Structural tests in additive regression

Wolfgang HÄRDLE, Stefan SPERLICH, and Vladimir SPOKOINY

We consider the component analysis problem for a regression model with an additive structure. The problem is to test if some of the additive components are of polynomial structure, e.g. linear, without specifying the structure of the remaining components. A particular case is the problem of selecting the significant covariates. The presented method is based on the wavelet transform using the Haar basis, which allows for applications under mild conditions on the design and smoothness of the regression function. The results demonstrate that each component of the model can be tested with the rate corresponding to the case if all the remaining components were known. The proposed procedure is also computationally straightforward. Simulation results and a real data example about female labor supply demonstrate the good performance of the test.

KEY WORDS: Additive models; Component analysis, Haar basis; Hypothesis testing; Nonparametric alternative; Regression

1. INTRODUCTION

In *multivariate regression problems* we study the structural relationship between the response variable Y and the vector of covariates $X = (X_1, \ldots, X_d)^\top$ via the regression curve

$$F(x) = E(Y|X = x)$$

with $x = (x_1, \ldots, x_d)^\top$. Purely nonparametric models do not make any assumption about the form of the d-variate function $F(x)$. The problem is then to fit a d-dimensional surface to the observed data $\{(X_i, Y_i) : i = 1, \ldots, n\}$. The obvious approach is to generalize the univariate smoothing techniques based on local 'averaging' to this multivariate situation. A serious problem arising here is that we need much more data material in higher dimensions in order to have enough data points in a local neighborhood of each point. Several approaches for dimensionality reduction have been proposed to deal with this so-called *curse of dimensionality*. A promising one is *additive modeling* as in economic theory it is a favorite structure anyway, see e.g. Deaton and Muellbauer (1980).

Such a nonparametric additive regression model has the form

$$
\begin{aligned}
y &= F(x) + \xi, & x = (x_1, \ldots, x_d) \in I\!\!R^d, & \quad (1) \\
F(x) &= f_1(x_1) + \ldots + f_d(x_d), & & \quad (2)
\end{aligned}
$$

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where y is a scalar variable, $\{f_m\}_{m=1}^d$ is a set of unknown component functions and ξ is a random error.

This class of models has been shown to be useful in statistical practice: it generalizes linear regression in a natural way and allows interpretation of marginal changes i.e. the effect of one variable on the mean function F holding all else constant. Additive models were considered first by Leontief (1947) for input-output analysis speaking of *separable* models. In the statistical literature the nonparametric additive regression has been introduced in the eighties, see Buja, Hastie and Tibshirani (1989) for a survey. An advantage of additive models is that they combine flexible modelling of many variables with statistical precision that is typical for just one explanatory variable, see Stone (1985, 1986). Algorithmic aspects of additive modelling by backfitting are discussed in Hastie and Tibshirani (1990) or Venables and Ripley (1994). Tjøstheim, Auestad (1994) and Linton, Nielsen (1995) proposed a method of analysis based on marginal integration.

An essential advantage in additive models is that they allow component-wise inferences. Important problems of component analysis in economics are the question of significance as well as of linearity, since nonlinearities often raise serious problems, e.g. of identification in equation or economic equilibrium systems. In nonparametrics, among others, Hastie and Tibshirani (1990) or Härdle and Korostelev (1996) considered also the problem of selection of significant covariates. In this paper we focus on the much more general problem of testing for component f_m the null hypothesis of it being of polynomial form, e.g., being constant or linear.

Theory for nonparametric hypothesis testing is well developed. So the problem of testing a simple null hypothesis versus a univariate nonparametric alternative is studied in

detail, see e.g. Ingster (1993), Härdle and Mammen (1993), Hart (1997), Stute (1997) for historical background and further references. Many tests have been shown to be sensitive against every directional local alternative, e.g. Bierens (1982), Eubank and Hart (1992), Stute (1997) and references therein. Unfortunately, the power of every particular test cannot be uniform w.r.t. the "direction" in the function space as shown in Burnashev (1979) or Ingster (1982). This particularly means that the finite sample power of every test could be better for some local alternatives and worse for the others. The same arguments apply to the so called "intermediate" efficiency approach of Inglot and Ledwina (1996).

Ingster (1982, 1993) has shown that a test could be uniformly consistent against a smooth alternative only if this alternative deviates from the null with the distance of order $n^{-2s/(4s+1)}$ with s being the degree of smoothness. The structure of the proposed rate-optimal tests also essentially relies on the smoothness properties of the underlying function though such kind of prior information about the underlying function is typically lacking in practical applications. Spokoiny (1996) offered an adaptive data-driven testing procedure which does not require knowledge of smoothness properties of the tested function and allow for a near optimal testing rate up to a $\log \log n$ factor. The latter can be viewed as the price for adaptation. Horowitz and Spokoiny (1999) proposed a similar test based on kernel smoothers with different bandwidths and showed that it is simultaneously consistent against any local "directional" alternative which approaches the null hypothesis at the rate $(n/\log\log n)^{-1/2}$.

It is worth noting that the adaptive testing procedure from Spokoiny (1996) is essentially a theoretical device rather than a practically applicable method since it is developed for the idealized "signal + white noise" model, simple null, known noise variance etc. Practically relevant procedures should address numerous issues arising in particular applications. In the context of multidimensional additive modeling, an additional challenge comes from the fact that the considered component f_1, even being completely specified, does not specify the whole model since nothing is assumed about the other components, $f_2, \ldots, f_d$ which can be viewed as an infinite-dimensional nuisance parameter. This particularly creates a serious problem with evaluating the critical value of the proposed test statistics which provides the prescribed type I level.

Therefore, the task is to develop a procedure which, independent of the functional form of the 'nuisance' components $f_2, \ldots, f_d$, leads to the given type I error α if f_1 is linear, and is sensitive against a smooth alternative with unknown degree of smoothness. In view of practical applications (see Section 4) we proceed with a deterministic

non-regular design allowing discrete components and with unknown noise variance.

In this paper we apply a Haar decomposition which is a particular and non-regular case of the wavelet transform. Nevertheless, for the hypothesis testing framework the application of the Haar basis leads not only to the desired optimal testing rate but also provides a test which is more stable w.r.t. the design non regularity. This is important for practical applications, allowing relaxation and simplification of the conditions on the design, reduction of computational burden and more.

Our approach is based on the simultaneous approximation of all components $f_1, \ldots, f_d$ by Haar sums: we first estimate the Haar coefficients for all components and then analyze the coefficients corresponding to the first one. The testing problem is formulated in the next section, the procedure is described in Section 2. The asymptotic properties are discussed in Section 3. The results demonstrate asymptotic optimality of the proposed procedure and they are stated under mild conditions on the design. Section 4 illustrates the finite sample performance of the test. In particular we present a comparative study of our test with the Eubank and Hart (1992) test in the one-dimensional case and with the ideal one designed for the case as if the other components and all smoothness properties were known. An application of the test to real data (study of the female labor supply in East Germany) is thoroughly discussed in Section 4.2. Extensions to more general problems including model check of additivity and multiple testing of several components simultaneously are shortly discussed in Section 3.3 and the proofs are postponed to Section 5.

1.1 Model and testing problem

We are given data (X_i, Y_i), $i = 1, \ldots, n$, with $X_i \in \mathbb{R}^d$, $Y_i \in \mathbb{R}^1$, obeying the regression equation

$$Y_i = F(X_i) + \xi_i \qquad (3)$$

where F is an unknown regression function with the additive structure

$$F(x) = f_1(x_1) + \ldots + f_d(x_d), \qquad (4)$$

and ξ_i are normal random errors with zero mean and known variance σ^2. We allow for deterministic non-regular design $X_1, \ldots, X_n$ with possible replications. It is only supposed that the design is rescaled to the unit cube $[0,1]^d$, that is, $X_{i,\ell} \in [0,1]$ for all $i \le n$ and $\ell \le d$.

Our aim is to analyze each component f_m, $m = 1, \ldots, d$. For simplicity we present the procedure focusing on the first component f_1, and on the problem of testing linearity, i.e. the hypothesis $H_0 : f_1(t) = a_1 + b_1 t$ for some constants a_1, b_1.

Let ϕ be a test, a measurable function of observations with values 0 (accept) and 1 (reject). Denote by $\boldsymbol{P}_F$

the distribution of the data $Y_1, \ldots, Y_n$ for a fixed model function F, see (3) and (4). Let now F_0 be a function with a linear first component. The type I error probability is the probability under F_0 to reject the hypothesis: $\alpha_{F_0}(\phi) = \boldsymbol{P}_{F_0}(\phi = 1)$. Similarly one defines the error probability $\beta_F(\phi)$ of the second type. If the first component f_1 is not linear, then $\beta_F(\phi) = \boldsymbol{P}_F(\phi = 0)$. Given $\alpha > 0$, we wish to construct such a test ϕ that $\alpha_{F_0}(\phi) \leq \alpha$ for all F_0 with a linear first component and, in addition, it is sensitive against a large class of alternatives F.

2. TESTING PROCEDURE

In order to illustrate the main ideas, we begin with the univariate case i.e. $d = 1$.

2.1 The case of $d = 1$

Consider the univariate regression model

$$Y_i = f(X_i) + \xi_i, \qquad i = 1, \ldots, n, \tag{5}$$

which corresponds to (3) with $d = 1$. We write here f instead of f_1 to minimize the notation. The problem consists in testing the hypothesis that the function f is linear.

Eubank and Hart (1992) nicely pointed out a common feature of many procedures for model checking. Let $\mathcal{F}_0$ be the set of regression functions considered under the null hypothesis (here the linear functions). Then f is written as $f(x, \theta_0) + \sum_j \theta_j \psi_j(x)$ with $f(x, \theta_0)$ a member of $\mathcal{F}_0$ and $\{\psi_j\}$ an orthonormal system. The testing problem reduces now to testing $\theta_j = 0$ for all j, cf. also Stute (1997).

The procedure proposed here follows this idea and relies on a special piecewise constant approximation (the Haar decomposition) of the function f.

Denote by I the multi-index $I = (j, k)$ with $j = 1, 2, \ldots$ and $k = 0, 1, \ldots, 2^j - 1$, and by $\mathcal{I}$, the set of all such multi-indices. Let now the function $\psi(t)$ (the mother wavelet) be defined by

$$\psi(t) = \begin{cases} 0, & t < 0, \, t \geq 1, \\ 1, & 0 \leq t < 1/2, \\ -1, & 1/2 \leq t < 1. \end{cases}$$

For every $I = (j, k)$ with $j \geq 0$ and $k = 0, \ldots, 2^j - 1$ set

$$h_I(t) = \psi(2^j t - k).$$

Clearly the function ψ_I with $I = (j, k)$ is supported on the interval $A_I = [2^{-j}k, 2^{-j}(k+1)]$. Denote also by $\mathcal{I}_j$ the index subset corresponding to the j-th resolution level:

$$\mathcal{I}_j = \{I = (j, k), \, k = 0, 1, \ldots, 2^j - 1\} \qquad j \geq 0.$$

The idea of the test is to estimate from the data the coefficients c_I of the approximation of the unknown regression

function f by the sum

$$c_0 + c_1 x + \sum_{\ell=0}^{j} \sum_{I \in \mathcal{I}_\ell} c_I h_I(x)$$

and then to check whether some of estimated coefficients c_I differ significantly from zero.

For a formal description, define with $I = (j, k) \in \mathcal{I}$

$$\begin{aligned} \mu_I^2 &= \sum_{i=1}^{n} h_I^2(X_i), \\ \psi_I(X_i) &= \mu_I^{-1} h_I(X_i). \end{aligned}$$

Clearly μ_I^2 is the number of design points in A_I, that is, $\mu_I^2 = \#\{i : X_i \in A_I\}$, $I \in \mathcal{I}$.

