# Universität Stuttgart 

## Fachbereich Mathematik

# Improved Classification Rates under Refined Margin 

 ConditionsIngrid Blaschzyk, Ingo Steinwart

Fachbereich Mathematik
Fakultät Mathematik und Physik
Universität Stuttgart
Pfaffenwaldring 57
D-70 569 Stuttgart

E-Mail: preprints@mathematik.uni-stuttgart.de
WWW: http://www.mathematik.uni-stuttgart.de/preprints
ISSN 1613-8309
(C) Alle Rechte vorbehalten. Nachdruck nur mit Genehmigung des Autors.

LATEX-Style: Winfried Geis, Thomas Merkle

# Improved Classification Rates under Refined Margin Conditions 

Ingrid Blaschzyk and Ingo Steinwart<br>Institute for Stochastics and Applications<br>University of Stuttgart Pfaffenwaldring 57 D-70569 Stuttgart<br>e-mail: ingrid.blaschzyk@mathematik.uni-stuttgart.de<br>e-mail: ingo.steinwart@mathematik.uni-stuttgart.de


#### Abstract

In this paper we present a simple partitioning based technique to refine the statistical analysis of classification algorithms. The core of our idea is to divide the input space into two parts such that the first part contains a suitable vicinity around the decision boundary, while the second part is sufficiently far away from the decision boundary. Using a set of margin conditions we are then able to control the classification error on both parts separately. By balancing out these two error terms we obtain a refined error analysis in a final step. We apply this general idea to the histogram rule and show that even for this simple method we obtain, under certain assumptions, better rates than the ones known for support vector machines, for certain plug-in classifiers, and for a recently analysed tree based adaptive-partitioning ansatz.


## Contents

1 Introduction . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . 1
2 General assumptions . . . . . . . . . . . . . . . . . . . . . . . . . . . . 2
3 Oracle inequality and learning rates . . . . . . . . . . . . . . . . . . . 5
4 Proofs . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . 11
A Appendix . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . 22
References . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . 23

## 1. Introduction

Given a dataset $D:=\left(\left(x_{i}, y_{i}\right), \ldots,\left(x_{n}, y_{n}\right)\right)$ of observations drawn in an i.i.d. fashion from a probability measure $P$ on $X \times Y$, where $X \subset \mathbb{R}^{d}$ and $Y:=$ $\{-1,1\}$, the learning goal of binary classification is to find a decision function $f_{D}: X \rightarrow\{-1,1\}$ such that for new data $(x, y)$ we have $f_{D}(x)=y$ with high probability.

The problem of classification is, apart from regression, one of the most considered problems in learning theory and many classical learning methods have been presented in the literature such as histogram rules, nearest neighbor methods or moving window rules. A general reference for these methods is [4]. Several
more recent methods use trees to build a classifier, for example the random forest algorithm, introduced in [3], makes a prediction by a majority vote over a collection of random forest trees. Another example is the tree based adaptivepartitioning algorithm, presented in [2]. Here, a classifier is picked by empirical risk minimization over a nested sequence $\left(S_{m}\right)_{m \geq 1}$ of families of sets which consists of dyadic or decorated trees. An example of a non-tree based algorithm is described in [1]. Here, the final classifier is found by empirical risk minimization over a suitable grid of plug-in rules. Another non-tree based algorithm is, for example, the support vector machine (SVM), which solves a regularized empirical risk minimization problem over a reproducing kernel Hilbert space $H$. For more details on statistical properties of SVM for classification we refer the reader to [7, Chapter 8].

In this paper we discuss a partitioning based technique to analyse the statistical properties of classification algorithms. In particular we show for the histogram rule that under certain assumptions this technique leads to rates, which are faster than the rates obtained in [2], [3], and [7]. To be more precise, we divide the input space $X$ into two overlapping regions that are adjustable by a parameter $r$ in such a way that one set, which we will denote by $A_{r}$, contains points near the decision boundary, whereas the other set $B_{r}$ contains those that are sufficiently bounded far away from the decision boundary. We examine the excess risks over these two sets separately by using an oracle inequality for empirical risk minimizers on both parts. It turns out that under a suitable assumption, which describes the location of critical noise, we have no approximation error as well as an optimal variance bound on $B_{r}$, which in turn leads to an $\mathcal{O}\left(n^{-1}\right)$ behaviour of the excess risk on $B_{r}$. However, this bound still depends on the parameter $r$, namely it increases for $r \rightarrow 0$. In contrast our bound on the risk on $A_{r}$ decreases for $r \rightarrow 0$. By balancing out these two risks with respect to $r$ we obtain a refined bound on $X$ under additional assumptions describing the concentration of mass around the decision boundary.

A more detailed discussion on this technique and the statistical result are presented in Section 3. Moreover a comparison of the resulting learning rates to the known ones for the SVM and the tree based adaptive-partitioning algorithm described in [2] can be find at the end of Section 3. We note that all proofs are deferred to Section 4.

## 2. General assumptions

To describe our learning goal we consider in the following the classification loss $L:=L_{\text {class }}(y, t): Y \times \mathbb{R} \rightarrow[0, \infty)$, defined by $L(y, t):=\mathbf{1}_{(-\infty, 0]}(y \cdot \operatorname{sign} t)$ for $y \in Y, t \in \mathbb{R}$, where $\mathbf{1}_{(-\infty, 0]}$ denotes the indicator function on $(-\infty, 0]$. We define the risk of a measurable estimator $f: X \rightarrow \mathbb{R}$ by

$$
\mathcal{R}_{L, P}(f):=\int_{X \times Y} L(y, f(x)) d P(x, y)
$$

and the empirical risk by

$$
\mathcal{R}_{L, D}(f):=\frac{1}{n} \sum_{i=1}^{n} L\left(y_{i}, f\left(x_{i}\right)\right),
$$

where $D:=\frac{1}{n} \sum_{i=1}^{n} \delta_{\left(x_{i}, y_{i}\right)}$ denotes the Dirac measure in $\left(x_{i}, y_{i}\right)$. The smallest possible risk

$$
\mathcal{R}_{L, P}^{*}:=\inf _{f: X \rightarrow \mathbb{R}} \mathcal{R}_{L, P}(f)
$$

is called the Bayes risk, and a measurable function $f_{L, P}^{*}: X \rightarrow \mathbb{R}$ so that $\mathcal{R}_{L, P}\left(f_{L, P}^{*}\right)=\mathcal{R}_{L, P}^{*}$ holds is called Bayes decision function. Recall that the Bayes decision function $f_{L, P}^{*}$ for the classification loss is given by $\operatorname{sign}(2 P(y=$ $1 \mid x)-1$ ) for $x \in X$, where $P(\cdot \mid x)$ is the conditional probability on $Y$ given $x$. Let us now briefly describe a particular histogram rule. To this end, let $\mathcal{A}=$ $\left(A_{j}\right)_{j \geq 1}$ be a partition of $\mathbb{R}^{d}$ into cubes of side length $s \in(0,1]$ and $X:=[-1,1]^{d}$. For $x \in X$ we denote by $A(x)$ the unique cell of $\mathcal{A}$ with $x \in A(x)$ and call the map $h_{P, s}: X \rightarrow Y$ defined by

$$
h_{P, s}(x):= \begin{cases}-1 & \text { if } f_{P, s}(x)<0,  \tag{1}\\ 1 & \text { if } f_{P, s}(x) \geq 0\end{cases}
$$

where $f_{P, s}(x):=P(A(x) \times\{1\})-P(A(x) \times\{-1\})$, infinite sample histogram rule. For a dataset $D$ we further write

$$
f_{D, s}:=\frac{1}{n} \sum_{i=1}^{n} \mathbf{1}_{\left\{y_{i}=+1\right\}} \mathbf{1}_{A(x)}\left(x_{i}\right)-\frac{1}{n} \sum_{i=1}^{n} \mathbf{1}_{\left\{y_{i}=-1\right\}} \mathbf{1}_{A(x)}\left(x_{i}\right) .
$$

Thus, the empirical histogram is defined by $h_{D, s}:=\operatorname{sign} f_{D, s}$. We define the set $\mathcal{F}$ by

$$
\mathcal{F}:=\left\{\sum_{A_{j} \cap[-1,1]^{d} \neq \emptyset} c_{j} \mathbf{1}_{A_{j}}: c_{j} \in\{-1,1\}\right\} .
$$

Then, it is easy to show that the empirical histogram rule $h_{D, s}$ is an empirical risk minimizer over $\mathcal{F}$ for the classification loss, that means

$$
\mathcal{R}_{L, D}\left(h_{D, s}\right)=\inf _{f \in \mathcal{F}} \mathcal{R}_{L, D}(f) .
$$

Since we aim in a further step to examine the risk on subsets of $X$ consisting of cells, we have to specify the loss on those subsets. Therefore, we define for an arbitrary index set $J \subset\{1, \ldots, m\}$ the set

$$
\begin{equation*}
T_{J}:=\bigcup_{j \in J} A_{j} \tag{2}
\end{equation*}
$$

and the related loss $L_{T_{J}}: X \times Y \times \mathbb{R} \rightarrow[0, \infty)$ by

$$
\begin{equation*}
L_{T_{J}}(x, y, t):=\mathbf{1}_{\bigcup_{j \in J} A_{j}}(x) L_{\text {class }}(y, t) \tag{3}
\end{equation*}
$$

Furthermore, we define the risk over $T$ by

$$
\mathcal{R}_{L_{T_{J}}, P}(f):=\int_{X \times Y} L_{T_{J}}(x, y, f(x)) d P(x, y)
$$

