h^{d}} \sum_{\mathbf{k} \in G_{\mathbf{n}-\mathbf{1}}} \int_{\left[x_{\mathbf{k}}, x_{\mathbf{k}+\mathbf{1}}\right]}\{ & (f * \psi)\left(x_{\mathbf{k}}\right) K_{n}\left(\frac{x-x_{\mathbf{k}}}{h}\right) \\
& \left.-(f * \psi)(z) K_{n}\left(\frac{x-z}{h}\right)\right\} \mathrm{d} z \tag{6.9}
\end{align*}
$$

denotes the "error" in the integral approximation and

$$
R_{n, 2}(x):=\frac{1}{(2 \pi)^{d} h^{d}} \int_{\left(\left[-1 / a_{n}, 1 / a_{n}\right]^{d}\right)^{C}}(f * \psi)(z) K_{n}\left((x-z) \frac{1}{h}\right) \mathrm{d} z .
$$

An application of the Plancherel identity (see for example Folland [15], Theorem 8.29) gives (observing Assumption A1 and A3)

$$
\begin{aligned}
\mathbb{E} \hat{f}_{n}(x)= & \frac{1}{(2 \pi)^{d / 2} h^{d}} \int_{\mathbb{R}^{d}} \mathcal{F} f\left(h^{-1} \xi\right) \mathcal{F} \psi\left(h^{-1} \xi\right) \frac{\mathcal{F} \eta(\xi)}{\mathcal{F} \psi(\xi / h)} \exp \left(\mathrm{i} h^{-1} x^{T} \xi\right) \mathrm{d} \xi \\
& +R_{n, 1}(x)+R_{n, 2}(x) \\
= & \frac{1}{(2 \pi)^{d / 2}} \int_{\mathbb{R}^{d}} \mathcal{F} f(\xi) \cdot \mathcal{F} \eta(\xi h) \exp \left(\mathrm{i} x^{T} \xi\right) \mathrm{d} \xi+R_{n, 1}(x)+R_{n, 2}(x) \\
= & f(x)+R_{n, 1}(x)+R_{n, 2}(x)+R_{n, 3}(x)+R_{n, 4}(x),
\end{aligned}
$$

where

$$
\begin{aligned}
& R_{n, 3}(x)=\frac{1}{(2 \pi)^{d / 2}} \int_{\left([-D / h, D / h]^{d}\right)^{C}} \mathcal{F} f(\xi) \exp (\mathrm{i} x \xi) \mathrm{d} \xi \\
& R_{n, 4}(x)=\frac{1}{(2 \pi)^{d / 2}} \int_{[-1 / h, 1 / h]^{d} \backslash[-D / h, D / h]^{d}} \mathcal{F} f(\xi) \cdot \mathcal{F} \eta(\xi h) \exp (\mathrm{i} x \xi) \mathrm{d} \xi
\end{aligned}
$$

We further obtain from Assumption C

$$
\left|\int_{\left\{\xi_{j}>D / h\right\}} \mathcal{F} f(\xi) \exp (-\mathrm{i} x \xi) \mathrm{d} \xi\right| \leq \frac{1}{D^{\gamma}} \int_{\left\{\xi_{j}>D / h\right\}}|\mathcal{F} f(\xi)|\left(h \xi_{j}\right)^{\gamma} \mathrm{d} \xi=\mathrm{o}\left(h^{\gamma}\right)
$$

and finally $\left|R_{n, 3}(x)\right| \leq \sum_{j=1}^{d} \int_{\left\{\xi_{j}>D / h\right\}}|\mathcal{F} f(\xi)| \mathrm{d} \xi=\mathrm{o}\left(h^{\gamma}\right)$. With the same arguments it follows $R_{n, 4}(x)=\mathrm{o}\left(h^{\gamma}\right)$, since $|\mathcal{F} \eta(\xi h)| \leq 1$ for all $\xi \in \mathbb{R}^{d}$. Define $\mathcal{A}_{n}=\left(\left[-\frac{1}{a_{n}}, \frac{1}{a_{n}}\right]^{d}\right)^{C}$, then we obtain from the representation (2.7) the estimate

$$
\begin{aligned}
\left|R_{n, 2}(x)\right| \leq & \frac{1}{(2 \pi)^{d} h^{\beta+d}}\left(\int_{\mathcal{A}_{n}}|(f * \psi)(z)|^{2} \mathrm{~d} z\right)^{1 / 2} \\
& \times\left[\left(\int_{\mathcal{A}_{n}}\left|K\left((x-z) \frac{1}{h}\right)\right|^{2} \mathrm{~d} z\right)^{1 / 2}\right. \\
& \left.\quad+\left(\int_{\mathcal{A}_{n}}\left|\left(h^{\beta} K_{n}-K\right)\left((x-z) \frac{1}{h}\right)\right|^{2} \mathrm{~d} z\right)^{1 / 2}\right] \\
= & \frac{1}{(2 \pi)^{d} h^{\beta}}\left(\int_{\mathcal{A}_{n}}|(f * \psi)(z)|^{2} \mathrm{~d} z\right)^{1 / 2}\left(\mathrm{O}\left(h^{d} a_{n}^{d / 2}\right)+\mathrm{O}\left(h^{\mu_{1}+d} a_{n}^{d / 2}\right)\right)=\mathrm{O}\left(\frac{a_{n}^{v+d / 2}}{h^{\beta}}\right)
\end{aligned}
$$