We also define two functions $\psi_0 \equiv \mu_0^{-1}$ and $\psi_1(t) = \mu_1^{-1} t$ with $\mu_0^2 = n$ and $\mu_1^2 = \sum_{i=1}^{n} X_i^2$ and introduce the index set

$$\mathcal{I}(j) = \{0, 1\} + \bigcup_{\ell=0}^{j} \mathcal{I}_\ell. \tag{6}$$

By $N(j)$ we denote the number of indices in $\mathcal{I}(j)$. Obviously $N(j) = 2^{j+1} + 1$. Let $\boldsymbol{\theta}(j)$ denote a vector in $\mathbb{R}^{N(j)}$ with entries θ_I, $I \in \mathcal{I}(j)$. Define the vector $\widehat{\boldsymbol{\theta}}(j)$ as solution to the quadratic problem

$$\widehat{\boldsymbol{\theta}}(j) = \operatorname*{arginf}_{\{\boldsymbol{\theta}(j) \in \mathbb{R}^{N(j)}\}} \sum_{i=1}^{n} \left(Y_i - \sum_{I \in \mathcal{I}(j)} \theta_I \psi_I(X_i) \right)^2.$$

To get an explicit expression for $\widehat{\boldsymbol{\theta}}(j)$ we introduce vector notation. Let g be a function observed at point $X_1, \ldots, X_n$. We identify every such function with the column-vector $\boldsymbol{g}$ in $\mathbb{R}^n$ with the entries $g(X_i)$ and define $\|\boldsymbol{g}\|_n$ by $\|\boldsymbol{g}\|_n^2 = \sum_{i=1}^{n} g^2(X_i)$. Let also $\boldsymbol{Y}$ stand for the column vector $(Y_1, \ldots, Y_n)^\top$. Introduce a $n \times N(j)$-matrix $\Psi(j)$ with entries $\psi_I(X_I)$:

$$\Psi(j) = \Big(\psi_I(X_i), \, i = 1, \ldots, n, \, I \in \mathcal{I}(j) \Big).$$

Then

$$\begin{aligned} \widehat{\boldsymbol{\theta}}(j) &= \operatorname*{arginf}_{\boldsymbol{\theta}(j) \in \mathbb{R}^{N(j)}} \|\boldsymbol{Y} - \Psi(j)\boldsymbol{\theta}(j)\|_n^2 \\ &= V(j)\Psi(j)^\top \boldsymbol{Y} = W(j)^\top \boldsymbol{Y}, \end{aligned}$$

where $V(j)$ is the pseudo-inverse of $\Psi(j)^\top \Psi(j)$, that is, $V(j) = \{\Psi(j)^\top \Psi(j)\}^-$ and $W(j) = \Psi(j)V(j)$ is a $n \times N(j)$-matrix.

Since the errors ξ_i are normal $\mathcal{N}(0, \sigma^2)$, we obtain by (5) that $\widehat{\boldsymbol{\theta}}(j)$ is a Gaussian vector with the mean $\boldsymbol{\theta}^*(j) = W(j)^\top \boldsymbol{f} = V(j)\Psi(j)^\top \boldsymbol{f}$ and the covariance matrix $\sigma^2 V(j)$,

$$\widehat{\boldsymbol{\theta}}(j) \sim \mathcal{N}\left\{ \boldsymbol{\theta}^*(j), \sigma^2 V(j) \right\}.$$

The entries of the matrix $V(j)$ (resp. $W(j)$) will be denoted by $v_{I,I'}$ (resp. $w_{i,I}$) where $I, I' \in \mathcal{I}$ and $i = 1, \ldots, n$. All these values depend on j, but do not indicate this dependence explicitly to simplify the notation.

By $\widehat{\boldsymbol{\theta}}_j$ we denote the part of the vector $\widehat{\boldsymbol{\theta}}(j)$ corresponding to j-th resolution level: $\widehat{\boldsymbol{\theta}}_j = (\widehat{\theta}_I, I \in \mathcal{I}_j)^\top$, so that $\widehat{\boldsymbol{\theta}}_j \in R^{2^j}$. Obviously $\widehat{\boldsymbol{\theta}}_j = W_j^\top \boldsymbol{Y}$ where W_j is the $n \times 2^j$-submatrix of $W(j)$ corresponding to the index set $\mathcal{I}_j$: $W_j = (w_{i,I}, i = 1, \ldots, n, I \in \mathcal{I}_j)$. Similarly we define the 2^j-vector $\boldsymbol{\theta}_j^*$ and $2^j \times 2^j$-submatrix V_j of $V(j)$:

$$\boldsymbol{\theta}_j^* = (\theta_I^*, I \in \mathcal{I}_j), \qquad V_j = (v_{I,I'}, I, I' \in \mathcal{I}_j).$$

Clearly $\widehat{\boldsymbol{\theta}}_j \sim \mathcal{N}(\boldsymbol{\theta}_j^*, \sigma^2 V_j)$ and $V_j = W_j^\top W_j$.

2.2   Level test statistic for $d = 1$

The proposed testing procedure is based on the fact that for f linear, all the empirical coefficients $\widehat{\theta}_I$, $I \neq 0, 1$, are zero mean Gaussian r.v.'s. We build for every j one test statistic corresponding to the hypothesis $\boldsymbol{\theta}_j^* = 0$.

By definition $\widehat{\boldsymbol{\theta}}_j = W_j^\top \boldsymbol{Y}$ which yields $\widehat{\boldsymbol{\theta}}_j \sim \mathcal{N}(\boldsymbol{\theta}_j^*, \sigma^2 V_j)$ with $V_j = W_j^\top W_j$. This naturally leads to the likelihood-based statistic $S_j = \widehat{\boldsymbol{\theta}}_j^\top V_j^- \widehat{\boldsymbol{\theta}}_j$ where V_j^- means the pseudo-inverse of V_j. Under the null hypothesis (that is, for a linear function f), it clearly holds $\boldsymbol{\theta}_j^* = 0$ and $\widehat{\boldsymbol{\theta}}_j = W_j^\top \boldsymbol{\xi}$, and hence,

$$S_j = \boldsymbol{\xi}^\top W_j V_j^- W_j^\top \boldsymbol{\xi} = \boldsymbol{\xi}^\top \mathcal{R}_j \boldsymbol{\xi} \tag{7}$$

where $\mathcal{R}_j = W_j V_j^- W_j^\top = W_j \left(W_j^\top W_j\right)^- W_j^\top$ is a projector in the space $I\!R^n$ (that is, $\mathcal{R}_j^2 = \mathcal{R}_j$). By N_j we denote the rank of $\mathcal{R}_j$. By definition $N_j \leq 2^j$. The definition (7) particularly yields that $\sigma^{-2} S_j$ follows the χ^2-distribution with N_j degrees of freedom.

The level test statistic T_j is defined via centering and standardization of S_j. The following simple properties are useful here:

$$\boldsymbol{E} S_j = \boldsymbol{E} \boldsymbol{\xi}^\top \mathcal{R}_j \boldsymbol{\xi} = \sigma^2 \, \mathrm{tr} \mathcal{R}_j = \sigma^2 N_j,$$
$$\boldsymbol{E} \left(S_j - \sigma^2 N_j\right)^2 = 2\sigma^4 N_j.$$

Since the noise variance σ^2 is usually unknown, we replace it by a pilot estimate $\widehat{\sigma}^2$, see Section 2.5 below. This leads to the test statistic T_j of the form:

$$T_j = \frac{\widehat{\boldsymbol{\theta}}_j^\top V_j^- \widehat{\boldsymbol{\theta}}_j - \widehat{\sigma}^2 N_j}{\widehat{\sigma}^2 \sqrt{2N_j}} \tag{8}$$

An important feature of this statistic is that under the null hypothesis, it has a nondegenerate distribution. (Which approaches the standard normal law as N_j grows.) Moreover, this distribution is known (see Section 2.6 for a closed form expression) which allows precise evaluation of the cor-

responding $(1 - \alpha)$-quantile $t_{j,\alpha}$ defined by

$$\boldsymbol{P}_0 \left(T_j > t_{j,\alpha}\right) = \alpha, \tag{9}$$

where $\boldsymbol{P}_0$ denotes the distribution of T_j under the null hypothesis.

2.3   A multiscale test for $d = 1$

The proposed test analyzes all statistics T_j for different j simultaneously. Similar ideas are discussed extensively in the literature. Eubank and Hart (1992) proposed the so called "order selection" test using a modified Mallows' criterion (Mallows, 1973) for selecting the number of considered terms of an orthogonal series expansion for the deviation of the underlying function f from the null hypothesis; see also Aerts, Claeskens and Hart (1999). This method leads to the maximum of $\sigma^{-2} S_j - (1 + C_n) N_j$ where $(1 + C_n) N_j$ is the penalty term for going to a more complicated model. A similar test, called the data-driven Neyman's smooth test is proposed in Ledwina (1994) and Kallenberg and Ledwina (1995). Fan (1996), Spokoiny (1996) and Fan and Huang (1998) considered the tests based on the maximum of centered and standardized sums like T_j. Such a test has a strong appeal: the results from Ingster (1982, 1993) show that the test T_j with a special choice of the index j is rate optimal against a smooth alternative with a smoothness degree s. The test based on the maximum of T_j is adaptive in the sense that it is near optimal against a smooth alternative with unknown degree of smoothness.

Here we apply the method based on the multiscaling idea (see Section 2.4 for further discussion) which is close to the proposal from Fan (1996) and Spokoiny (1996): the test statistic T^* is defined as the maximum of $T_j - t_{j,\alpha}$ over all considered levels j with $t_{j,\alpha}$ from (9). Namely, we consider all j from $j = 0$ until the finest resolution level j_n defined as

$$j_n = [\log_2(n/3)]$$

where $[a]$ means the integer part of a. We now set

$$T^* = \max_{0 \leq j \leq j_n} \left(T_j - t_{j,\alpha}\right).$$

A choice of the critical value for this test is discussed in Section 2.6.

2.4   A multiscale test for $d > 1$

The basic idea of testing is similar to the univariate case and it is based on the approximation of each component f_m from (3) by the sum

$$c_{1,m} x_m + \sum_{j=0}^{j_n} \sum_{I \in \mathcal{I}_j} c_{I,m} h_I(x_m), \qquad m = 1, \ldots, d.$$

(We skip here the constant term to provide identifiability of each component.) Let us fix a level j for the first component and a level j_n for the remaining ones, and let $\mathcal{I}(j) = \{0, 1\} + \bigcup_{0 < \ell \le j} \mathcal{I}_\ell$. We also define $\mathcal{I}'(j) = \{1\} + \bigcup_{0 \le \ell \le j} \mathcal{I}_\ell$. To define the level test, we approximate $F(x)$ by

$$\sum_{I \in \mathcal{I}(j)} c_{I,1} h_I(x_1) + \sum_{m=2}^{d} \sum_{I \in \mathcal{I}'(j_n)} c_{I,m} h_I(x_m).$$

Here $N = 2^{j_n+1}$ coefficients are used for each component f_m, $m \ge 2$, and, assuming that $j \le j_n$, the total number of coefficients is at most $Nd + 1$. We modify now the definition of j_n from the one-dimensional case to provide $Nd + 1 \le 2n/3$ that leads to the choice

$$j_n = \left[\log_2 \left(\frac{n}{3d} \right) \right]. \tag{10}$$

To define the test, we first standardize each basis function:

$$\psi_{I,m}(t) = \mu_{I,m}^{-1} h_I(t) \text{ with } \mu_{I,m}^2 = \sum_{i=1}^{n} h_I^2(X_{i,m})$$

Here $(X_{i,1}, \dots, X_{i,d})$ is the coordinate representation of X_i.