As mentioned in the introduction, we have to make assumptions on $P$ to obtain rates. Therefore we recall some notions from [7, Chapter 8], which describe the behaviour of $P$ in the vicinity of the decision boundary. To this end, let $\eta: X \rightarrow[0,1]$, defined by $\eta(x):=P(y=1 \mid x), x \in X$ be a version of the posterior probability of $P$, that means that the probability measures $P(\cdot \mid x)$ form a regular conditional probability of P. We write

$$
\begin{aligned}
X_{1} & :=\{x \in X: \eta(x)>1 / 2\} \\
X_{-1} & :=\{x \in X: \eta(x)<1 / 2\} .
\end{aligned}
$$

Then, the function $\Delta_{\eta}: X \rightarrow[0, \infty]$ defined by

$$
\Delta_{\eta}(x):= \begin{cases}d\left(x, X_{1}\right) & \text { if } x \in X_{-1}  \tag{4}\\ d\left(x, X_{-1}\right) & \text { if } x \in X_{1} \\ 0 & \text { otherwise }\end{cases}
$$

where $d(x, A):=\inf _{x^{\prime} \in A} d\left(x, x^{\prime}\right)$, is called distance to the decision boundary. This helps us to describe the mass of the marginal distribution $P_{X}$ of $P$ around the decision boundary by the following exponents. We say that $P$ has strong margin exponent (SME) $\alpha \in(0, \infty]$ if there exists a constant $c_{\text {SME }}>0$ such that

$$
P_{X}\left(\left\{\Delta_{\eta}(x)<t\right\}\right) \leq\left(c_{\mathrm{SME}} t\right)^{\alpha}
$$

for all $t>0$. Descriptively, the strong margin exponent $\alpha$ measures the amount of mass close to the decision boundary. Therefore, large values of $\alpha$ are better since they reflect a low concentration of mass in this region, which makes the classification easier. Furthermore, we say that $P$ has margin-noise exponent (MNE) $\beta \in(0, \infty]$ if there exists a constant $c_{\text {MNE }}>0$ such that

$$
\int_{\left\{\Delta_{\eta}(x)<t\right\}}|2 \eta(x)-1| d P_{X}(x) \leq\left(c_{\mathrm{MNE}} t\right)^{\beta}
$$

for all $t>0$. The margin-noise exponent $\beta$ measures the mass and the noise, that means the amount of points $x \in X$ with $\eta(x) \approx 1 / 2$, around the decision boundary. That is, we have high margin-noise exponent if we have low mass and/or high noise around the decision boundary. Next, we say that the distance to the decision boundary $\Delta_{\eta}$ controls the noise from below by the exponent $\gamma$ if there exist a $\gamma \in[0, \infty)$ and a constant $c_{\mathrm{LC}}>0$ with

$$
\begin{equation*}
\Delta_{\eta}^{\gamma}(x) \leq c_{\mathrm{LC}}|2 \eta(x)-1| \tag{5}
\end{equation*}
$$

for $P_{X}$-almost all $x \in X$. That means, if $\eta(x)$ is close to $1 / 2$ for some $x \in X$, this $x$ is close to the decision boundary. For examples of typical values of these exponents and relations between them we refer the reader to [7, Chapter 8].

Finally, in order to describe the region of the decision boundary in a more geometrical way, we say according to [5, 3.2.14(1)] that a general set $T \subset X$ is $m$-rectifiable for an integer $m>0$ if there exists a Lipschitzian function mapping some bounded subset of $\mathbb{R}^{m}$ onto $T$. Moreover, we denote by $\mathcal{H}^{d-1}$ the ( $d-1$ )-dimensional Hausdorff measure on $\mathbb{R}^{d}$.

The following lemma, which is based on [6, Lemma A.10.4], describes the Lebesgue measure of the decision boundary in terms of the Hausdorff measure. Its result will be necessary for the analysis of the main theorem in Section 3.
Lemma 2.1. Let $X:=[-1,1]^{d}$ and $P$ be a probability measure on $X \times\{-1,1\}$ with fixed version $\eta: X \rightarrow[0,1]$ of its posterior probability. Moreover let $\lambda^{d}$ be the $d$-dimensional Lebesgue measure and $\mathcal{H}^{d-1}$ be the $(d-1)$-dimensional Hausdorff measure on $\mathbb{R}^{d}$. Furthermore, let $X_{0}$ equal the relative boundary of $X_{1}$ in $X$, that means $X_{0}=\delta_{X} X_{1}$, with $\mathcal{H}^{d-1}\left(X_{0}\right)>0$ and let $X_{0}$ be $(d-1)$-rectifiable. Then, there exists a $\delta^{*}>0$ such that for all $\delta \in\left(0, \delta^{*}\right]$ we have

$$
\lambda^{d}(\{x \in X \mid \Delta(x) \leq \delta\}) \leq 4 \mathcal{H}^{d-1}(\{x \in X \mid \eta(x)=1 / 2\}) \cdot \delta
$$

## 3. Oracle inequality and learning rates

Our goal is to find an upper bound for the excess risk $\mathcal{R}_{L, P}\left(h_{D, s}\right)-\mathcal{R}_{L, P}^{*}$. The idea is to split $X$ into two overlapping sets and to find a bound on the risks over these sets by using information on $P$. To this end, we denote the set of indices of cubes that intersect $X$ by

$$
J:=\left\{j \geq 1 \mid A_{j} \cap[-1,1]^{d} \neq \emptyset\right\} .
$$

Next, we split this set into cubes that lie near the decision boundary and into cubes that are bounded away from the decision boundary. To be more precisely, we define, for $r>0$ and a version $\eta$ for which the assumptions at the end of Section 2 hold, the set of indices of cubes near the decision boundary by

$$
J_{A}^{r}:=\left\{j \in J \mid \forall x \in A_{j}: \Delta_{\eta}(x) \leq 3 r\right\}
$$

and the set of indices of cubes that are sufficiently bounded away by

$$
J_{B}^{r}:=\left\{j \in J \mid \forall x \in A_{j}: \Delta_{\eta}(x) \geq r\right\} .
$$

Moreover, we write

$$
\begin{align*}
A_{r} & :=\bigcup_{j \in J_{A}^{r}} A_{j},  \tag{6}\\
B_{r} & :=\bigcup_{j \in J_{B}^{r}} A_{j} . \tag{7}
\end{align*}
$$

As the following lemma shows, we need to define requirements on the side length of the cells to ensure that $X \subset A_{r} \cup B_{r}$. Besides that, it shows that we are able to assign all $x \in A_{j}$, where $j \in J_{B}^{r}$, either to the class $X_{-1}$ or to $X_{1}$.

Lemma 3.1. Let $\mathcal{A}=\left(A_{j}\right)_{j \geq 1}$ be a partition of $\mathbb{R}^{d}$ into cubes of side length $s \in(0,1]$ and let $X:=[-1,1]^{d}$. For $r \geq s / 2$ define the sets $A_{r}$ and $B_{r}$ by (6) and (7). Then,
i) we have $X \subset A_{r} \cup B_{r}$,
ii) we have either $A_{j} \cap X_{1}=\emptyset$ or $A_{j} \cap X_{-1}=\emptyset$ for $j \in J_{B}^{r}$.

Since the excess risk is non-negative, we obtain under the assumption of Lemma 3.1(i) that

$$
\begin{align*}
& \mathcal{R}_{L, P}\left(h_{D, s}\right)-\mathcal{R}_{L, P}^{*} \\
& \leq\left(\mathcal{R}_{L_{A_{r}}, P}\left(h_{D, s}\right)-\mathcal{R}_{L_{A_{r}}, P}^{*}\right)+\left(\mathcal{R}_{L_{B_{r}}, P}\left(h_{D, s}\right)-\mathcal{R}_{L_{B_{r}}, P}^{*}\right) \tag{8}
\end{align*}
$$

That means, we can bound the excess risk $\mathcal{R}_{L, P}\left(h_{D, s}\right)-\mathcal{R}_{L, P}^{*}$ if we find bounds on the excess risks over the sets $A_{r}$ and $B_{r}$. For that purpose, we use an oracle inequality for empirical risk minimizer separately on both error terms, see [7, Theorem 7.2]. This is possible, since the following lemma shows that, considering the loss $L_{T_{J}}$ for any set $T_{J}$ constructed as in (2), the empirical histogram rule $h_{D, s}$ is still an empirical risk minimizer over $\mathcal{F}$.

Lemma 3.2. Consider for an arbitrary index set $J \subset\{1, \ldots, m\}$ the set $T_{J}:=$ $\bigcup_{j \in J} A_{j}$ and the related loss $L_{T_{J}}: X \times Y \times \mathbb{R} \rightarrow[0, \infty)$ defined in (3). Then, the empirical histogram rule $h_{D, s}$ is an empirical risk minimizer over $\mathcal{F}$ for the loss $L_{T_{J}}$, that means

$$
\mathcal{R}_{L_{T_{J}, D}}\left(h_{D, s}\right)=\inf _{f \in \mathcal{F}} \mathcal{R}_{L_{T_{J}, D}}(f)
$$

Before we state our oracle inequality, we discuss in a more detailed way the improvement that we gained by our separation technique described above. First, we make no approximation error on the set $B_{r}$, which consists of cells that are sufficiently bounded away from the decision boundary. This follows from the circumstance that $h_{D, s}$ learns correctly on those cells and follows even intuitively inasmuch as the noise concentration is rather low in this region. We refer the reader to Part 1 of the proof of Lemma 3.4 for details. Second, the main refinement arises from the fact that we achieve, under the condition that the decision boundary controls the noise from above, in the variance bound on $B_{r}$, a bound of the form

$$
\mathbb{E}_{P}\left(L \circ f-L \circ f_{L, P}^{*}\right)^{2} \leq V \cdot \mathbb{E}_{P}\left(L \circ f-L \circ f_{L, P}^{*}\right)^{\theta}
$$

where $V>0$, the best possible exponent, $\theta=1$. This bound plays an important part in the analysis of the risk terms, since we have small variance if the righthand side of the latter bound is small, as the next lemma shows.