uniformly with respect to $x \in[0,1]^{d}$. Note that by Assumption C we have $f \in \mathcal{W}^{\lfloor m\rfloor}\left(\mathbb{R}^{d}\right)$ and since $m>2+\frac{d}{2}$ Sobolev's Embedding Theorem (Folland [15], Theorem 8.54) implies the ex-
istence of a function $\tilde{f} \in C^{2}\left(\mathbb{R}^{d}\right)$ with $f=\tilde{f}$ almost everywhere. Observing that the convolution function $\psi$ is integrable gives $\partial^{\alpha}(f * \psi)=\left(\partial^{\alpha} f\right) * \psi \in C\left(\mathbb{R}^{d}\right)$ for all $\alpha \in\{0,1,2\}$ with $|\alpha| \leq 2$ (see for example Folland [15], Proposition 8.10), which justifies the application of Taylor's Theorem. Straightforward but tedious calculations give for the remaining term (6.9) $R_{n, 1}(x)=\mathrm{O}\left(\frac{1}{n^{3} a_{n}^{3} h^{\beta+2}}\right)+\mathrm{O}\left(\frac{a_{n}^{v}}{n h^{\beta}}\right)$ uniformly with respect to $x \in[0,1]^{d}$.

### 6.4. Proofs of Theorem 2 and Corollary 2

First we will show that the kernel $\check{K}_{n}$ satisfies conditions B1 and B2, with the kernel $\check{K}$ defined (4.3). If Assumption B̌ holds we have

$$
\begin{aligned}
& \int_{\mathbb{R}^{d}}\left(\frac{h^{\beta}}{\mathcal{F}_{d} \psi(\xi / h)}-\Psi(\xi)\right) \mathcal{F}_{d+1} \check{\eta}(\xi, \tau) \exp \left(\mathrm{i} x^{T} \xi\right) \mathrm{d} \xi \\
& =\sum_{p=1}^{M-1} h^{\mu_{p}} \int_{\mathbb{R}^{d}} \Psi_{p}(\xi) \mathcal{F}_{d+1} \check{\eta}(\xi, \tau) \exp \left(\mathrm{i} x^{T} \xi\right) \mathrm{d} \xi \\
& \quad+h^{\mu_{M}} \int_{\mathbb{R}^{d}} \Psi_{M, n}(\xi) \mathcal{F}_{d+1} \check{\eta}(\xi, \tau) \exp \left(\mathrm{i} x^{T} \xi\right) \mathrm{d} \xi
\end{aligned}
$$

which implies

$$
\begin{aligned}
h^{\beta} \check{K}_{n}(x, t)-\check{K}(x, t)= & \sum_{p=1}^{M-1} h^{\mu_{p}} \int_{\mathbb{R}^{d+1}} \Psi_{p}(\xi) \mathcal{F}_{d+1} \check{\eta}(\xi, \tau) \exp \left(\mathrm{i} x^{T} \xi+\mathrm{i} t \tau\right) \mathrm{d}(\xi, \tau) \\
& +h^{\mu_{M}} \int_{\mathbb{R}^{d+1}} \Psi_{M, n}(\xi) \mathcal{F}_{d+1} \check{\eta}(\xi, \tau) \exp \left(\mathrm{i} x^{T} \xi+\mathrm{i} t \tau\right) \mathrm{d}(\xi, \tau)
\end{aligned}
$$

A careful inspection of the proofs of Theorem 1 and Corollary 1 shows that the arguments can be transferred to the time-dependent case if the increase of $n$ and $m(n)$ as well as the decrease of $a_{n}, b_{m(n)}, h$ and $h_{t}$ are balanced as given in the assumptions of the theorem. The details are omitted for the sake of brevity.

## 7. Proof of auxiliary results

### 7.1. Proof Lemma 2

Define $S_{\mathbf{k}}:=\sum_{\mathbf{j} \in G_{\mathbf{k}}^{+}} \varepsilon_{\mathbf{j}}$, set $S_{\mathbf{j}} \equiv 0$ if $\min \left\{j_{1}, \ldots, j_{d}\right\}=0$ and recall the definition of $Y_{n}^{(+)}$and $\tau_{x}$ in (6.4) and before Lemma 2, respectively. In a first step we will replace the errors $\varepsilon_{k}$ by increments given in terms of partial sums $S_{\mathbf{k}-\alpha}$ for $\alpha \in\{0,1\}^{d}$. To be precise, we use the representation

$$
\varepsilon_{\mathbf{k}}=\sum_{\underline{\alpha} \in\{0,1\}^{d}}(-1)^{|\underline{\alpha}|} S_{(\mathbf{k}-\underline{\alpha})}=\sum_{\underline{\alpha} \in\{0,1\}^{d}}(-1)^{|\underline{\alpha}|} S_{((\mathbf{k}-\mathbf{1}): \mathbf{k})_{\underline{\alpha}}} .
$$