Now let some $j \le j_n$ be fixed. Denote by $\mathcal{I}(d, j)$ the index set

$$\mathcal{I}(d, j) = \left\{ (I, 1), I \in \mathcal{I}(j) \right\} \times \prod_{m=2}^{d} \left\{ (I, m), I \in \mathcal{I}'(j_n) \right\}$$

and let

$$N(d, j) = N(j) + (d-1)N = 2^{j+1} + (d-1)2^{j_n+1} + 1$$

be the number of elements in $\mathcal{I}(d, j)$.

Set $\Psi(d, j)$ for the $n \times N(d, j)$ matrix with entries $\psi_{I,m}(X_i) = \mu_{I,m}^{-1} h_I(X_{i,m})$, $i = 1, \dots, n$, $(I, m) \in \mathcal{I}(d, j)$, and define the vector $\widehat{\boldsymbol{\theta}}(d, j)$ in $\mathbb{R}^{N(d,j)}$ as a solution to the quadratic problem:

$$\begin{aligned} \widehat{\boldsymbol{\theta}}(d, j) &= \underset{\boldsymbol{\theta}(d,j) \in \mathbb{R}^{N(d,j)}}{\text{arginf}} \| \boldsymbol{Y} - \Psi(d, j)\boldsymbol{\theta}(d, j) \|_n^2 \\ &= \underset{\boldsymbol{\theta}(d,j) \in \mathbb{R}^{N(d,j)}}{\text{arginf}} \sum_{i=1}^{n} \Bigg(Y_i - \sum_{I \in \mathcal{I}(j)} \theta_{I,1} \psi_{I,1}(X_{i,1}) \\ &\quad - \sum_{m=2}^{d} \sum_{I \in \mathcal{I}'(j_n)} \theta_{I,m} \psi_{I,m}(X_{i,m}) \Bigg)^2. \end{aligned}$$

As in the univariate case, we derive

$$\widehat{\boldsymbol{\theta}}(d, j) = V(d, j)\Psi(d, j)^\top \boldsymbol{Y} = W(d, j)^\top \boldsymbol{Y} \tag{11}$$

where the matrix $V(d, j)$ is the pseudo-inverse of $\Psi(d, j)^\top \Psi(d, j)$, i.e. $V(d, j) = \left\{ \Psi(d, j)^\top \Psi(d, j) \right\}^-$ and $W(d, j) = \Psi(d, j)V(d, j)$. The entries of the matrix

$V(d, j)$ (resp. $W(d, j)$) will be denoted by $v_{(I,m),(I',m')}$ (resp. $w_{i,(I,m)}$).

Similarly to the univariate case, we define the level test making use of the subvector $\widehat{\boldsymbol{\theta}}_j = (\widehat{\theta}_{I,1}, I \in \mathcal{I}_j)$ and the submatrix $V_j = (v_{(I,1),(I',1)}, I, I' \in \mathcal{I}_j)$ of the covariance matrix $V(d, j)$. Let W_j again denote the submatrix of $W(d, j)$ corresponding to the level j of the first component: $W_j = (w_{i,(I,1)}, i = 1, \dots, n, I \in \mathcal{I}_j)$. Then clearly $\widehat{\boldsymbol{\theta}}_j = W_j^\top \boldsymbol{Y}$ and $V_j = W_j^\top W_j$. The test statistic T_j is defined as follows, cf. (8):

$$T_j = \frac{\widehat{\boldsymbol{\theta}}_j^\top V_j^- \widehat{\boldsymbol{\theta}}_j - \widehat{\sigma}^2 N_j}{\widehat{\sigma}^2 \sqrt{2N_j}} = \frac{\boldsymbol{Y}^\top \mathcal{R}_j \boldsymbol{Y} - \widehat{\sigma}^2 N_j}{\widehat{\sigma}^2 \sqrt{2N_j}},$$

where $\mathcal{R}_j = W_j V_j^- W_j^\top$ and N_j is the rank of $\mathcal{R}_j$ (or equivalently of V_j), $N_j \le 2^j$. With $t_{j,\alpha}$ fulfilling (9), the final test statistic is again of the form:

$$T^* = \max_{0 \le j \le j_n} (T_j - t_{j,\alpha}). \tag{12}$$

Remark 2.1. In some practical applications, see e.g. our example in Section 4, one or more explanatory variables X_m can be discrete with only a few possible values, say two or three. In that case the corresponding component function f_m is completely determined by its values at these points and it can be precisely expanded by a finite Haar sum with very few Haar levels. Of course, for such situations it is not reasonable to consider all j_n Haar levels for those components and the required number of levels for every particular component f_m should be determined by the identifiability reasons, see Section 4 for an example.

2.5 Estimation of the noise variance

Here we indicate how the noise variance σ^2 can be estimated from the data. One may apply two different approaches for variance estimation. One way is based on residuals from locally polynomial fitting, see e.g. Rice (1984) or Gasser et al (1986) for the univariate case or Hall et al (1991) and Spokoiny (1999b) and references therein for a detailed discussion of the multivariate case. Another approach is to retrieve the residuals from the same orthogonal series expansion which is used for model checking. Here we follow the latter proposal.

Let j_n be defined in (10). Due to this definition we have $n/3 \le d2^{j_n+1} \le 2n/3$.

Let $\widehat{\boldsymbol{\theta}}(d, j_n)$ be the least squares estimator from (11) with $j = j_n$, that is, the maximal number of Haar coefficients are used for all components f_m. This vector is Gaussian with the mean $\boldsymbol{\theta}^*(d, j_n) = W(d, j_n)\boldsymbol{F}$ and the covariance matrix $\sigma^2 V(d, j_n)$. Moreover, $\Psi(d, j_n)\widehat{\boldsymbol{\theta}}(d, j_n) = \Pi_n Y$ where

$$\Pi_n = \Psi(d, j_n)\left(\Psi(d, j_n)^\top \Psi(d, j_n)\right)^- \Psi^T(d, j_n)$$

is the projector in $I\!R^n$ on the subspace generated by additive functions of the form

$$\theta_0 + \sum_{m=1}^{d} \sum_{I \in \mathcal{I}'(j_n)} \theta_{I,m} \psi_{I,m}(x_m).$$

One can easily check that

$$\begin{aligned}
\boldsymbol{E}\|\boldsymbol{Y} - \Psi(d, j_n)\widehat{\boldsymbol{\theta}}(d, j_n)\|_n^2 \\
&= \|\boldsymbol{F} - \Pi_n \boldsymbol{F}\|_n^2 + \boldsymbol{E}\|\boldsymbol{\xi} - \Pi_n \boldsymbol{\xi}\|_2^2 \\
&= \|\boldsymbol{F} - \Pi_n \boldsymbol{F}\|_n^2 + \sigma^2 \operatorname{tr}(I_n - \Pi_n) \\
&= \|\boldsymbol{F} - \Pi_n \boldsymbol{F}\|_n^2 + \sigma^2(n - r_n)
\end{aligned}$$

where I_n denotes the identity $n \times n$-matrix and r_n is the rank of Π_n. By definition $r_n \leq 2n/3$.

Under regularity conditions on the function F, see e.g. Lemma 1 in the next section, the accuracy of approximating F by such an expansion tends to zero as n tends to infinity in the sense that

$$n^{-1}\|\boldsymbol{F} - \Pi_n \boldsymbol{F}\|_n^2 \to 0, \qquad n \to \infty.$$

This consideration prompts one to use the value

$$\widehat{\sigma}^2 = \frac{1}{n - r_n}\|\boldsymbol{Y} - \Psi(d, j_n)\widehat{\boldsymbol{\theta}}(d, j_n)\|_n^2$$

for estimating σ^2. It is important to mention that if $F \equiv 0$, then $(n - r_n)\widehat{\sigma}^2 = \|\boldsymbol{\xi} - \Pi_n \boldsymbol{\xi}\|_n^2$ follows the χ^2-distribution with $n - r_n$ degrees of freedom and $\widehat{\sigma}^2$ and $\widehat{\boldsymbol{\theta}}(d, j_n)$ are independent.

2.6   Critical level of the test

First we again discuss the univariate situation with $d = 1$. In that case the function F coincides with the first component f_1 and its structure is known under the null hypothesis. Moreover, in view of the method of approximation, the linear trend in f_1 has no influence on the remaining coefficients and we may assume that the function f_1 is exactly zero. The same applies to the variance estimate $\widehat{\sigma}^2$. This reduces the linear hypothesis to the case of a simple null hypothesis $f_1 \equiv 0$, that is, the observations Y_i coincide with the noise ξ_i. In this situation one has $S_j = \boldsymbol{\xi}^\top \mathcal{R}_j \boldsymbol{\xi}$, $\widehat{\sigma}^2 = (n-r_n)^{-1}\|(1\!\!1 - \Pi_n)\boldsymbol{\xi}\|_n^2 = (n - r_n)^{-1}\boldsymbol{\xi}^\top(1\!\!1 - \Pi_n)\boldsymbol{\xi}$, where $\mathcal{R}_j = W_j V_j^- W_j^\top$ and $1\!\!1$ denotes the unit operator in $I\!R^n$ and the test statistics T_j can be represented in the form

$$\begin{aligned}
T_j &= \frac{\boldsymbol{\xi}^\top \mathcal{R}_j \boldsymbol{\xi} - \widehat{\sigma}^2 N_j}{\widehat{\sigma}^2 \sqrt{2N_j}} \\
&= \frac{\boldsymbol{\xi}^\top \mathcal{R}_j \boldsymbol{\xi}(n - r_n)}{\boldsymbol{\xi}^\top(1\!\!1 - \Pi_n)\boldsymbol{\xi}\sqrt{2N_j}} - \sqrt{N_j/2}.
\end{aligned} \tag{13}$$

Therefore, each T_j is the ratio of two quadratic forms of ξ_i's and as a consequence, it does not depend on the noise variance and its distribution can be precisely described via

the Fisher distribution $F_{N_j, n - r_n}$ with N_j and $n - r_n$ degrees of freedom. The values $t_{j,\alpha}$ defined in (9) can therefore be calculated using the proper quantile of this Fisher distribution and they depend only on N_j, $n - r_n$ and α. Since all the T_j's are constructed on the base of the same data, they are dependent in a rather complicated way and hence, the closed form expression for the distribution of the maximum $T^* = \max_{j \leq j_n}(T_j - t_{j,\alpha})$ is difficult to obtain. Therefore, some Monte-Carlo experiments can be used for finding a proper quantile λ satisfying $\boldsymbol{P}_0(T^* > \lambda) = \alpha$, where $\boldsymbol{P}_0$ means that each T_j follows (13) with a standard Gaussian vector $\boldsymbol{\xi}$. Having done this, we define the test ϕ^* as

$$\phi^* = \mathbf{1}(T^* > \lambda). \tag{14}$$

For the general multivariate case one can show that under some regularity conditions (see Condition (D) in the next section) the influence of the remaining components $f_2, \ldots, f_m$ on the test statistic T^* is asymptotically negligible and we therefore determine the critical value λ in the same way using simulated data from the d-variate model with the same regression design and with the vanishing regression function and standard Gaussian errors. For further discussion and details concerning this Monte-Carlo method, see Section 4 or Spokoiny (1999a).

Remark 2.2.   Note that the adjustment of T_j by $t_{j,\alpha}$ is more of theoretical nature allowing for the unified exposition. Indeed, all the test statistics T_j have non-degenerate distribution with the variance 1 and moreover, for large j, this distribution is close to the standard normal CDF. This means that all the $t_{j,\alpha}$'s are of the same order and the effect of this adjustment is negligible. This issue is confirmed by our simulation results, see Section 4.