Lemma 3.3. Let $X:=[-1,1]^{d}$ and $P$ be a probability measure on $X \times\{-1,1\}$ with fixed version $\eta: X \rightarrow[0,1]$ of its posterior probability. Assume that the associated distance to the decision boundary $\Delta_{\eta}$ controls the noise from below by the exponent $\gamma \in[0, \infty)$ and consider for some fixed $r>0$ the set $B_{r}$, defined in (7). Furthermore let $L:=L_{\text {class }}$ be the classification loss and let $f_{L, P}^{*}$ be a fixed Bayes decision function. Then, for all measurable $f: X \rightarrow\{-1,1\}$ we have

$$
\mathbb{E}_{P}\left(L_{B_{r}} \circ f-L_{B_{r}} \circ f_{L, P}^{*}\right)^{2} \leq \frac{c_{L C}}{r^{\gamma}} \mathbb{E}_{P}\left(L_{B_{r}} \circ f-L_{B_{r}} \circ f_{L, P}^{*}\right)
$$

We remark that the right-hand side of the variance bound on $B_{r}$ depends on the separation parameter $r$. This dependence is also reflected in the risk term on $B_{r}$. In particular, we show in the proof of our main theorem by applying [7, Theorem 7.2] on the risk term on the set $B_{r}$ that the improvements mentioned above lead to

$$
\mathcal{R}_{L_{B_{r}}, P}\left(h_{D, s}\right)-\mathcal{R}_{L_{B_{r}}, P}^{*} \leq \frac{32 c_{1}\left(8^{d+1} s^{-d}+\tau\right)}{r^{\gamma} n}
$$

with probability $P^{n} \geq 1-e^{-\tau}$, where $\tau \geq 1$ and $c_{1}$ is a positive constant. Whereas this error term increases for $r \rightarrow 0$, the error term on the set $A_{r}$ behaves exactly the opposite way, that is, it decreases for $r \rightarrow 0$. In fact, bounding the risk on $A_{r}$ requires additional knowledge of the behaviour of $P$ in the vicinity of the decision boundary. By applying [7, Theorem 7.2] on the risk on the set $A_{r}$ we show under the assumption that $P$ has strong margin exponent $\alpha$ and margin-noise exponent $\beta$ that

$$
\mathcal{R}_{L_{A}, P}\left(h_{D}\right)-\mathcal{R}_{L_{A}, P}^{*} \leq 6 c_{\mathrm{MNE}} s^{\beta}+4\left(\frac{8 V\left(c_{5} r s^{-d}+\tau\right)}{n}\right)^{\frac{\alpha+\gamma}{\alpha+2 \gamma}}
$$

holds with probability $P^{n} \geq 1-e^{-\tau}$. Here, $c_{5}$ is a positive constant, $\tau \geq 1$ and $V$ is the prefactor of the variance bound on $A_{r}$, shown in Part 2 of the proof of Lemma 3.4. If we balance the obtained risk terms over $A_{r}$ and $B_{r}$ with respect to $r$, we obtain the oracle inequality presented in our main theorem. For this purpose, we define the positive constant

$$
\tilde{c}_{\alpha, \gamma, d}:=\left(\frac{16 \gamma(\alpha+2 \gamma) \cdot 8^{d+1} \max \left\{c_{\mathrm{LC}}, 2^{\gamma}\right\} \cdot(\alpha+\gamma)^{-1}}{\hat{c}^{\frac{\alpha+\gamma}{\alpha+2 \gamma}}}\right)^{\frac{\alpha+\gamma}{\alpha+\gamma+\gamma(\alpha+2 \gamma)}}
$$

which depends on $\alpha, \gamma$ and $d$ and where $\hat{c}:=24 \max \left\{12 \mathcal{H}^{d-1}(\{\eta=1 / 2\}), 1\right\}$. $\max \left\{1, \frac{\alpha+\gamma}{\gamma} c_{\text {SME }}^{\frac{\alpha \gamma}{\alpha+\gamma}}\left(\frac{\gamma c_{\mathrm{LC}}}{\alpha}\right)^{\frac{\alpha}{\alpha+\gamma}}\right\}$.
Theorem 3.4. Let $\mathcal{A}=\left(A_{j}\right)_{j \geq 1}$ be a partition of $\mathbb{R}^{d}$ into cubes of side length $s \in(0,1]$. Let $X:=[-1,1]^{d}$ and $P$ be a probability measure on $X \times\{-1,1\}$ with fixed version $\eta: X \rightarrow[0,1]$ of its posterior probability. Assume that the associated distance to the decision boundary $\Delta_{\eta}$ controls the noise from below by the exponent $\gamma \in[0, \infty)$ and assume as well that $P$ has $M N E \beta \in[0, \infty)$ and

SME $\alpha \in(0, \infty]$. Furthermore, let $X_{0}$ equal the relative boundary of $X_{1}$ in $X$, that means $X_{0}=\delta_{X} X_{1}$, with $\mathcal{H}^{d-1}\left(X_{0}\right)>0$ and let $X_{0}$ be $(d-1)$-rectifiable. Let $L$ be the classification loss and let for fixed $n \geq 1$ and $\tau \geq 1$ the bounds

$$
\begin{equation*}
s \leq \tilde{c}_{\alpha, \gamma, d}^{\frac{(1+\gamma)(\alpha+\gamma)+\gamma^{2}}{(1+\gamma)(\alpha+\gamma)+\gamma^{2}+d \gamma}}\left(\frac{\tau}{n}\right)^{\frac{\gamma}{(1+\gamma)(\alpha+\gamma)+\gamma^{2}+d \gamma}} \tag{9}
\end{equation*}
$$

and

$$
\begin{equation*}
s^{d} n \geq \tau\left(\frac{\tilde{c}_{\alpha, \gamma, d}}{\min \left\{\frac{\delta^{*}}{3}, 1\right\}}\right)^{\frac{(1+\gamma)(\alpha+\gamma)+\gamma^{2}}{\gamma}} \tag{10}
\end{equation*}
$$

be satisfied, where the constant $\tilde{c}_{\alpha, \gamma, d}>0$ depends on $\alpha, \gamma, d$ and the constant $\delta^{*}>0$ is the one of Lemma 2.1. Then, there exists a constant $c_{\alpha, \gamma, d}>0$ such that

$$
\begin{equation*}
\mathcal{R}_{L, P}\left(h_{D, s}\right)-\mathcal{R}_{L, P}^{*} \leq 6\left(c_{M N E} s\right)^{\beta}+c_{\alpha, \gamma, d}\left(\frac{\tau}{s^{d} n}\right)^{\frac{(1+\gamma)(\alpha+\gamma)}{(1+\gamma)(\alpha+\gamma)+\gamma^{2}}} \tag{11}
\end{equation*}
$$

holds with probability $P^{n} \geq 1-2 e^{-\tau}$, where the constant $c_{\alpha, \gamma, d}$ only depends on $\alpha, \gamma$ and $d$.

The proof shows that the constants $c_{\alpha, \gamma, d}$ is given by

$$
\begin{equation*}
c_{\alpha, \gamma, d}:=128 \cdot 8^{d+1} \max \left\{c_{\mathrm{LC}}, 2^{\gamma}\right\} \cdot \max \left\{\frac{\gamma(\alpha+2 \gamma)}{\alpha+\gamma}, 1\right\} \cdot \tilde{c}_{\alpha, \gamma, d}^{-\gamma} \tag{12}
\end{equation*}
$$

Note that the assumptions (9) and (10) on the side length $s$ of the cubes are natural assumptions, since $s$ has to be small enough given a specific number of observations, but yet should not shrink too fast for grown observations. By choosing an appropriate sequence of $s_{n}$ in dependence of our data length $n$ and setting a constraint on the margin-noise exponent $\beta$ we state learning rates in the next theorem. Prior to that, we define the positive constant

$$
\begin{aligned}
& \tilde{c}_{\alpha, \beta, \gamma, \tau, d} \\
& :=\left(\frac{d(1+\gamma)(\alpha+\gamma) \cdot c_{\alpha, \gamma, d} \cdot \tau^{\frac{(1+\gamma)(\alpha+\gamma)}{(1+\gamma)(\alpha+\gamma)+\gamma^{2}}}}{6 \beta c_{\mathrm{MNE}}^{\beta}\left((1+\gamma)(\alpha+\gamma)+\gamma^{2}\right)}\right)^{\frac{(1+\gamma)(\alpha+\gamma)+\gamma^{2}}{\beta\left((1+\gamma)(\alpha+\gamma)+\gamma^{2}\right)+d(1+\gamma)(\alpha+\gamma)}}
\end{aligned}
$$

that depends on $\alpha, \beta, \gamma, \tau$ and $d$ and where $c_{\alpha, \gamma, d}$ is the constant from (12).
Theorem 3.5. Assume that $X$ and $P$ satisfy the assumptions of Theorem 3.4 for $\beta \leq \gamma^{-1} \kappa$, where $\kappa:=(1+\gamma)(\alpha+\gamma)$. In addition assume that the side length $s_{n}$ in Theorem 3.4 is given by

$$
s_{n}=\tilde{c}_{\alpha, \beta, \gamma, \tau, d} n^{-\frac{\kappa}{\beta\left(\kappa+\gamma^{2}\right)+d \kappa}} .
$$

Then, there exists a constant $c_{\alpha, \beta, \gamma, \tau, d}>0$ such that for all $n \geq n_{0}$

$$
\mathcal{R}_{L, P}\left(h_{D, s_{n}}\right)-\mathcal{R}_{L, P}^{*} \leq c_{\alpha, \beta, \gamma, \tau, d} n^{-\frac{\beta \kappa}{\beta\left(\kappa+\gamma^{2}\right)+d \kappa}}
$$

holds with probability $P^{n} \geq 1-2 e^{-\tau}$, where $n_{0}$ and the constant $c_{\alpha, \beta, \gamma, \tau, d}$ only depend on $\tau, \alpha, \beta, \gamma$ and $d$.