A straightforward calculation gives

$$
\begin{aligned}
Y_{n}^{(+)}(x):= & \frac{h^{\beta}}{\sigma\|K\|_{2} \sqrt{n^{d} h^{d} a_{n}^{d}}} \sum_{\mathbf{k} \in G_{n}^{+}} K_{n} \circ \tau_{x}(\mathbf{k}-\mathbf{1}) \sum_{\alpha \in\{0,1\}^{d}}(-1)^{|\alpha|} S_{\mathbf{k}-\alpha} \\
= & \frac{h^{\beta}}{\sigma\|K\|_{2} \sqrt{n^{d} h^{d} a_{n}^{d}}} \sum_{\alpha \in\{0,1\}^{d}}(-1)^{|\alpha|} \sum_{\mathbf{k} \in G_{n}^{+}} K_{n} \circ \tau_{x}(\mathbf{k}-\mathbf{1}) S_{((\mathbf{k}-\mathbf{1}): \mathbf{k}) \alpha} \\
= & \frac{h^{\beta}}{\sigma\|K\|_{2} \sqrt{n^{d} h^{d} a_{n}^{d}}} \\
& \times\left(\sum_{\alpha \in\{0,1\}^{d}}(-1)^{|\alpha|} \sum_{\mathbf{k} \in G_{n}^{+}}\left(K_{n} \circ \tau_{x}(\mathbf{k}-\mathbf{1})-K_{n} \circ \tau_{x}\left(((\mathbf{k}-\mathbf{1}): \mathbf{k})_{\alpha}\right)\right) S_{((\mathbf{k}-\mathbf{1}): \mathbf{k}) \alpha}\right. \\
& \left.\quad+\sum_{\alpha \in\{0,1\}^{d}}(-1)^{|\alpha|} \sum_{\mathbf{k} \in G_{n}^{+}} K_{n} \circ \tau_{x}\left(((\mathbf{k}-\mathbf{1}): \mathbf{k})_{\alpha}\right) S_{((\mathbf{k}-\mathbf{1}): \mathbf{k}) \alpha}\right) .
\end{aligned}
$$

Now we can make use of Proposition 6 and Proposition 3 of Owen [25] to rewrite the sums, such that the increments given in terms of partial sums can be expressed by increments given in terms of the kernel $K_{n}$. We obtain

$$
\begin{aligned}
Y_{n}^{(+)}(x)= & \frac{h^{\beta}}{\sigma\|K\|_{2} \sqrt{n^{d} h^{d} a_{n}^{d}}} \\
& \times\left[K_{n} \circ \tau_{x}(\mathbf{n}) S_{(\mathbf{n})}\right. \\
& \quad+\sum_{\alpha \in\{0,1\}^{d}}(-1)^{|\alpha|} \sum_{\mathbf{k} \in G_{n}^{+}} \sum_{\beta \in\{0,1\}^{d} \backslash\{\mathbf{0}\}}(-1)^{|\beta|} \Delta_{\beta}\left(K_{n} \circ \tau_{x} ; \mathbf{k}-\mathbf{1},\right. \\
& \left.\left.((\mathbf{k}-\mathbf{1}): \mathbf{k})_{\alpha}\right) S_{((\mathbf{k}-\mathbf{1}): \mathbf{k}) \alpha}\right] .
\end{aligned}
$$

The quantity $\Delta_{\beta}\left(K_{n} \circ \tau_{x} ; \mathbf{k}-\mathbf{1},((\mathbf{k}-\mathbf{1}): \mathbf{k})_{\alpha}\right)$ can only take values different from zero if $\alpha \leq$ $\mathbf{1}-\beta$. Note that for $\alpha \leq \mathbf{1}-\beta$ the equality $\left.(\mathbf{k})_{\beta}=((\mathbf{k}-\mathbf{1}): \mathbf{k})_{\alpha}\right)_{\beta}$ holds which implies that in this case we also have $\left[(\mathbf{k}-\mathbf{1})_{\beta},\left(((\mathbf{k}-\mathbf{1}): \mathbf{k})_{\alpha}\right)_{\beta}\right]=\left[(\mathbf{k}-\mathbf{1})_{\beta},(\mathbf{k})_{\beta}\right]$. We further obtain

$$
\begin{aligned}
Y_{n}^{(+)}(x)= & \frac{h^{\beta}}{\sigma\|K\|_{2} \sqrt{n^{d} h^{d} a_{n}^{d}}} \\
& \times\left[K_{n} \circ \tau_{x}(\mathbf{n}) S_{(\mathbf{n})}\right. \\
& \quad+\sum_{\beta \in\{0,1\}^{d} \backslash\{\mathbf{0}\}}(-1)^{|\beta|} \sum_{\mathbf{k} \in G_{\mathbf{n}}^{+}} \sum_{\tilde{\alpha} \in\{0,1\}^{d-|\beta|}}(-1)^{|\tilde{\alpha}|}
\end{aligned}
$$

$$
\begin{aligned}
& \times \Delta_{\beta}\left(K_{n} \circ \tau_{x} ; \mathbf{k}-\mathbf{1},(\mathbf{k})_{\beta}:\right. \\
& \left.\quad\left((\mathbf{k}-\mathbf{1})_{\mathbf{1}-\beta}:(\mathbf{k})_{\mathbf{1}-\beta}\right)_{\tilde{\alpha}}\right) \\
& \times S_{\left.(\mathbf{k})_{\beta}:\left((\mathbf{k}-\mathbf{1})_{1-\beta}:(\mathbf{k})_{1-\beta}\right)_{\tilde{\alpha}}\right] .} .
\end{aligned}
$$

The alternating sum with respect to the index $\tilde{\alpha}$ can be written as an increment $\Delta$ as defined in (6.1) which then defines a telescope sum according to Owen [25], Proposition 2. Taking into account that $S(\mathbf{k}) \equiv 0$ if $k_{j}=0$ for at least one $j \in\{1, \ldots, d\}$ gives

$$
Y_{n}^{(+)}(x)=\frac{h^{\beta}}{\sigma\|K\|_{2} \sqrt{n^{d} h^{d} a_{n}^{d}}} \sum_{\beta \in\{0,1\}^{d}}(-1)^{|\beta|} \sum_{\mathbf{j} \in G_{\mathbf{n}}^{|\beta|,+}} \Delta_{\beta}\left(K_{n} \circ \tau_{x} ; I_{\mathbf{j}}\right) \cdot S_{\mathbf{j}:(\mathbf{n})_{1-\beta}} .
$$