3.   MAIN RESULTS

In this section we present asymptotic properties of the proposed testing procedure. We state the results on the type one and type two error probabilities separately since we evaluate them under different assumptions on the design variables. The result on the type I error probabilities $\alpha_{F_0}(\phi^*)$ is valid under mild assumptions on the design. But for high sensitivity of the test, we need slightly stronger regularity conditions on the design variables. We begin by describing the required assumptions on the model.

3.1   Assumptions

When testing the first component of the function F from (4), the remaining components $f_2, \ldots, f_d$ can be viewed as nonparametrically specified nuisance parameters which are to be estimated by a pilot estimator. In order to ensure the required accuracy of estimation, we need some

conditions on the rate of approximation of each function f_m with $2 \leq m \leq d$ by the Haar series. We formulate these conditions exactly in the required form. Later we show that these conditions are met, for instance, under mild conditions on smoothness of f_m and on the design $X_1, \ldots, X_n$.

Recall that we identify every function g on $\mathbb{R}^d$ with the vector $\boldsymbol{g} = (g(X_i), i = 1, \ldots, n)^\top$ in $\mathbb{R}^n$. In particular, each component f_m is identified with the vector $\boldsymbol{f}_m = (f_m(X_{i,m}), i = 1, \ldots, n)^\top$ and $\boldsymbol{\psi}_{I,m}$ is understood as the vector with the elements $\psi_{I,m}(X_{i,m}) = \lambda_{I,m}^{-1} h_I(X_{i,m})$. Recall also the notation $\|\boldsymbol{g}\|_n^2 = \sum_{i=1}^n g^2(X_i)$.

Denote by $\mathcal{L}_m(j)$ the linear subspace in $\mathbb{R}^n$ generated by the functions (vectors) $\{\boldsymbol{\psi}_{I,m}\}$, $I \in \mathcal{I}_\ell$, $0 \leq \ell \leq j$,

$$\mathcal{L}_m(j) = \left\{ \theta_{0,m} + \theta_{1,m}\boldsymbol{\psi}_{1,m} + \sum_{\ell=0}^j \sum_{I \in \mathcal{I}_\ell} \theta_{I,m}\boldsymbol{\psi}_{I,m} \right\}.$$

Clearly all the functions (or vectors) from $\mathcal{L}_m(j)$ depend only on m-th coordinates $X_{i,m}$ of design points X_i, $i = 1, \ldots, n$. By $\Pi_{m,n}\boldsymbol{f}_m$ we denote the projection of $\boldsymbol{f}_m$ onto $\mathcal{L}_m(j_n)$ w.r.t. the distance $\|\cdot\|_n$,

$$\Pi_{m,n}\boldsymbol{f}_m = \operatorname*{arginf}_{\boldsymbol{g} \in \mathcal{L}_m(j_n)} \|\boldsymbol{f}_m - \boldsymbol{g}\|_n$$

$$= \operatorname*{arginf}_{\boldsymbol{g} \in \mathcal{L}_m(j_n)} \sum_{i=1}^n |f_m(X_{i,m}) - g(X_{i,m})|^2.$$

In our results we impose the following condition:

Condition (D) *For some fixed constant C and n large*

$$\sum_{m=1}^d \|\boldsymbol{f}_m - \Pi_{m,n}\boldsymbol{f}_m\|_n \leq C\sigma n^{-1/2}.$$

The following lemma shows that condition (D) is satisfied under mild smoothness conditions on each component f_m.

Lemma 1. *Let $\mu_{n,m}$ be the m-th marginal of the empirical design measure μ_n,*

$$\mu_{n,m}(A) = n^{-1} \sum_{i=1}^n \mathbf{1}(X_{i,m} \in A), \qquad m = 1, \ldots, d.$$

Let further C_1 be a constant such that for every $0 \leq a < b \leq 1$ with $b - a > 1/n$, it holds

$$\mu_{n,m}[a,b] \leq C_1(b-a).$$

If each f_m, $m = 2, \ldots, d$, is a Lipschitz function i.e.

$$|f_m(x) - f_m(x')| \leq C_2|x - x'|, \qquad \forall x, x' \in [0,1],$$

then condition (D) is fulfilled with C depending on σ, C_1 and C_2 only.

Another situation in which the difference $\|\boldsymbol{f}_m - \Pi_{m,n}\boldsymbol{f}_m\|_n$ can be easily controlled, is the case of a discrete m-th component (i.e. when all $X_{i,m}$ belong to some

finite set). In that case, the value $\|\boldsymbol{f}_m - \Pi_{m,n}\boldsymbol{f}_m\|_n$ is zero provided that n is large enough.

3.2 Asymptotic properties of the test

Let ϕ^* be the test introduced above in (14).

Theorem 1. *Suppose that the observations (X_i, Y_i), $i = 1, \ldots, n$, obey the regression model (3) and (4), and let condition (D) hold. If the first component f_1 of the function F is linear, then*

$$\boldsymbol{P}_F(\phi^* = 1) \leq \alpha + \delta_1(n),$$

where $\delta_1(n)$ tends to 0 as $n \to \infty$ and depends only on n and constant C arising in condition (D).

The proof of the theorem is given in Section 5.

We now turn to the results concerning the sensitivity of test ϕ^*.

The first assertion shows under which conditions we reject an alternative with a high probability.

Theorem 2. *Let the function F in model (3) be of the form (4). Let also $\boldsymbol{\theta}_j^* = (\theta_{I,1}^*, I \in \mathcal{I}_j)$ be the subvector of the vector $\boldsymbol{\theta}^*(d, j)$ corresponding to j-th resolution level of the first component and let $V_j = (v_{(I,1),(I',1)}, I, I' \in \mathcal{I}_j)$ be the covariance submatrix corresponding this index set. If, for some $j \leq j_n$, $\beta > 0$ and $c > 0$, we have*

$$T_j^* \equiv (2N_j)^{-1/2}\sigma^{-2}\boldsymbol{\theta}_j^{*\top}V_j^-\boldsymbol{\theta}_j^* > t_{j,\alpha} + \lambda + t_{j,\beta}' + c\sqrt{\log j_n},$$

with $t_{j,\beta}'$ defined by the equality $\boldsymbol{P}(T_j^0 < -t_{j,\beta}') = \beta$, then

$$\boldsymbol{P}_F(\phi^* = 0) \leq \beta + \delta_1(n)$$

where $\delta_1(n)$ depends on α, β and c only and $\delta_1(n) \to 0$, for $n \to \infty$.

Remark 3.1. This result claims that the test ϕ^* rejects with a probability close to 1 any alternative for which at least one of T_j^* exceeds $c'\sqrt{\log j_n}$ with some fixed constant c'. Therefore, we may suppose that the error of the second type may occur only if

$$T_j^* \leq c'\sqrt{\log j_n}, \qquad 0 \leq j \leq j_n. \tag{15}$$

Now we discuss how this statement can be transformed into the result about the uniform rate of testing. Following Ingster (1982, 1993) we consider the class of alternatives with the first component f_1 separated from the null (the set of the linear functions) with distance at least ϱ,

$$\inf_{a,b} \|f_1(\cdot) - a - b \cdot\| \geq \varrho$$

where $\|\cdot\|$ means the usual L_2-norm, and in addition we assume that f_1 is smooth in the sense that f_1 belongs to

some class of functions $\mathcal{F}$. Ingster (1982) established his results assuming that the underlying function belongs to a Hölder or L_2-Sobolev ball $\mathcal{F}$, Spokoiny (1998) studied the case of a more general L_p-Sobolev ball with any $p \geq 1$.

We are interested in a minimal separation distance ρ which still allows for a uniform testing. To state the result we need some regularity conditions on the design and smoothness conditions on the first component f_1. The reason why stronger conditions on the design are required can be explained by the fact that a degenerate design leads to an identification problem: the components cannot be separated and therefore it is impossible to make any inference about them. Set

$$
\begin{aligned}
u_*(j) &= \inf_{I \in \mathcal{I}_j} 2^j M_I / n, \\
u^*(j) &= \sup_{I \in \mathcal{I}_j} 2^j M_I / n,
\end{aligned}
$$

where, given $I = (\ell, k)$, the value M_I stands for the number of design points X_i whose first component belongs to the interval $A_I = [k2^{-\ell}, (k+1)2^{-\ell})$, that is, $M_I = \#\{i : X_{i,1} \in A_I\}$. *Design regularity* means in particular that $u_*(j)$ is bounded away from zero i.e. each interval A_I contains enough design points $X_{i,1}$, cf. the condition in Lemma 1.

Recall the notation $V_j = \left(v_{(I,1),(I',1)}, I, I' \in \mathcal{I}_j\right)$ and N_j denotes its rank, $N_j \leq 2^j$. Set

$$
v^*(j) = \|V_j\|.
$$

Here, the norm $\|A\|$ of a matrix A is understood as the maximal eigenvalue of this matrix. We understand design regularity in the sense that V_j is non-degenerate and all the $v^*(j)$'s are bounded.

Finally, given an integer s, suppose that the function f_1 is s times differentiable and the value

$$
r_s^2 = \int_0^1 |f_1^{(s)}(x)|^2 dx
$$

is finite, where $f_1^{(s)}$ means the s-th derivative of f_1.

Theorem 3.   Let condition (D) hold. Suppose there exists an integer s and for some $j \leq j_n$, the first component f_1 of the model function F satisfies the following inequality

$$
\inf_{a,b} \|f_1 - a - b\psi_{1,1}\|_n^2 \geq
$$
$$
C_1 \, r_s^2 \, n \, 2^{-js} + C_2 \frac{u^*(j)}{u_*(j)} v^*(j) \, 2^{j/2} \sigma^2 \sqrt{\log j_n}
$$

with $\psi_{1,1}(x) = x_1$ and constants C_1 and C_2 depending on s only, then

$$
\boldsymbol{P}_F(\phi^* = 0) \leq \delta_1(n) \to 0, \qquad n \to \infty,
$$

with $\delta_1(n)$ as in Theorem 2.

The proof of this assertion is based only on (15) and can be found in Härdle et al. (1997) or Spokoiny (1999a).

Remark 3.2.   By minimizing the sum of the form $B_1 n 2^{-js} + B_2 2^{j/2} \sigma^2 \sqrt{\log \log n}$ with fixed B_1 and B_2 with respect to j we find that a smooth alternative will be rejected with a high probability if

$$
\inf_{a,b} n^{-1} \|f_1 - a - b\psi_{1,1}\|_n^2 \geq B_3 \left(\frac{n}{\sigma^2 \sqrt{\log \log n}}\right)^{-\frac{2s}{4s+1}}
$$

for a constant B_3 depending on B_1 and B_2 only. Spokoiny (1996) has shown that this rate is optimal in the problem of testing against a smooth alternative with an unknown degree of smoothness s.

3.3 Extensions

Here we briefly discuss possible extensions of the test which we introduced previously.

3.3.1   Testing additivity

Though our test was constructed for testing functional forms of the additive components, it can also be useful when the presence of interaction is at question. Often, the additive structure is given or wanted by the economic theory the particular model is based on, see e.g. Deaton and Muellbauer (1980) or also our application in Section 4. However, not only from a statistical point of view it is interesting to scrutinize this assumption in some cases. Several approaches of testing additivity are discussed in Hart (1997), but nonparametric theory for this problem is quite recent, see e.g. Sperlich, Tjøstheim and Yang (1999), also for more references.