The proof shows that the constant $c_{\alpha, \beta, \gamma, \tau, d}$ is given by

$$
c_{\alpha, \beta, \gamma, \tau, d}:=2 \max \left\{\frac{d \cdot \kappa}{\beta\left(\kappa+\gamma^{2}\right)}, 1\right\} c_{\alpha, \gamma, \delta} \cdot \tau^{\frac{\kappa}{\kappa+\gamma^{2}}} \cdot \tilde{c}_{\alpha, \beta, \gamma, \tau, d}^{-\frac{d \kappa}{\kappa+\gamma^{2}}} .
$$

To obtain the rates we have to know the $P$ describing parameters. However, it is possible to yield the rates in Theorem 3.5 by a training validation ansatz, that is by splitting the dataset into two parts and considering a suitable set of candidates $s_{n}$. In order to compare our rate obtained in Theorem 3.5, we now consider, besides our geometric assumption on $X$, namely
(i) $X_{0}$ is $(d-1)$-rectifiable with $\mathcal{H}^{d-1}\left(X_{0}\right)>0$ and $X_{0}$ equals the relative boundary of $X_{1}$ in $X$,
the following two assumptions on $P$ :
(ii) $P$ has SME $\alpha \in(0, \infty]$,
(iii) there exists a $\gamma \in[0, \infty)$ and constants $c_{1}, c_{2}, c_{\mathrm{UC}}>0$ such that
a) $c_{1}|2 \eta(x)-1| \geq c_{\mathrm{LC}} \Delta_{\eta}^{\gamma}(x)$,
b) $c_{2}|2 \eta(x)-1| \leq c_{\mathrm{UC}} \Delta_{\eta}^{\gamma}(x)$.

Here, assumption $(i i i)_{a)}$ coincides up to the constant $c_{1}$ with the definition in (5). Furthermore, assumption $(i i i)_{b}$ indicates that we have an upper control by $\Delta_{\eta}$ on the noise, which is a kind of inverse to $(i i i)_{a}$. Then, [7, Lemma 8.17] shows under the assumptions (ii) and (iii) that P has MNE $\beta=\alpha+\gamma$. Hence, we find in Theorem 3.5 with $\kappa:=(1+\gamma)(\alpha+\gamma)$ and a suitable cell-width that $h_{D, s_{n}}$ learns with a rate with exponent

$$
\begin{aligned}
\frac{\beta(1+\gamma)(\alpha+\gamma)}{\beta\left[(1+\gamma)(\alpha+\gamma)+\gamma^{2}\right]+d(1+\gamma)(\alpha+\gamma)} & =\frac{(\alpha+\gamma)(1+\gamma)(\alpha+\gamma)}{(\alpha+\gamma)\left[(1+\gamma)(\alpha+\gamma)+\gamma^{2}\right]+d(1+\gamma)(\alpha+\gamma)} \\
& =\frac{(1+\gamma)(\alpha+\gamma)}{(1+\gamma)(\alpha+\gamma)+\gamma^{2}+d(1+\gamma)}
\end{aligned}
$$

A simple transformation shows that this exponent equals

$$
\begin{equation*}
\frac{(1+\gamma)(\alpha+\gamma)}{(1+\gamma)(\alpha+\gamma)+\gamma^{2}+d(1+\gamma)}=\frac{\alpha+\gamma}{\alpha+\gamma+\frac{\gamma^{2}}{1+\gamma}+d}=\frac{\alpha+\gamma}{\alpha+2 \gamma+\frac{\gamma^{2}}{1+\gamma}+d-\gamma}=\frac{\alpha+\gamma}{\alpha+2 \gamma+d-\frac{\gamma}{1+\gamma}} \tag{13}
\end{equation*}
$$

First, we compare the rate with exponent (13) with the rate achieved by support vector machines (SVM) for the hinge loss by considering the assumptions (i), (ii) and (iii). For this purpose, [7, Chapter 8.3(8.18)] shows that the best possible rate for the SVM is obtained by

$$
n^{-\frac{\alpha+\gamma}{\alpha+2 \gamma+d}+\rho}
$$

where $\rho>0$ is an arbitrary small number. Hence, our rate in (13) is better by $-\frac{\gamma}{1+\gamma}$ in the denominator. For the typical value of $\gamma=1$, indicating a moderate control of noise by the decision boundary, our rate is better by $-1 / 2$ in the denominator. Finally, we remark that for both results less assumptions are sufficient.

Second, we compare our rate with the ones for certain plug-in classifiers, see [1], and with the rates of a classification scheme, described in [2]. In both cases, the authors use the so called margin assumption, which is comparable to the Definition of the noise exponent in [7]. Hence, we find under assumptions (ii) and (iii) with [7, Exercise 8.5] that the margin condition is fulfilled for $\alpha / \gamma$. In addition to ( $i$ ) and (iii) we impose the following two conditions:
(iv) $\eta$ is Hoelder-continuous with exponent $\gamma$,
(v) $P_{X}$ is the Lebesgue measure.

Under condition $(i)$ and $(v)$ we find with Lemma 2.1 that assumption (ii) is fulfilled for $\alpha=1$. Furthermore, we find under condition (iv) with Lemma A. 1 that assumption (iii) is fulfilled with exponent $\gamma$. Hence, the conditions ( $i$ ) and $(i i i)-(v)$ yield in (13) a rate with exponent

$$
\frac{1+\gamma}{1+2 \gamma+d-\frac{\gamma}{1+\gamma}}
$$

Furthermore, [1, Theorem 4.3] shows that certain plug-in classifier yield under the same conditions $(i)$ and $(i i i)-(v)$ the rate

$$
\begin{equation*}
n^{-\frac{1+\gamma}{1+2 \gamma+d}} \tag{14}
\end{equation*}
$$

and we find that our rate is better by $-\frac{\gamma}{1+\gamma}$ in the denominator. Under the same assumptions, [2, Corollary 5.2 (ii)] shows that the described classification scheme obtains the rate

$$
\left(\frac{(\log n)^{\frac{1}{2+d}}}{n}\right)^{-\frac{1+\gamma}{2 \gamma+d}}
$$

Hence, our rate is worse by $\frac{1}{1+\gamma}$. However, the results from [2] are also comparable under another set of assumptions. Indeed, if we assume that the conditions $(i)-(i v)$ hold, then, our rate given in (13) holds and [2, Corollary 5.2(i)] shows that the described classification scheme yields the rate

$$
\left(\frac{\log n}{n}\right)^{-\frac{\alpha+\gamma}{\alpha+2 \gamma+d}}
$$

and our rate is again better by $-\frac{\gamma}{1+\gamma}$ in the denominator. Note that for $\alpha=$ 1 this rate equals (14) up to the logarithm. Finally, we remark that for our results as well as for the results from [1] and [2] less assumptions are sufficient and in the comparisons above we tried to formulate reasonable sets of common assumptions.

## 4. Proofs

Proof of Lemma 2.1: For a set $T \subset X$ and $\delta>0$ we define as in [6] the sets

$$
\begin{aligned}
& T^{+\delta}:=\{x \in X \mid d(x, T) \leq \delta\} \\
& T^{-\delta}:=X \backslash(X \backslash T)^{+\delta}
\end{aligned}
$$

Since $X_{1}:=\{x \in X \mid \eta(x) \leq 1 / 2\}$ is bounded and measurable, we find with [6, Lemma A.10.3] and the proof of [6, Lemma+A.10.4(ii)] that there exists a $\delta^{*}>0$, such that for all $\delta \in\left(0, \delta^{*}\right]$ we have

$$
\begin{equation*}
\lambda^{d}\left(X_{1}^{+\delta} \backslash X_{1}^{-\delta}\right) \leq 4 \mathcal{H}^{d-1}\left(\partial X_{1}\right) \cdot \delta=4 \mathcal{H}^{d-1}(\{x \in X \mid \eta(x)=1 / 2\}) \cdot \delta \tag{15}
\end{equation*}
$$

Next, we show that

$$
\begin{equation*}
\{x \in X \mid \Delta(x) \leq \delta\} \subset X_{1}^{+\delta} \backslash X_{1}^{-\delta} \cup X_{0} . \tag{16}
\end{equation*}
$$

For this purpose, we remark that according to (4) we have

$$
\begin{aligned}
& \{x \in X \mid \Delta(x) \leq \delta\} \\
& =\left\{x \in X_{1} \mid d\left(x, X_{-1}\right) \leq \delta\right\} \cup\left\{x \in X_{-1} \mid d\left(x, X_{1}\right) \leq \delta\right\} \cup X_{0}
\end{aligned}
$$