With the definitions $X(A):=\sum_{\mathbf{k} \in A \subset \mathbb{Z}^{d}} X_{\mathbf{k}}$ for any subset $A \in \mathbb{Z}^{d}$, we can rewrite these partial sums as set-indexed partial sums with index class $n \cdot \mathscr{S}$, where $\mathscr{S}:=\left\{(0, \gamma] \mid 0<\gamma_{j} \leq 1,1 \leq\right.$ $j \leq d\}$ and $n \cdot \mathscr{S}:=\{n \cdot S \mid S \in \mathscr{S}\}$. It follows directly that $\mathscr{S}$ is a sufficiently smooth VC-class of sets, which justifies the application of Theorem 1 in Rio [28]. Therefore there exists a version of a Brownian sheet on $[0, \infty)^{d}$, say $B_{1}$, such that

$$
\begin{equation*}
\sup _{\mathbf{k} \in G_{\mathbf{n}}^{+}}\left|\frac{S_{\mathbf{k}}}{\sigma}-B_{1}(\mathbf{k})\right|=\mathrm{O}\left((\log (n))^{1 / 2} n^{(d-\delta) / 2}\right) \quad \text { a.s. } \tag{7.1}
\end{equation*}
$$

Recalling the definition of $I_{\mathbf{j}}$ in (6.6) we further obtain

$$
\begin{aligned}
Y_{n}^{(+)} & (x)-Y_{n, 1}^{(+)}(x) \\
= & \frac{h^{\beta}}{\|K\|_{2} \sqrt{n^{d} h^{d} a_{n}^{d}}} \\
& \times \sum_{\beta \in\{0,1\}^{d}}(-1)^{|\beta|} \sum_{\mathbf{j} \in G_{\mathbf{n}}^{|\beta|,+}} \Delta_{\beta}\left(K_{n} \circ \tau_{x} ; I_{\mathbf{j}}\right) \cdot\left(\frac{1}{\sigma} S_{\mathbf{j}:(\mathbf{n})_{1-\beta}}-B_{1}\left(\mathbf{j}:(\mathbf{n})_{\mathbf{1}-\beta}\right)\right) .
\end{aligned}
$$

The estimate (7.1) implies the existence of a constant $C \in \mathbb{R}_{+}$such that

$$
\begin{aligned}
\left|Y_{n}^{(+)}(x)-Y_{n, 1}^{(+)}(x)\right| & \\
\leq C \cdot \sqrt{\frac{\log (n)}{n^{\delta} h^{\delta} a_{n}^{\delta}}} h^{\beta} & \sum_{\gamma \in\{0,1\}^{d},|\gamma|=1} \int_{\left[0,\left(a_{n} h\right)^{-1}\right]^{d}}(u)_{\gamma}^{(d-\delta) / 2}\left|\partial^{\mathbf{1}} K_{n}(x-u)\right| \mathrm{d} u \\
& +\sum_{\beta \in\{0,1\}^{d} \backslash\{\mathbf{0}, \mathbf{1}\}} \int_{\left[0,\left(a_{n} h\right)^{-1}\right]|\beta|}\left|\partial^{\beta} K_{n}\left(\left(x-\left(u:\left(a_{n} h\right)^{-1} \mathbf{1}\right)\right)_{\beta}\right)\right|(\mathrm{d} u)_{\beta} \\
& \left.+\left|K_{n}\left(x-\left(a_{n} h\right)^{-1} \mathbf{1}\right)\right|\right] \quad \text { a.s. }
\end{aligned}
$$

It follows from Assumption B that the function $u \mapsto(u)_{\gamma}^{|\alpha| / 2} \partial^{\alpha} K(u)$ is integrable on $\mathbb{R}^{d}$ for all $\alpha \in\{0,1\}^{d}$ such that

$$
\int_{\left[0,\left(a_{n} h\right)^{-1}\right]^{d}}(u)_{\gamma}^{(d-\delta) / 2}\left|\partial^{1} K_{n}(x-u)\right| \mathrm{d} u=\mathrm{O}\left(h^{(\delta-d) / 2-\beta}\right)
$$

and

$$
\int_{\left[0,\left(a_{n} h\right)^{-1}\right]|\beta|}\left|\partial^{\beta} K_{n}\left(\left(x-\left(u:\left(a_{n} h\right)^{-1} \mathbf{1}\right)\right)_{\beta}\right)\right|(\mathrm{d} u)_{\beta}+\left|K_{n}\left(x-\left(a_{n} h\right)^{-1} \mathbf{1}\right)\right|=\mathrm{O}\left(\left(a_{n} h\right)^{d / 2} h^{-\beta}\right) .
$$

Note that for sufficiently large $n$ such that $a_{n}<\frac{1}{2}$ we obtain $-\frac{1}{2 a_{n} h} \geq x_{j}-\left(a_{n} h\right)^{-1}=\frac{a_{n}-1}{a_{n} h}$ uniformly with respect to $j$ (note that $x_{j} \in\left[0, h^{-1}\right]$ ). Let $\tilde{B}$ be a continuous version of $B_{1}$. We set $\tilde{B}(t) \equiv 0$ if $t_{j}<0$ for at least one index $j \in\{1, \ldots, d\}$ and let $\left\{\tilde{B}_{\alpha} \mid \alpha \in\{0,1\}^{d}\right\}$ be $2^{d}$ mutually independent copies of $\tilde{B}$. For $t \in \mathbb{R}^{d}$ define

$$
B_{\alpha}(t):=\tilde{B}_{\alpha}\left((-1)^{\alpha_{1}} t_{1},(-1)^{\alpha_{2}} t_{2}, \ldots,(-1)^{\alpha_{d}} t_{d}\right),
$$

then the process $\left\{B(t):=\sum_{\alpha \in\{0,1\}^{d}} B_{\alpha}(t) \mid t \in \mathbb{R}^{d}\right\}$ is a Wiener field on $\mathbb{R}^{d}$.

