Vol. 15 (2021) 1743–1782

ISSN: 1935-7524

https://doi.org/10.1214/21-EJS1806

Consistent regression using data-dependent coverings

Vincent Margot, Jean-Patrick Baudry, Frederic Guilloux and Olivier Wintenberger

Sorbonne Université, Campus Pierre et Marie Curie, Laboratoire de Probabilités, Statistique et Modélisation, Tour 16-26, 1er étage, 4, Place Jussieu, 75005 Paris

Abstract: We introduce a procedure to generate an estimator of the regression function based on a data-dependent quasi-covering of the feature space. A quasi-partition is generated from the quasi-covering and the estimator predicts the conditional empirical expectation over the cells of the quasi-partition. We provide sufficient conditions to ensure the consure of the estimator. Each element of the quasi-covering is labeled as significant or insignificant. We avoid the condition of cell shrinkage commonly found in the literature for data-dependent partitioning estimators. This reduces the number of elements in the quasi-covering. An important feature of our estimator is that it is interpretable.

The proof of the consistency is based on a control of the convergence rate of the empirical estimation of conditional expectations, which is interesting in itself.

MSC2020 subject classifications: Primary 62G05, 62G08; secondary

Keywords and phrases: Consistency, nonparametric regression, rule-based algorithm, data-dependent covering, interpretable learning.

Received September 2019.

1. Introduction

We consider the following regression setting: let (\mathbf{X}, Y) be a pair of random variables in $\mathbb{R}^d \times \mathbb{R}$ of unknown distribution \mathbb{Q} such that

$$Y = g^*(\mathbf{X}) + Z,$$

where $\mathbb{E}[Z] = 0$, $\mathbb{V}(Z) = \sigma^2 < \infty$ and g^* is a measurable function from \mathbb{R}^d to \mathbb{R} . We make the following common assumptions:

•
$$Z$$
 is independent of X ; (H1)

Y is bounded:
$$\mathbb{Q}(S) = 1$$
 where $S = \mathbb{R}^d \times [-L, L]$ for some unknown $L > 0$. (H2)

The accuracy of $g:\mathbb{R}^d\to\mathbb{R}$ such that $\mathbb{E}[g(\mathbf{X})^2]<\infty$ is measured by its quadratic risk defined as

$$\mathcal{L}(g) := \mathbb{E}\left[(g(\mathbf{X}) - Y)^2 \right].$$

Then we have

$$g^{*}(\mathbf{X}) = \mathbb{E}\left[Y|\mathbf{X}\right] = \underset{g}{\operatorname{arg\,min}} \mathcal{L}\left(g\right) \text{ a.s.}$$

where the arg min ranges over all the measurable functions g with $\mathbb{E}[g(\mathbf{X})^2] < \infty$. Given a sample $\mathbf{D}_n = ((\mathbf{X}_1, Y_1), \dots, (\mathbf{X}_n, Y_n))$, we aim at predicting Y conditionally on \mathbf{X} . The observations (\mathbf{X}_i, Y_i) are assumed independent and identically distributed (i.i.d.) from the distribution \mathbb{Q} .

Let us denote the empirical measure for any $\mathbf{r} \subseteq \mathbb{R}^d$ and $I \subseteq \mathbb{R}$

$$\mathbb{Q}_n(\mathbf{r} \times I) := \frac{1}{n} \sum_{i=1}^n \mathbf{1}_{\mathbf{X}_i \in \mathbf{r}} \mathbf{1}_{Y_i \in I}.$$

For simplicity of notation, for any $\mathbf{r} \subseteq \mathbb{R}^d$ we write $\mathbb{Q}_n(\mathbf{r})$ instead of $\mathbb{Q}_n(\mathbf{r} \times \mathbb{R})$ as well as $\mathbb{Q}(\mathbf{r})$ instead of $\mathbb{Q}(\mathbf{r} \times \mathbb{R})$.

We consider a set of functions \mathcal{G}_n . Any $\tilde{g}_n \in \arg\min_{g \in \mathcal{G}_n} \mathcal{L}(g)$ minimizes the risk over \mathcal{G}_n but is not an estimator since it depends on the knowledge of \mathbb{Q} . Following the Empirical Risk Minimization (ERM) principle (Vapnik, 1995, Section 1.5), we define the empirical risk and an empirical risk minimizer over \mathcal{G}_n as, respectively,

$$\mathcal{L}_n(g) := \frac{1}{n} \sum_{i=1}^n \left(g(\mathbf{X}_i) - Y_i \right)^2 \text{ and } g_n \in \underset{g \in \mathcal{G}_n}{\arg \min} \mathcal{L}_n(g).$$
 (1.1)

The aim of this paper is to provide $interpretable^1$ learning algorithms that generate weakly consistent empirical risk minimizers g_n , i.e. such that their excess of risk $\ell(g^*, g_n)$ fulfill

$$\ell(g^*, g_n) := \mathcal{L}(g_n) - \mathcal{L}(g^*) = \mathbb{E}[(g_n(\mathbf{X}) - g^*(\mathbf{X}))^2] = o_{\mathbb{P}}(1).$$

Note that \mathbb{P} refers to an underlying probability measure of reference that is unique (on the contrary to \mathbb{Q} that is arbitrary) and that \mathbb{E} corresponds to the expectation under \mathbb{P} .

1.1. Rule-based algorithms using partitions and quasi-coverings

In this paper, we consider algorithms generating interpretable models that are rule-based, such as *CART* (Breiman et al., 1984), *ID3* (Quinlan, 1986), *C4.5* (Quinlan, 1993), *FORS* (Karalič and Bratko, 1997), *M5 Rules* (Holmes, Hall and Prank, 1999). In these models, the estimator is explained by the realization

¹A formal definition of interpretability convenient to our framework is proposed and discussed in Section 1.2.

of a simple condition, an *If-Then* statement of the form:

IF
$$(\mathbf{X} \in c_1)$$
 And $(\mathbf{X} \in c_2)$ And ... And $(\mathbf{X} \in c_k)$ (1.2)
THEN $g_n(\mathbf{X}) = p$

where $p \in \mathbb{R}$ and each $c_i \subseteq \mathbb{R}^d$ is expressed in its simplest shape. For instance, the subset $[a_1, b_1] \times [a_2, b_2] \times \mathbb{R}^{d-2}$ will be expressed as two c_i 's:

$$c_1 = [a_1, b_1] \times \mathbb{R}^{d-1}$$
 and $c_2 = \mathbb{R} \times [a_2, b_2] \times \mathbb{R}^{d-2}$.

The If part, called the condition of the rule, or simply the rule, is composed of the conjunction of k basic tests which forms a subset of \mathbb{R}^d and k is called the length of the rule. The Then part, called the conclusion of the rule, is the estimated value when the rule is activated (i.e., when the condition in the If part is satisfied). The rules are easy to understand and allow an interpretable decision process when k is small. For a review of the best-known algorithms for descriptive and predictive rule learning, see Zhao and Bhowmick (2003) and Fürnkranz and Kliegr (2015).

Formally, models generated by such algorithms are defined by a corresponding data-dependent partition \mathcal{P}_n of \mathbb{R}^d . Each element of the partition is named a cell. Let us define for any $\mathbf{r} \subseteq \mathbb{R}^d$ such that $\mathbb{Q}(\mathbf{r}) > 0$ and for any h measurable

$$\mathbb{E}[h(Y) \mid \mathbf{X} \in \mathbf{r}] := \frac{\mathbb{E}[h(Y)\mathbf{1}_{\mathbf{X} \in \mathbf{r}}]}{\mathbb{P}(\mathbf{X} \in \mathbf{r})}$$

and

$$\mathbb{V}(Y\mid \mathbf{X}\in\mathbf{r}):=\mathbb{E}\left[Y^2\mid \mathbf{X}\in\mathbf{r}\right]-\mathbb{E}\left[Y\mid \mathbf{X}\in\mathbf{r}\right]^2$$

and their empirical counterparts for any $\mathbf{r} \subseteq \mathbb{R}^d$ such that $\mathbb{Q}_n(\mathbf{r}) > 0$:

$$\mathbb{E}_n[h(Y) \mid \mathbf{X} \in \mathbf{r}] := \frac{\sum_{i=1}^n h(Y_i) \mathbf{1}_{\mathbf{X}_i \in \mathbf{r}}}{\sum_{i=1}^n \mathbf{1}_{\mathbf{X}_i \in \mathbf{r}}}$$

and

$$\mathbb{V}_n(Y \mid \mathbf{X} \in \mathbf{r}) := \mathbb{E}_n \left[Y^2 \mid \mathbf{X} \in \mathbf{r} \right] - \mathbb{E}_n \left[Y \mid \mathbf{X} \in \mathbf{r} \right]^2.$$

Any such algorithm following the ERM principle associates to \mathcal{P}_n an estimator such that

$$g_n \mathbf{1}_{\bigcup \{A \in \mathcal{P}_n : \mathbb{Q}_n(A) > 0\}} = \sum_{\substack{A \in \mathcal{P}_n \\ \mathbb{Q}_n(A) > 0}} \mathbb{E}_n[Y \mid \mathbf{X} \in A] \mathbf{1}_A$$
 (1.3)

which minimizes the empirical risk among the class of piecewise constant functions over \mathcal{P}_n .

Those algorithms use the dataset \mathbf{D}_n twice: The partition $\mathcal{P}_n = \mathcal{P}_n(\mathbf{D}_n)$ is constructed according to the dataset, then the dataset is used again to calculate the value of g_n on each element of \mathcal{P}_n as in (1.3).

In this paper, we present a novel family of algorithms by relaxing the condition that the rules constitute a partition of \mathbb{R}^d . We consider collections of sets $\mathcal{C}_n = \mathcal{C}_n(\mathbf{D}_n)$ which asymptotically cover \mathbb{R}^d (see (H4)). We call such a collection a quasi-covering, and a quasi-partition if the sets do not intersect, which

we do not need to assume in general. Such a quasi-covering is turned into a quasi-partition $\mathcal{P}(\mathcal{C}_n)$ using the definition below. Then according to the ERM principle we define a quasi-covering estimator g_n by (1.3) where $\mathcal{P}_n = \mathcal{P}(\mathcal{C}_n)$. The main interest of quasi-covering estimators compared to partitioning estimators is its interpretability. This will be discussed in Section 1.2.

Definition 1.1. Let C be a finite collection of sets of \mathbb{R}^d and let $\mathbf{c} = \bigcup_{\mathbf{r} \in C} \mathbf{r}$. We define the activation function as

$$\varphi_{\mathcal{C}}: \mathbb{R}^d \mapsto 2^{\mathcal{C}}; \qquad \varphi_{\mathcal{C}}(\mathbf{x}) = \{ \mathbf{r} \in \mathcal{C}: \mathbf{x} \in \mathbf{r} \}.$$

Then $\mathcal{P}(\mathcal{C})$, the partition of **c** generated from \mathcal{C} , is defined as

$$\mathcal{P}(\mathcal{C}) := \varphi_{\mathcal{C}}^{-1}(Im(\varphi_{\mathcal{C}})).$$

We introduce the maximal (resp. minimal) redundancy of C as

$$M(\mathcal{C}) := \max_{\mathbf{x} \in \mathbf{c}} \#\varphi_{\mathcal{C}}(\mathbf{x}), \qquad m(\mathcal{C}) := \min_{\mathbf{x} \in \mathbf{c}} \#\varphi_{\mathcal{C}}(\mathbf{x}).$$

Moreover, the cells of $\mathcal{P}(\mathcal{C})$ which are included in an element of $\mathcal{C}' \subset \mathcal{C}$ are gathered in $\mathcal{P}_{\mathcal{C}}(\mathcal{C}')$:

$$\mathcal{P}_{\mathcal{C}}(\mathcal{C}') := \{ A \in \mathcal{P}(\mathcal{C}) : \exists \mathbf{r} \in \mathcal{C}', A \subseteq \mathbf{r} \}.$$

We provide a more explicit characterization of $\mathcal{P}(\mathcal{C})$ in Proposition 1.1.

Proposition 1.1. Let C be a finite collection of sets of \mathbb{R}^d . Then

$$\mathcal{P}(\mathcal{C}) = \left\{ \bigcap_{\mathbf{r} \in \tilde{\mathcal{C}}} \mathbf{r} \setminus \bigcup_{\mathbf{r} \in \mathcal{C} \setminus \tilde{\mathcal{C}}} \mathbf{r} : \tilde{\mathcal{C}} \subseteq \mathcal{C} \right\} \setminus \{\emptyset\}.$$

The proof of Proposition 1.1 is deferred to the appendix. Figures 1 and 2 illustrate the fact that $\mathcal{P}(\mathcal{C})$ can contain (much) more sets than \mathcal{C} itself.

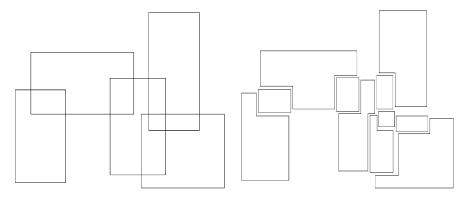


Fig 1. The 5 overlapping elements of C. Fig 2. The 11 cells of the partition $\mathcal{P}(C)$.

As the construction of $\mathcal{P}(\mathcal{C})$ can be time consuming, it is noteworthy that there is no need to fully describe $\mathcal{P}(\mathcal{C})$ to compute the value of $g_n(\mathbf{x})$ for some $\mathbf{x} \in \mathbb{R}^d$. The trick is to identify the unique cell of $\mathcal{P}(\mathcal{C})$ which contains \mathbf{x} . By creating binary vectors of size $\#\mathcal{C}$, whose value is 1 if \mathbf{x} fulfills the rule's condition and 0 otherwise, this cell identification is a simple sequence of vectorial operations. Figure 3 is an illustration of this process (we refer to Margot et al. (2018) for more details).

As compared to ours, the predictors designed by Random Forests (RF) and other tree ensembles algorithms do not satisfy the ERM principle (1.1): they are based on the averaging of the predictions of the activated rules at the point \mathbf{x} . This averaging is justified empirically and theoretically as it lowers the variance of the prediction thanks to the independence of the random rules as explained in Breiman (2001). The averaging does not decrease the approximation error of the prediction and there are few results about their consistency, we may cite Denil, Matheson and Freitas (2013); Scornet et al. (2015) for RF only. On the opposite, our estimator defined by (1.3) is the empirical conditional expectation on the cell of $\mathcal{P}(\mathcal{C}_n)$. Under some conditions, the approximation error of the prediction is low thanks to the ERM principle. The variance will be controlled through the significance condition (see Definition 2.1 below) as in the RIPE algorithm of Margot et al. (2018).

1.2. Interpretability

In many sensitive areas, such as health care, justice, defense or asset management the importance of interpretability in the decision-making process has been underlined. As explained in Lipton (2018), there are different meanings of interpretability depending on the desiderata of the user and the expected properties of the algorithms. In this paper, we embrace the intuitive definition of model interpretability of Biran and Cotton (2017): Interpretability is the degree to which an observer can understand the cause of a decision. Interpretable models should provide a parsimonious characterization of an estimator of g^* . Nowadays, the most popular and accurate algorithms for regression, such as Support Vector Machines, Neural networks, Random Forests,... are no interpretable. This lack of interpretability comes from the complexity of the models they generate. We call them black box models. We state that the novel family of (quasi-)covering algorithms described in Section 4 can achieve different interpretability-accuracy trade-off by reducing the complexity of the generated models while being accurate enough, maintaining weak consistency guarantees.