Let us first show that $\left\{x \in X_{1} \mid d\left(x, X_{-1}\right) \leq \delta\right\} \subset X_{1}^{+\delta} \backslash X_{1}^{-\delta}$. To this end, consider an $x \in X_{1}$ with $d\left(x, X_{-1}\right) \leq \delta$, where we check at once that $x \in X_{1}^{+\delta}$. Now, assume that $x \in X_{1}^{-\delta}=X \backslash\left(X \backslash X_{1}\right)^{+\delta}$. Then, we find that $x \notin\left(X \backslash X_{1}\right)^{+\delta}$ such that $d\left(x, X \backslash X_{1}\right)=d\left(x, X_{-1} \cup X_{0}\right)>\delta$. Hence, $x \notin X_{1}^{-\delta}$. Next, let us show that $\left\{x \in X_{-1} \mid d\left(x, X_{1}\right) \leq \delta\right\} \subset X_{1}^{+\delta} \backslash X_{1}^{-\delta}$. To this end, consider an $x \in X_{-1}$ with $d\left(x, X_{1}\right) \leq \delta$. Then, it is clear that $x \in X_{1}^{+\delta}$ by definition of $X_{1}^{+\delta}$. Furthermore, $x \notin X_{1}^{-\delta}$ since $X_{1}^{-\delta}=X \backslash\left(X_{-1}\right)^{+\delta} \subset X_{1}$. Having showed (16), we find together with the fact that $\lambda^{d}\left(X_{0}\right)=0$ since $X_{0}=\partial X_{1}$ is $(d-1)$ rectifiable that

$$
\lambda^{d}(\{x \in X \mid \Delta(x) \leq \delta\}) \leq \lambda^{d}\left(X_{1}^{+\delta} \backslash X_{1}^{-\delta}\right)
$$

Finally, with (15) we obtain that

$$
\lambda^{d}(\{x \in X \mid \Delta(x) \leq \delta\}) \leq \lambda^{d}\left(X_{1}^{+\delta} \backslash X_{1}^{-\delta}\right) \leq 4 \mathcal{H}^{d-1}(\{x \in X \mid \eta(x)=1 / 2\}) \cdot \delta
$$

for all $\delta \in\left(0, \delta^{*}\right]$.

## Proof of Lemma 3.1:

i) We define the set of indices

$$
J_{C}^{r}:=\left\{j \in J \mid \exists x \in A_{j}: \Delta_{\eta}(x)<r\right\}
$$

and define the set

$$
C_{r}:=\bigcup_{j \in J_{C}^{r}} A_{j} .
$$

Since $X \subset B_{r} \cup C_{r}$, it suffices to show that $C_{r} \subset A_{r}$. To show the latter we fix an $x \in C_{r}$. If $x \in X_{0}$ we immediately have $\Delta_{\eta}(x)=0 \leq 3 r$, hence we assume w.l.o.g. that $x \in X_{1}$. Then, there exists a $j \in J_{C}^{r}$ such that $x \in A_{j}$. Furthermore, there exists an $x^{*} \in A_{j}$ with $\Delta_{\eta}\left(x^{*}\right)<r$ and we find

$$
\begin{aligned}
\Delta_{\eta}(x) & =\inf _{x^{\prime} \in X_{-1}}\left\|x-x^{\prime}\right\|_{\infty} \\
& \leq \inf _{x^{\prime} \in X_{-1}}\left(\left\|x-x^{*}\right\|_{\infty}+\left\|x^{*}-x^{\prime}\right\|_{\infty}\right) \\
& \leq s+\Delta_{\eta}\left(x^{*}\right) \\
& <s+r
\end{aligned}
$$

where $\|\cdot\|_{\infty}$ is the supremum norm in $\mathbb{R}^{d}$. Since $s \leq 2 r$, it follows that $\Delta_{\eta}(x) \leq 3 r$ and therefore $x \in A_{r}$.
ii) We assume for $A_{j}$ with $j \in J_{B}^{r}$ that we have an $x_{1} \in A_{j} \cap X_{1} \neq \emptyset$ and an $x_{-1} \in A_{j} \cap X_{-1} \neq \emptyset$. Then, the connecting line $\overline{x_{-1} x_{1}}$ from $x_{-1}$ to $x_{1}$ is contained in $A_{j}$ since $A_{j}$ is convex and we have $\left\|x_{-1}-x_{1}\right\|_{\infty} \leq s$. Moreover, since $\Delta_{\eta}(x) \geq r$ for all $x \in B_{r}$ we have $X_{0} \cap B_{r}=\emptyset$. Next, pick an $m>1$ such that

$$
t_{0}=0, \quad t_{m}=1, \quad t_{i}=\frac{i}{m}
$$

and

$$
x_{i}:=t_{i} x_{-1}+\left(1-t_{i}\right) x_{1}
$$

for $i=0, \ldots, m$. Clearly, $x_{i} \in \overline{x_{-1} x_{1}}$ and $x_{i} \in X_{-1} \cup X_{1}$. Since $x_{0} \in X_{1}$ and $x_{m} \in X_{-1}$, there exists an $i$ with $x_{i} \in X_{1}$ and $x_{i+1} \in X_{-1}$ and we find that

$$
\left\|x_{i}-x_{i+1}\right\|_{\infty} \geq \Delta_{\eta}\left(x_{i}\right) \geq r
$$

On the other hand,

$$
\left\|x_{i}-x_{i+1}\right\|_{\infty}=\frac{1}{m}\left\|x_{-1}-x_{1}\right\|_{\infty} \leq \frac{s}{m} \leq \frac{2 r}{m}
$$

such that $r \leq \frac{2 r}{m}$, which is not true for $m \geq 3$. Hence, we can not have an $x_{1} \in A_{j} \cap X_{1} \neq \emptyset$ and an $x_{-1} \in A_{j} \cap X_{-1} \neq \emptyset$ for $j \in J_{B}^{r}$.

Proof of Lemma 3.2: For $f \in \mathcal{F}$ we have

$$
\begin{aligned}
& \mathcal{R}_{L_{T_{J}}, D}(f) \\
& =\int_{X \times Y} L_{T_{J}}(x, y, f(x)) d D(x, y) \\
& =\sum_{j \in J} \int_{A_{j} \times Y} L_{\text {class }}(y, f(x)) d D(x, y) .
\end{aligned}
$$

Next, we take a closer look at the risk on a single cell $A_{j}$ for a $j \in J$. That is,

$$
\int_{A_{j} \times Y} L_{\text {class }}(y, f(x)) d D(x, y)=\frac{1}{n} \sum_{i=1}^{n} \mathbf{1}_{A_{j}}\left(x_{i}\right) \mathbf{1}_{y_{i} \neq c_{i}} .
$$

The risk on a cell is the smaller the less often we have $y_{i} \neq c_{i}$ such that the best classifier on a cell is the one which decides by majority. This is true for the histogram rule by definition. Since the risk is zero on $A_{j}$ with $j \notin J$, the histogram rule minimizes the risk with respect to $L_{T_{J}}$.

Proof of Lemma 3.3: We define $h_{f}:=L_{B_{r}} \circ f-L_{B_{r}} \circ f_{L, P}^{*}$ for a measurable $f: X \rightarrow\{-1,1\}$. Since $\left(L_{B_{r}} \circ f-L_{B_{r}} \circ f_{L, P}^{*}\right)^{2}=\mathbf{1}_{B_{r}} \frac{\left|f-f_{L, P}^{*}\right|}{2}$ we obtain

$$
\begin{aligned}
& \mathbb{E}_{P}\left(h_{f}-\mathbb{E}_{P} h_{f}\right)^{2} \\
& \leq \mathbb{E}_{P}\left(h_{f}\right)^{2} \\
& =\mathbb{E}_{P}\left(L_{B_{r}} \circ f-L_{B_{r}} \circ f_{L, P}^{*}\right)^{2} \\
& =\frac{1}{2} \mathbb{E}_{P} \mathbf{1}_{B_{r}}\left|f-f_{L, P}^{*}\right|
\end{aligned}
$$

For $x \in B_{r}$ we have $\Delta_{\eta}(x) \geq r$ and thus we find with our lower-control assumption that

$$
r^{\gamma} \leq \Delta_{\eta}^{\gamma}(x) \leq c_{L C}|2 \eta(x)-1|
$$

and therefore

$$
1 \leq c_{L C} r^{-\gamma}|2 \eta(x)-1|
$$

By using $\mathbf{1}_{B_{r}} \frac{\left|f-f^{*}\right|}{2}=\mathbf{1}_{\left(X_{-1} \triangle\{f<0\}\right) \cap B_{r}}$, where $\triangle$ denotes the symmetric difference defined by $C \triangle D:=(C \backslash D) \cup(D \backslash C)$ for sets $C, D \subset X$ and by using Lemma A. 1 we obtain for the variance bound

$$
\begin{aligned}
\mathbb{E}_{P}\left(h_{f}-\mathbb{E}_{P} h_{f}\right)^{2} & \leq \frac{1}{2} \int \mathbf{1}_{B_{r}}(x)\left|f(x)-f_{L, P}^{*}(x)\right| d P_{X}(x) \\
& \leq \frac{c_{L C}}{2 r^{\gamma}} \int \mathbf{1}_{B_{r}}(x)|2 \eta(x)-1|\left|f(x)-f_{L, P}^{*}(x)\right| d P_{X}(x) \\
& =\frac{c_{L C}}{r^{\gamma}} \int_{\left(X_{-1} \triangle\{f<0\}\right) \cap B_{r}}|2 \eta(x)-1| d P_{X}(x) \\
& =\frac{c_{L C}}{r^{\gamma}}\left(\mathcal{R}_{L_{B_{r}}, P}(f)-\mathcal{R}_{L_{B_{r}}, P}^{*}\right) \\
& =\frac{c_{L C}}{r^{\gamma}} \mathbb{E}_{P} h_{f}
\end{aligned}
$$