### 7.2. Proof of Lemma 4

Note that $\partial^{\alpha} K$ exists and is integrable for each $\alpha \in\{0,1\}^{d}$. Consequently, the kernel $K$ is of bounded variation on $\left[0,\left(a_{n} h\right)^{-1}\right]^{d}$ in the sense of Hardy Krause for each fixed $n$ (see Owen [25], Definition 2). Therefore an application of integration by parts for the Wiener integral (note that the kernel $K$ has not necessarily a compact support) and rescaling of the Brownian sheet $Y_{n, 3}^{(+)}$yields

$$
\begin{aligned}
& Y_{n, 3}^{(+)}(x) \stackrel{\mathcal{D}}{=} \sum_{\alpha \in\{0,1\}^{d} \backslash\{\mathbf{0}\}}(-1)^{|\alpha|} \int_{\left[0,\left(a_{n} h\right)^{-1}\right]^{|\alpha|}} B\left(\left(u:\left(a_{n} h\right)^{-1} \mathbf{1}\right)_{\alpha}\right) \mathrm{d} K\left(x-\left(u:\left(a_{n} h\right)^{-1} \mathbf{1}\right)_{\alpha}\right) \\
&+\Delta\left(K(x-\cdot) \cdot B(\cdot),\left[0,\left(a_{n} h\right)^{-1}\right]^{d}\right) \\
&= \sum_{\alpha \in\{0,1\}^{d} \backslash\{\boldsymbol{0}\}}(-1)^{|\alpha|} \int_{\left[0,\left(a_{n} h\right)^{-1}\right]^{|\alpha|}} B\left(\left(u:\left(a_{n} h\right)^{-1} \mathbf{1}\right)_{\alpha}\right) \partial^{\alpha} K\left(x-\left(u:\left(a_{n} h\right)^{-1} \mathbf{1}\right)_{\alpha}\right)(\mathrm{d} u)_{\alpha} \\
&+\Delta\left(K(x-\cdot) \cdot B(\cdot),\left[0,\left(a_{n} h\right)^{-1}\right]^{d}\right) .
\end{aligned}
$$

Recalling the definition of $Y_{n, 2}^{(+)}(x)$ and identity (15) in Owen [25] we can replace the increments by the corresponding integrals, that is

$$
\begin{aligned}
Y_{n, 2}^{(+)}(x) \stackrel{\mathcal{D}}{=} & \sum_{\alpha \in\{0,1\}^{d} \backslash\{\mathbf{0}\}}(-1)^{|\alpha|} \\
& \times \sum_{\mathbf{k} \in G_{\mathbf{n}-\mathbf{1}}^{|\alpha|,+}} \int_{\left[\left(n a_{n} h\right)^{-1}(\mathbf{k}-\mathbf{1})_{\alpha},\left(n a_{n} h\right)^{-1}(\mathbf{k})_{\alpha}\right]} \partial^{\alpha} K\left(x-\left(u:\left(a_{n} h\right)^{-1} \mathbf{1}\right)_{\alpha}\right)(\mathrm{d} u)_{\alpha}
\end{aligned}
$$

$$
\begin{aligned}
& \times B\left(\left(\left(n a_{n} h\right)^{-1} \mathbf{k}:\left(a_{n} h\right)^{-1} \mathbf{1}\right)_{\alpha}\right) \\
& +\Delta\left(K(x-\cdot) \cdot B(\cdot),\left[0,\left(a_{n} h\right)^{-1}\right]^{d}\right)
\end{aligned}
$$

$$
=Y_{n, 3}^{+}(x)+R_{n, \mathrm{SI}}(x)
$$

where the remainder $R_{n, \mathrm{SI}}(x)$ is defined in an obvious manner. From the modulus of continuity for the Brownian Sheet (see Khoshnevisan [20], Theorem 3.2.1) it follows that for $a, b \in \mathbb{R}^{d}$

$$
\begin{equation*}
\limsup _{\delta \rightarrow 0+} \sup _{s, t \in[a, b],\|s-t\|_{\infty}<\delta} \frac{|B(s)-B(t)|}{\sqrt{\delta \log (1 / \delta)}} \leq 24 \cdot d\|b\|_{\infty}^{d / 2} \tag{7.2}
\end{equation*}
$$

which yields

$$
\begin{aligned}
\left|Y_{n, 2}^{(+)}(x)-Y_{n, 3}^{(+)}(x)\right|= & \left|R_{n, \mathrm{SI}}(x)\right| \\
\leq & \sup _{\delta<1 / n} \sup _{s, t \in[0,2]^{d}:\|s-t\|_{\infty} \leq \delta}|B(s)-B(t)| \\
& \times \sqrt{\frac{\log (n)}{n}}\left[\int_{\left[0,(a-n h)^{-1}\right]^{d}}\left((u)^{1}\right)^{1 / 2}\left|\partial^{\mathbf{1}} K(x-u)\right| \mathrm{d} u+\mathrm{O}\left(h^{-(d-1) / 2}\right)\right]
\end{aligned}
$$