We distinguish two main approaches to generate interpretable prediction models. The first approach is to create black-box models and then to summarize them in a so-called *post-hoc* interpretable model. Recent researches propose to use explanation models such as *LIME* (Ribeiro, Singh and Guestrin, 2016), *DeepLIFT* (Shrikumar, Greenside and Kundaje, 2019) or *SHAP* (Lundberg and Lee, 2017) to interpret black-box models. These explanation models measure the importance of a feature on the prediction process (see Guidotti et al. (2018) for

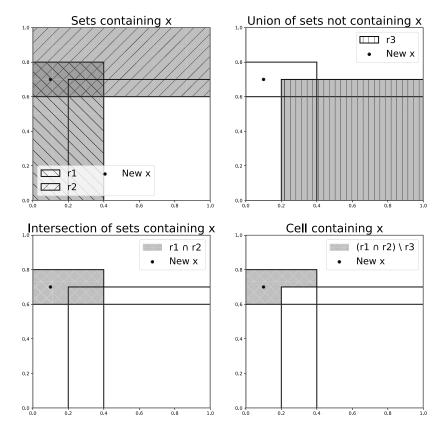


FIG 3. Evaluation steps of the cell containing $\mathbf{x} = (0.1, 0.7)$ of the partition generated from the covering of $[0, 1]^2$, $C = \{\mathbf{r}_1, \mathbf{r}_2, \mathbf{r}_3\}$. Using partition from a covering allows to generate complex cells with a simple interpretation $(\mathbf{r}_1 \text{ And } \mathbf{r}_2)$, where a classical partitioning algorithm cannot. Note that the condition \mathbf{x} satisfies $(\mathbf{r}_1 \text{ And } \mathbf{r}_2)$ implicitly implies that \mathbf{x} does not satisfy \mathbf{r}_3 .

a survey of existing methods). The other possibility is to use *intrinsic* interpretable algorithms that directly generate interpretable models.

The interpretability of a rule-based estimator is achieved when the length k of each rule and the number of rules are small, i.e. when the corresponding collection of rules is simple. Considering that a collection $\mathcal C$ with one rule of length k is just as interpretable as a collection with k rules of length 1 since they consists in the same number of tests, interpretability gets a certain additivity property. With this in mind, we are able to quantify the interpretability index of a collection $\mathcal C$.

Definition 1.2. The interpretability index of a collection of rules C is defined by

$$Int(\mathcal{C}) := \sum_{\mathbf{r} \in \mathcal{C}} length(\mathbf{r}),$$
 (1.4)

By convention, we define the interpretability index of the estimator g_n generated by a set of rules C_n as $Int(g_n) = Int(C_n)$.

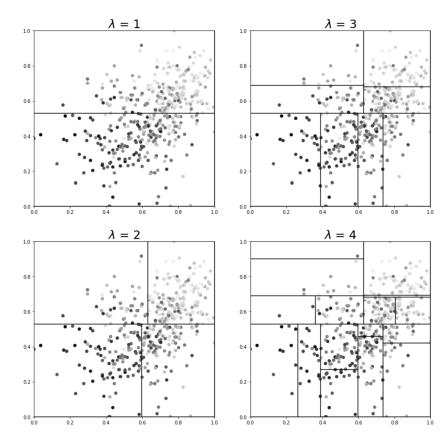


FIG 4. Partitions generated by fully deployed decision tree algorithm without pruning for a maximal depth $\lambda \in \{1, 2, 3, 4\}$. Here the number of rules satisfies 2^{λ} , $\lambda \in \{1, 2, 3, 4\}$.

Remark 1. The interpretability index $Int(g_n)$ does not take into account the notion of interaction described in Friedman and Popescu (2008). Indeed, a low interpretability index does not inform whether g_n exhibits interactions between variables or not. Many interactions may exhibit from a set of k rules of length 1 whereas it is less the case for a unique rule of length k. Due to the additivity property, both cases have the same interpretability index.

We aim at designing estimators which combine consistency guarantees keeping an interpretability index as small as possible. To demonstrate the consistency of a partition-based estimator g_n , results such as Theorem 13.1 in Györfi et al. (2006), which is a generalization of Stone's Theorem to data-dependent partitioning estimates, are usually applied. Then, "it is required that the diameter of the cells of the data-dependent partition [...] converge in some sense to zero" (Condition (13.10) in Györfi et al., 2006). However, as illustrated in Figure 4, such approaches based on fine partitions have a high interpretability index. A cell shrinkage condition such as Condition (13.10) in Györfi et al. (2006) neces-

sarily implies that the number of elements in the partition \mathcal{P}_n tends to infinity and that the interpretability index also tends to infinity.

For an estimator defined on a data-dependent quasi-covering, we can get rid of this condition and avoid the condition of shrinkage of the sets in the collection C_n . Any prediction can potentially be explained by a small set of rules that are interpretable: See Table 5 in Section 4 for an example. The generated partition may be finer and more complex than a classical data-dependent partition, while the explanation of the quasi-covering is easily understood by humans.

Despite the fact that we have not yet proven the parsimony of the selected set of rules in general, the use of quasi-covering instead of partition and conditions other than the usual ones drastically reduce the number of rules in simple cases. This is discussed in detail in Section 2.1.2.

We prove the consistency of a quasi-covering based estimator g_n by carefully designing the covering elements. The key concept is that of suitable data-dependent quasi-covering introduced in Section 2. Proposition 3.2 provides rates of convergence of empirical conditional expectations. This preliminary result of independent interest is crucial to prove the weak consistency of suitable data-dependent covering estimators stated in Theorem 2.1. Our main result Theorem 2.1 is proven in Section 3 while the proofs of intermediate results are postponed to the appendix. We apply our quasi-covering approach on artificial data and real data using different rule generators, and we compare it to usual existing algorithms in Section 4.

2. Consistent prediction based on suitable quasi-coverings

In order to ensure the consistency of quasi-covering based algorithms we have to introduce the notion of *suitable* sequence of data-dependent quasi-coverings.

2.1. Suitable data-dependent quasi-coverings

2.1.1. Definition

We use the classical notation $x_+ = \max\{x, 0\}$ for any $x \in \mathbb{R}$.

Definition 2.1. Let $\alpha \in [0, 1/2)$. We call a sequence $(C_n)_{n \geq 1}$ of data-dependent finite collections of sets \mathbf{r} of \mathbb{R}^d α -suitable if it satisfies:

$$\mathbb{Q}_n(\mathbf{r}) > n^{-\alpha}, \quad \mathbf{r} \in \mathcal{C}_n, n \ge 1.$$

2. The collection coverage condition (H4)

$$\mathbb{Q}_n(\mathbf{c}_n^c) = o_{\mathbb{P}}(1)$$

where $\mathbf{c}_n = \bigcup_{\mathbf{r} \in \mathcal{C}_n} \mathbf{r}$.

3. The significance condition (H5) For some consistent estimator (σ_n^2) of σ^2 and sequences $\beta_n = o_{\mathbb{P}}(1)$ and $\varepsilon_n = o_{\mathbb{P}}(1)$, it holds for all $n \geq 1$

$$C_n = C_n^s \cup C_n^i$$

where the significant sets C_n^s are defined by

$$C_n^s := \left\{ \mathbf{r} \in C_n : \beta_n \big| \mathbb{E}_n[Y|\mathbf{X} \in \mathbf{r}] - \mathbb{E}_n[Y] \big| \ge \sqrt{(\mathbb{V}_n(Y|\mathbf{X} \in \mathbf{r}) - \sigma_n^2)_+} \right\}$$
(2.1)

and the insignificant sets C_n^i are defined by

$$C_n^i := \left\{ \mathbf{r} \in C_n \setminus C_n^s : \varepsilon_n \ge \sqrt{(\mathbb{V}_n(Y|\mathbf{X} \in \mathbf{r}) - \sigma_n^2)_+} \right\}.$$
 (2.2)

4. The redundancy condition (H6)

The redundancies of C_n^s and C_n^i satisfy

$$\frac{M(\mathcal{C}_n^s)}{m(\mathcal{C}_n^s)} = o_{\mathbb{P}}(\beta_n^{-2} \wedge n^{1/2 - \alpha})$$
(2.3)

and

$$\frac{M(\mathcal{C}_n^i)}{m(\mathcal{C}_n^i)} = o_{\mathbb{P}}(\varepsilon_n^{-2} \wedge n^{1/2 - \alpha}). \tag{2.4}$$

A few remarks on the conditions of Definition 2.1 are in order. The set coverage condition (H3) controls the minimal number of observations in a set of the collection C_n , ensuring that the empirical conditional expectation is close to the conditional expectation given any of the sets in C_n . The collection coverage condition (H4) guarantees that C_n covers a sufficient part of \mathbb{R}^d . The significance conditions (H5) control the variance of the rules in two different ways; The significant rules are predicting significantly differently than the average and may have a large variance, the insignificant rules are not significant and have a small variance. The redundancy condition (H6) controls how far from a partition the collections C_n^s and C_n^i can be. This condition is fulfilled if the sets in these collections are mutually disjoint.

Remark 2. The significant condition (2.1) is not related to a condition on the diameter of the set \mathbf{r} . It is intended to detect areas of the feature space where the conditional expectation of Y is noticeably different from its unconditional expectation over the whole feature space. See Example 1. The insignificant condition (2.2) can be proved to hold under a shrinkage condition on the diameter of the set, see Proposition 3.3. This is not a necessary condition.

Remark 3. We may assume in the following that the noise variance σ^2 is known or that a good enough estimator σ_n^2 exists. Theorem 2.1 proves the consistency of the procedure under the assumption that $|\sigma_n^2 - \sigma^2| = O_{\mathbb{P}}(n^{\alpha - 1/2})$. One refers to Liitiäinen et al. (2009) as a reference on the difficult question of the bias-variance

trade-off of some residual variance estimators. Most of the existing estimators are based on Nearest-Neighbor methods. Random forests were not used for residual variance estimation until recently, see Ramosaj and Pauly (2019). None of these results is satisfactory in our setting of rule-based procedure. We let for a future work the study of a good enough estimator σ_n^2 based on a covering algorithm. We will use on real data in Section 4 the smallest of the empirical variances of the rules.

Remark 4. An easy way to ensure (H6) is to constrain the proportion of any element of C_n which can be included in the union of the others to be small. Let (C_n) be a sequence of collections of sets \mathbf{r} of \mathbb{R}^d that fulfills (H3) and $\gamma \in (0,1)$. We consider $(\mathbf{r}_i)_{1 \leq i \leq \#C_n}$ any ordering of C_n . If

$$\mathbb{Q}_n\Big(\mathbf{r}_{i+1}\bigcap\Big\{\bigcup_{1\leq j\leq i}\mathbf{r}_j\Big\}\Big)\leq \gamma\,\mathbb{Q}_n(\mathbf{r}_{i+1})\,,\qquad 1\leq i\leq \#\mathcal{C}_n-1\,,$$

then the cardinality of C_n is upper bounded by $\frac{n^{\alpha}}{1-\gamma}$ for n sufficiently large. Indeed, by the inclusion-exclusion principle we get

$$1 \ge \mathbb{Q}_n(\cup_{1 \le i \le \#\mathcal{C}_n} \mathbf{r}_i) = \sum_{i=1}^{\#\mathcal{C}_n} \mathbb{Q}_n(\mathbf{r}_i \setminus \cup_{1 \le j \le i-1} \mathbf{r}_j) \ge \#\mathcal{C}_n (1-\gamma) n^{-\alpha}.$$

Thus (H6) can be checked for any $\alpha \in [0, 1/4)$, using the fact that $M(\mathcal{C}_n^s)$ and $M(\mathcal{C}_n^i)$ are smaller than $\frac{n^{\alpha}}{1-\gamma}$ and setting $\beta_n = o_{\mathbb{P}}(n^{-\alpha/2})$ and $\varepsilon_n = o_{\mathbb{P}}(n^{-\alpha/2})$.

2.1.2. Advantages of suitable quasi-coverings

We discuss how our approach can provide estimators with low interpretability indices by considering quasi-coverings instead of partitions on the one hand and by introducing conditions on the elements of the quasi-coverings which are completely different from Condition (13.10) in Györfi et al. (2006) on the partition cells on the other hand.

We already discussed and illustrated in Figures 1 and 2 that a quasi-covering C_n contains less elements than the quasi-partition $\mathcal{P}(C_n)$. Then its interpretability index is lower.

Moreover, $\mathcal{P}(\mathcal{C})$ contain more complex elements than those of a partition-based algorithm: the cells in $\mathcal{P}(\mathcal{C})$ in Figure 2 are not necessarily conjunctions of simple tests on coordinates as the sets in Figure 1. Our covering-based approach generates complex cells without sacrificing interpretability since the algorithm keeps small the number of elements of \mathcal{C} which have to be designed.

Let us illustrate this point in the following example considering rules designed as hyperrectangles which is usual, see Remark 6.

Proposition 2.1. Let
$$C = [0, 1]^d \subset \mathbb{R}^d$$
.

1. The minimal cardinality of a partition \mathcal{P} of \mathbb{R}^d by hyperrectangles such that $C \in \mathcal{P}$ is 2d + 1.

2. The minimal cardinality of a covering C of \mathbb{R}^d by hyperrectangles such that $C \in \mathcal{P}(C)$ is 2.

The proof of Proposition 2.1 is deferred to the appendix and is a consequence of the fact that $\mathbb{R}^d \setminus [0,1]^d$ is a more complex set than an hyperrectangle.

On the other hand, our Condition (H5) suits better the requirement of interpretability than Condition (13.10) of Györfi et al. (2006) as they do not imply any shrinkage of the sets in the collection C_n and the number of sets in C_n does not necessarily have to grow with n. We illustrate in the following example this point when g^* itself is rule-based with a low interpretability index: in this important situation, an estimator with low interpretability index is expected.

Example 1. The condition involved in (2.1) can hold for a set \mathbf{r} with arbitrary diameter that does not satisfy Condition (13.10) of Györfi et al. (2006). Consider the case $g^* = \mathbf{1}_A$ for some Borel set A such that $0 < \mathbb{P}(\mathbf{X} \in A) < 1$ and assume that σ^2 is known. Then $\mathbf{r} = A$ is a significant set as it satisfies the condition involved in (2.1) for some well-chosen $\beta_n = o_{\mathbb{P}}(1)$. Indeed, from the Strong Law of Large Numbers $k_n := \#\{i : \mathbf{X}_i \in A\} \sim n\mathbb{P}(\mathbf{X} \in A)$ a.s. as $n \to \infty$. Then repeated application of the Central Limit Theorem yields

$$\begin{aligned} \left| \mathbb{E}_n[Y|\mathbf{X} \in A] - \mathbb{E}_n[Y] \right| &\geq \mathbb{E}_n[Y|\mathbf{X} \in A] - \mathbb{E}_n[Y] \\ &= 1 - \frac{k_n}{n} + \frac{1}{k_n} \sum_{i=1}^n Z_i \mathbf{1}_{\mathbf{X}_i \in A} - \frac{1}{n} \sum_{i=1}^n Z_i \\ &= 1 - \mathbb{P}(\mathbf{X} \in A) + O_{\mathbb{P}}(n^{-1/2}) \end{aligned}$$

and

$$(\mathbb{V}_n(Y|\mathbf{X} \in A) - \sigma^2)_+ \le |\mathbb{V}_n(Y|\mathbf{X} \in A) - \sigma^2|$$

$$= \left| \frac{1}{k_n} \sum_{i=1}^n Z_i^2 \mathbf{1}_{\mathbf{X}_i \in A} - \left(\frac{1}{k_n} \sum_{i=1}^n Z_i \mathbf{1}_{\mathbf{X}_i \in A} \right)^2 - \sigma^2 \right|$$

$$= O_{\mathbb{P}}(n^{-1/2}).$$

It follows that (2.1) holds for $\mathbf{r} = A$ with

$$\beta_n = \frac{\sqrt{(\mathbb{V}_n(Y|\mathbf{X}\in A) - \sigma^2)_+}}{|\mathbb{E}_n[Y|\mathbf{X}\in A] - \mathbb{E}_n[Y]|} = O_{\mathbb{P}}(n^{-\frac{1}{4}}).$$

For similar reasons (2.1) also holds for $\mathbf{r} = A^c$ with

$$\beta_n' = \beta_n \vee \frac{\sqrt{(\mathbb{V}_n(Y|\mathbf{X} \in A^c) - \sigma^2)_+}}{|\mathbb{E}_n[Y|\mathbf{X} \in A^c] - \mathbb{E}_n[Y]|} = O_{\mathbb{P}}(n^{-\frac{1}{4}}).$$

Finally it can easily be checked that $C_n = \{\mathbf{r} \in \{A, A^c\} : \mathbb{Q}_n(\mathbf{r}) > n^{-\alpha}\}$ defines an α -suitable sequence of quasi-partitions for any $\alpha \in (0, \frac{1}{2})$. The number of rules in C_n tends to 2.

Remark 5. In Section 4 and for many usual procedures, only hyperrectangles are considered as rules. In Example 1, if A is assumed to be an hyperrectangle itself, A^c can be written as the union of a set of hyperrectangles as in the proof of Proposition 2.1 and similar arguments as above apply to prove that an α -suitable sequence of quasi-coverings which are all subsets of a small set of hyperrectangles can thus be defined accordingly.