As was pointed out at the beginning, our procedure can test for presence of a component. Thus, for testing of no interaction one can proceed as follows. Introduce artificial covariates $X_{m,m'} = X_m X_{m'}$ for $m \neq m'$. No interaction between X_m and $X_{m'}$ means that the covariate $X_{m,m'}$ has no effect, which is a particular case of the problem we considered before.

3.3.2   Non-Gaussian errors

In our results we suppose Gaussian homoskedastic noise with unknown dispersion σ^2. This assumption allows simplification of the calculations and highlights the main ideas, skipping a lot of technical details which appear when considering non-Gaussian noise. However, the results from Section 3 apply to i.i.d. errors with unknown distribution under some moment conditions. We refer to Spokoiny (1999a) for the analysis of non-Gaussian noise in the univariate case. An extension to the multivariate situation is straightforward.

3.3.3   Multiple testing

The above test was developed for testing one component of an additive model. In practice one could also be

interested in testing all the components of the model simultaneously. This leads to a multiple testing problem which requires a more careful evaluation of the corresponding critical values. Following the rule proposed in Section 2 one can construct for every component f_m the corresponding test statistic T_m^* and calculate the corresponding critical level λ_m. Now we apply the same idea of multiscale testing as one used for construction of every component test. Namely, to provide a prescribed nominal level α of the multiple test, which checks all components f_m simultaneously, all these critical values λ_m should be slightly increased, e.g. by the same value $\Delta\lambda$ such that

$$\boldsymbol{P}_0\left(\max_{m=1,\dots,d}(T_m^* - \lambda_m) > \Delta\lambda\right) \leq \alpha$$

where $\boldsymbol{P}_0$ means the distribution on the space of observations under the model with $F \equiv 0$ and with the standard Gaussian errors (that is, Y_i are i.i.d. standard Gaussian).

3.3.4 Local test

In parallel to the test T^* based on the maximum of some quadratic forms of the empirical Haar coefficients $\widehat{\theta}_I$, one may consider another test which is called the "local" test in Härdle, Sperlich and Spokoiny (1997). This test is based on the maximum of the standardized empirical coefficients $\widehat{\theta}_I$ over all $I \in \mathcal{I}_j$. More precisely, for every $j \leq j_n$, we define

$$T_{j,\mathrm{loc}} = \max_{I \in \mathcal{I}_j} \frac{\widehat{\theta}_I^2}{\widehat{\sigma}_I^2} - \tau_j$$

where $\widehat{\sigma}_I^2 = \widehat{\sigma}^2 v_{(I,1),(I,1)}$ and τ_j are such that

$$\boldsymbol{P}_0\left(T_{j,\mathrm{loc}} > \tau_j\right) = \alpha_{\mathrm{loc}}$$

with $\boldsymbol{P}_0$ being again the distribution under the no-response model with standard normal errors. The multi-level "local" test ϕ_{loc}^* is defined by

$$\phi_{\mathrm{loc}}^* = \boldsymbol{1}\left(\max_{j \leq j_n} T_{j,\mathrm{loc}} > \tau^*\right)$$

where τ^* fulfills

$$\boldsymbol{P}_0\left(\max_{j \leq j_n} T_{j,\mathrm{loc}} > \tau^*\right) = \alpha_{\mathrm{loc}}.$$

For applications one can use an approximation $\tau_j \approx 2\log N_j - 2\log\log N_j + 2\log\alpha_{\mathrm{loc}}^{-1}$. Such defined "local" test has been shown to be sensitive against a "non-smooth" alternative (e.g. an alternative with jumps), see Härdle et al. (1997). In practical applications one would be willing to apply both tests T^* and T_{loc}^* simultaneously which requires some additional adjustment of the critical levels for both tests. Taking into account the specific structure of the test ϕ_{loc}^*, our recommendation is to perform this "local"

test at a very small significance level, e.g. $\alpha_{\mathrm{loc}} = 0.005$ or even smaller which does not require an additional adjustment of the test ϕ^*.

Also the theoretical properties of such defined test are presented and discussed in Härdle et al. (1997).

4. SIMULATION STUDIES AND AN APPLICATION

The performance of the suggested test procedure for finite samples was examined in a simulation study. Then we apply the procedure to the analysis of female labor supply data. The goal of the simulation study was to illustrate the performance of the test for different smoothness properties of the investigated function, impact of non-normally distributed error terms, and to observe the (relative) power of the test against smooth alternatives.

4.1 Some simulated examples

We considered 3-dimensional regression problems

$$Y = m(x) + \xi \tag{16}$$

$x = (x_1, x_2, x_3)^T$, with a function m having additive components taken from the following set of functions with different smoothness properties:

$$f_1(x) = 2\sin(\pi x) \quad , \quad f_2(x) = 2\sin(2\pi x),$$
$$f_3(x) = 2\sin(3\pi x), \quad \text{and} \quad f_4(x) = x^2 .$$

Note that the indices $\gamma = 1, 2, 3, 4$ of the f_γ in this section refer to the functional form and not to their ordering. For investigating the level and the power of the test, we consider the following three specific models:

$$\begin{aligned} m_1(x;v) &= (1-v)x_1 + vf_2(x_1) + f_1(x_2) + f_4(x_3), \\ m_2(x;v) &= (1-v)x_1 + vf_3(x_1) + f_1(x_2) + f_4(x_3), \\ m_3(x;v) &= (1-v)x_1 + vf_4(x_1) + f_1(x_2) + f_3(x_3), \end{aligned}$$

each time testing the linearity of the first component $(1-v)x_1 + vf_\gamma(x_1)$ with v running from zero to one. This parameter v has the same meaning as the separation distance between the null and the alternative.

The explanatory variables were always uniformly distributed on the cube $[-2, 2]^3$. Unless stated otherwise, the sample size was set to $n = 150$ and the error term standard normal. We did not assume to know the standard deviation but estimated σ as suggested in Section 2.5 and got, as expected, only slightly overestimated $\widehat{\sigma}_\xi$ (5 to 15%). For getting the critical values we applied 249 Monte-Carlo replications to economize on the computational time. However, some examples were conducted with a larger number of replicates and the results are very similar. For practical applications, more precise critical values can be expected when resampling 499 or even 999 times.

Further, as discussed in Remark 2.6, we have to decide how to choose $t_{j,\alpha}$. We present all results for the two most

natural choices. First, we set $t_{j,\alpha}$ equal to the $F_{N_j, n-r_n}(\alpha)$-quantile with $\alpha = .01, .05, .1$ being the significance level (in tables indicated by $F(\alpha)$); second, we tried our procedure with simply $t_{j,\alpha} \equiv 0$ (in tables indicated by $'0'$).

All calculations are done in GAUSS, graphics in XploRe. The results after 500 runs can be found in Table 1, together with the average over the resolution levels at which our null is rejected. The latter delivers some qualitative information to the practitioner about the frequency where the violation from the linear null occurs.

The tested (first) component of the mean regression function has different smoothness in these three examples.

The results demonstrate a lower power of the test for less smooth first component which confirms the theoretical issues. Smoothness of the x_2 and x_3 functions has no strong influence on the results. One also can see that the resolution level at which the procedure rejects the null, clearly depends on the smoothness of the first component as well as on the distance between null and alternative. It can be seen, that only looking at one special level would reduce a lot the power of our procedure. All the numerical results are completely in agreement with the theoretical investigations from Section 3.

Table 1. Percentage of rejections and average of active resolution level j_1 (underlined) for functions $(1 - v)x_1 + vf_\gamma(x_1)$, $\gamma = 2, 3, 4$ in model $m_k(x; v)$, $k = 1, 2, 3$. In left column v running from 0 to 0.8. Results given for test with $t_{j,\alpha} = F(\alpha)$ and for test with $'0'$.

$\alpha =$	f_2, m_1						f_3, m_2						f_4, m_3					
	.01		.05		.10		.01		.05		.10		.01		.05		.10	
$t_{j,\alpha} =$	$'0'$	$F(\alpha)$	$'0'$	$F(\alpha)$	$'0'$	$F(\alpha)$	$'0'$	$F(\alpha)$	$'0'$	$F(\alpha)$	$'0'$	$F(\alpha)$	$'0'$	$F(\alpha)$	$'0'$	$F(\alpha)$	$'0'$	$F(\alpha)$
0.0	.012	.010	.040	.044	.082	.080	.004	.010	.048	.050	.106	.102	.008	.010	.050	.062	.108	.094
	1.00	2.60	2.25	2.59	1.88	2.55	2.00	3.40	1.88	3.08	2.11	2.49	1.50	3.60	2.48	3.06	2.13	2.72
0.1	.016	.014	.072	.050	.126	.106	.004	.006	.058	.060	.110	.104	.022	.020	.078	.062	.132	.128
	1.38	2.71	1.06	2.00	1.38	1.83	2.00	2.00	2.14	2.67	2.04	2.44	1.73	3.20	1.49	1.94	1.53	1.94
0.2	.060	.028	.212	.160	.338	.278	.028	.024	.128	.128	.224	.232	.058	.044	.170	.140	.288	.244
	1.00	1.79	1.11	1.49	1.10	1.47	2.07	2.83	1.72	2.58	1.97	2.36	1.28	2.05	1.31	1.69	1.31	1.52
0.3	.208	.130	.474	.394	.596	.550	.066	.060	.246	.290	.372	.396	.188	.108	.426	.352	.556	.510
	0.83	1.20	0.97	1.14	0.99	1.13	2.18	2.77	2.35	2.72	2.42	2.73	1.04	1.13	1.14	1.20	1.43	1.23
0.4	.476	.362	.748	.704	.824	.786	.210	.236	.464	.502	.620	.642	.376	.284	.634	.582	.728	.704
	0.92	1.05	0.96	1.07	0.99	1.06	2.26	2.80	2.33	2.69	2.41	2.63	1.04	1.16	1.08	1.15	1.11	1.14
0.5	.756	.638	.926	.898	.972	.960	.398	.454	.712	.738	.820	.838	.628	.542	.860	.814	.918	.896
	0.90	1.03	0.92	0.97	0.93	0.97	2.34	2.89	2.40	2.75	2.42	2.63	1.04	1.09	1.08	1.10	1.11	1.08
0.6	.932	.864	.984	.984	.998	.996	.596	.692	.878	.912	.950	.958	.864	.810	.970	.956	.982	.980
	0.88	0.95	0.89	0.93	0.89	0.92	2.64	2.94	2.67	2.82	2.68	2.77	1.02	1.05	1.02	1.04	1.02	1.04
0.7	.970	.962	.998	.994	1.00	1.00	.832	.878	.970	.974	.988	.986	.964	.936	.990	.984	1.00	.998
	0.91	0.97	0.92	0.95	0.92	0.93	2.79	2.97	2.77	2.89	2.77	2.85	1.02	1.03	1.02	1.02	1.02	1.02
0.8	.994	.992	1.00	1.00	1.00	1.00	.912	.946	.990	.990	.996	.998	.998	.988	1.00	1.00	1.00	1.00
	0.92	0.95	0.92	0.94	0.92	0.93	2.72	2.94	2.68	2.85	2.68	2.78	1.01	1.02	1.01	1.01	1.01	1.01

The presented results in Table 1 for two different tests are very similar giving a slight advantage to the choice $t_{j,\alpha} \equiv 0$. We therefore consider only this choice in the sequel.

So far all simulations were done generating the data with standard normal errors. As we use also the normal distribution in the Monte-Carlo method for estimating the critical value of the test, it is of interest to check for a notable loss of power if the underlying error distribution is non-normal. We therefore examined the test performance when the errors are from the centered and standardized χ^2_{df}-distribution with $df = 5, 10$ (and for comparison ∞). The simulations were done for the model (16) with $m(x) = (1 - v)x_1 + vf_1(x_1) + f_3(x_2) + f_4(x_3)$. The

results are given in Table 2 and show an astonishing stable performance of the test w.r.t the different error distributions.