(note that the dominating term in $R_{n, \mathrm{SI}}(x)$ is given by the summand where $|\alpha|=d$ ). With the same arguments as in the proof of Lemma 2 we finally obtain

$$
\left|Y_{n, 2}^{(+)}(x)-Y_{n, 3}^{(+)}(x)\right|=\mathrm{O}_{P}\left(\sqrt{\frac{\ln \left(n a_{n} h\right)}{n h^{d}}}\right)
$$

where we used the estimate (7.2) for the modulus of continuity of the Brownian sheet (note that this estimate is independent of $x$ ).

### 7.3. Proof of Lemma 5

Integration by parts gives

$$
\begin{align*}
\Delta_{n, 3}:= & \left|Y_{n, 3}^{(+)}(x)-Y^{(+)}(x)\right| \\
\leq & \left|\int_{[0, \infty)^{d} \backslash\left[0,1 /\left(a_{n} h\right)\right]^{d}} B(u) \partial^{\mathbf{1}} K(x-u) \mathrm{d} u\right| \\
& +\mid \sum_{\alpha \in\{0,1\}^{d} \backslash\{\mathbf{0}, \mathbf{1}\}}(-1)^{|\alpha|}  \tag{7.3}\\
& \quad \times \int_{\left[0,1 /\left(a_{n} h\right)\right]^{|\alpha|}} B\left(\left(u:\left(a_{n} h\right)^{-1} \mathbf{1}\right)_{\alpha}\right) \partial^{\alpha} K\left(x-\left(u:\left(a_{n} h\right)^{-1} \mathbf{1}\right)_{\alpha}\right)(\mathrm{d} u)_{\alpha} \mid
\end{align*}
$$

$$
\begin{aligned}
& +\left|\Delta\left(K(x-\cdot) B(\cdot) ;\left[0,\left(a_{n} h\right)^{-1}\right]^{d}\right)\right| \\
:= & \left|\Delta_{n, 3}^{(1)}(x)\right|+\left|\Delta_{n, 3}^{(2)}(x)\right|+\left|\Delta_{n, 3}^{(3)}(x)\right|,
\end{aligned}
$$

where the processes $\Delta_{n, 3}^{(j)}(x), j=1,2,3$ are defined in an obvious manner. Let $n$ be sufficiently large such that $\frac{1}{a_{n} h} \geq 1$ and $a_{n}<\frac{1}{2}$. Since $B(u)=0$ if $t_{j}=0$ for at least one index $j \in\{1, \ldots, d\}$ we have

$$
\begin{aligned}
\left|\Delta_{n, 3}^{(3)}(x)\right| & =\left|K\left(x-\left(a_{n} h\right)^{-1} \mathbf{1}\right) \cdot B\left(\left(a_{n} h\right)^{-1} \mathbf{1}\right)\right| \\
& =\sqrt{2 d\left(a_{n} h\right)^{-d} \ln \left(d \ln \left(\left(a_{n} h\right)^{-1}\right)\right)} \frac{\left|K\left(x-\left(a_{n} h\right)^{-1} \mathbf{1}\right)\right|\left|B\left(\left(a_{n} h\right)^{-1} \mathbf{1}\right)\right|}{\sqrt{2 d\left(a_{n} h\right)^{-d} \ln \left(d \ln \left(\left(a_{n} h\right)^{-1}\right)\right)}} .
\end{aligned}
$$

An application of the version of a law of the iterated logarithm given in Theorem 3 of Paranjape and Park [26] yields the estimate

$$
\begin{aligned}
\sup _{x \in\left[0, h^{-1}\right]}\left|\Delta_{n, 3}^{(3)}(x)\right| & =\mathrm{O}(1) \cdot \sqrt{2 d\left(a_{n} h\right)^{-d} \ln \left(d \ln \left(\left(a_{n} h\right)^{-1}\right)\right)} \sup _{x \in\left[0, h^{-1}\right]}\left|K\left(x-\left(a_{n} h\right)^{-1} \mathbf{1}\right)\right| \\
& \leq \mathrm{O}(1) \cdot \sqrt{2 d\left(a_{n} h\right)^{-d} \ln \left(d \ln \left(\left(a_{n} h\right)^{-1}\right)\right)} \sup _{v \leq a_{n}-1 /\left(a_{n} h\right)}|K(v)| \\
& =\mathrm{o}\left(\frac{1}{\log (n)}\right) \quad \text { a.s. }
\end{aligned}
$$

uniformly with respect to $x$.
To show that $\Delta_{n, 3}^{(2)}(x)$ and $\Delta_{n, 3}^{(1)}(x)$ are asymptotically negligible we also apply the LIL for the Brownian sheet. For each summand, say $\Delta_{n, 3, \alpha}^{(2)}$, in $\Delta_{n, 3}^{(2)}(x)(|\alpha|<d)$ we have

$$
\begin{aligned}
\Delta_{n, 3, \alpha}^{(2)}(x) & :=\left|\int_{\left[0,1 /\left(a_{n} h\right)\right]^{|\alpha|}} B\left(\left(u:\left(a_{n} h\right)^{-1} \mathbf{1}\right)_{\alpha}\right) \partial^{\alpha} K\left(x-\left(u:\left(a_{n} h\right)^{-1} \mathbf{1}\right)_{\alpha}\right)(\mathrm{d} u)_{\alpha}\right| \\
& =\left(a_{n} h\right)^{-|\alpha|}\left|\int_{[0,1]^{|\alpha|}} B\left((u: \mathbf{1})_{\alpha}\left(a_{n} h\right)^{-1}\right) \partial^{\alpha} K\left(x-(u: \mathbf{1})_{\alpha}\left(a_{n} h\right)^{-1}\right)(\mathrm{d} u)_{\alpha}\right|
\end{aligned}
$$