Proof of Theorem 3.4: We define the set of cubes $A_{r}$ and $B_{r}$ as in (6), (7) for the choice of

$$
\begin{equation*}
r:=\tilde{c}_{\alpha, \gamma, d}\left(\frac{\tau}{s^{d} n}\right)^{\frac{1-\theta}{1+\gamma(2-\theta)}} \tag{17}
\end{equation*}
$$

where

$$
\begin{equation*}
\theta:=\frac{\alpha}{\alpha+\gamma} \tag{18}
\end{equation*}
$$

To estimate the excess risk $\mathcal{R}_{L, P}\left(h_{D, s}\right)-\mathcal{R}_{L, P}^{*}$, we split the risk as in (8) by

$$
\begin{aligned}
& \mathcal{R}_{L, P}\left(h_{D, s}\right)-\mathcal{R}_{L, P}^{*} \\
& \leq\left(\mathcal{R}_{L_{A_{r}}, P}\left(h_{D, s}\right)-\mathcal{R}_{L_{A_{r}}, P}^{*}\right)+\left(\mathcal{R}_{L_{B_{r}}, P}\left(h_{D, s}\right)-\mathcal{R}_{L_{B_{r}}, P}^{*}\right)
\end{aligned}
$$

This separation is valid by Lemma 3.1 (i), since $s \leq r$. To see that, we remark that

$$
\begin{aligned}
s \leq \tilde{c}_{\alpha, \gamma, d}\left(\frac{\tau}{s^{d} n}\right)^{\frac{1-\theta}{1+\gamma(2-\theta)}} & \Longleftrightarrow s^{\frac{1+\gamma(2-\theta)+d(1-\theta)}{1+\gamma(2-\theta)}} \leq \tilde{c}_{\alpha, \gamma, d}\left(\frac{\tau}{n}\right)^{\frac{1-\theta}{1+\gamma(2-\theta)}} \\
& \Longleftrightarrow s \leq\left(\tilde{c}_{\alpha, \gamma, d}\left(\frac{\tau}{n}\right)^{\frac{1-\theta}{1+\gamma(2-\theta)}}\right)^{\frac{1+\gamma(2-\theta)}{1+\gamma(2-\theta)+d(1-\theta)}}
\end{aligned}
$$

and conclude by replacing $\theta$ by (18) that $s \leq r$ holds if

$$
s \leq \tilde{c}_{\alpha, \gamma, d}^{\frac{(1+\gamma)(\alpha+\gamma)+\gamma^{2}}{(1+\gamma)(\alpha+\gamma)+\gamma^{2}+d \gamma}}\left(\frac{\tau}{n}\right)^{\frac{\gamma}{(1+\gamma)(\alpha+\gamma)+\gamma^{2}+d \gamma}}
$$

which equals (9). The rest of the proof is structured in three parts, where we establish error bounds on $A_{r}$ and $B_{r}$ in the first two parts and combine the results obtained in the third and last part of the proof. In the following we write $A:=A_{r}$ and $B:=B_{r}$ and keep in mind, that these sets depend on a parameter $r$. Furthermore we write $h_{D}:=h_{D, s}$.

Part 1: In the first part we establish an oracle inequality for $\mathcal{R}_{L_{B}, P}\left(h_{D, s}\right)-$ $\mathcal{R}_{L_{B}, P}^{*}$. Therefore we define $h_{f}^{B}:=L_{B} \circ f-L_{B} \circ f_{L_{B}, P}^{*}$ and find that

$$
\left\|h_{f}^{B}\right\|_{\infty}=\left\|L_{B} \circ f-L_{B} \circ f_{L_{B}, P}^{*}\right\|_{\infty} \leq 1
$$

for all $f \in \mathcal{F}$. Furthermore with Lemma 3.3 we obtain

$$
\begin{equation*}
\mathbb{E}_{P}\left(h_{f}^{B}\right)^{2} \leq \frac{c_{\mathrm{LC}}}{r^{\gamma}} \mathbb{E}_{P} h_{f}^{B} \leq \frac{c_{1}}{r^{\gamma}} \mathbb{E}_{P} h_{f}^{B} \tag{19}
\end{equation*}
$$

where $c_{1}:=\max \left\{c_{\mathrm{LC}}, 2^{\gamma}\right\}$. We observe that $r^{\gamma} \leq c_{1}$, since with assumption (10), where we rewrite the exponent by $\frac{(1+\gamma)(\alpha+\gamma)+\gamma^{2}}{\gamma}=\frac{1-\theta}{1+\gamma(2-\theta)}$, we find

$$
\begin{aligned}
r & =\tilde{c}_{\alpha, \gamma, d}\left(\frac{\tau}{s^{d} n}\right)^{\frac{1-\theta}{1+\gamma(2-\theta)}} \\
& \leq \tilde{c}_{\alpha, \gamma, d}\left(\left(\frac{\min \left\{\frac{\delta^{*}}{3}, 1\right\}}{\tilde{c}_{\alpha, \gamma, d}}\right)^{\frac{1+\gamma(2-\theta)}{1-\theta}}\right)^{\frac{1-\theta}{1+\gamma(2-\theta)}} \\
& =\min \left\{\frac{\delta^{*}}{3}, 1\right\} \\
& \leq 1
\end{aligned}
$$

and therefore $r^{\gamma} \leq 2^{\gamma} \leq c_{1}$. As we conclude from Lemma 3.2 that $h_{D}$ is an empirical risk minimizer over $\mathcal{F}$ for the loss $L_{B}$, we are able to use [7, Theorem 7.2], an improved oracle inequality for ERM. We obtain for all fixed $\tau \geq 1$ and $n \geq 1$ that

$$
\mathcal{R}_{L_{B}, P}\left(h_{D}\right)-\mathcal{R}_{L_{B}, P}^{*}<6\left(\mathcal{R}_{L_{B}, P, \mathcal{F}}^{*}-\mathcal{R}_{L_{B}, P}^{*}\right)+\frac{32 c_{1}(\log (|\mathcal{F}|+1)+\tau)}{r^{\gamma} n}
$$

holds with probability $P^{n} \geq 1-e^{-\tau}$, where $\mathcal{R}_{L_{B}, P, \mathcal{F}}^{*}:=\inf _{f \in \mathcal{F}} \mathcal{R}_{L_{B}, P}(f)$. Next, we refine the right-hand side of this oracle inequality. Obviously we have $|\mathcal{F}| \leq 2^{|J|}$. We bound the the cardinality $|J|$ by using a volume comparison argument. To this end, we define the set $\tilde{J}:=\left\{j \geq 1 \mid A_{j} \cap 2[-1,1]^{d} \neq \emptyset\right\}$ and observe that $\bigcup_{j \in J} A_{j} \subset \bigcup_{j \in J} A_{j} \subset 2 B_{\ell_{\propto}^{d}}$. Then,

$$
|J| s^{d}=\lambda^{d}\left(\bigcup_{j \in J} A_{j}\right) \leq \lambda^{d}\left(\bigcup_{j \in \tilde{J}} A_{j}\right) \leq \lambda^{d}\left(4 B_{\ell_{\infty}^{d}}\right)=8^{d},
$$

such that we deduce with $|J| \leq 8^{d} s^{-d}$ that

$$
\begin{aligned}
\log (|\mathcal{F}|+1) & \leq \log \left(2^{8^{d} s^{-d}}+1\right) \\
& \leq \log \left(2 \cdot 2^{8^{d} s^{-d}}\right) \\
& =\log \left(2^{8^{d} s^{-d}+1}\right) \\
& =\left(8^{d} s^{-d}+1\right) \log (2) \\
& \leq 8^{d} s^{-d}+1 \\
& \leq 8^{d+1} s^{-d} .
\end{aligned}
$$