Table 2. Percentage of rejections for the function $(1 - v)x_1 + vf_1(x_1) + f_3(x_2) + f_4(x_3)$ (v in left column) when the errors are χ^2, respectively normal distributed.

$\xi \sim$	χ_5^2			χ_{10}^2			χ_∞^2		
$\alpha =$	.01	.05	.10	.01	.05	.10	.01	.05	.10
0.0	.008	.058	.100	.006	.042	.102	.011	.057	.104
0.1	.028	.102	.152	.016	.082	.144	.019	.093	.150
0.2	.102	.210	.312	.104	.222	.314	.108	.263	.356
0.3	.290	.534	.640	.302	.530	.648	.282	.506	.620
0.4	.546	.806	.874	.530	.774	.830	.580	.800	.855
0.5	.808	.918	.958	.802	.934	.976	.798	.930	.962
0.6	.944	.984	.994	.930	.984	.992	.925	.985	.992
0.7	.980	.994	1.00	.980	.996	.998	.985	1.00	1.00
0.8	.996	1.00	1.00	.998	1.00	1.00	.998	1.00	1.00

Next we compared the performance of our procedure with the ideal ("oracle") parametric t-test (or Neymann-Pearson NP), see below, for the sample sizes 150 and 300. This gives us an idea about the relative efficiency of the test. Here, t-test means testing the hypothesis $H_0 : \beta_2 = 0$ in the model

$$Y = \beta_0 + \beta_1 x_1 + \beta_2 f_1(x_1) + \beta_3 f_3(x_2) + \beta_4 f_4(x_3) + \xi$$

but with known functions f_1, f_3, f_4.

Figure 1 is giving the power functions of our wavelet test and the ideal t-test for the 5% significance level. They demonstrate how fast the power of our procedure increases and the separation distance between the null and the alternative decreases for an increasing number of observations.
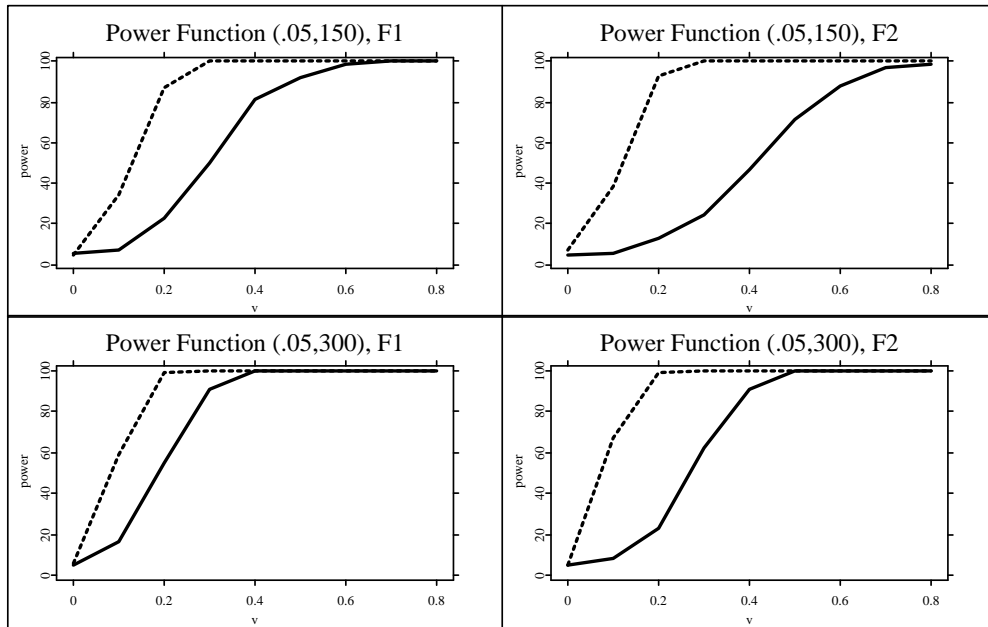

Figure 1. The power functions for $n = 150$ (upper) and $n = 300$ (lower), testing linearity of the first component of $m_k(x; v)$, for $k = 1$ on the left, $k = 2$ on the right with 5% significance level. Solid line is for the wavelet test, dotted line if using t-test with all information about functional forms.

For one dimensional models there exists a huge variety of nonparametric methods to test functional forms such as linearity, see Section 2.3. Although our procedure has been developed for component analysis in additive models, it could be of interest to know how well it does in the one dimensional case compared to existing methods. We chose the linearity test of Eubank and Hart (1992), see their Example 2 (p.1416) for exactly the problem of testing linearity. Along their recommendation we applied the variance estimator of Hall, Kay and Titterington (1990)

and used a polynomial basis to compute the residuals. We considered the model $m(x) = (1-v)x+vf_1(x)$. The results, given in Table 3 indicate that our method is competitive even in the special one dimensional case. Note, however, that our test is using Monte-Carlo-Methods whereas the Eubank/Hart test is based on asymptotic results. Consequently, the latter has strong computational advantages but can underperform in such small samples as $n = 150$.

Table 3. Percentage of rejections for the function $m(x) = (1 - v)x + vf_1(x)$ (v in left column) for our wavelet test with $t_{j,\alpha} = F(\alpha)$ and Eubank/Hart test with $n = 150$.

$\alpha =$	Wavelet test			Eubank/Hart test		
	.01	.05	.10	.01	.05	.10
0.0	.005	.033	.086	.024	.064	.120
0.1	.030	.167	.237	.016	.078	.190
0.2	.235	.449	.564	.068	.344	.530
0.3	.656	.876	.938	.402	.802	.908
0.4	.945	.991	.998	.822	.982	.994
0.5	.993	1.00	1.00	.986	1.00	1.00
0.6	1.00	1.00	1.00	1.00	1.00	1.00

4.2 Applications

We now turn to an application to illustrate the method with real data. The data set is a subsample of the Socio Economic Panel of Germany from 1992. To study the female labor supply in East Germany, 607 women with job and living together with a partner in East Germany have been asked for their weekly number of working hours, Y_i. The following values have been chosen as explanatory variables: the age of the woman X_1, her earnings per hour X_2, the prestige index of her kind of profession X_3 (called "Treimann Prestige Index", see Treimann, 1978), the monthly rent or redemption X_4 for their apartment or house, the monthly net income of her partner (in most cases her husband) X_5, her education X_6 measured in years, the unemployment rate X_7 of the particular country of the Federal Republic of Germany where the woman is living in and the number of children younger than 16 years, X_8. These data have already been analyzed nonparametrically, see e.g. Sperlich (1998) for functional forms and further discussion.

The study of female labor supply is quite common in economic empirical research and usually done with (log-) linear additive models using indicators for which the above mentioned set of variables is typical, see e.g. Mroz (1987), Blundell, Duncan, Meghir (1998), Strøm and Wagenhals (1991) or any Handbook of Labor Economics. Among them, especially Mroz (1987) investigated the sensitivity against model specification in this context and found tremendous differences in results depending on the particular specification. A first natural generalization would be to allow the additive components to be non-(log-)linear. Exactly this we now want to investigate. Later, we will additionally have a look at the additivity assumption.

Since some of these explanatory variables are not only discrete but have even less then 10 different values observed (e.g. for children and unemployment rate of coun-

try - there are only five countries in East Germany), we have to choose respectively low $j_k(n)$ (the highest considered wavelet level for k-th component) to avoid overparametrization in this component, see Remark 3. Therefore we chose $j_6(n) = 2$ for X_6, $j_7(n) = 2$ for X_7, and for X_8 (number of children) only $j_8(n) = 1$. For X_1, X_2, X_3, X_4 and X_5 we chose $j_n = 4$. These are also the functions of interest we want to analyze.

An analysis of the residuals from our nonparametric fit indicates that the variance of the residual does not really differ over the range of every explanatory variables which justifies the assumption of a homogeneous noise for this particular application.

In Figure 2 we have displayed the wavelet coefficient estimates used for the test statistics. They are standardized, i.e. divided by $v_I \hat{\sigma}$, but not corrected for the correlation inside the levels $j_1 = 0, 1, 2, 3, 4$. The length is indicating their absolute value.

Often, the earnings per hour (X_2) are modeled log-linear rather than linear by some reasoning from economic theory. So in a second run we also want to test the influence of $\ln(X_2)$ instead of using X_2 against linearity and gave the coefficients in the lower right.

Taking into account the construction of test statistics T_j, T^*, Figure 2 gives some ideas where we would expect the test to reject the null: e.g. for "earnings per hour" at $j_1 = 0$, "prestige" at $j_1 = 4$, and "log(earnings per hour)" at $j_1 = 2$, whereas it is not that clear for "age", "income of partner" or "rent/redemption". For the latter one we even would guess that there is no significance in the coefficients.

The wavelet test delivers the following results: the linearity hypothesis was rejected for X_1 ("age") at only 10% (at $j_1 = 2$), X_2 ("earnings per hour") at 1% (at $j_1 = 0$), for X_3 ("prestige") at (almost) 5% (p-value$\approx$ 0.052, at $j_1 = 4$), and for $\ln(X_2)$ also at the 1% significance level (at $j_1 = 2$).

Though the additivity assumption is not of prime interest for us, we finally also looked for possible second order interactions between the regressors. We applied the procedure described in Section 3.1 to all combinations $X_j X_k$, $k \neq j$, $j, k = 1, \ldots 8$. As before, we chose $j_6(n) = j_7(n) = 2$, $j_8(n) = 1$. It turned out, that the null hypothesis no interaction between "age" and "prestige" is rejected at 1%, "prestige" and "years of education" exactly at 5%, and between "age" and "u-rate" and "earnings of husband" and "u-rate" at the 10% level. Hence, our testing procedure enabled us to detect that the underlying data are inconsistent with the classic female labor supply model assumptions concerning the function form. Including now the two interactions "age"/"prestige" and "prestige"/"years of education" repeated testing all one dimensional components on linearity, we now got the following p-values: for "age" 0.68, "hourly earnings" 0.004, "ln(hourly earnings)" 0.008, and "prestige" 0.208.