Scaling of the Brownian sheet yields

$$
\begin{aligned}
\Delta_{n, 3, \alpha}^{(2)}(x) & \stackrel{\mathcal{D}}{=}\left(a_{n} h\right)^{-(2|\alpha|+d) / 2}\left|\int_{[0,1]^{|\alpha|}} B\left((u: \mathbf{1})_{\alpha}\right) \partial^{\alpha} K\left(x-(u: \mathbf{1})_{\alpha}\left(a_{n} h\right)^{-1}\right)(\mathrm{d} u)_{\alpha}\right| \\
& =\mathrm{O}\left(\left(a_{n} h\right)^{-d / 2}\right)\left|\int_{\left[0,1 /\left(a_{n} h\right)\right]^{|\alpha|}} \partial^{\alpha} K\left(x-(u: \mathbf{1})_{\alpha}\left(a_{n} h\right)^{-1}\right)(\mathrm{d} u)_{\alpha}\right| \quad \text { a.s. }
\end{aligned}
$$

With the same arguments as in the proof of Lemma 2 we conclude that the leading contributions are given by the quantities $\Delta_{n, 3, \alpha}^{(2)}(x)$, where $|\alpha|=d-1$. For $\alpha=(0,1, \ldots, 1)$ obtain

$$
\sup _{x \in\left[0, h^{-1}\right]^{d}}\left|\Delta_{n, 3, \alpha}^{(2)}(x)\right|=\mathrm{O}_{P}\left(\left(a_{n} h\right)^{-d / 2}\right) \sup _{v \leq-1 /\left(2 a_{n} h\right)} \int_{\mathbb{R}^{d-1}}\left|\partial^{\alpha} K\left(v, u_{2}, \ldots, u_{d}\right)\right| \mathrm{d}\left(u_{2}, \ldots, u_{d}\right)
$$

This gives $\sup _{x \in\left[0, h^{-1}\right]^{d}}\left|\Delta_{n, 3,(0,1, \ldots, 1)}^{(2)}(x)\right|=\mathrm{o}\left(\frac{1}{\log (n)}\right)$. Applying the same argument to the other terms yields $\Delta_{n, 3}^{(2)}(x)=\mathrm{o}_{P}\left(\frac{1}{\log (n)}\right)$ uniformly with respect to $x \in[0,1 / h]^{d}$. Finally, a similar $\operatorname{argument}$ gives for the remaining term in (7.3) $\Delta_{n, 3}^{(1)}(x)=\mathrm{o}_{P}\left(\frac{1}{\log (n)}\right)$, which completes the proof of Lemma 5.

### 7.4. Proof of Lemma 3

Note that we have

$$
Y_{n, 1}^{(+)}(x) \stackrel{\mathcal{D}}{=} \frac{h^{\beta}}{\sqrt{n^{d} a_{n}^{d} h^{d}}\|K\|_{L^{2}\left(\mathbb{R}^{d}\right)}} \sum_{\alpha \in\{0,1\}^{d}}(-1)^{|\alpha|} \sum_{\mathbf{j} \in G_{n}^{|\alpha|,+}} \Delta_{\alpha}\left(\overline{\mathcal{F} f_{p}} \circ \tau_{x} ; I_{\mathbf{j}}\right) B\left(\mathbf{j}:(\mathbf{n})_{1-\alpha}\right)
$$

The representation (2.7) and the definition (6.6) yield

$$
\begin{aligned}
& \left|Y_{n, 1}^{(+)}(x)-Y_{n, 2}^{+}(x)\right| \\
& \quad=\sum_{p=1}^{M-1} \frac{h^{\mu_{p}}}{\sqrt{n^{d} a_{n}^{d} h^{d}}}\left|\sum_{\alpha \in\{0,1\}^{d}}(-1)^{|\alpha|} \sum_{\mathbf{j} \in G_{\mathbf{n}}^{|\alpha|,+}} \Delta_{\alpha}\left(\overline{\mathcal{F} f_{p}} \circ \tau_{x} ; I_{\mathbf{j}}\right) B\left(\mathbf{j}:(\mathbf{n})_{\mathbf{1}-\alpha}\right)\right|+\mathrm{o}_{P}\left(\frac{1}{\log (n)}\right) .
\end{aligned}
$$

For each fixed $p$ we can now perform the approximation steps of the previous lemmas and obtain

$$
\begin{aligned}
& \log (n) \sup _{x \in\left[0, h^{-1}\right]^{d}} \left\lvert\, \frac{1}{\sqrt{n^{d} a_{n}^{d} h^{d}}} \sum_{\alpha \in\{0,1\}^{d}}(-1)^{|\alpha|} \sum_{\mathbf{j} \in G_{\mathbf{n}}^{|\alpha|,+}} \Delta_{\alpha}\left(\overline{\mathcal{F} f_{p}} \circ \tau_{x} ; I_{\mathbf{j}}\right) B\left(\mathbf{j}:(\mathbf{n})_{\mathbf{1}-\alpha}\right)\right. \\
&-\int_{\mathbb{R}_{+}^{d}} \overline{\mathcal{F} f_{p}}(x-u) \mathrm{d} B(u) \mid=\mathrm{o}_{P}(1)
\end{aligned}
$$