Note that in Proposition 2.1, Example 1, and Remark 5, the number of elements of the quasi-coverings and the interpretability index are bounded as d or n grows whereas they would need to tend to infinity for any method partition-based or relying on Condition (13.10) of Györfi et al. (2006).

Besides, another advantage is that the check of the conditions involved in (2.1) and (2.2) for elements of a quasi-covering can be done simultaneously with parallel computing on the contrary to any condition on partition cells that depend on each other.

2.2. Consistency of data-dependent quasi-covering algorithms

2.2.1. Generalization of the partitioning number

Conditions are required to control the complexity of the family of partitions $\mathcal{P}(\mathcal{C}_n)$ that the algorithm can generate. These conditions use some concepts introduced in (Nobel, 1996, Sec. 1.2) (see also (Györfi et al., 2006, Def 13.1)). The standard definitions have to be adapted since we consider collections which do not necessarily cover \mathbb{R}^d . To discriminate with the standard definitions for partitions we denote these quantities with tildes.

Definition 2.2. Let Π be a family of finite collections of disjoint sets of \mathbb{R}^d .

1. The maximal number of sets in an element of Π is denoted by

$$\widetilde{\mathcal{M}}(\Pi) := \sup \left\{ \# \mathcal{P} : \mathcal{P} \in \Pi \right\}.$$

2. For a set $\mathbf{x}_1^n = {\mathbf{x}_1, \dots, \mathbf{x}_n} \in (\mathbb{R}^d)^n$, let

$$\widetilde{\Delta}(\mathbf{x}_1^n,\Pi):=\#\big\{\{\mathbf{x}_1^n\cap A:A\in\mathcal{P}\}\setminus\{\emptyset\}:\mathcal{P}\in\Pi\big\}.$$

3. The partitioning number $\widetilde{\Delta}_n(\Pi)$ of Π is defined by

$$\widetilde{\Delta}_n(\Pi) := \max_{\mathbf{x}_1^n \in (\mathbb{R}^d)^n} \widetilde{\Delta}(\mathbf{x}_1^n, \Pi).$$

The partitioning number is the maximal number of different collections of disjoint non-empty subsets of any n-point set that can be induced by elements of Π and is smaller than or equal to $(n+1)^n$, which is an upper bound of the number of collections of disjoint subsets of \mathbf{x}_1^n as the number of maps from $\{1,\ldots,n\}$ to $\{0,\ldots,n\}$.

Classical consistency theorems require to control the partitioning number of the family of partitions Π_n over which the estimators are defined (see (13.7) and (13.8) in Györfi et al., 2006, Theorem 13.1). As illustrated in Figure 2, the complexity of the sets in $\mathcal{P}(\mathcal{C}_n)$ may be high even if the complexity of the sets in \mathcal{C}_n is low. This fact makes the partitioning number of the family of partitions tough, if not impossible, to bound because no control on the shape of the cells is available.

Usually Π_n is chosen as $\{\mathcal{P}(\mathcal{C}_n(\mathbf{d}_n)) : \mathbf{d}_n \in \mathcal{S}^n\}$. It can be useful to take into account the building process of the family of sets $\mathcal{C}_n(\mathbf{d}_n)$ in the evaluation of the partitioning number when the partitioning number of $\{\mathcal{P}(\mathcal{C}_n(\mathbf{d}_n)) : \mathbf{d}_n \in \mathcal{S}^n\}$ cannot be controlled easily. This is what the definition below enables by considering Π_n as a refinement of $\{\mathcal{P}(\mathcal{C}_n(\mathbf{d}_n)) : \mathbf{d}_n \in \mathcal{S}^n\}$.

Definition 2.3. Let \mathcal{F} and \mathcal{F}' be two families of finite collections of disjoint sets. \mathcal{F}' is called a refinement of \mathcal{F} if

$$\forall F \in \mathcal{F}, \quad \exists F' \in \mathcal{F}', \quad \forall A \in F, \quad A = \bigcup_{A' \in F': A' \subseteq A} A'.$$

2.2.2. Main result

Our main result provides the consistency of estimators based on suitable quasicovering sequences.

Theorem 2.1. Assume that g^* and \mathbb{Q} satisfy (H1), (H2). Let $\alpha \in [0, \frac{1}{2})$ and (\mathcal{C}_n) be an α -suitable sequence of finite collections of sets of \mathbb{R}^d such that:

$$\mathbb{Q}(\mathbf{c}_n^c) = O_{\mathbb{P}}(n^{\alpha - 1/2}) \tag{H7}$$

where $\mathbf{c}_n = \cup_{\mathbf{r} \in \mathcal{C}_n} \mathbf{r}$;

$$\widetilde{\mathcal{M}}(\Pi_n) = o(n)$$
 and $\log(\widetilde{\Delta}_n(\Pi_n)) = o(n)$ (H8)

where Π_n is a refinement of $\{\mathcal{P}(\mathcal{C}_n(\mathbf{d}_n)) : \mathbf{d}_n \in \mathcal{S}^n\};$

$$\forall n \in \mathbb{N}^*, \forall \mathbf{d}_n \in \mathcal{S}^n, \{\mathbf{r} \times \mathbb{R}, \mathbf{r} \in \mathcal{C}_n(\mathbf{d}_n)\} \subseteq \mathcal{B}$$
(H9)

where \mathcal{B} is a \mathbb{Q} -Donsker class.

Assume moreover that the estimator (σ_n^2) involved in Definition 2.1 is such that

$$|\sigma_n^2 - \sigma^2| = O_{\mathbb{P}}(n^{\alpha - 1/2}).$$

Then

$$g_n = \sum_{\substack{A \in \mathcal{P}(\mathcal{C}_n) \\ \mathbb{Q}_n(A) > 0}} \mathbb{E}_n[Y \mid \mathbf{X} \in A] \mathbf{1}_A$$

is weakly consistent:

$$\ell\left(q^*, q_n\right) = o_{\mathbb{P}}(1).$$

Remark 6. Many algorithms as CART generate rules as in (1.2) where each simple condition ($\mathbf{X} \in c$) consists in some distinct coordinate of \mathbf{X} belonging to an interval of \mathbb{R} . In Section 4 we use such rule generators for which (H9) is automatically satisfied. Indeed let \mathcal{H} be the set of all hyperrectangles of \mathbb{R}^d :

$$\mathcal{H} := \left\{ I_1 \times \ldots \times I_d \subset \mathbb{R}^d : I_i \text{ is an interval of } \mathbb{R} \right\} \cup \{\emptyset\}.$$

Any set of rules C_n such that $C_n \subseteq \mathcal{H}$, which is always the case of the sets of rules considered in Section 4, fulfills (H9) since it is a VC class (see Wenocur and Dudley, 1981) and VC classes are \mathbb{Q} -Donsker (see for example Van der Vaart, 2000, Lemma 19.15 and comments).

3. Proof of Theorem 2.1

In order to prove the main theorem, we need some preliminary results based on the notions of \mathbb{Q} -Donsker class and outer probability. The outer probability, defined for $A\subseteq \Omega$ by $\mathbb{P}^*(A):=\inf\left\{\mathbb{P}(\tilde{A}):A\subset \tilde{A},\tilde{A}\in \mathcal{A}\right\}$, is introduced to handle functions which are not necessarily measurable. The usual notion of boundedness in probability for sequences of random variables is generalized because sequences of maps are considered with values in metric spaces which are not Euclidean spaces (thus bounded and closed sets need not be compact) and which are not guaranteed to be measurable. See (Van der Vaart, 2000, Chapter 18) for details.

Definition 3.1 (Van der Vaart, 2000, Chapter 18). A sequence $(M_n)_{n\in\mathbb{N}}$ of maps defined on Ω and with values in a metric space (\mathbb{D},d) is said to be asymptotically tight if

$$\forall \varepsilon > 0, \exists K \subset \mathbb{D} \ compact, \forall \delta > 0, \limsup_{n \to \infty} \mathbb{P}^*(M_n \notin K^{\delta}) < \varepsilon,$$

with
$$K^{\delta} = \{ y \in \mathbb{D} : d(y, K) < \delta \}.$$

Remark 7. If $\mathbb{D} = \mathbb{R}$, (M_n) is asymptotically tight if and only if

$$\forall \varepsilon > 0, \exists M > 0 \text{ such that } \limsup_{n \to \infty} \mathbb{P}^*(|M_n| > M) < \varepsilon.$$

The notation $O_{\mathbb{P}^*}(1)$ stands for asymptotically tight instead of the usual $O_{\mathbb{P}}(1)$ (bounded in probability).

For $f: \mathcal{S} \to \mathbb{R}$ in $\mathcal{L}^1(\mathbb{Q})$ we define $v_n f := \sqrt{n} (\int f d\mathbb{Q}_n - \int f d\mathbb{Q})$ and consider the empirical process indexed by a set \mathcal{F} of such functions: $\{v_n f: f \in \mathcal{F}\}$.

Definition 3.2 (Van der Vaart, 2000, Section 19.2). A class of functions \mathcal{F} is called \mathbb{Q} -Donsker if the sequence of processes $\{v_n f : f \in \mathcal{F}\}$ converges in distribution to a tight limit process in the space $\ell^{\infty}(\mathcal{F})$.

The limit process is then a Q-Brownian bridge.

Definition 3.3. A class of sets $\mathcal{B} \subseteq \mathcal{B}_{\mathcal{S}}$ is called \mathbb{Q} -Donsker if $\mathcal{I}_{\mathcal{B}} := \{\mathbf{1}_A : A \in \mathcal{B}\}$ is a \mathbb{Q} -Donsker class of functions.

If \mathcal{F} is a \mathbb{Q} -Donsker class of functions, then the empirical process $((v_n f)_{f \in \mathcal{F}})_{n \in \mathbb{N}}$ is asymptotically tight as a sequence of maps with values in $\ell^{\infty}(\mathcal{F})$ (this is a consequence of Prohorov's Theorem adapted to this framework – see Van der Vaart, 2000, Theorem 18.12). Keeping in mind that a compact set in $\ell^{\infty}(\mathcal{F})$ is bounded, we have:

Proposition 3.1. Let \mathcal{F} be a \mathbb{Q} -Donsker class of functions. Then

$$\|\mathbb{Q}_n - \mathbb{Q}\|_{\mathcal{F}} = O_{\mathbb{P}^*}(n^{-1/2}),$$

where for any $v: \mathcal{F} \to \mathbb{R}$, $||v||_{\mathcal{F}} = \sup_{f \in \mathcal{F}} |v(f)|$.

Remark 8. If $\mathcal{B} \subseteq \mathcal{B}_{\mathcal{S}}$ is a \mathbb{Q} -Donsker class of sets, where $\mathcal{B}_{\mathcal{S}}$ is the Borel set on \mathcal{S} , then

$$\|\mathbb{Q}_n - \mathbb{Q}\|_{\mathcal{B}} = O_{\mathbb{P}^*}(n^{-1/2}),$$

where for any $v: \mathcal{B} \to \mathbb{R}$, $||v||_{\mathcal{B}} = \sup_{A \in \mathcal{B}} |v(A)|$.

Remark 9. It can be checked that if $(Z_n)_{n\in\mathbb{N}}$ is a sequence of non-negative random variables, $(a_n)_{n\in\mathbb{N}} \in (\mathbb{R}^+)^{\mathbb{N}}$ such that $a_n = o_{\mathbb{P}}(1)$ and $(M_n)_{n\in\mathbb{N}}$ is a sequence of maps (non necessarily measurable) such that $M_n = O_{\mathbb{P}^*}(1)$ and $Z_n \leq a_n M_n$ for any n, then $Z_n \xrightarrow[n \to +\infty]{\mathbb{P}} 0$.

3.1. Empirical estimation of conditional expectations

We shall also prove and use the following proposition, which is inspired by the work of Grunewalder (2018) (Proposition 3.2).

Proposition 3.2. Let $\mathcal{B} \subseteq \mathcal{B}_{\mathcal{S}}$ and let $\mathcal{F}_{\mathcal{B}} := \{ f\mathbf{1}_A : f \in \mathcal{F}, A \in \mathcal{B} \}$ where \mathcal{F} is a set of functions in $\mathcal{L}^1(\mathbb{Q})$ uniformly bounded. If \mathcal{B} and $\mathcal{F}_{\mathcal{B}}$ are \mathbb{Q} -Donsker classes then for any $\alpha \in [0, 1/2)$ and with $\mathcal{B}_n := \{ A \in \mathcal{B}, \mathbb{Q}_n(A) \geq n^{-\alpha} \}$ we have

$$\sup_{f \in \mathcal{F}} \sup_{A \in \mathcal{B}_n} \left| \mathbb{E}_n \left[f \mid A \right] - \mathbb{E} \left[f \mid A \right] \right| = O_{\mathbb{P}^*}(n^{\alpha - 1/2}).$$

Corollary 3.1. Let $\mathcal{B} \subseteq \mathcal{B}_{\mathcal{S}}$ be a \mathbb{Q} -Donsker class. If Y is bounded then for any $i \in \mathbb{N}$ and any $\alpha \in [0, 1/2)$, with $\mathcal{B}_n := \{A \in \mathcal{B}, \mathbb{Q}_n(A) \geq n^{-\alpha}\}$ we have

$$\sup_{A \in \mathcal{B}_n} \left| \mathbb{E}_n \left[Y^i \mid (\mathbf{X}, Y) \in A \right] - \mathbb{E} \left[Y^i \mid (\mathbf{X}, Y) \in A \right] \right| = O_{\mathbb{P}^*}(n^{\alpha - 1/2}), \quad (3.1)$$

and

$$\sup_{A \in \mathcal{B}_n} |\mathbb{V}_n \left[Y \mid (\mathbf{X}, Y) \in A \right] - \mathbb{V} \left[Y \mid (\mathbf{X}, Y) \in A \right] | = O_{\mathbb{P}^*}(n^{\alpha - 1/2}). \tag{3.2}$$

Proofs of these results are deferred to the appendix.

It seems that the result of Corollary 3.1, which is of independent interest, does not appear as such in the existing literature. As a first application of

Corollary 3.1, we show that any sequence of partitions with shrinking cell diameters is a suitable covering. We define the diameter of a cell \mathbf{r} as $\mathrm{Diam}(\mathbf{r}) = \sup_{\mathbf{x} \in \mathbf{r}, \mathbf{x}' \in \mathbf{r}} \|\mathbf{x} - \mathbf{x}'\|$, where $\|\cdot\|$ is any norm of \mathbb{R}^d .

Proposition 3.3. Consider a sequence $(\mathcal{P}_n)_{n\in\mathbb{N}}$ of data-dependent partitions that satisfies the coverage condition (H3) with $\alpha \in [0, \frac{1}{2})$ and such that

$$\bigcup_{n\in\mathbb{N}^*}\bigcup_{\mathbf{d}_n\in(\mathbb{R}^d)^n}\{\mathbf{r}\times\mathbb{R}:\mathbf{r}\in\mathcal{P}_n(\mathbf{d}_n)\}$$

is a.s. a \mathbb{Q} -Donsker class. Suppose that σ^2 is known, g^* is uniformly continuous and

$$\max_{\mathbf{r}\in\mathcal{P}_n}Diam(\mathbf{r}) = o_{\mathbb{P}}(1). \tag{3.3}$$

Then the sequence (\mathcal{P}_n) is α -suitable.

Proof. Let us show that each cell is either significant or insignificant. Thanks to Condition (H3), Corollary 3.1 (Eq. (3.2)) and Remark 9,

$$\max_{\mathbf{r} \in \mathcal{P}_n} |\mathbb{V}_n(Y | \mathbf{X} \in \mathbf{r}) - \mathbb{V}(Y | \mathbf{X} \in \mathbf{r})| = O_{\mathbb{P}}(n^{\alpha - 1/2}).$$
 (3.4)

Moreover $\mathbb{V}(Y \mid \mathbf{X} \in \mathbf{r}) = \mathbb{V}(g^*(\mathbf{X}) \mid \mathbf{X} \in \mathbf{r}) + \sigma^2$. Thus, as the redundancy condition (H6) is automatically satisfied for cells of a partition, the desired result will follow if we check that $\varepsilon_n = o_{\mathbb{P}}(1)$ with

$$\varepsilon_n := \max_{\mathbf{r} \in \mathcal{P}_n} \sqrt{(\mathbb{V}_n(Y|\mathbf{X} \in \mathbf{r}) - \sigma^2)_+}.$$

From (3.4) we remark that

$$\varepsilon_n \le \max_{\mathbf{r} \in \mathcal{P}_n} \sqrt{\mathbb{V}(g^*(\mathbf{X}) \mid \mathbf{X} \in \mathbf{r})} + O_{\mathbb{P}}(n^{\alpha/2 - 1/4})$$
.