Certainly, the same study could be done for any higher order interaction, or it could be applied to constructing a general test of additivity. This, however, lies beyond the scope of our illustration.
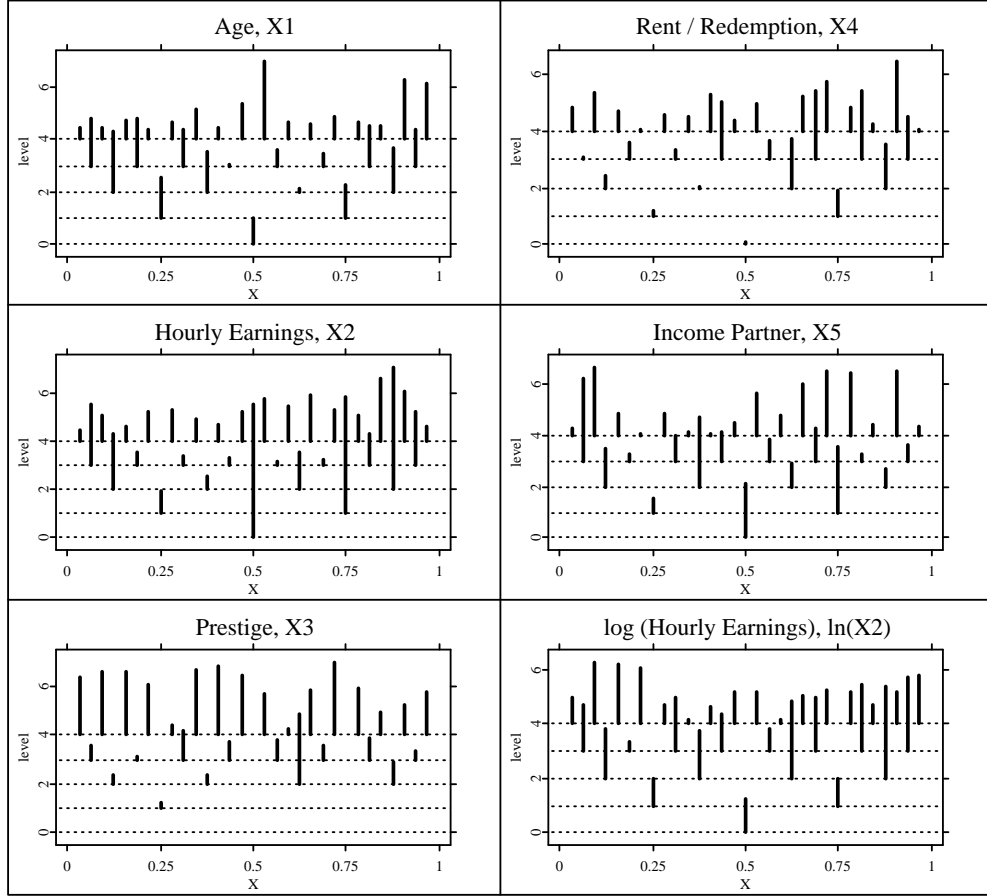

Figure 2. The estimated wavelet coefficients for some additive component functions. The coefficients θ_I are first standardized by $v_I \widehat{\sigma}$; the length indicates size. Coefficients with no length are zero. Resolution levels j_1 are given at vertical axes. The range of the regressors is normed to $[0,1]$ in which the coefficients are positioned in the center of the support of its corresponding mother wavelet ψ, e.g. for $j = 0$ on 0.5. Not given are the coefficients for the constant nor for the linear term since they do not enter to the test statistic.

5. PROOFS

In this section we collect the proofs of Theorems 1 through 2 and of the other statements presented in Section 3.

5.1 Proof of Lemma 1

By definition of j_n it holds $3dn^{-1}/2 \leq 2^{-j_n-1} \leq 3dn^{-1}$. Next, define $\Pi'_{m,n}$ as the projector on the space of piecewise constant functions of the m-th component x_m with the piece length 2^{-j_n-1}. Since $\Pi_{m,n}$ projects on a larger space generated by piecewise constant functions and the linear function $\psi_{1,m}(x) = x_m$, it clearly holds $\|\boldsymbol{f}_m - \Pi_{m,n}\boldsymbol{f}_m\|_n \leq \|\boldsymbol{f}_m - \Pi'_{m,n}\boldsymbol{f}_m\|_n$.

Let A be an interval of the form $A = [k2^{-j_n-1}, (k+1)2^{-j_n-1})$ and let $N_{A,m}$ denote the number of design points X_i with $X_{i,m} \in A$. The condition of the lemma on the marginals $\mu_{m,n}$ of the empirical measure implies that $N_{A,m} \leq C_1 2^{-j_n-1}n$. Denote also by $f_{m,A}$ the arithmetic mean of the values $f_m(X_i)$ over all X_i with $X_{i,m} \in A$. Then $\Pi'_{m,n}f_m(X_i) = f_{m,A}$ and the Lipschitz condition on the component functions f_m yields $|f_m(X_{i,m}) - f_{m,A}| \leq C_2 2^{-j_n-1}$ for $X_{i,m} \in A$ and hence

$$\sum_{i:X_{i,m}\in A} |f_m(X_i) - f_{m,A}|^2 \leq N_{A,m} \left| C_2 2^{-j_n-1} \right|^2$$
$$\leq C_1 C_2^2 n 2^{-3j_n-3}.$$

We have 2^{j_n+1} such intervals and therefore

$$\|\boldsymbol{f}_m - \Pi_{m,n}\boldsymbol{f}_m\|_n^2 \leq C_1 C_2^2 n 2^{-2j_n-2} \leq C_1 C_2^2 d^2 n^{-1}$$

and the assertion follows.

5.2   Some properties of the variance estimate

It is well known that under mild regularity conditions, the unknown variance σ^2 can be estimated at the rate $n^{-1/2}$. We now show that the proposed estimate $\widehat{\sigma}^2$ is also root-n consistent under the condition (D).

The estimate $\widehat{\sigma}^2$ can be represented in the form

$$
\begin{aligned}
\widehat{\sigma}^2 &= (n - r_n)^{-1} \boldsymbol{Y}^\top (\mathbb{I} - \Pi_n) \boldsymbol{Y} \\
&= (n - r_n)^{-1} (\boldsymbol{\xi} + \boldsymbol{F})^\top (\mathbb{I} - \Pi_n)(\boldsymbol{\xi} + \boldsymbol{F}).
\end{aligned}
$$

where r_n was the rank of Π_n. Condition (D) provides $\|(\mathbb{I} - \Pi_n)\boldsymbol{F}\|_n \leq C\sigma n^{-1/2}$, see the proof of Lemma 1.

Lemma 2.   Under the condition (D) it holds

$$
\boldsymbol{P} \left(\left| \sigma^{-2} \widehat{\sigma}^2 - 1 \right| > \sqrt{(n - r_n)^{-1} \log n} \right) = o_n(1)
$$

where $o_n(1)$ denotes a numerical sequence tending to zero as $n \to \infty$. Moreover,

$$
\boldsymbol{P} \left(\sigma^{-2} \left| \widehat{\sigma}^2 - \widehat{\sigma}_0^2 \right| > n^{-1} \right) = o_n(1)
$$

where

$$
\widehat{\sigma}_0^2 = (n - r_n)^{-1} \boldsymbol{\xi}^\top (\mathbb{I} - \Pi_n) \boldsymbol{\xi}.
$$

Proof. By definition

$$
\begin{aligned}
&\sigma^{-2} (n - r_n) \left(\widehat{\sigma}^2 - \widehat{\sigma}_0^2 \right) \\
&= \sigma^{-2} \boldsymbol{F}^\top (\mathbb{I} - \Pi_n) \boldsymbol{F} + 2\sigma^{-2} \boldsymbol{\xi}^\top (\mathbb{I} - \Pi_n) \boldsymbol{F}.
\end{aligned}
$$

Condition (D) provides

$$
\boldsymbol{F}^\top (\mathbb{I} - \Pi_n) \boldsymbol{F} = \|(\mathbb{I} - \Pi_n)\boldsymbol{F}\|_n^2 \leq C^2 n^{-1}.
$$

Next, since $\sigma^{-2} \boldsymbol{\xi}^\top (\mathbb{I} - \Pi_n) \boldsymbol{F}$ is the linear combination of Gaussian errors ξ_i, it is also a Gaussian random variable with zero mean and the variance

$$
\begin{aligned}
\sigma^{-4} \boldsymbol{E} \left| \boldsymbol{\xi}^\top (\mathbb{I} - \Pi_n) \boldsymbol{F} \right|^2 &= \sigma^{-4} \boldsymbol{E} \boldsymbol{F}^\top (\mathbb{I} - \Pi_n) \boldsymbol{\xi} \boldsymbol{\xi}^\top (\mathbb{I} - \Pi_n) \boldsymbol{F} \\
&= \sigma^{-2} \boldsymbol{F}^\top (\mathbb{I} - \Pi_n) \boldsymbol{F} \\
&\leq C^2 n^{-1}.
\end{aligned}
$$

This implies for every $t \geq 1$

$$
\boldsymbol{P} \left(\sigma^{-2} \left(\widehat{\sigma}^2 - \widehat{\sigma}_0^2 \right) \frac{C^2}{n(n - r_n)} + \frac{C n^{-1/2} t}{n - r_n} \right) \leq e^{-t^2/2}
$$

and the second assertion of the lemma follows in view of $n - r_n \leq n/3$.

For the first one, it remains to estimate $\sigma^{-2} \widehat{\sigma}_0^2 - 1 = \sigma^{-2} (n - r_n)^{-1} \boldsymbol{\xi}^\top (\mathbb{I} - \Pi_n) \boldsymbol{\xi}$. Since $\sigma^{-2} \boldsymbol{\xi}^\top (\mathbb{I} - \Pi_n) \boldsymbol{\xi}$ follows the χ^2-distribution with $n - r_n$ degrees of freedom, the expression $\sqrt{(n - r_n)/2} \left(\sigma^{-2} \widehat{\sigma}_0^2 - 1 \right)$ is asymptotically standard normal and the required assertion follows. ∎

5.3   Proof of Theorem 1

Let j_n be due to (10), then with condition (D),

$$
\sum_{m=2}^d \|\boldsymbol{f}_m - \Pi_{m,n} \boldsymbol{f}_m\|_n \leq C\sigma n^{-1/2}
$$

and each $\Pi_{m,n} \boldsymbol{f}_m$ can be represented in the form

$$
\Pi_{m,n} \boldsymbol{f}_m = \sum_{I \in \mathcal{I}'(j_n)} \theta_{I,m} \boldsymbol{\psi}_{I,m}, \qquad m = 2, \dots, d,
$$

with some coefficients $\theta_{I,m}$, $I \in \mathcal{I} \in \mathcal{I}'(j_n)$.

Under the null hypothesis, $\boldsymbol{f}_1 = \theta_{0,1} + \theta_{1,1} \boldsymbol{\psi}_{1,1}$. Define

$$
\boldsymbol{F}' = \boldsymbol{F} - \boldsymbol{f}_1 - \Pi_{2,n} \boldsymbol{f}_2 - \dots - \Pi_{d,n} \boldsymbol{f}_d.
$$

Then the above bound yields

$$
\begin{aligned}
\|\boldsymbol{F}'\|_n &= \|\boldsymbol{F} - \boldsymbol{f}_1 - \Pi_{2,n} \boldsymbol{f}_2 - \dots - \Pi_{d,n} \boldsymbol{f}_d\|_n \\
&\leq C\sigma n^{-1/2}. \qquad (17)
\end{aligned}
$$

Now we show that the original regression function F can be replaced by F'.

Lemma 3.   The change F for F' does not affect the test statistics $T^* = \max_{j \leq j_n} T_j$.