Thus,

$$
\begin{equation*}
\mathcal{R}_{L_{B}, P}\left(h_{D}\right)-\mathcal{R}_{L_{B}, P}^{*}<6\left(\mathcal{R}_{L_{B}, P, \mathcal{F}}^{*}-\mathcal{R}_{L_{B}, P}^{*}\right)+\frac{32 c_{1}\left(8^{d+1} s^{-d}+\tau\right)}{r^{\gamma} n} \tag{20}
\end{equation*}
$$

holds with probability $P^{n} \geq 1-e^{-\tau}$.
Finally, we have to bound the approximation error $\mathcal{R}_{L_{B}, P, \mathcal{F}}^{*}-\mathcal{R}_{L_{B}, P}^{*}=$ $\inf _{f \in \mathcal{F}} \mathcal{R}_{L_{B}, P}(f)-\mathcal{R}_{L_{B}, P}^{*}$. We find with $h_{P, s} \in \mathcal{F}$ and Lemma A. 1 that

$$
\begin{align*}
\mathcal{R}_{L_{B}, P, \mathcal{F}}^{*}-\mathcal{R}_{L_{B}, P}^{*} & \leq \mathcal{R}_{L_{B}, P}\left(h_{P, s}\right)-\mathcal{R}_{L_{B}, P}^{*} \\
& =\int_{\left(X_{1} \Delta\left\{h_{P, s} \geq 0\right\}\right) \cap B}|2 \eta-1| d P_{X} \\
& =\sum_{j \in J_{B}^{r}} \int_{\left(X_{1} \Delta\left\{h_{P, s} \geq 0\right\}\right) \cap A_{j}}|2 \eta-1| d P_{X}  \tag{21}\\
& =0,
\end{align*}
$$

since $\left(X_{1} \triangle\left\{h_{P, s} \geq 0\right\}\right) \cap A_{j}=\emptyset$ for each $j \in J_{B}^{r}$. To see the latter, we first remark that the latter set contains those $x \in A_{j}$ for that either $h_{P, s}(x) \geq 0$ and
$\eta(x) \leq 1 / 2$ or $h_{P, s}(x)<0$ and $\eta(x)>1 / 2$. Since we have $A_{j} \subset X_{-1} \cup X_{1}$ we can ignore the case $\eta(x)=1 / 2$. Furthermore, we know by Lemma 3.1(ii) that either $A_{j} \cap X_{-1}=\emptyset$ or $A_{j} \cap X_{1}=\emptyset$. Let us first consider the case $A_{j} \cap X_{-1}=\emptyset$ and thus $A_{j} \subset X_{1}$. According to the definition of the histogram rule, c.f. (1), we find for all $x \in A_{j}$ that $h_{P, s}(x)=1$, since

$$
\begin{aligned}
& f_{P, s}(x) \\
& =P\left(A_{j}(x) \times\{1\}\right)-P\left(A_{j}(x) \times\{-1\}\right) \\
& =\int_{A_{j}} \int_{Y} \mathbf{1}_{A_{j} \times\{1\}}(x, y) P(d y \mid x) d P_{X}(x) \\
& \quad-\int_{A_{j}} \int_{Y} \mathbf{1}_{A_{j} \times\{-1\}}(x, y) P(d y \mid x) d P_{X}(x) \\
& =\int_{A_{j}} \mathbf{1}_{A_{j} \times\{1\}}(x, 1) \eta(x) d P_{X}(x)-\int_{A_{j}} \mathbf{1}_{A_{j} \times\{-1\}}(x,-1)(1-\eta(x)) d P_{X}(x) \\
& =\int_{A_{j}} 2 \eta(x)-1 d P_{X}(x) \\
& \geq 0
\end{aligned}
$$

Obviously we have $\eta(x) \geq 1 / 2$ and $h_{P, s}(x)=1$ for all $x \in A_{j}$. Analogously we can show for cells with $A_{j} \cap X_{1}=\emptyset$ for $j \in J_{B}^{r}$ that $\eta(x) \leq 1 / 2$ and $h_{P, s}(x)=-1$ for all $x \in A_{j}$. Hence, $\left(X_{1} \triangle\left\{h_{P, s} \geq 0\right\}\right) \cap A_{j}=\emptyset$ for all $j \in J_{B}^{r}$ and the approximation error vanishes on the set $B$.

Altogether, for the oracle inequality on B we obtain with (20) and (21) that

$$
\begin{equation*}
\mathcal{R}_{L_{B}, P}\left(h_{D}\right)-\mathcal{R}_{L_{B}, P}^{*}<\frac{32 c_{1}\left(8^{d+1} s^{-d}+\tau\right)}{r^{\gamma} n} \tag{22}
\end{equation*}
$$

holds with probability $P^{n} \geq 1-e^{-\tau}$.
Part 2: In the second part we establish an oracle inequality for $\mathcal{R}_{L_{A}, P}\left(h_{D}\right)-$ $\mathcal{R}_{L_{A}, P}^{*}$, again by using [7, Theorem 7.2]. Analogously to Part 1 we define $h_{f}^{A}:=$ $L_{A} \circ f-L_{A} \circ f_{L_{A}, P}^{*}$ for $f \in \mathcal{F}$ and find $\left\|h_{f}^{A}\right\|_{\infty} \leq 1$. Since $\left(h_{f_{0}}^{A}\right)^{2}=\mathbf{1}_{A} \frac{\left|f-f_{L, P}^{*}\right|}{2}=$
$\mathbf{1}_{\left(X_{-1} \triangle\{f<0\}\right) \cap A}$ we find with [Appendix, Lemma A.1] that

$$
\begin{align*}
& \mathbb{E}_{P}\left(h_{f_{0}}^{A}\right)^{2} \\
&= \frac{1}{2} \int_{A}\left|f_{0}(x)-f_{L_{A}, P}^{*}(x)\right| d P_{X}(x) \\
&= \frac{1}{2} \int_{A \cap\{|2 \eta-1| \geq t\}}\left|f_{0}(x)-f_{L_{A}, P}^{*}(x)\right| d P_{X}(x) \\
& \quad+\frac{1}{2} \int_{A \cap\{|2 \eta-1|<t\}}\left|f_{0}(x)-f_{L_{A}, P}^{*}(x)\right| d P_{X}(x) \\
& \leq \frac{1}{2 t} \int_{A \cap\{|2 \eta-1| \geq t\}}|2 \eta(x)-1|\left|f_{0}(x)-f_{L_{A}, P}^{*}(x)\right| d P_{X}(x)  \tag{23}\\
& \quad \quad P_{X}(\{x \in A:|2 \eta(x)-1|<t\}) \\
& \leq \frac{1}{2 t} \int_{A}|2 \eta(x)-1|\left|f_{0}(x)-f_{L_{A}, P}^{*}(x)\right| d P_{X}(x) \\
& \quad+P_{X}(\{x \in A:|2 \eta(x)-1|<t\}) \\
& \leq t^{-1} \mathbb{E}_{P} h_{f_{0}}^{A}+\min \left\{P_{X}(A), P_{X}(\{x \in X:|2 \eta(x)-1|<t\})\right\}
\end{align*}
$$

for all $t>0$. We turn our attention to the minimum and note, that by the definition of $A$ we have

$$
\begin{equation*}
P_{X}(A) \leq P_{X}\left(\left\{\Delta_{\eta}(x) \leq 3 r\right\}\right) \tag{24}
\end{equation*}
$$

For $x \in X$ with $|2 \eta(x)-1|<t$ by the definition of the lower control we conclude from

$$
\frac{\Delta_{\eta}^{\gamma}(x)}{c_{\mathrm{LC}}} \leq|2 \eta(x)-1|<t
$$

that

$$
\Delta_{\eta}(x) \leq\left(c_{\mathrm{LC}} t\right)^{\frac{1}{\gamma}}
$$

and consequently

$$
\begin{equation*}
\{x \in X:|2 \eta(x)-1|<t\} \subset\left\{x \in X: \Delta_{\eta}(x) \leq\left(c_{\mathrm{LC}} t\right)^{\frac{1}{\gamma}}\right\} \tag{25}
\end{equation*}
$$

Then we find by (24), (25) and by the definition of the strong margin exponent that

$$
\begin{align*}
& \min \left\{P_{X}(A), P_{X}(\{x \in X:|2 \eta(x)-1|<t\})\right\} \\
& \leq \min \left\{P_{X}\left(\left\{\Delta_{\eta}(x) \leq 3 r\right\}\right), P_{X}\left(\left\{x \in X: \Delta_{\eta}(x) \leq\left(c_{\mathrm{LC}} t\right)^{\frac{1}{\gamma}}\right\}\right)\right\}  \tag{26}\\
& \leq \min \left\{\left(c_{\mathrm{SME}} 3 r\right)^{\alpha}, c_{\mathrm{SME}}^{\alpha}\left(c_{\mathrm{LC}} t\right)^{\frac{\alpha}{\gamma}}\right\}
\end{align*}
$$

Combining (26) with (23) we obtain

$$
\begin{align*}
\mathbb{E}_{P}\left(h_{f_{0}}^{A}-\mathbb{E}_{P} h_{f_{0}}^{A}\right)^{2} & \leq t^{-1} \mathbb{E}_{P} h_{f_{0}}^{A}+\min \left\{\left(c_{\mathrm{SME}} 3 r\right)^{\alpha}, c_{\mathrm{SME}}^{\alpha}\left(c_{\mathrm{LC}} t\right)^{\frac{\alpha}{\gamma}}\right\}  \tag{27}\\
& \leq t^{-1} \mathbb{E}_{P} h_{f_{0}}^{A}+c_{\mathrm{SME}}^{\alpha}\left(c_{\mathrm{LC}} t\right)^{\frac{\alpha}{\gamma}}
\end{align*}
$$

Minimizing the right-hand side of (27) yields

$$
\min _{t>0}\left(t^{-1} \mathbb{E}_{P} h_{f_{0}}^{A}+c_{\mathrm{SME}}^{\alpha}\left(c_{\mathrm{LC}} t\right)^{\frac{\alpha}{\gamma}}\right)=c_{2}\left(\mathbb{E}_{P} h_{f_{0}}^{A}\right)^{\frac{\alpha}{\alpha+\gamma}}
$$

where $c_{2}:=\frac{\alpha+\gamma}{\gamma} c_{\text {SME }}^{\frac{\alpha \gamma}{\alpha+\gamma}}\left(\frac{\gamma c_{\mathrm{LC}}}{\alpha}\right)^{\frac{\alpha}{\alpha+\gamma}}$, such that with

$$
\begin{equation*}
V:=\max \left\{1, c_{2}\right\} \tag{28}
\end{equation*}
$$

and (18) we have

$$
\begin{equation*}
\mathbb{E}_{P}\left(h_{f_{0}}^{A}\right)^{2} \leq t^{-1} \mathbb{E}_{P} h_{f_{0}}^{A}+c_{\mathrm{SME}}\left(c_{\gamma} t^{\frac{1}{\gamma}}\right)^{\alpha}=c_{2}\left(\mathbb{E}_{P} h_{f_{0}}^{A}\right)^{\frac{\alpha}{\alpha+\gamma}} \leq V\left(\mathbb{E}_{P} h_{f_{0}}^{A}\right)^{\theta} \tag{29}
\end{equation*}
$$