For all n, if $\mathbf{r} \in \mathcal{P}_n$, then $\mathbf{r} \times \mathbb{R} \in \mathcal{B}_{\mathcal{S}}$. We denote $\mathbf{X}_{\mathbf{r}}$ and $\mathbf{X}'_{\mathbf{r}}$ two independent variables distributed as \mathbf{X} given that $\mathbf{X} \in \mathbf{r}$. We obtain

$$V(g^*(\mathbf{X}) \mid \mathbf{X} \in \mathbf{r}) = V(g^*(\mathbf{X}_{\mathbf{r}}))$$

$$= \frac{1}{2}V(g^*(\mathbf{X}_{\mathbf{r}}) - g^*(\mathbf{X}'_{\mathbf{r}}))$$

$$\leq \frac{1}{2}E\left[(g^*(\mathbf{X}_{\mathbf{r}}) - g^*(\mathbf{X}'_{\mathbf{r}}))^2\right].$$

Thus, if we denote w the modulus of continuity of g^* , we get

$$\sqrt{\mathbb{V}(g^*(\mathbf{X}) \mid \mathbf{X} \in \mathbf{r})} \le 2^{-1/2} w(\operatorname{Diam}(\mathbf{r})).$$

By uniform continuity, the condition (3.3) implies that

$$\varepsilon_n \le 2^{-1/2} \max_{\mathbf{r} \in \mathcal{P}_n} \left(w(\mathrm{Diam}(\mathbf{r})) \right) + O_{\mathbb{P}}(n^{\alpha/2 - 1/4}) = o_{\mathbb{P}}(1).$$

Thus, from (2.2), each cell which is not significant is insignificant and the corresponding covering sequence is α -suitable.

3.2. Estimation-approximation decomposition

Let \mathcal{G}_n denote the set of piecewise constant functions on the partition $\mathcal{P}(\mathcal{C}_n) \cup \{\mathbf{c}_n^c\}$ such that $\forall g \in \mathcal{G}_n, \forall \mathbf{x} \in \mathbb{R}^d, |g(\mathbf{x})| \leq L$ and $\forall g \in \mathcal{G}_n, \forall \mathbf{x} \in \mathbf{c}_n^c, g(\mathbf{x}) = 0$.

The excess risk $\ell(g^*, g_n)$ can be decomposed into two terms mimicking Lemma 10.1 of Györfi et al. (2006). First notice that under the conditions of Theorem 2.1 on g_n we have

$$\frac{1}{n}\sum_{i=1}^{n}(g_n(\mathbf{X}_i) - Y_i)^2 \mathbf{1}_{\mathbf{c}_n}(\mathbf{X}_i) \le \frac{1}{n}\sum_{i=1}^{n}(g(\mathbf{X}_i) - Y_i)^2 \mathbf{1}_{\mathbf{c}_n}(\mathbf{X}_i), \quad \forall g \in \mathcal{G}_n. \quad (3.5)$$

Now, since moreover $g_n \in \mathcal{G}_n$,

$$\ell(g^*, g_n) = \mathbb{E}\left[\left(g_n(\mathbf{X}) - g^*(\mathbf{X})\right)^2\right]$$

$$= \mathbb{E}\left[\left(g_n(\mathbf{X}) - Y\right)^2\right] - \mathbb{E}\left[\left(g^*(\mathbf{X}) - Y\right)^2\right]$$

$$= \mathbb{E}\left[\left(g_n(\mathbf{X}) - Y\right)^2\right] - \inf_{g \in \mathcal{G}_n} \mathbb{E}\left[\left(g(\mathbf{X}) - Y\right)^2\right]$$

$$+ \inf_{g \in \mathcal{G}_n} \mathbb{E}\left[\left(g(\mathbf{X}) - Y\right)^2\right] - \mathbb{E}\left[\left(g^*(\mathbf{X}) - Y\right)^2\right]$$

$$= \sup_{g \in \mathcal{G}_n} \left\{ \mathbb{E}\left[\left(g_n(\mathbf{X}) - Y\right)^2\right] - \frac{1}{n} \sum_{i=1}^n (g_n(\mathbf{X}_i) - Y_i)^2$$

$$+ \frac{1}{n} \sum_{i=1}^n (g_n(\mathbf{X}_i) - Y_i)^2 - \frac{1}{n} \sum_{i=1}^n (g(\mathbf{X}_i) - Y_i)^2$$

$$+ \frac{1}{n} \sum_{i=1}^n (g(\mathbf{X}_i) - Y_i)^2 - \mathbb{E}\left[\left(g(\mathbf{X}) - Y\right)^2\right] \right\}$$

$$+ \inf_{g \in \mathcal{G}_n} \mathbb{E}\left[\left(g(\mathbf{X}) - g^*(\mathbf{X})\right)^2\right]$$

$$= 2 \sup_{g \in \mathcal{G}_n} \left| \mathbb{E}\left[\left(g(\mathbf{X}) - Y\right)^2\right] - \frac{1}{n} \sum_{i=1}^n (g(\mathbf{X}_i) - Y_i)^2\right] \right|$$

$$+ \underbrace{\frac{1}{n} \sum_{i=1}^n (g_n(\mathbf{X}_i) - Y_i)^2 \mathbf{1}_{\mathbf{c}_n^c}(\mathbf{X}_i)}_{\leq 4L^2 \mathbb{Q}_n(\mathbf{c}_n^c)}$$

$$+ \inf_{g \in \mathcal{G}_n} \mathbb{E}\left[\left(g(\mathbf{X}) - g^*(\mathbf{X})\right)^2\right].$$

Up to the term $4L^2\mathbb{Q}_n(\mathbf{c}_n^c)$ which converges to zero by assumption, this is the standard decomposition of the risk into the estimation error and the approximation error. To prove the theorem it is sufficient to prove that

$$\inf_{g \in \mathcal{G}_n} \mathbb{E}\left[\left(g(\mathbf{X}) - g^*(\mathbf{X}) \right)^2 \right] = o_{\mathbb{P}}(1)$$
(3.6)

1760

and

$$\sup_{g \in \mathcal{G}_n} \left| \mathbb{E} \left[\left(g(\mathbf{X}) - Y \right)^2 \right] - \frac{1}{n} \sum_{i=1}^n (g(\mathbf{X}_i) - Y_i)^2 \right] \right| = o_{\mathbb{P}}(1). \tag{3.7}$$

3.3. Approximation error

Let us define $C_n^{i \setminus s} = \{ \mathbf{r} \setminus \mathbf{c}_n^s : \mathbf{r} \in C_n^i \}$ where $\mathbf{c}_n^s = \bigcup_{\mathbf{r} \in C_n^s} \mathbf{r}$. Notice that $\mathcal{P}(C_n^{i \setminus s}) = C_n^s$ $\{A \setminus \mathbf{c}_n^s : A \in \mathcal{P}(\mathcal{C}_n^i)\}$. Indeed, by Proposition 1.1:

$$\mathcal{P}(\mathcal{C}_{n}^{i \setminus s}) = \left\{ \bigcap_{\mathbf{r} \in \tilde{\mathcal{C}}} \mathbf{r} \setminus \bigcup_{\mathbf{r} \in \mathcal{C}_{n}^{i \setminus s} \setminus \tilde{\mathcal{C}}} \mathbf{r} : \tilde{\mathcal{C}} \subseteq \mathcal{C}_{n}^{i \setminus s} \right\}$$

$$= \left\{ \bigcap_{\mathbf{r} \in \tilde{\mathcal{C}}^{i}} (\mathbf{r} \setminus \mathbf{c}_{n}^{s}) \setminus \bigcup_{\mathbf{r} \in \mathcal{C}_{n}^{i} \setminus \tilde{\mathcal{C}}^{i}} (\mathbf{r} \setminus \mathbf{c}_{n}^{s}) : \tilde{\mathcal{C}}^{i} \subseteq \mathcal{C}_{n}^{i} \right\}$$

$$= \left\{ \left(\bigcap_{\mathbf{r} \in \tilde{\mathcal{C}}^{i}} \mathbf{r} \setminus \bigcup_{\mathbf{r} \in \mathcal{C}_{n}^{i} \setminus \tilde{\mathcal{C}}^{i}} \mathbf{r} \right) \setminus \mathbf{c}_{n}^{s} : \tilde{\mathcal{C}}^{i} \subseteq \mathcal{C}_{n}^{i} \right\}$$

$$= \left\{ A \setminus \mathbf{c}_{n}^{s} : A \in \mathcal{P}(\mathcal{C}_{n}^{i}) \right\}. \tag{3.8}$$

So that, by Proposition 1.1 again, $\mathcal{P}(\mathcal{C}_n)$ is a partition of \mathbf{c}_n finer than $\mathcal{P}(\mathcal{C}_n^s) \cup$ $\mathcal{P}(\mathcal{C}_n^{i \setminus s})$. Hence, with

$$\tilde{g}_n = \sum_{\substack{A \in \mathcal{P}(\mathcal{C}_n^s) \cup \mathcal{P}(\mathcal{C}_n^{i \setminus s}) \\ \mathbb{O}(A) > 0}} \mathbb{E}[Y | \mathbf{X} \in A] \mathbf{1}_A,$$

we have $\tilde{g}_n \in \mathcal{G}_n$ and

$$\inf_{g \in \mathcal{G}_n} \sqrt{\mathbb{E}\left[\left(g(\mathbf{X}) - g^*(\mathbf{X})\right)^2\right]} \leq \sqrt{\mathbb{E}\left[\left(\tilde{g}_n(\mathbf{X}) - g^*(\mathbf{X})\right)^2\right]} \\
\leq \sqrt{\mathbb{E}\left[\left(\tilde{g}_n(\mathbf{X}) - g^*(\mathbf{X})\mathbf{1}_{\mathbf{c}_n}(\mathbf{X})\right)^2\right]} \\
+ \underbrace{\sqrt{\mathbb{E}\left[g^*(\mathbf{X})^2\mathbf{1}_{\mathbf{c}_n^c}(\mathbf{X})\right]}}_{\leq L\sqrt{\mathbb{Q}(\mathbf{c}^c)}}.$$

Thus, to prove (3.6), it suffices to show that $\mathbb{W}_n = o_{\mathbb{P}}(1)$ where

$$\mathbb{W}_n := \mathbb{E}\left[\left(\tilde{g}_n(\mathbf{X}) - g^*(\mathbf{X}) \mathbf{1}_{\mathbf{c}_n}(\mathbf{X}) \right)^2 \right].$$

We will repeatedly make use of the following equalities:

$$\begin{cases}
\mathbb{E}[g^*(\mathbf{X})\mathbf{1}_{\mathbf{c}_n}(\mathbf{X}) \mid \mathbf{X} \in \mathbf{r}] = \mathbb{E}[g^*(\mathbf{X}) \mid \mathbf{X} \in \mathbf{r}] = \mathbb{E}[Y \mid \mathbf{X} \in \mathbf{r}], \\
\mathbb{V}(g^*(\mathbf{X})\mathbf{1}_{\mathbf{c}_n}(\mathbf{X}) \mid \mathbf{X} \in \mathbf{r}) = \mathbb{V}(g^*(\mathbf{X}) \mid \mathbf{X} \in \mathbf{r}), \quad \forall \mathbf{r} \subseteq \mathbf{c}_n
\end{cases}$$
(3.9)

and

$$\mathbb{V}(Y|\mathbf{X} \in \mathbf{r}) = \mathbb{V}(g^*(\mathbf{X})|\mathbf{X} \in \mathbf{r}) + \sigma^2, \qquad \forall \mathbf{r} \subseteq \mathbb{R}^d.$$
 (3.10)

Let us first remark that

$$\mathbb{W}_{n} = \mathbb{E}\left[\left(\sum_{A \in \mathcal{P}(C_{n}^{s}) \cup \mathcal{P}(C_{n}^{i \setminus s})} \mathbb{E}\left[Y \mid \mathbf{X} \in A\right] \mathbf{1}_{A}(\mathbf{X}) - g^{*}(\mathbf{X}) \mathbf{1}_{\mathbf{c}_{n}}(\mathbf{X})\right)^{2}\right] \\
= \sum_{A' \in \mathcal{P}(C_{n}^{s}) \cup \mathcal{P}(C_{n}^{i \setminus s}) \cup \{\mathbf{c}_{n}^{c}\}} \mathbb{E}\left[\left(\sum_{A \in \mathcal{P}(C_{n}^{s}) \cup \mathcal{P}(C_{n}^{i \setminus s})} \mathbb{E}\left[g^{*}(\mathbf{X}) \mathbf{1}_{\mathbf{c}_{n}}(\mathbf{X}) \mid \mathbf{X} \in A\right] \mathbf{1}_{A}(\mathbf{X}) - g^{*}(\mathbf{X}) \mathbf{1}_{\mathbf{c}_{n}}(\mathbf{X})\right]^{2} \mathbf{1}_{A'}(\mathbf{X})\right] \\
= \sum_{A' \in \mathcal{P}(C_{n}^{s}) \cup \mathcal{P}(C_{n}^{i \setminus s}) \cup \{\mathbf{c}_{n}^{c}\}} \mathbb{E}\left[\left(\mathbb{E}\left[g^{*}(\mathbf{X}) \mathbf{1}_{\mathbf{c}_{n}}(\mathbf{X}) \mid \mathbf{X} \in A'\right] - g^{*}(\mathbf{X}) \mathbf{1}_{\mathbf{c}_{n}}(\mathbf{X})\right)^{2} \mathbf{1}_{A'}(\mathbf{X})\right] \\
= \sum_{A' \in \mathcal{P}(C_{n}^{s}) \cup \mathcal{P}(C_{n}^{i \setminus s}) \cup \{\mathbf{c}_{n}^{c}\}} \mathbb{E}\left[\left(\mathbb{E}\left[g^{*}(\mathbf{X}) \mathbf{1}_{\mathbf{c}_{n}}(\mathbf{X}) \mid \mathbf{X} \in A'\right] - g^{*}(\mathbf{X}) \mathbf{1}_{\mathbf{c}_{n}}(\mathbf{X})\right)^{2} \mid \mathbf{X} \in A'\right] \\
\times \mathbb{P}\left(\mathbf{X} \in A'\right) \tag{3.11}$$

which shows that \mathbb{W}_n is a within-group variance for the variable $g^*(\mathbf{X})\mathbf{1}_{\mathbf{c}_n}(\mathbf{X})$ and the groups $\mathcal{P}(\mathcal{C}_n^s) \cup \mathcal{P}(\mathcal{C}_n^{i \setminus s}) \cup \{\mathbf{c}_n^c\}$.