It can easily be shown that for all $p=1, \ldots, M-1$

$$
\lim _{n \rightarrow \infty} n^{d} a_{n}^{d} h^{d} \operatorname{Var}\left(\frac{1}{n^{d} a_{n}^{d} h^{d}} \sum_{\mathbf{k} \in G_{\mathbf{n}}} Y_{\mathbf{k}} \overline{\mathcal{F} f_{p}}\left(\left(x-x_{\mathbf{k}}\right) \frac{1}{h}\right)\right)=\sigma^{2}\left\|f_{p}\right\|_{2}^{2}
$$

where the limit does not depend on $x$. We finally obtain, repeating the approximation steps given in the previous lemmas for each of the $2^{d}-1$ remaining orthants

$$
\begin{aligned}
& \log (n) \sup _{x \in\left[0, h^{-1}\right]^{d}} \left\lvert\, \frac{1}{\sqrt{n^{d} a_{n}^{d} h^{d}}} \sum_{\alpha, \gamma \in\{0,1\}^{d}}(-1)^{|\alpha|} \sum_{\mathbf{j} \in G_{\mathbf{n}}^{|\alpha|,+}} \Delta_{\alpha}\left(\overline{\mathcal{F} f_{p}} \circ \tau_{x} ;(-1)^{\gamma} I_{\mathbf{j}}\right) B\left((-1)^{\gamma} \mathbf{j}:(\mathbf{n})_{\mathbf{1}-\alpha}\right)\right. \\
&-\int_{\mathbb{R}^{d}} \overline{\mathcal{F} f_{p}}(x-u) \mathrm{d} B(u) \mid=\mathrm{o}_{P}(1)
\end{aligned}
$$

Note that
$r(x-z):=\mathbb{E}\left(\int_{\mathbb{R}^{d}} \overline{\mathcal{F} f_{p}}(x-u) \mathrm{d} B(u) \int_{\mathbb{R}^{d}} \mathcal{F} f_{p}(z-u) \mathrm{d} B(u)\right)=\int_{\mathbb{R}^{d}} f_{p}(x-z+u) f_{p}(u) \mathrm{d} u$
and $\|r\|_{1} \leq\left\|f_{p}\right\|^{2}<\infty$. The system of sets $\left\{\left[0, h^{-1}\right]^{d} \mid n \in \mathbb{N}\right\}$ is a blowing up system of sets in the sense of Definition 14.1 in Piterbarg [27]. If we define

$$
Z_{p}(x)=\frac{1}{\left\|f_{p}\right\|_{L^{2}\left(\mathbb{R}^{d}\right)}} \int_{\mathbb{R}^{d}} \overline{\mathcal{F} f_{p}}(x-u) \mathrm{d} B(u)
$$

then Theorem 1 in Bickel and Rosenblatt [3] gives the asymptotic independence of the scaled minimum and maximum of the process $Z_{p}$, which, with the observation that $Z_{p}$ and $-Z_{p}$ have the same distribution and an application of Theorem 14.1 in Piterbarg [27] yields that for $G \sim \operatorname{Gumbel}(\ln (2), 1)$

$$
\sup _{x \in\left[0, h^{-1}\right]^{d}}\left(\left(\left|Z_{p}(x)\right|-\tilde{C}_{n, 3}\right) \tilde{C}_{n, 3}\right) \xrightarrow{\mathcal{D}} G \quad \text { for } n \rightarrow \infty
$$

where the constants $\tilde{C}_{1}, \tilde{C}_{n, 2}$ and $\tilde{C}_{n, 3}$ are given by

$$
\begin{aligned}
\tilde{C}_{1} & =\operatorname{det}\left(\left[\frac{1}{\left\|f_{p}\right\|_{L^{2}\left(\mathbb{R}^{d}\right)}^{2}} \int_{\mathbb{R}^{d}}\left|f_{p}(v)\right|^{2} v_{i} v_{j} \mathrm{~d} v\right], i, j=1, \ldots, d\right), \\
\tilde{C}_{n, 2} & =\sqrt{\frac{\tilde{C}_{1}}{(2 \pi)^{d+1}}} \frac{1}{h^{d}}, \\
\tilde{C}_{n, 3} & =\sqrt{2 \ln \left(\tilde{C}_{n, 2}\right)}+\frac{(d-1) \ln \left(2 \ln \left(\tilde{C}_{n, 2}\right)\right)}{2 \sqrt{2 \ln \left(\tilde{C}_{n, 2}\right)}} .
\end{aligned}
$$

Since $h^{\mu_{p}}=\mathrm{o}\left(\frac{1}{\log n}\right)$ we obtain $h^{\mu_{p}} \sup _{x \in\left[0, h^{-1}\right]}\left|Z_{p}(x)\right|=\mathrm{o}_{P}\left((\log n)^{-1 / 2}\right)$ for each $p=$ $1, \ldots, M-1$, which justifies the replacement of $h^{\beta} K_{n}$ by $K$. Since the outer sum does not depend on $n$ this gives the desired result.

### 7.5. Proof of Lemma 1

With the same arguments as in the proof of the previous lemmas we can replace the errors by combinations of partial sums and perform the same approximation steps. In each replacement we obtain at most a $d-1$-fold sum which yields the desired result right away.

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