Proof. Let some $j \leq j_n$ be fixed. Denote by $\boldsymbol{\theta}(d, j) = (\theta_{I,m}, (I, m) \in \mathcal{I}(d, j))$ the vector with $\theta_{I,1} = 0$ for $I \in \mathcal{I}_j$, $0 \leq j \leq j_n$, and with the above defined $\theta_{I,m}$ for $m \geq 2$. Then $\boldsymbol{F}' = \boldsymbol{F} - \Psi(d, j)\boldsymbol{\theta}(d, j)$ and the inequality in (17) can be rewritten in the form $\|\boldsymbol{F} - \Psi(d, j)\boldsymbol{\theta}(d, j)\|_n \leq \sigma n^{-1/4}$. Define also $\boldsymbol{\theta}^*(d, j) = W(d, j)^\top \boldsymbol{F}$ and let $\boldsymbol{\theta}_j^*$ be the subvector of $\boldsymbol{\theta}^*(d, j)$ corresponding to the jth resolution level of the first component. This vector can be written in the form $\boldsymbol{\theta}_j^* = \mathcal{E}_j \boldsymbol{\theta}^*(d, j)$ with $\mathcal{E}_j$ being the projector from $\mathbb{R}^{N(d,j)}$ onto $\mathbb{R}^{2^j}$ keeping the entries $\theta_{I,1}^*$, $I \in \mathcal{I}_j$, of the vector $\boldsymbol{\theta}^*(d, j)$ corresponding to the jth resolution level of the first component. Then it holds

$$
\boldsymbol{\theta}_j^* = \mathcal{E}_j W(d, j)^\top \boldsymbol{F} = W_j^\top \boldsymbol{F}.
$$

Since the test statistic T^* is calculated via the estimates $\widehat{\boldsymbol{\theta}}_j = W_j^\top \boldsymbol{Y}$ for $j \leq j_n$, and since $W_j^\top \boldsymbol{Y} = W_j^\top \boldsymbol{\xi} + W_j^\top \boldsymbol{F}$, it only remains to check that $W_j^\top \Psi(d, j)\boldsymbol{\theta}(d, j) = 0$ for all $j \leq j_n$. The definition of $\boldsymbol{\theta}(d, j)$ provides $\mathcal{E}_j \boldsymbol{\theta}(d, j) = 0$, and hence,

$$
\begin{aligned}
&W_j^\top \Psi(d, j)\boldsymbol{\theta}(d, j) \\
&= \mathcal{E}_j W(d, j)^\top \Psi(d, j)\boldsymbol{\theta}(d, j) \\
&= \mathcal{E}_j \left(\Psi(d, j)^\top \Psi(d, j) \right)^- \Psi(d, j)^\top \Psi(d, j)\boldsymbol{\theta}(d, j) \\
&= 0
\end{aligned}
$$

as required. ∎

This lemma allows to reduce the statement of the theorem to the case with $\|\boldsymbol{F}\|_n \leq C\sigma n^{-1/2}$.

Recall that the critical value of the test is evaluated under the condition $F \equiv 0$. Now we intend to show that $\boldsymbol{P}_F(\phi^* = 1) = \alpha + o_n(1)$ for every regression function F satisfying $\|\boldsymbol{F}\|_n \leq C\sigma n^{-1/2}$. The test ϕ^* is based on the test statistic $T^* = \max_{j \leq j_n}(T_j - t_{j,\alpha})$ with

$$T_j = \frac{\boldsymbol{Y}^\top W_j V_j^- W_j^\top \boldsymbol{Y}}{\widehat{\sigma}^2\sqrt{2N_j}} - \sqrt{N_j/2} = \frac{\boldsymbol{Y}^\top \mathcal{R}_j \boldsymbol{Y}}{\widehat{\sigma}^2\sqrt{2N_j}} - \sqrt{N_j/2}.$$

Here W_j is the submatrix of the matrix $W(d, j)$ corresponding to the jth resolution level of the first component, $W_j = \mathcal{E}_j W(d, j)$, and $V_j = W_j^\top W_j$, so that $\mathcal{R}_j = W_j V_j^- W_j^\top$ is a projector in $\mathbb{R}^n$ on the N_j-dimensional subspace. The model $\boldsymbol{Y} = \boldsymbol{F} + \boldsymbol{\xi}$ implies

$$T_j = \frac{\boldsymbol{\xi}^\top \mathcal{R}_j \boldsymbol{\xi}}{\widehat{\sigma}^2\sqrt{2N_j}} - \sqrt{N_j/2} + \frac{2\boldsymbol{\xi}^\top \mathcal{R}_j \boldsymbol{F} + \boldsymbol{F}^\top \mathcal{R}_j \boldsymbol{F}}{\widehat{\sigma}^2\sqrt{2N_j}}.$$

Define

$$T_j^0 = \frac{\boldsymbol{\xi}^\top \mathcal{R}_j \boldsymbol{\xi}}{\widehat{\sigma}_0^2\sqrt{2N_j}} - \sqrt{N_j/2}. \tag{18}$$

We intend to bound the difference $T_j - T_j^0$.

Lemma 4. *Let condition* (D) *be fulfilled and the component* f_1 *be a linear functions. Then it holds*

$$\sum_{j=0}^{j_n} \boldsymbol{P}\left(|T_j - T_j^0| > \epsilon_j\right) = o_n(1). \tag{19}$$

where $\epsilon_j = 3C\sqrt{\frac{\log j_n}{nN_j}}$.

Proof. Clearly we have

$$T_j - T_j^0 = \frac{\boldsymbol{\xi}^\top \mathcal{R}_j \boldsymbol{\xi}}{\sqrt{2N_j}}\left(\widehat{\sigma}^{-2} - \widehat{\sigma}_0^{-2}\right) + \frac{2\boldsymbol{\xi}^\top \mathcal{R}_j \boldsymbol{F} + \boldsymbol{F}^\top \mathcal{R}_j \boldsymbol{F}}{\widehat{\sigma}^2\sqrt{2N_j}}.$$

Similarly to the proof of Lemma 2 one can show that

$$\frac{\boldsymbol{\xi}^\top \mathcal{R}_j \boldsymbol{\xi}}{\sqrt{2N_j}}\left(\widehat{\sigma}^{-2} - \widehat{\sigma}_0^{-2}\right) = o(n^{-1})$$

and for every $t \geq 1$,

$$\boldsymbol{P}\left(\frac{2\boldsymbol{\xi}^\top \mathcal{R}_j \boldsymbol{F}}{\sigma^2\sqrt{2N_j}} > \frac{2\|\mathcal{R}_j \boldsymbol{F}\|_n t}{\sigma^2\sqrt{2N_j}}\right) \leq e^{-t^2/2}. \tag{20}$$

Since $\|\mathcal{R}_j \boldsymbol{F}\|_n \leq \|\boldsymbol{F}\|_n \leq Cn^{-1/2}$, this inequality applied with $t = \sqrt{3\log j_n}$ yields

$$\sum_{j=0}^{j_n} \boldsymbol{P}\left(\frac{2\boldsymbol{\xi}^\top \mathcal{R}_j \boldsymbol{F} + \boldsymbol{F}^\top \mathcal{R}_j \boldsymbol{F}}{\sigma^2\sqrt{2N_j}} > Cn^{-1} + 2C\sqrt{\frac{3\log j_n}{2nN_j}}\right)$$
$$\leq (1 + j_n)e^{3/2\log j_n} = o_n(1).$$

Now the required assertion follows in view of the first statement of Lemma 2. ∎

By construction, it holds

$$\boldsymbol{P}\left(\max_{0 \leq j \leq j_n}(T_j^0 - t_{j,\alpha}) > \lambda\right) = \alpha.$$

The idea is to show that this equality remains valid in the asymptotic sense if we replace here T_j^0 by T_j. Assertion (19) yields

$$\boldsymbol{P}\left(\max_{0 \leq j \leq j_n}(T_j^0 - t_{j,\alpha} - \epsilon_j) > \lambda\right) - o_n(1)$$
$$\leq \boldsymbol{P}\left(\max_{0 \leq j \leq j_n}(T_j - t_{j,\alpha}) > \lambda\right)$$
$$\leq \boldsymbol{P}\left(\max_{0 \leq j \leq j_n}(T_j^0 - t_{j,\alpha} + \epsilon_j) > \lambda\right) + o_n(1).$$

Now it suffices to check that

$$\sum_{j=0}^{j_n} \boldsymbol{P}\left(|T_j^0 - t_{j,\alpha} - \lambda| \leq \epsilon_j\right) = o_n(1).$$

The distribution of T_j^0 is precisely known and for sufficiently large n it is very close to the centered and standardized χ^2-distribution with N_j degrees of freedom. This particularly yields that the density of this distribution with respect to the Lebesgue measure is bounded by 1 and therefore,

$$\sum_{j=0}^{j_n} \boldsymbol{P}\left(|T_j^0 - t_{j,\alpha} - \lambda| \leq 3C\sqrt{\frac{\log j_n}{nN_j}}\right)$$
$$\leq \sum_{j=0}^{j_n} 6C\sqrt{\frac{\log j_n}{nN_j}} = o_n(1)$$

and the theorem is proved.

5.4 Proof of Theorem 2

The proof utilizes the following technical

Lemma 5. *For sufficiently large* n, *it holds*

$$\max_{0 \leq j \leq j_n} t_{j,\alpha} + \lambda \leq 2\sqrt{\log j_n}.$$

Proof. The statement obviously follows from the fact that

$$\sum_{j=0}^{j_n} \boldsymbol{P}\left(T_j^0 > 2\sqrt{\log j_n}\right) = o_n(1)$$

where every T_j^0 is defined by centering and standardization of a χ^2-sum with N_j degrees of freedom, see Spokoiny (1999a) for more details. ∎

Let, for some $j \le j_n$, it holds

$$
\begin{aligned}
T_j^* &= (2N_j)^{-1/2}\sigma^{-2}\boldsymbol{\theta}_j^{*T}V_j^{-}\boldsymbol{\theta}_j^* \\
&\ge (t_{j,\alpha} + \lambda + t_{j,\beta}') + c\sqrt{\log j_n}
\end{aligned}
$$

with some $c > 0$. This inequality can be represented in the form

$$
\frac{\boldsymbol{F}^{\top}\mathcal{R}_j\boldsymbol{F}}{\sigma^2\sqrt{2N_j}} > 2t_{j,\alpha} + \lambda + t_{j,\beta}' + a\sqrt{\log j_n}. \tag{21}
$$

We will show that under the above assumption,

$$
\boldsymbol{P}_F(T_j < t_{j,\alpha} + \lambda) \le \beta + o_n(1),
$$

which obviously implies the assertion.

Similarly to the proof of Theorem 1 we derive

$$
\begin{aligned}
T_j - T_j^0 &= \frac{\boldsymbol{\xi}^{\top}\mathcal{R}_j\boldsymbol{\xi}}{\sqrt{2N_j}}\left(\frac{1}{\widehat{\sigma}^2} - \frac{1}{\widehat{\sigma}_0^2}\right) + \frac{2\boldsymbol{\xi}^{\top}\mathcal{R}_j\boldsymbol{F} + \boldsymbol{F}^{\top}\mathcal{R}_j\boldsymbol{F}}{\widehat{\sigma}^2\sqrt{2N_j}} \\
&= o(n^{-1}) + \frac{2\boldsymbol{\xi}^{\top}\mathcal{R}_j\boldsymbol{F} + \boldsymbol{F}^{\top}\mathcal{R}_j\boldsymbol{F}}{\widehat{\sigma}^2\sqrt{2N_j}} \\
&= o(n^{-1}) + \frac{2\boldsymbol{\xi}^{\top}\mathcal{R}_j\boldsymbol{F} + \boldsymbol{F}^{\top}\mathcal{R}_j\boldsymbol{F}}{\sigma^2\sqrt{2N_j}}
\end{aligned}
$$

with T_j^0 from (18). The definition of $t_{j,\beta}'$ provides

$$
\boldsymbol{P}\left(-T_j^0 < -t_{j,\beta}'\right) = \beta.
$$

Now (20) implies

$$
\begin{aligned}
\boldsymbol{P}_F(T_j < t_{j,\alpha} + \lambda) &\le \boldsymbol{P}\left(T_j^0 < -t_{j,\beta}'\right) \\
&\quad + \boldsymbol{P}\left(-\gamma_j > T_j^* - t_{j,\alpha} - \lambda - t_{j,\beta}'\right) + o_n(1)
\end{aligned}
$$

where $\gamma_j = \frac{2\boldsymbol{\xi}^{\top}\mathcal{R}_j\boldsymbol{F}}{\sigma^2\sqrt{2N_j}}$ is a Gaussian r.v. with zero mean and $\boldsymbol{E}\gamma_j^2 = 4T_j^*(2N_j)^{-1/2}$. It remains to check that condition (21) and Lemma 5 imply

$$
\frac{4T_j^*(2N_j)^{-1/2}}{(T_j^* - t_{j,\alpha} - \lambda - t_{j,\beta}')^2} = o_n(1).
$$

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