According to the decomposition

$$\mathbb{W}_{n} = \underbrace{\sum_{A \in \mathcal{P}(\mathcal{C}_{n}^{s})} \mathbb{E}\left[\left(\mathbb{E}\left[Y \mid \mathbf{X} \in A\right] - g^{*}(\mathbf{X})\mathbf{1}_{\mathbf{c}_{n}}(\mathbf{X})\right)^{2} \middle| \mathbf{X} \in A\right] \mathbb{P}\left(\mathbf{X} \in A\right)}_{\mathbb{W}_{n}^{s}} + \underbrace{\sum_{A \in \mathcal{P}(\mathcal{C}_{n}^{i \setminus s})} \mathbb{E}\left[\left(\mathbb{E}\left[Y \mid \mathbf{X} \in A\right] - g^{*}(\mathbf{X})\mathbf{1}_{\mathbf{c}_{n}}(\mathbf{X})\right)^{2} \middle| \mathbf{X} \in A\right] \mathbb{P}\left(\mathbf{X} \in A\right)}_{\mathbb{W}_{n}^{i \setminus s}} + \underbrace{\mathbb{E}\left[\left(\mathbb{E}\left[g^{*}(\mathbf{X})\mathbf{1}_{\mathbf{c}_{n}}(\mathbf{X}) \mid \mathbf{X} \in \mathbf{c}_{n}^{c}\right] - g^{*}(\mathbf{X})\mathbf{1}_{\mathbf{c}_{n}}(\mathbf{X})\right)^{2} \middle| \mathbf{X} \in \mathbf{c}_{n}^{c}\right] \mathbb{P}\left(\mathbf{X} \in \mathbf{c}_{n}^{c}\right)}_{0},$$

it is sufficient to prove that $\mathbb{W}_n^s \xrightarrow[n \to \infty]{\mathbb{P}} 0$ and $\mathbb{W}_n^{i \setminus s} \xrightarrow[n \to \infty]{\mathbb{P}} 0$. To deal with \mathbb{W}_n^s , we start from the decomposition of the total variance into

the within-group and the between-group variances:

$$W_n^s = \mathbb{E}[(g^*(\mathbf{X})\mathbf{1}_{\mathbf{c}_n}(\mathbf{X}) - \mathbb{E}[g^*(\mathbf{X})\mathbf{1}_{\mathbf{c}_n}(\mathbf{X})])^2\mathbf{1}_{\mathbf{c}_n^s}(\mathbf{X})] - \mathbb{B}_n^s$$

$$\leq \mathbb{V}(g^*(\mathbf{X})\mathbf{1}_{\mathbf{c}_n}(\mathbf{X})) - \mathbb{B}_n^s$$
(3.12)

where

$$\mathbb{B}_n^s := \sum_{A \in \mathcal{P}(\mathcal{C}_n^s)} \left(\mathbb{E}\left[g^*(\mathbf{X})\mathbf{1}_{\mathbf{c}_n}(\mathbf{X}) \mid \mathbf{X} \in A\right] - \mathbb{E}\left[g^*(\mathbf{X})\mathbf{1}_{\mathbf{c}_n}(\mathbf{X})\right] \right)^2 \mathbb{P}\left(\mathbf{X} \in A\right).$$

The between group variance \mathbb{B}_n^s will be lower estimated by a function of \mathbb{W}_n^s from which an upper bound of \mathbb{W}_n^s will follow. The key point of this lower bound is to use the definition of \mathcal{C}_n^s in (H5). To make the terms controlled by this significant condition appear, we lower bound \mathbb{B}_n^s by a sum over the elements of \mathcal{C}_n^s instead of elements of $\mathcal{P}(\mathcal{C}_n^s)$ and the expectations is replaced by empirical expectations. After applying the inequality provided by the significant condition we'll have to do the reverse operation to make \mathbb{W}_n^s appear again.

Let us first lower bound \mathbb{B}_n^s by a sum over the elements of \mathcal{C}_n^s instead of elements of $\mathcal{P}(\mathcal{C}_n^s)$ and replace the expectations by empirical expectations:

$$\mathbb{B}_{n}^{s} = \sum_{\mathbf{r} \in \mathcal{C}_{n}^{s}} \sum_{A \in \mathcal{P}_{\mathcal{C}_{n}^{s}}(\mathbf{r})} \frac{1}{\#\varphi_{\mathcal{C}_{n}^{s}}(A)} \left(\mathbb{E}\left[Y \mid \mathbf{X} \in A\right] - \mathbb{E}\left[g^{*}(\mathbf{X})\mathbf{1}_{\mathbf{c}_{n}(\mathbf{X})}\right] \right)^{2} \mathbb{P}\left(\mathbf{X} \in A\right) \\
\geq \sum_{\mathbf{r} \in \mathcal{C}_{n}^{s}} \frac{1}{M(\mathcal{C}_{n}^{s}, \mathbf{r})} \sum_{A \in \mathcal{P}_{\mathcal{C}_{n}^{s}}(\mathbf{r})} \left(\mathbb{E}\left[Y \mid \mathbf{X} \in A\right] - \mathbb{E}\left[g^{*}(\mathbf{X})\mathbf{1}_{\mathbf{c}_{n}(\mathbf{X})}\right] \right)^{2} \mathbb{P}\left(\mathbf{X} \in A\right) \\
\geq \frac{1}{M(\mathcal{C}_{n}^{s})} \sum_{\mathbf{r} \in \mathcal{C}_{n}^{s}} \sum_{A \in \mathcal{P}_{\mathcal{C}_{n}^{s}}(\mathbf{r})} \left(\mathbb{E}\left[Y \mid \mathbf{X} \in A\right] - \mathbb{E}\left[g^{*}(\mathbf{X})\mathbf{1}_{\mathbf{c}_{n}(\mathbf{X})}\right] \right)^{2} \mathbb{P}\left(\mathbf{X} \in A \mid \mathbf{X} \in \mathbf{r}\right) \mathbb{P}\left(\mathbf{X} \in \mathbf{r}\right) \\
\geq \frac{1}{M(\mathcal{C}_{n}^{s})} \sum_{\mathbf{r} \in \mathcal{C}_{n}^{s}} \left(\sum_{A \in \mathcal{P}_{\mathcal{C}_{n}^{s}}(\mathbf{r})} \mathbb{E}\left[Y \mid \mathbf{X} \in A\right] \mathbb{P}\left(\mathbf{X} \in A \mid \mathbf{X} \in \mathbf{r}\right) - \mathbb{E}\left[g^{*}(\mathbf{X})\mathbf{1}_{\mathbf{c}_{n}(\mathbf{X})}\right] \right)^{2} \mathbb{P}\left(\mathbf{X} \in \mathbf{r}\right) \\
= \frac{1}{M(\mathcal{C}_{n}^{s})} \sum_{\mathbf{r} \in \mathcal{C}_{n}^{s}} \left(\mathbb{E}\left[Y \mid \mathbf{X} \in \mathbf{r}\right] - \mathbb{E}\left[g^{*}(\mathbf{X})\mathbf{1}_{\mathbf{c}_{n}(\mathbf{X})}\right] \right)^{2} \mathbb{P}\left(\mathbf{X} \in \mathbf{r}\right) \\
\geq \frac{1}{M(\mathcal{C}_{n}^{s})} \sum_{\mathbf{r} \in \mathcal{C}_{n}^{s}} \left(V_{n,\mathbf{r} \times \mathbb{R}}^{2} - \Delta_{n} \right) \mathbb{P}(\mathbf{X} \in \mathbf{r}) \tag{3.13}$$

where we applied Jensen's inequality for the third to last inequality and where

$$\Delta_n := \sup_{A \in \mathcal{B}_n} \{ V_{n,A}^2 - U_{n,A}^2 \}$$

with for any $A \in \mathcal{B}$,

$$U_{n,A} := \mathbb{E}\left[Y \mid (\mathbf{X}, Y) \in A\right] - \mathbb{E}\left[g^*(\mathbf{X})\mathbf{1}_{\mathbf{c}_n}(\mathbf{X})\right]$$
$$V_{n,A} := \mathbb{E}_n\left[Y \mid (\mathbf{X}, Y) \in A\right] - \mathbb{E}_n\left[Y\right]$$

and $\mathcal{B}_n = \{A \in \mathcal{B} \text{ s.t. } \mathbb{Q}_n(A) > n^{-\alpha}\}.$ Continuing (3.13) with the definition of \mathcal{C}_n^s in (H5) in mind,

$$\mathbb{B}_{n}^{s} \geq \frac{1}{M(\mathcal{C}_{n}^{s})} \sum_{\mathbf{r} \in \mathcal{C}_{n}^{s}} \left(\beta_{n}^{-2} (\mathbb{V}_{n}(Y | \mathbf{X} \in \mathbf{r}) - \sigma_{n}^{2}) - \Delta_{n} \right) \mathbb{P}(\mathbf{X} \in \mathbf{r})$$

$$\geq \frac{1}{M(\mathcal{C}_{n}^{s})} \sum_{\mathbf{r} \in \mathcal{C}_{n}^{s}} \left(\beta_{n}^{-2} (\mathbb{V}_{n}(Y | \mathbf{X} \in \mathbf{r}) - \sigma^{2}) - \beta_{n}^{-2} | \sigma_{n}^{2} - \sigma^{2} | - \Delta_{n} \right) \mathbb{P}(\mathbf{X} \in \mathbf{r}).$$

Since $\mathbb{W}_n^s = \sum_{A \in \mathcal{P}_{\mathcal{C}_n}(\mathcal{C}_n^s)} \mathbb{V}(g^*(\mathbf{X})\mathbf{1}_{\mathbf{c}_n}(\mathbf{X})|\mathbf{X} \in A)\mathbb{P}(\mathbf{X} \in A)$, this last term can be lower bounded by \mathbb{W}_n^s if the empirical variances are replaced by variances and the sum over \mathcal{C}_n^s by a sum over $\mathcal{P}(\mathcal{C}_n^s)$. Let us then define

$$\Delta'_n := \sup_{A \in \mathcal{B}_n} \left\{ \left| \mathbb{V}(Y | (\mathbf{X}, Y) \in A) - \mathbb{V}_n(Y | (\mathbf{X}, Y) \in A) \right| \right\}$$

and write

Together with (3.12), this yields

$$\mathbb{W}_n^s \leq \frac{\mathbb{V}(g^*(\mathbf{X})\mathbf{1}_{\mathbf{c}_n}(\mathbf{X})) + \beta_n^{-2}\Delta_n' + \beta_n^{-2}|\sigma_n^2 - \sigma^2| + \Delta_n}{1 + \beta_n^{-2}\frac{m(\mathcal{C}_n^s)}{M(\mathcal{C}_n^s)}}$$

Under (H3) and (H9), Corollary 3.1 applies and we obtain

$$\Delta_{n} = \sup_{A \in \mathcal{B}_{n}} \left\{ (V_{n,A} - U_{n,A})(V_{n,A} + U_{n,A}) \right\}$$

$$\leq \sup_{A \in \mathcal{B}_{n}} \left\{ \left| \mathbb{E}_{n}[Y \mid (\mathbf{X}, Y) \in A] - \mathbb{E}[Y \mid (\mathbf{X}, Y) \in A] \right| + \left| \mathbb{E}[g^{*}(\mathbf{X}) \mathbf{1}_{\mathbf{c}_{n}(\mathbf{X})}] - \mathbb{E}_{n}[Y] \right| \right\} \times 4L$$

$$\leq 4L \times \left(\sup_{A \in \mathcal{B}_{n}} \left| \mathbb{E}_{n}[Y \mid (\mathbf{X}, Y) \in A] - \mathbb{E}[Y \mid (\mathbf{X}, Y) \in A] \right| + \left| \mathbb{E}[g^{*}(\mathbf{X}) \mathbf{1}_{\mathbf{c}_{n}}(\mathbf{X})] - \mathbb{E}[g^{*}(\mathbf{X})] \right| + \left| \mathbb{E}[g^{*}(\mathbf{X})] - \mathbb{E}_{n}[Y] \right| \right)$$

$$\leq 4L \times \left(O_{\mathbb{P}^{*}}(n^{\alpha - 1/2}) + L\mathbb{Q}(\mathbf{c}_{n}^{c}) + O_{\mathbb{P}}(n^{-1/2}) \right)$$

$$= O_{\mathbb{P}^{*}}(n^{\alpha - 1/2})$$

since $\mathbb{Q}(\mathbf{c}_n^c) = O_{\mathbb{P}^*}(n^{\alpha-1/2})$. Corollary 3.1 also yields

$$\Delta_n' = O_{\mathbb{P}^*}(n^{\alpha - 1/2}).$$

Since it is moreover assumed that $|\sigma_n^2 - \sigma^2| = O_{\mathbb{P}}(n^{\alpha - \frac{1}{2}})$, (2.3) of Condition (H6) and Remark 9 lead to

$$\mathbb{W}_n^s \xrightarrow[n \to \infty]{\mathbb{P}} 0.$$

To deal with $\mathbb{W}_n^{i \setminus s}$, remember (3.8) and write

$$\begin{split} & \mathbb{W}_{n}^{i \backslash s} = \sum_{A \in \mathcal{P}(\mathcal{C}_{n}^{i \backslash s})} \mathbb{E}\left[\left(\mathbb{E}\left[Y \mid \mathbf{X} \in A \right] - g^{*}(\mathbf{X}) \mathbf{1}_{\mathbf{c}_{n}}(\mathbf{X}) \right)^{2} \mathbf{1}_{A}(\mathbf{X}) \right] \\ & = \sum_{A \in \mathcal{P}(\mathcal{C}_{n}^{i})} \mathbb{E}\left[\left(\mathbb{E}\left[Y \mid \mathbf{X} \in A \setminus \mathbf{c}_{n}^{s} \right] - g^{*}(\mathbf{X}) \mathbf{1}_{\mathbf{c}_{n}}(\mathbf{X}) \right)^{2} \mathbf{1}_{A \setminus \mathbf{c}_{n}^{s}}(\mathbf{X}) \right] \\ & \leq \sum_{A \in \mathcal{P}(\mathcal{C}_{n}^{i})} \mathbb{E}\left[\left(\mathbb{E}\left[Y \mid \mathbf{X} \in A \right] - g^{*}(\mathbf{X}) \mathbf{1}_{\mathbf{c}_{n}}(\mathbf{X}) \right)^{2} \mathbf{1}_{A \setminus \mathbf{c}_{n}^{s}}(\mathbf{X}) \right] \\ & \leq \sum_{A \in \mathcal{P}(\mathcal{C}_{n}^{i})} \mathbb{E}\left[\left(\mathbb{E}\left[Y \mid \mathbf{X} \in A \right] - g^{*}(\mathbf{X}) \mathbf{1}_{\mathbf{c}_{n}}(\mathbf{X}) \right)^{2} \mathbf{1}_{A}(\mathbf{X}) \right] \\ & \leq \sum_{A \in \mathcal{P}(\mathcal{C}_{n}^{i})} \mathbb{E}\left[\left(\mathbb{E}\left[Y \mid \mathbf{X} \in A \right] - g^{*}(\mathbf{X}) \mathbf{1}_{\mathbf{c}_{n}}(\mathbf{X}) \right)^{2} \mathbf{1}_{A}(\mathbf{X}) \right] \\ & \leq \frac{1}{m(\mathcal{C}_{n}^{i})} \sum_{\mathbf{r} \in \mathcal{C}_{n}^{i}} \sum_{A \in \mathcal{P}_{\mathcal{C}_{n}^{i}}} \mathbb{E}\left[\left(\mathbb{E}\left[Y \mid \mathbf{X} \in \mathbf{r} \right] - g^{*}(\mathbf{X}) \mathbf{1}_{\mathbf{c}_{n}}(\mathbf{X}) \right)^{2} \mathbf{1}_{A}(\mathbf{X}) \right] \\ & \leq \frac{1}{m(\mathcal{C}_{n}^{i})} \sum_{\mathbf{r} \in \mathcal{C}_{n}^{i}} \mathbb{E}\left[\mathbb{E}\left[Y \mid \mathbf{X} \in \mathbf{r} \right] - g^{*}(\mathbf{X}) \mathbf{1}_{\mathbf{c}_{n}}(\mathbf{X}) \right)^{2} \mathbf{1}_{\mathbf{r}}(\mathbf{X}) \right] \\ & = \frac{1}{m(\mathcal{C}_{n}^{i})} \sum_{\mathbf{r} \in \mathcal{C}_{n}^{i}} \mathbb{E}\left[\mathbb{E}\left[Y \mid \mathbf{X} \in \mathbf{r} \right] - g^{*}(\mathbf{X}) \mathbf{1}_{\mathbf{c}_{n}}(\mathbf{X}) \right)^{2} \mathbf{1}_{\mathbf{r}}(\mathbf{X}) \right] \\ & \leq \frac{1}{m(\mathcal{C}_{n}^{i})} \sum_{\mathbf{r} \in \mathcal{C}_{n}^{i}} (\mathbb{V}(Y \mid \mathbf{X} \in \mathbf{r}) - \sigma^{2}) \mathbb{P}\left(\mathbf{X} \in \mathbf{r}\right) \qquad \text{by (3.9) and (3.10)} \\ & \leq \frac{1}{m(\mathcal{C}_{n}^{i})} \sum_{\mathbf{r} \in \mathcal{C}_{n}^{i}} (\mathbb{V}_{n}(Y \mid \mathbf{X} \in \mathbf{r}) - \sigma^{2} + |\sigma_{n}^{2} - \sigma^{2}| + \Delta_{n}^{\prime}) \mathbb{P}\left(\mathbf{X} \in \mathbf{r}\right) \qquad \text{under (H5)} \\ & \leq \frac{M(\mathcal{C}_{n}^{i})}{m(\mathcal{C}_{n}^{i})} \left(\varepsilon_{n}^{2} + |\sigma_{n}^{2} - \sigma^{2}| + \Delta_{n}^{\prime} \right) \frac{\mathbb{P}}{n \to \infty} 0 \\ & \text{under (2.4) of Condition (H6) and with Remark 9}. \end{aligned}$$

3.4. Estimation error

The proof of (3.7) is inspired by (Györfi et al., 2006, Theorem 13.1) and its proof. It is not needed here to truncate the functions in \mathcal{G}_n nor Y since they are all bounded by assumption. Thus the assumptions (H8) do not need to involve the truncation constant.

Recall that \mathcal{G}_n is the set of piecewise constant functions with values in [-L, L] on the elements of the partition $\mathcal{P}(\mathcal{C}_n) \cup \{\mathbf{c}_n^c\}$ such that $\forall g \in \mathcal{G}_n, \forall \mathbf{x} \in \mathbf{c}_n^c, g(\mathbf{x}) = 0$. Then \mathcal{G}_n is a subset of

$$\mathcal{G}_c \circ \Pi_n := \left\{g: \mathbb{R}^d o \mathbb{R}: g = \sum_{A \in \mathcal{P}} f_A \mathbf{1}_A, \mathcal{P} \in \Pi_n, f_A \in \mathcal{G}_c
ight\},$$

where Π_n is a refinement of $\{\mathcal{P}(\mathcal{C}_n(\mathbf{d}_n)) : \mathbf{d}_n \in \mathcal{S}^n\}$ (see Defintion 2.3) and \mathcal{G}_c is the set of constant functions $\mathbb{R}^d \to [-L, L]$. Then

$$\sup_{g \in \mathcal{G}_n} \left| \frac{1}{n} \sum_{i=1}^n (g(\mathbf{X}_i) - Y_i)^2 - \mathbb{E} [(g(\mathbf{X}) - Y)^2] \right|$$

$$\leq \sup_{g \in \mathcal{G}_c \circ \Pi_n} \left| \frac{1}{n} \sum_{i=1}^n (g(\mathbf{X}_i) - Y_i)^2 - \mathbb{E} [(g(\mathbf{X}) - Y)^2] \right|,$$

According to (Györfi et al., 2006, Theorem 9.1 and Problem 10.4) we have,

$$\mathbb{P}\left\{ \sup_{g \in \mathcal{G}_c \circ \Pi_n} \left| \frac{1}{n} \sum_{i=1}^n (g(\mathbf{X}_i) - Y_i)^2 - \mathbb{E}\left[(g(\mathbf{X}) - Y)^2 \right] \right| > \varepsilon \right\} \\
\leq 8\mathbb{E}\left[\mathcal{N}_1 \left(\frac{\varepsilon}{32L}, \mathcal{G}_c \circ \Pi_n, \mathbf{X}_1^n \right) \right] \exp\left\{ \frac{-n\varepsilon^2}{128.(4L^2)^2} \right\}, \quad (3.14)$$

where $\mathbf{X}_{1}^{n} = {\{\mathbf{X}_{1}, \dots, \mathbf{X}_{n}\}}.$

Here \mathcal{N}_1 $(\varepsilon, \mathcal{G}_c \circ \Pi_n, \mathbf{X}_1^n)$ is the random variable corresponding to the minimal number $N \in \mathbb{N}$ such that there exist functions $g_1, \ldots, g_N : \mathbb{R}^d \to [-L, L]$ with the property that for every $g \in \mathcal{G}_c \circ \Pi_n$ there is a $j \in \{1, ..., N\}$ such that

$$\frac{1}{n}\sum_{i=1}^{n}|g(\mathbf{X}_i)-g_j(\mathbf{X}_i)|\leq \varepsilon.$$

This number is called the ε -covering number of $\mathcal{G}_c \circ \Pi_n$. It can be interpreted as the complexity of the class. Then using (Györfi et al., 2006, Lemma 13.1) (it can be checked to hold in the current situation, in particular for a family of finite collections of disjoint sets instead of a family of partitions of \mathbb{R}^d) we have

$$\mathcal{N}_{1}\left(\frac{\varepsilon}{32L}, \mathcal{G}_{c} \circ \Pi_{n}, \mathbf{X}_{1}^{n}\right) \\
\leq \widetilde{\Delta}_{n}(\Pi_{n}) \left\{ \sup_{z_{1}, \dots, z_{m} \in \{\mathbf{X}_{1}, \dots, \mathbf{X}_{n}\}, m \leq n} \mathcal{N}_{1}\left(\frac{\varepsilon}{32L}, \mathcal{G}_{c}, z_{1}^{m}\right) \right\}^{\widetilde{\mathcal{M}}(\Pi_{n})},$$

According to (Györfi et al., 2006, Lemma 9.2) for any set of functions \mathcal{G} and any sample z_1^m we have

$$\mathcal{N}_1\left(\frac{\varepsilon}{32L},\mathcal{G},z_1^m\right) \leq \mathcal{M}_1\left(\frac{\varepsilon}{32L},\mathcal{G},z_1^m\right),$$

where $\mathcal{M}_1(\varepsilon, \mathcal{G}, z_1^m)$ is the maximal $N \in \mathbb{N}$ such that there exist functions $g_1, \ldots, g_N \in \mathcal{G}$ with

$$\frac{1}{n}\sum_{i=1}^{m}|g_j(z_i)-g_k(z_i)|\geq \varepsilon,$$

for all $1 \leq j < k \leq N$. It is called L_1 ε -packing of \mathcal{G} on z_1^m . See (Györfi et al., 2006, Definition 9.4 (c)).

Now, from the definition of \mathcal{G}_c ,

$$\sup_{z_1,\dots,z_m\in\{\mathbf{X}_1,\dots,\mathbf{X}_n\},m\leq n} \mathcal{M}_1\left(\varepsilon,\mathcal{G}_c,z_1^m\right) \leq \left\lceil \frac{2L}{\varepsilon} \right\rceil + 1.$$

Finally,

$$\sup_{z_1, \dots, z_m \in \{\mathbf{X}_1, \dots, \mathbf{X}_n\}, m \le n} \mathcal{N}_1\left(\frac{\varepsilon}{32L}, \mathcal{G}_c \circ \Pi_n, z_1^m\right) \le \widetilde{\Delta}(\Pi_n) \left(\left\lceil \frac{64L^2}{\varepsilon} \right\rceil + 1\right)^{\widetilde{\mathcal{M}}(\Pi_n)}.$$
(3.15)

According to (3.14) and (3.15) we have:

$$\mathbb{P}\left\{ \sup_{g \in \mathcal{G}_n} \left| \frac{1}{n} \sum_{i=1}^n |g(\mathbf{X}_i) - Y_i|^2 - \mathbb{E}\left[|g(\mathbf{X}) - Y|^2 \right] \right| > \varepsilon \right\} \\
\leq 8\widetilde{\Delta}_n(\Pi_n) \left(\left\lceil \frac{64L^2}{\varepsilon} \right\rceil + 1 \right)^{\widetilde{\mathcal{M}}(\Pi_n)} \exp\left(-\frac{n\varepsilon^2}{128.(4L^2)^2} \right)$$

and since

$$8\widetilde{\Delta}_{n}(\Pi_{n})\left(\left\lceil\frac{64L^{2}}{\varepsilon}\right\rceil+1\right)^{\widetilde{\mathcal{M}}(\Pi_{n})}\exp\left(-\frac{n\varepsilon^{2}}{2048L^{4}}\right)$$

$$=8\exp\left(-\frac{n}{L^{4}}\left(\frac{\varepsilon^{2}}{2048}-\frac{\log\widetilde{\Delta}_{n}(\Pi_{n})L^{4}}{n}-\frac{\widetilde{\mathcal{M}}(\Pi_{n})L^{4}\log\left(\left\lceil\frac{64L^{2}}{\varepsilon}\right\rceil+1\right)}{n}\right)\right),$$

this concludes the proof of (3.7) and of Theorem 2.1.

4. Illustrations

In this section we propose a simple algorithm to generate data-dependent coverings using either Random Forests (RF, Breiman (2001)) or Gradient Boosted trees (GB, Friedman (2001)) or Stochastic Gradient Boosting trees (SGB, Friedman (2002)) as rule generator. The interest is twofold; First, it exhibits examples of data-dependent quasi-coverings that are likely to be suitable as in Definition 2.1. Second, we observe that the interpretability indices are small when applying our algorithm predicting real data selecting suitable sets of rules. Note that all rules considered here are based on intervals of $\mathbb R$ as in Remark 6 and then their length is automatically smaller than or equal to d.

4.1. Covering algorithm

The proposed algorithm generates an estimator based on a data-dependent quasi-covering:

- 1. Generate trees with a given method among RF, GB or SGB with a maximal tree size, *tree_size*, and a maximal number of generated rules (all the nodes and leaves of the trees). *max_rules*.
- nodes and leaves of the trees), max_rules . 2. For a chosen $\alpha \in (0, 1/2)$, set $\beta_n = n^{\alpha/2-1/4}$ and $\varepsilon_n = \beta_n s_n$, where s_n is the empirical standard deviation of Y. Keep all rules which length is less than or equal to $l_max \in \{1, \ldots, d\}$ and which fulfill (H3).
- 3. Split this set of rules into two sets: The set of significant rules S_n (i.e., rules \mathbf{r} that fulfill $\beta_n \big| \mathbb{E}_n[Y|\mathbf{X} \in \mathbf{r}] \mathbb{E}_n[Y] \big| \ge \sqrt{(\mathbb{V}_n(Y|\mathbf{X} \in \mathbf{r}) \sigma_n^2)_+}$ and the set of insignificant rules I_n (i.e., rules \mathbf{r} that are not in S_n and that fulfill $\varepsilon_n \ge \sqrt{(\mathbb{V}_n(Y|\mathbf{X} \in \mathbf{r}) \sigma_n^2)_+}$).
- 4. Select a minimum set of rules C_n using Algorithm 1 given in the appendix. If at any step the set of selected rules forms a covering, the selection process is stopped. A rule is added to the currently selected set of rules if and only if it has at least a rate $1 \gamma \in (0, 1)$ of points not covered by the current set of rules². The set S_n , sorted by decreasing empirical coverage rate, is browsed first. Then, if necessary, the set I_n , sorted by increasing empirical variance, is browsed.

As explained in Definition 2.1 the set C_n is the union of the significant set $C_n^s = C_n \cap S_n$ and the insignificant set $C_n^i = C_n \cap I_n$. In practice the set C_n^s is the most interesting one because it identifies the rules where the conditional mean is prominent.

The datasets and the code for the illustrations are available on GitHub. The code is written in both the Python and R languages.

 $^{^{2}}$ This step was already described in Remark 4 in order to fulfill the redundancy condition (H6).

4.2. Artificial data

We consider here the same model as in Friedman and Popescu (2008). We generate n = 5000 observations following the regression setting

$$Y = g^*(\mathbf{X}) + Z,$$

where d = 100 (the dimension of **X**) and

$$g^*(\mathbf{X}) = 9 \prod_{j=1}^{3} \exp\left(-3\left(1 - X_j\right)^2\right) - 0.8 \exp\left(-2\left(X_4 - X_5\right)\right) + 2\sin^2(\pi \cdot X_6) - 2.5\left(X_7 - X_8\right), \quad (4.1)$$

and $Z \sim \mathcal{N}(0, \sigma^2)$. The value of $\sigma > 0$ was chosen to produce a two-to-one signal-to-noise ratio. The variables were generated from a uniform distribution on $\{0/10, \ldots, 9/10\}$. It is important to note that only the eight first variables are informative; the 92 others are just noise. The coefficients that multiply each of the terms in g^* have been chosen to ensure that the variables have approximately the same influence.

We evaluate the accuracy of the estimators with the mean squared error defined by

$$MSE = \frac{\mathbb{E}_{\mathbb{Q}}\left[(Y - g_n(\mathbf{X}))^2\right]}{\mathbb{V}(Y)}.$$

In order to evaluate the error without the noise variance, we also consider the following criterion:

$$MSE^* = \frac{\mathbb{E}_{\mathbb{Q}}\left[\left(g^*(\mathbf{X}) - g_n(\mathbf{X})\right)\right]}{\mathbb{V}(g^*(\mathbf{X}))}.$$

We approximate the criteria MSE and MSE^* with 50000 test observations sampled independently from \mathbb{Q} .

4.2.1. Execution

We run M=100 simulations with 5000 independent observations. For each simulation we compare the Covering Algorithm with RF as rules generator (Covering) with a classical Random Forest (RF) and with RuleFit (Friedman and Popescu, 2008) with rules only.

RuleFit is a very accurate rule-based algorithm. First it generates a list of rules by considering all nodes and leaves of a boosted tree ensemble ISLE (Friedman et al., 2003). Then rules are used as features in a sparse linear regression model obtained by Lasso (Tibshirani, 1996).

For RF, we set the number of trees at 100. For RuleFit and Covering we set the maximal number of generated rules at 4000 and the maximal length of a rule is fixed at $k_{max} = 3$. And we set $\alpha = 1/2 - 1/100$ and $\gamma = 0.90$ for Covering.

4.2.2. Results

The MSE and MSE^* and the interpretability index (1.4) for each algorithm across the experiments are summarized in Table 1 and plotted in Figure 5.

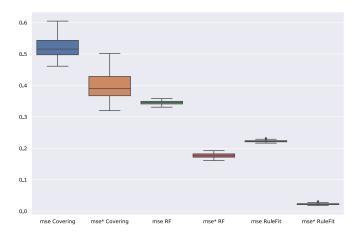


Fig 5. MSE on 100 realizations from model (4.1) for Covering Algorithm (Covering), Random Forest (RF) and RuleFit.

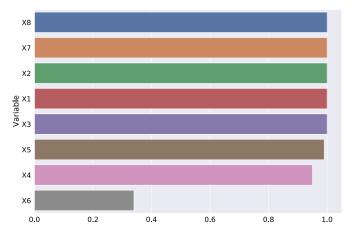


FIG 6. For each informative variable, frequency of occurrence in at least one rule selected by the Covering Algorithm in 100 independent simulations.

4.2.3. Comments

Figure 6 shows the frequency of occurrence of a variable in at least one rule of the selected set of rules. Note that only informative variables are involved in the selected rules. Moreover, in one thirds of the experiments, all informative

Table 1 Number of rules generated by the algorithm (Nb Rules), interpretability index (1.4) (Interpretability) and mean squared errors MSE and MSE* for each algorithm.

Random Forest	Nb rules	Interpretability	MSE	MSE*
mean	630718.78	6871475.52	0.35	0.18
std	429.56	28124.49	0.01	0.01
min	629458.00	6815078.00	0.33	0.16
25%	630434.00	6848260.00	0.34	0.17
50%	630701.00	6871049.00	0.35	0.18
75%	631020.50	6890434.00	0.35	0.18
max	631790.00	6947934.00	0.36	0.19

Covering Algorithm	Nb rules	Interpretability	MSE	MSE*
mean	15.89	43.16	0.52	0.40
std	2.53	8.65	0.03	0.04
min	10.00	23.00	0.46	0.32
25%	14.00	36.00	0.50	0.37
50%	16.00	42.00	0.52	0.39
75%	17.00	48.00	0.54	0.43
max	24.00	70.00	0.60	0.50

RuleFit	Nb rules	Interpretability	MSE	MSE*
mean	360.65	1155.63	0.22	0.02
std	55.30	232.09	0.00	0.00
\min	253.00	747.00	0.22	0.02
25%	314.75	965.25	0.22	0.02
50%	358.50	1143.00	0.22	0.02
75%	397.00	1297.25	0.22	0.02
max	507.00	1792.00	0.23	0.03

variables are identified. It means the support of g^* is well identified in these cases. In the most of the experiences, the variable X6 is the only variable that is not involved. One reason could be that RF is a random generator of rules that may not capture the importance of X6 at every run. The problem could be solved by considering a deterministic rule generator algorithm.

Table 1 emphasizes that Covering Algorithm generates more interpretable models than RF and RuleFit according to the interpretability index defined in (1.4). Indeed, for the same constraints on the maximal number of rules and the maximal length of the rules, Covering Algorithm selects far fewer rules than RuleFit. Nevertheless, RuleFit is much more accurate than Covering Algorithm. The accuracy of RuleFit is high even if many noise variables are included in its generated model. Therefore, a posteriori analysis of the importance of the variables and rules of the generated model of RuleFit is crucial, see Section 9.1.2 of Friedman and Popescu (2008) for more information.

4.3. Real data

For this application, we consider seven public datasets presented in Table 2. with different dimension and number of observations. For the dataset *Student*

Name	$(n \times d)$	Short description
Diabetes	443 × 10	Prediction of quantitative measure of disease progression one year after baseline (Efron et al., 2004).
Prostate	97 × 8	Prediction of the level of prostate-specific antigen based on clinical measures in men who were about to receive a radical prostatectomy (Hastie, Friedman and Tibshirani, 2001).
Ozone	330 × 9	Prediction of atmospheric ozone concentration from daily meteorological measurements (Hastie, Friedman and Tibshirani, 2001).
Machine	209 × 8	Prediction of published relative performance (Dua and Graff, 2017).
MPG	398 × 8	Prediction of city-cycle fuel consumption in miles per gallon (Dua and Graff, 2017).
Boston	506×13	Prediction of the median price of neighborhoods, (Harrison Jr and Rubinfeld, 1978).
Student	649 × 32	Prediction of the final grade of students based on attributes collected by reports and questionnaires (Cortez and Silva, 2008).

we have removed the variables G1 and G2 that are the first and the second grade respectively because the target attribute G3 has a strong correlation with G2 and G1. In Cortez and Silva (2008) the authors state that despite the complexity of the regression model, the prediction of G3 without G2 and G1 is much more useful in practice.

4.3.1. Execution

For each dataset we run 20 executions. For each execution, the data are randomly split into a training set and a test set, with 70% and 30% ratios respectively. As a baseline of accuracy we consider RF (Breiman, 2001) and as a baseline of interpretability we consider the CART algorithm (Breiman et al., 1984). We compare Covering Algorithm with RF as rule generator (CA RF), GB trees as rule generator (CA GB), SGB trees as rule generator (CA SGB) with RuleFit (Friedman and Popescu, 2008) (with rules only), Node harvest (Meinshausen, 2010) (Nh) and SIRUS (Bénard et al., 2020). These algorithms are the main existing interpretable rule-based algorithms for regression settings that are based on tree ensembles.

Node harvest uses also a tree ensemble as rule generator. The algorithm considers all nodes and leaves of RF as rules and solves a linear quadratic problem to fit a weight to each rule. So the prediction is a convex combination of rules.

SIRUS (Stable and Interpretable RUle Set) is designed as a stable predictive algorithm. It uses a custom RF algorithm to generate many rules (considering nodes and leaves). Then it selects the rules with a rate of occurrence greater than a tuning parameter p_0 . To ensure a large number of occurrences on rules, the features are discretized.

The maximal number of leaves of CART is set at $tree_size = 10$ and the number of tree for RF is set at $nb_tree = 1000$. For CA and RuleFit we set similarly the maximal number of rules generated by the rule generator $max_rules = 4000$ and the size of a tree $tree_size = 8$. For CA, Node harvest and SIRUS we set identically the maximum rule length $l_max = 3$. For CA we also set $\alpha = 1/2 - 1/100$ and $\gamma = 0.90$. In these applications the real value of σ^2 is unknown. We estimate it by σ_n^2 the minimal variance of the generated rules fulfilling the covering condition (H3). We have no guarantee that this estimator is good enough (see Remark 3). For SIRUS, the hyperparameter p0 is estimated by 10-fold cross-validation and the maximal number of rules is set at 25. The parameters setting is summarized in Table 3.

Table 3
Algorithm parameter settings.

Algorithm	Parameters
CART	$tree_size = 10.$
RF	$nb_tree = 1000.$
RuleFit	$tree_size = 8,$
	$max_rules = 4000.$
Node harvest (Nh)	$l_max = 3.$
SIRUS	$l_max = 3,$
	$max_selected_rules = 25,$
	ncv = 10.
CA	$max_rules = 4000,$
	$tree_size = 8,$
	alpha = 1/2 - 1/100,
	gamma = 0.90,
	l- $max = 3$.

4.3.2. Results

Experimental results are gathered in Table 4. For each dataset and algorithm we present an average over the 20 executions of the number of selected rules (Nb rules), the empirical coverage rate on the training set (Cov), the interpretability index (1.4) (Int) and the mean squared error on the test set divided by the empirical variance of the target (MSE).

These results emphasize that suitable data-dependent quasi-coverings (see Definition 2.1) are very efficient to generate an interpretable rule-based model, see also Section 4.3.3. They generate very simple models with an interpretability index (1.4) much lower than the other algorithms. Moreover the sets of selected rules C_n form in most cases a covering of the training set. CA achieves a good

Table 4

Average of the number of rules (NB rules), the empirical coverage rate on the training set (Cov), the interpretability index (Int) and the ratio mean squared error/empirical variance of the target (MSE) over 20 executions of usual interpretable algorithms for various regression public datasets. Best values are in bold, as well as values within 10% of the best for each dataset.

Dataset	Dataset Regression Tree (CART)				Random Forest (RF)			
Dataset	Nb rules	Int	MSE	Nb rules	Int	MSE		
Diabetes	10	30.2	1.14	191190.6	1070454.75	0.58		
Prostate	10	25.3	1.11	42047.95	169328.6	0.56		
Ozone	10	32	0.5	108082.95	572067.75	0.28		
Machine	10	24.2	0.18	73754.25	319376.45	0.06		
MPG	10	29.0	0.27	148954.4	623266.0	0.14		
Boston	10	26.3	0.29	212961.05	1270252.25	0.14		
Student	10	37.75	1.55	199278.65	1838107.25	0.74		

Dataset	CA RF			CA GB				
Dataset	Nb rules	Cov	Int	MSE	Nb rules	Cov	Int	MSE
Diabetes	7.6	0.99	17.95	0.70	9.1	0.99	19.7	0.74
Prostate	5.8	1.0	13.95	0.85	6.05	1.0	13.35	0.86
Ozone	3.3	0.99	5.85	0.42	3.95	0.99	6.55	0.39
Machine	2.85	0.99	5.05	0.40	3.25	0.99	6.05	0.50
MPG	3.1	0.99	4.3	0.28	3.65	0.99	6.85	0.30
Boston	4.9	0.99	10.9	0.45	6.8	0.99	15.15	0.45
Student	18.45	0.99	51.4	0.86	32	0.98	86.75	1.01

Dataset CA SGB				RuleFit				
Dataset	Nb rules	Cov	Int	MSE	Nb rules	Cov	Int	MSE
Diabetes	8.75	0.99	19.8	0.67	249.35	1.0	1168.7	0.71
Prostate	6.5	1.0	15.05	0.76	78.4	1.0	268.8	0.66
Ozone	4.05	0.99	6.5	0.38	187.45	1.0	851	0.28
Machine	3.2	0.99	6.1	0.52	99.85	1.0	278.15	0.08
MPG	3.75	0.99	7.25	0.29	174.7	1.0	745.15	0.13
Boston	6.35	0.99	13.85	0.46	264.95	1.0	1207.65	0.12
Student	37.3	0.98	100.05	1.01	277.95	1.0	1280.4	0.90

Dataset		NodeHarvest			SIRUS			
Dataset	Nb rules	Cov	Int	MSE	Nb rules	Cov	Int	MSE
Diabetes	227.15	1.0	613.05	0.59	21.25	0.99	32.35	0.57
Prostate	66	1.0	160.35	0.57	23.25	0.99	33.7	0.59
Ozone	181.45	1.0	472.9	0.29	24.9	1.0	39.4	0.32
Machine	67.15	1.0	159.0	0.32	13.9	1.0	21.85	0.26
MPG	100.15	1.0	257.85	0.15	25.0	1.0	45.2	0.21
Boston	139.85	1.0	345.4	0.20	25.0	0.99	38.5	0.25
Student	78.5	1.0	228.45	0.76	15.35	1.0	26.2	0.77

interpretability-accuracy trade-off with emphasis on interpretability except for the dataset *Student* that might be too complex to be interpretable. SIRUS appears as a good challenger achieving a more balanced trade-off. Finally, these results show that the choice of the rule generator for the CA has an impact on the interpretability index and on the accuracy.

4.3.3. Examples of interpretation

Any model generated by a CA can be summarized in a table of rules. We present one set of rules selected by CA RF for the dataset *Ozone* in Table 5. The measure Δ_n is the empirical mean deviation ratio of the prediction from the mean of Y.

Table 5

Summary of the rules selected by CA RF for the dataset Ozone. All rules are significant. The temp variable is the Sandburg Air Force base temperature in degrees Fahrenheit, the ibt variable is the inversion base temperature at LAX in degrees Fahrenheit and the humidity variable is the humidity in percent at LAX.

Rule	Conditions	Coverage	Prediction	Std	Δ_n
R1	$temp \in [30, 67.5]$	0.65	7.99	4.57	-0.2
R2	$ibt \in [-15, 191]$	0.60	8.02	4.5	-0.2
R3	$ibt \in [172, 326]$ $humidity \in [35.5, 92]$	0.41	18.61	7	0.8
R4	$ibt \in [172, 326]$ $humidity \in [19, 35.5]$	0.08	6.11	2.64	-0.4

Table 5 is a description of the model. Together with Table 6 we are able to translate in natural language the significant rules in Table 5, according to the definition of interpretability of Biran and Cotton (2017).

Table 6

Description of the variables selected by CA RF for the dataset Ozone. The Y variable is the daily maximum of the hourly average ozone concentration in Upland.

	Y	temp	ibt	humidity
mean	11.78	61.75	161.16	58.13
std	8.01	14.46	76.68	19.87
min	1.00	25.00	-25.00	19.00
25%	5.00	51.00	107.00	47.00
50%	10.00	62.00	167.50	64.00
75%	17.00	72.00	214.00	73.00
max	38.00	93.00	332.00	93.00

- Rule R1 has a length equal to 1 and suggests that if the *temp* variable is reasonably low then the daily maximum of the hourly average ozone concentrations in Upland tends to be low.
- Rule R2 has a length equal to 1 and suggests that if the *ibt* variable is not too high then the target tends to be low.
- Rules R3 and R4 have a length equal to 2 and suggest that CA has detected an interaction between high *ibt* and the *humidity* variable. Indeed, high value of *ibt* and high *humidity* indicate a high hourly average ozone concentrations and high value of *ibt* and low *humidity* indicate a low hourly average ozone concentrations.

4.3.4. Comments

In Figure 7 and Figure 8 we present boxplots of the MSE over the 20 executions for each algorithm for the *Machine* and *Ozone* datasets. These two datasets has

been chosen to show that CA can have both a very variable accuracy (Figure 7) and a stable and a good accuracy (Figure 8) for two datasets with the same dimension.

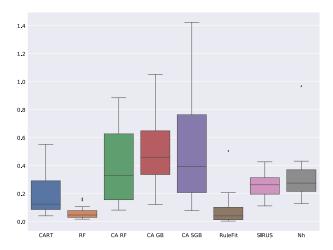


FIG 7. Boxplots of the ratio mean squared error/empirical variance of the target (MSE) over 20 executions of algorithms for the Machine dataset.

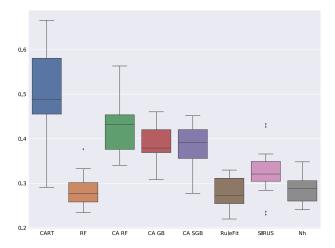


FIG 8. Boxplots of the ratio mean squared error/empirical variance of the target (MSE) over 20 executions of algorithms for the Ozone dataset.

One reason could be that rule generators of CA are random generators of rules that may not capture the good features at every run. In particular we have identified the problem of predicting the empty cells of the partition generated by the quasi-coverings. Then, by convention, $g_n(\mathbf{x}) = 0$ for $\mathbf{x} \in A$ with $\mathbb{Q}_n(A) = 0$. Thus, the accuracy of CA strongly depends on the coverage $\mathbb{Q}_n(A)$

of the sets $A \in \mathcal{P}(\mathcal{C}_n)$ which can be too small to be stable. In Definition 2.1 of suitable data-dependent quasi-coverings only the coverage of the elements of \mathcal{C}_n is controlled. The problem can be solved by considering a deterministic rule generator.

5. Conclusion and perspectives

In this paper, we provide a general framework for studying the consistency of rule-based interpretable estimators. We introduce the definition of a suitable quasi-covering. It is composed of two types of sets, namely the significant sets and the insignificant sets. The significant sets are considered as interpretable sets by construction. The insignificant ones are sets whose variance tends to zero. We provide a Covering Algorithm that extracts a suitable data-dependent quasi-covering from any rule generator.

In Section 4.2, we run a Monte Carlo experiment on Covering Algorithm applied to Random Forest. We compare its results with those of Random Forest (Breiman, 2001) and RuleFit (Friedman and Popescu, 2008). This experiment shows that Covering Algorithm, which strives for interpretability, also identifies the support of the regression function.

In Section 4.3, we apply Covering Algorithm to Random Forest (Breiman, 2001), Gradient Boosting (Friedman, 2001) and Stochastic Gradient Boosting (Friedman, 2002) and we compare their results with those of CART (Breiman et al., 1984), Random Forest (Breiman, 2001), RuleFit (Friedman and Popescu, 2008), Node harvest (Meinshausen, 2010) and SIRUS (Bénard et al., 2020). The loss of accuracy in the prediction is the cost of having an interpretable model according to our definition of interpretability. We broaden the accuracy-interpretability trade-off of classical algorithms by providing a much more interpretable method that remains consistent.

Our methodology based on quasi-coverings is very effective in generating interpretable models. The use of tree ensembles such as RF, GB or SGB as rule generator is questionable; In Section 4.3.4, we pointed out the possible negative effect of the randomization procedure in combination with the ERM principle instead of averaging. We also noted that the choice of the rule generator has an important effect on accuracy, interpretability and stability. We are now looking for an algorithm that satisfies (H8) and that deterministically generates significant and insignificant rules that form a suitable sequence of data-dependent quasi-coverings.

In practice the variance σ^2 is unknown and has to be estimated efficiently enough. We let for future work the difficult question of defining a rule-based estimator of the variance with rate of convergence $O_{\mathbb{P}}(n^{\alpha-1/2})$ as required in Theorem 2.1. Our setting could also be broadened; unbounded Y may be considered by introducing a truncation operator as in Györfi et al. (2006); strong consistency and rates of convergence of the data-dependent covering estimators may be established under regularity conditions on g^* . Finally, the scope could be easily adapted from the regression setting to the classification setting by adapting the significant condition accordingly.

Appendix

We gather here some proofs and provide the pseudo-code of the Covering Algorithm's selection process.

Proof of Proposition 1.1

Proof. By the definition of $\varphi_{\mathcal{C}}$, for any $\mathbf{x} \in \mathbb{R}^d$, $\varphi_{\mathcal{C}}^{-1}(\varphi_{\mathcal{C}}(\mathbf{x})) = \bigcap_{\mathbf{r} \in \mathcal{C}: \mathbf{x} \in \mathbf{r}} \mathbf{r} \setminus \bigcup_{\mathbf{r} \in \mathcal{C}: \mathbf{x} \notin \mathbf{r}} \mathbf{r}$. Thus

$$A \in \mathcal{P}(\mathcal{C}) = \varphi_{\mathcal{C}}^{-1}(Im(\varphi_{\mathcal{C}})) \iff \exists \mathbf{x} \in \mathbb{R}^d / A = \varphi_{\mathcal{C}}^{-1}(\varphi_{\mathcal{C}}(\mathbf{x})) = \bigcap_{\substack{\mathbf{r} \in \mathcal{C} \\ \mathbf{x} \in \mathbf{r} \\ \mathbf{x} \notin \mathbf{r}}} \mathbf{r} \setminus \bigcup_{\substack{\mathbf{r} \in \mathcal{C} \\ \mathbf{x} \notin \mathbf{r} \\ \mathbf{r} \neq \mathbf{r}}} \mathbf{r}$$
$$\iff \exists \tilde{\mathcal{C}} \subseteq \mathcal{C} / A = \bigcap_{\mathbf{r} \in \tilde{\mathcal{C}}} \mathbf{r} \setminus \bigcup_{\mathbf{r} \in \mathcal{C} \setminus \tilde{\mathcal{C}}} \mathbf{r} \text{ and } A \neq \emptyset. \quad \Box$$

Proof of Proposition 2.1

• Lower bound: Consider the set \mathcal{E}_1 of the points $x \in \mathbb{R}^d$ whose Proof. coordinates are all zero except x[i] = 2, for $1 \le i \le d$ and the set \mathcal{E}_2 of the points $x \in \mathbb{R}^d$ whose coordinates are all zero except x[i] = 1and x[i+1] = -1, for $1 \le i \le d$ (with the convention x[d+1] = x[1]). If \mathcal{P} is a partition by hyperrectangles and $C \in \mathcal{P}$, then no pair of the 2d points of $\mathcal{E}_1 \cup \mathcal{E}_2$ can be in the same element of $\mathcal{P} \setminus \{C\}$. Indeed, suppose that $(x,y) \in \mathcal{E}_1^2$ (or analogously $\in \mathcal{E}_1^2$) lies in some $D \in \mathcal{P} \setminus \{C\}$, with x[i] = 2 and y[j] = 2. Since R is an hyperrectangle, the point z whose coordinates are all zero except $z[i] = \frac{x[i] + y[i]}{2} = 1$ and $z[j] = \frac{x[j] + y[j]}{2} = 1$ is in R too. Thus C and R intersect. In the same way, suppose that $(x, y) \in \mathcal{E}_1 \times \mathcal{E}_2$ lies in some $D \in \mathcal{P} \setminus \{C\}$, with x[i] = 2, y[j] = 1 and y[j+1] = -1. Then, the point z whose coordinates are all zero except z[i] = 1 (while z[j] = x[j] = 0 and z[j+1] = x[j+1] = 0 is in R too. Thus C and R intersect. Thus, the 2d points of $\mathcal{E}_1 \cup \mathcal{E}_2$ induce 2d distinct elements of \mathcal{P} that are added to C.

• Upper bound: The partition

$$\mathcal{P} = \{ (-\infty, 0) \times \mathbb{R}^{d-1}; (1, +\infty) \times \mathbb{R}^{d-1};$$

$$[0, 1] \times (-\infty, 0) \times \mathbb{R}^{d-2}; [0, 1] \times (1, +\infty) \times \mathbb{R}^{d-2};$$

$$\dots;$$

$$[0, 1]^{d-1} \times (-\infty, 0); [0, 1]^{d-1} \times (1, +\infty);$$

$$[0, 1]^{d} \}$$

fulfills the conditions of the proposition.

2. The covering $\mathcal{C} = \{[0,1]^d ; \mathbb{R}^d\}$ fulfills the conditions of the proposition.

Proof of Proposition 3.2

For any $f: \mathcal{S} \to \mathbb{R}$ in $\mathcal{L}^1(\mathbb{Q})$ we note classically $\mathbb{Q}f := \int f d\mathbb{Q}$, $\mathbb{Q}_n f := \int f d\mathbb{Q}_n$. Proof. Let $\varepsilon > 0$. First, for any $f \in \mathcal{F}$ and $A \in \mathcal{B}_n$, since $\mathbb{Q}_n(A) > 0$ and then $\mathbb{Q}(A) > 0$,

$$\begin{aligned}
|\mathbb{E}_{n}\left[f\mid A\right] - \mathbb{E}\left[f\mid A\right]| \\
&= \left|\frac{\int_{A} f d\mathbb{Q}_{n}}{\mathbb{Q}_{n}(A)} - \frac{\int_{A} f d\mathbb{Q}}{\mathbb{Q}(A)}\right| \\
&= \left|\frac{\mathbb{Q}(A)\left(\int_{A} f d\mathbb{Q}_{n} - \int_{A} f d\mathbb{Q}\right) + (\mathbb{Q}(A) - \mathbb{Q}_{n}(A))\int_{A} f d\mathbb{Q}}{\mathbb{Q}(A)\mathbb{Q}_{n}(A)}\right| \\
&\leq \left|\frac{\int_{A} f d\mathbb{Q}_{n} - \int_{A} f d\mathbb{Q}}{\mathbb{Q}_{n}(A)}\right| + \left|(\mathbb{Q}(A) - \mathbb{Q}_{n}(A))\frac{\int_{A} f d\mathbb{Q}}{\mathbb{Q}(A)\mathbb{Q}_{n}(A)}\right|.
\end{aligned} (A.1)$$

Now, according to Proposition 3.1,

$$\sup_{\tilde{f}\in\mathcal{F},\tilde{A}\in B}\left|\int_{\tilde{A}}\tilde{f}d\mathbb{Q}_n-\int_{\tilde{A}}\tilde{f}d\mathbb{Q}\right|=O_{\mathbb{P}^*}(n^{-1/2}),$$

and

$$\sup_{\tilde{A} \in \mathcal{B}} \left| \mathbb{Q}_n(\tilde{A}) - \mathbb{Q}(\tilde{A}) \right| = O_{\mathbb{P}^*}(n^{-1/2}).$$

Thus, According to Remark 7, there exists M>0 such that for any n large enough,

$$\mathbb{P}^* \left\{ \sup_{\tilde{f} \in \mathcal{F}, \tilde{A} \in B} \left| \int_{\tilde{A}} \tilde{f} d\mathbb{Q}_n - \int_{\tilde{A}} \tilde{f} d\mathbb{Q} \right| > M n^{-1/2} \right\} < \frac{\varepsilon}{2},$$

and

$$\mathbb{P}^* \left\{ \sup_{\tilde{A} \in \mathcal{B}} \left| \mathbb{Q}_n(\tilde{A}) - \mathbb{Q}(\tilde{A}) \right| > M n^{-1/2} \right\} < \frac{\varepsilon}{2},$$

so that $\mathbb{P}^*(\Omega_n) \geq 1 - \varepsilon$ with

$$\Omega_n := \left\{ \sup_{\tilde{f} \in \mathcal{F}, \tilde{A} \in B} \left| \int_{\tilde{A}} \tilde{f} d\mathbb{Q}_n - \int_{\tilde{A}} \tilde{f} d\mathbb{Q} \right| \le M n^{-1/2} \right\} \bigcap \left\{ \sup_{\tilde{A} \in B} \left| \mathbb{Q}_n(\tilde{A}) - \mathbb{Q}(\tilde{A}) \right| \le M n^{-1/2} \right\}.$$

Then (A.1) yields, with $c := \sup_{f \in \mathcal{F}, \mathbf{x} \in \mathcal{S}} |f(\mathbf{x})| < \infty$ and since $\mathbb{Q}_n(A) \ge n^{-\alpha}$, for n large enough, in the event Ω_n ,

$$\sup_{f \in \mathcal{F}, A \in \mathcal{B}_n} \left| \mathbb{E}_n \left[f \mid A \right] - \mathbb{E} \left[f \mid A \right] \right| \le M n^{\alpha - 1/2} (1 + c),$$

since $\frac{\int_A f d\mathbb{Q}}{\mathbb{Q}(A)} \leq c$. Finally, it has been proved that $\forall \varepsilon > 0, \exists M > 0, \exists N \in \mathbb{N}^* / \forall n \geq N$,

$$\mathbb{P}^* \left\{ \sup_{f \in \mathcal{F}, A \in \mathcal{B}_n} \left| \mathbb{E}_n \left[f \mid A \right] - \mathbb{E} \left[f \mid A \right] \right| > M n^{\alpha - 1/2} \right\} < \varepsilon,$$

and then $\forall \varepsilon > 0, \exists M > 0$ such that

$$\limsup_{n \to \infty} \mathbb{P}^* \left\{ \sup_{f \in \mathcal{F}, A \in \mathcal{B}_n} \left| \mathbb{E}_n \left[f \mid A \right] - \mathbb{E} \left[f \mid A \right] \right| > M n^{\alpha - 1/2} \right\} \le \varepsilon,$$

which, together with Remark 7 again, proves the proposition.

Proof of Corollary 3.1

Proof of (3.1). Let $L = \operatorname{ess\,sup} Y$, $i \in \mathbb{N}$, and $f_i \in \mathcal{L}^1(\mathbb{Q})$ be defined by

$$f_i: \mathbb{R}^d \times [-L, L] \to [-L^i, L^i]$$

 $(\mathbf{x}, y) \mapsto y^i.$

 f_i is bounded and $\{f_i\}$ is finite thus Donsker. The result is then a straightforward application of Proposition 3.2.

Proof of (3.2). This part follows from Proposition 3.2 since Y is bounded and

$$\mathbb{V}_n\left[Y\mid (\mathbf{X},Y)\in A\right]:=\mathbb{E}_n\left[Y^2\mid (\mathbf{X},Y)\in A\right]-\mathbb{E}_n\left[Y\mid (\mathbf{X},Y)\in A\right]^2.$$

Covering algorithm selection process

Algorithm 1: Selection of minimal set of rules

```
Input:
               • the rate 0 < \gamma < 1;
               • a set of significant rules S;
               • a set of insignificant rules I;
       Output:
• a minimal set of rules C_n;
  1 C_n \leftarrow \arg\max_{\mathbf{r} \in S} \mathbb{Q}_n(\mathbf{r});
  2 S \leftarrow S \setminus C_n;
  з while \sum_{\mathbf{r} \in \mathcal{C}_n} \mathbb{Q}_n(\mathbf{r}) < 1 do
               \mathbf{r}^* \leftarrow \arg\max_{\mathbf{r} \in S} \mathbb{Q}_n(\mathbf{r});
               if \mathbb{Q}_n(\mathbf{r}^* \cap \{ \cup_{\mathbf{r} \in \mathcal{C}_n} \mathbf{r} \}) \leq \gamma \, \mathbb{Q}_n(\mathbf{r}^*) then
                 C_n \leftarrow C_n \cup \mathbf{r}^*;
  6
  7
  8
               S \leftarrow S \setminus \mathbf{r}^*;
               if \#S = 0 then
  9
10
                       Break;
11
12 end
13
      while \sum_{\mathbf{r}\in\mathcal{C}_n} \mathbb{Q}_n(\mathbf{r}) < 1 do
                     \leftarrow \operatorname{arg\,min}_{\mathbf{r}\in I} \mathbb{V}_n(Y|\mathbf{X}\in\mathbf{r});
14
               if \mathbb{Q}_n(\mathbf{r}^* \cap \{ \cup_{\mathbf{r} \in \mathcal{C}_n} \mathbf{r} \}) \leq \gamma \, \mathbb{Q}_n(\mathbf{r}^*) then
15
                      C_n \leftarrow C_n \cup \mathbf{r}^*;
16
17
                I \leftarrow I \setminus \mathbf{r}^*;
18
               if \#I = 0 then
19
                       Break;
20
21
22 end
23 return C_n;
```

References

BÉNARD, C., BIAU, G., DA VEIGA, S. and SCORNET, E. (2020). Interpretable Random Forests via Rule Extraction. arXiv preprint arXiv:2004.14841.

BIRAN, O. and COTTON, C. (2017). Explanation and justification in machine learning: A survey. In *IJCAI-17 workshop on explainable AI (XAI)* 8 1.

Breiman, L. (2001). Random forests. *Machine learning* **45** 5–32. MR3874153 Breiman, L., Friedman, J. H., Olshen, R. and Stone, C. (1984). *Classification and Regression Trees*. CRC press. MR0726392

CORTEZ, P. and SILVA, A. M. G. (2008). Using data mining to predict secondary school student performance. In *Proceedings of 5th FUture BUsiness TEChnology Conference (FUBUTEC 2008)*.

Denil, M., Matheson, D. and Freitas, N. (2013). Consistency of online random forests. In *International Conference on Machine Learning* 1256–1264. Dua, D. and Graff, C. (2017). UCI Machine Learning Repository.

- EFRON, B., HASTIE, T., JOHNSTONE, I., TIBSHIRANI, R. et al. (2004). Least angle regression. *The Annals of statistics* **32** 407–499. MR2060166
- FRIEDMAN, J. H. (2001). Greedy function approximation: a gradient boosting machine. *Annals of statistics* 1189–1232. MR1873328
- FRIEDMAN, J. H. (2002). Stochastic gradient boosting. Computational statistics & data analysis 38 367–378. MR1884869
- FRIEDMAN, J. H., POPESCU, B. E. et al. (2003). Importance sampled learning ensembles. *Journal of Machine Learning Research* **94305**.
- FRIEDMAN, J. H. and POPESCU, B. E. (2008). Predective Learning via Rule Ensembles. *The Annals of Applied Statistics* 916–954. MR2522175
- FÜRNKRANZ, J. and KLIEGR, T. (2015). A brief overview of rule learning. In *International Symposium on Rules and Rule Markup Languages for the Semantic Web* 54–69. Springer.
- Grunewalder, S. (2018). Plug-in Estimators for Conditional Expectations and Probabilities. In *International Conference on Artificial Intelligence and Statistics* 1513–1521.
- Guidotti, R., Monreale, A., Ruggieri, S., Turini, F., Giannotti, F. and Pedreschi, D. (2018). A survey of methods for explaining black box models. *ACM computing surveys (CSUR)* **51** 1–42.
- GYÖRFI, L., KOHLER, M., KRZYZAK, A. and WALK, H. (2006). A Distribution-Free Theory of Nonparametric Regression. Springer Science & Business Media. MR1920390
- Harrison Jr, D. and Rubinfeld, D. L. (1978). Hedonic housing prices and the demand for clean air. *Journal of environmental economics and management* 5 81–102.
- Hastie, T., Friedman, J. H. and Tibshirani, R. (2001). *The Elements of Statistical Learning* 1. Springer series in statistics Springer, Berlin. MR2722294
- Holmes, G., Hall, M. and Prank, E. (1999). Generating rule sets from model trees. In *Australasian Joint Conference on Artificial Intelligence* 1–12. Springer.
- Karalič, A. and Bratko, I. (1997). First order regression. *Machine Learning* **26** 147–176.
- LIITIÄINEN, E., VERLEYSEN, M., CORONA, F. and LENDASSE, A. (2009). Residual variance estimation in machine learning. *Neurocomputing* **72** 3692–3703.
- LIPTON, Z. C. (2018). The mythos of model interpretability. Queue 16 31–57. LUNDBERG, S. M. and LEE, S.-I. (2017). A unified approach to interpreting
- model predictions. In Advances in Neural Information Processing Systems 4765–4774.
- MARGOT, V., BAUDRY, J.-P., GUILLOUX, F. and WINTENBERGER, O. (2018). Rule Induction Partitioning Estimator. In *International Conference on Machine Learning and Data Mining in Pattern Recognition* 288–301. Springer.
- MEINSHAUSEN, N. (2010). Node harvest. The Annals of Applied Statistics 4 2049–2072. MR2829946
- NOBEL, A. (1996). Histogram Regression Estimation using Data-Dependent

- Partitions. The Annals of Statistics 24 1084–1105. MR1401839
- Quinlan, J. R. (1986). Induction of decision trees. *Machine learning* 1 81–106. MR3874153
- Quinlan, J. R. (1993). C4.5: Programs for Machine Learning.
- RAMOSAJ, B. and PAULY, M. (2019). Consistent estimation of residual variance with random forest Out-Of-Bag errors. *Statistics & Probability Letters* **151** 49–57. MR3937367
- RIBEIRO, M. T., SINGH, S. and GUESTRIN, C. (2016). Why should I trust you?: Explaining the predictions of any classifier. In *Proceedings of the 22nd ACM SIGKDD international conference on knowledge discovery and data mining* 1135–1144. ACM.
- Scornet, E., Biau, G., Vert, J.-P. et al. (2015). Consistency of random forests. *The Annals of Statistics* **43** 1716–1741. MR3357876
- Shrikumar, A., Greenside, P. and Kundaje, A. (2019). Learning important features through propagating activation differences. arXiv preprint arXiv:1704.02685.
- TIBSHIRANI, R. (1996). Regression shrinkage and selection via the lasso. Journal of the Royal Statistical Society. Series B (Methodological) 267–288. MR1379242
- Van der Vaart, A. W. (2000). Asymptotic statistics 3. Cambridge university press. MR1652247
- Vapnik, V. N. (1995). The Nature of Statistical Learning Theory. Springer-Verlag. MR1367965
- Wenocur, R. S. and Dudley, R. M. (1981). Some special Vapnik-Chervonenkis classes. *Discrete Mathematics* **33** 313–318. MR0602047
- Zhao, Q. and Bhowmick, S. S. (2003). Association rule mining: A survey. Nanyang Technological University, Singapore.