Note, that the definition of $V$ yields $V^{\frac{1}{2-\theta}} \geq 1$. Since $h_{D}$ is an ERM over $\mathcal{F}$ for the loss $L_{A}$ due to Lemma 3.2, by using [7, Theorem 7.2] we obtain for fixed $\tau \geq 1$ and $n \geq 1$ that

$$
\begin{align*}
& \mathcal{R}_{L_{A}, P}\left(h_{D}\right)-\mathcal{R}_{L_{A}, P}^{*} \\
& <6\left(\mathcal{R}_{L_{A}, P, \mathcal{F}}^{*}-\mathcal{R}_{L_{A}, P}^{*}\right)+4\left(\frac{8 V(\log (|\mathcal{F}|+1)+\tau)}{n}\right)^{\frac{1}{2-\theta}} \tag{30}
\end{align*}
$$

holds with probability $P^{n} \geq 1-e^{-\tau}$. In order to refine the right-hand side in (30), we establish a bound on the cardinality $|\mathcal{F}|=2^{\left|J_{A}\right|}$ and on the approximation error. To bound the mentioned cardinality we use the fact that A lies in a tube around the decision line, that is $\bigcup_{j \in J_{A}} A_{j} \subset\left\{\Delta_{\eta}(x) \leq 3 r\right\}$, see (6). We remark that $3 r \leq \delta^{*}$ holds, where $\delta^{*}$ is the constant from Lemma 2.1, since with assumption (10) we have

$$
3 r=3 \tilde{c}_{\alpha, \gamma, d}\left(\frac{\tau}{s^{d} n}\right)^{\frac{1-\theta}{1+\gamma(2-\theta)}} \leq 3 \min \left\{\frac{\delta^{*}}{3}, 1\right\} \leq \delta^{*}
$$

Then, with Lemma 2.1 we find that

$$
\lambda^{d}\left(\left\{\Delta_{\eta}(x) \leq 3 r\right\}\right) \leq 12 \mathcal{H}^{d-1}(\{\eta=1 / 2\}) r
$$

and we obtain

$$
\left|J_{A}\right| s^{d}=\lambda^{d}\left(\bigcup_{j \in J_{A}} A_{j}\right) \leq \lambda^{d}\left(\left\{\Delta_{\eta}(x) \leq 3 r\right\}\right) \leq 12 \mathcal{H}^{d-1}(\{\eta=1 / 2\}) r
$$

This yields to

$$
\left|J_{A}\right| \leq 12 \mathcal{H}^{d-1}(\{\eta=1 / 2\}) r s^{-d}=c_{3} r s^{-d}
$$

where $c_{3}:=12 \mathcal{H}^{d-1}(\{\eta=1 / 2\})$. By $r \geq s \geq s^{d}$ we hence conclude that

$$
\begin{align*}
\log (|\mathcal{F}|+1) & =\log \left(2^{c_{3} r s^{-d}}+1\right) \\
& \leq \log \left(2 \cdot 2^{c_{3} r s^{-d}}\right) \\
& =\log \left(2^{c_{3} r s^{-d}+1}\right)  \tag{31}\\
& =\left(c_{3} r s^{-d}+1\right) \log (2) \\
& \leq c_{3} r s^{-d}+r s^{-d} \\
& \leq c_{4} r s^{-d}
\end{align*}
$$

where $c_{4}:=2 \max \left\{12 \mathcal{H}^{d-1}(\{\eta=1 / 2\}), 1\right\}$. Thus (30) changes to

$$
\begin{equation*}
\mathcal{R}_{L_{A}, P}\left(h_{D}\right)-\mathcal{R}_{L_{A}, P}^{*} \leq 6\left(\mathcal{R}_{L_{A}, P, \mathcal{F}}^{*}-\mathcal{R}_{L_{A}, P}^{*}\right)+4\left(\frac{8 V\left(c_{4} r s^{-d}+\tau\right)}{n}\right)^{\frac{1}{2-\theta}} \tag{32}
\end{equation*}
$$

with probability $P^{n} \geq 1-e^{-\tau}$.
Finally, we have to bound the approximation error $\mathcal{R}_{L_{A}, P, \mathcal{F}}^{*}-\mathcal{R}_{L_{A}, P}^{*}$ in (32). For $f_{0}=h_{P, s}$ we have with Lemma A. 1 that

$$
\begin{aligned}
\mathcal{R}_{L_{A}, P}\left(h_{P, s}\right)-\mathcal{R}_{L_{A}, P}^{*} & =\int_{\left(X_{1} \triangle\left\{h_{P, s} \geq 0\right\}\right) \cap A}|2 \eta-1| d P_{X} \\
& =\sum_{j \in J_{A}^{r}} \int_{\left(X_{1} \triangle\left\{h_{P, s} \geq 0\right\}\right) \cap A_{j}}|2 \eta-1| d P_{X}
\end{aligned}
$$

We split $J_{A}^{r}$ in indices where cells do not intersect the decision line and those which do by

$$
\begin{aligned}
& J_{A_{1}}^{r}:=\left\{j \in J_{A}^{r} \mid P_{X}\left(A_{j} \cap X_{1}\right)=0 \vee P_{X}\left(A_{j} \cap X_{-1}\right)=0\right\} \\
& J_{A_{2}}^{r}:=\left\{j \in J_{A}^{r} \mid P_{X}\left(A_{j} \cap X_{1}\right)>0 \wedge P_{X}\left(A_{j} \cap X_{-1}\right)>0\right\} .
\end{aligned}
$$

such that

$$
\begin{aligned}
& \sum_{j \in J_{A}^{r}} \int_{\left(X_{1} \triangle\left\{h_{P, s} \geq 0\right\}\right) \cap A_{j}}|2 \eta-1| d P_{X} \\
& =\sum_{j \in J_{A_{1}}^{r}} \int_{\left(X_{1} \triangle\left\{h_{P, s} \geq 0\right\}\right) \cap A_{j}}|2 \eta-1| d P_{X} \\
& \quad+\sum_{j \in J_{A_{2}}^{r}} \int_{\left(X_{1} \triangle\left\{h_{P, s} \geq 0\right\}\right) \cap A_{j}}|2 \eta-1| d P_{X}
\end{aligned}
$$

We notice that, as in the calculation of the approximation error in Part 1, the first sum vanishes, since $\left(X_{1} \Delta\left\{h_{P, s} \geq 0\right\}\right) \cap A_{j}=\emptyset$ for all $j \in J_{A_{1}}^{r}$. Moreover, we remark that $J_{A_{2}}^{r}$ only contains cells of width $s$ that intersect the decision
boundary. Hence, by using the margin-noise assumption we find

$$
\begin{align*}
\mathcal{R}_{L_{A}, P}\left(h_{P, s}\right)-\mathcal{R}_{L_{A}, P}^{*} & =\sum_{j \in J_{A_{2}}^{r}} \int_{\left(X_{1} \triangle\left\{h_{P, s} \geq 0\right\}\right) \cap A_{j}}|2 \eta-1| d P_{X} \\
& \leq \int_{\{\Delta(x) \leq s\}}|2 \eta-1| d P_{X}  \tag{33}\\
& \leq\left(c_{\mathrm{MNE}} s\right)^{\beta}
\end{align*}
$$

Altogether for the oracle inequality on A with (32) we find that

$$
\begin{equation*}
\mathcal{R}_{L_{A}, P}\left(h_{D}\right)-\mathcal{R}_{L_{A}, P}^{*} \leq 6\left(c_{\mathrm{MNE}}\right)^{\beta}+4\left(\frac{8 V\left(c_{4} r s^{-d}+\tau\right)}{n}\right)^{\frac{1}{2-\theta}} \tag{34}
\end{equation*}
$$

holds with with probability $P^{n} \geq 1-e^{-\tau}$.
Part 3: In the last part we combine the results obtained in Part 1, the oracle inequality on B and Part 2, the oracle inequality on A. That means, with the separation in (8) we obtain with (22) and (34) for the oracle inequality on $X$ that

$$
\begin{align*}
& \mathcal{R}_{L, P}\left(h_{D, s}\right)-\mathcal{R}_{L, P}^{*} \\
& \leq\left(\mathcal{R}_{L_{A}, P}\left(h_{D, s}\right)-\mathcal{R}_{L_{A}, P}^{*}\right)+\left(\mathcal{R}_{L_{B}, P}\left(h_{D, s}\right)-\mathcal{R}_{L_{B}, P}^{*}\right) \\
& \leq 6\left(c_{\mathrm{MNE}} s\right)^{\beta}+4\left(\frac{8 V\left(c_{4} r s^{-d}+\tau\right)}{n}\right)^{\frac{1}{2-\theta}}+\frac{32 c_{1}\left(8^{d+1} s^{-d}+\tau\right)}{r^{\gamma} n} \tag{35}
\end{align*}
$$

holds with probability $P^{n} \geq 1-2 e^{-\tau}$. Since $s \in(0,1]$ and $r \geq s$, we find that $r s^{-d} \geq 1$. Together with the fact $s^{-d}, \tau \geq 1$ and $c_{4} \geq 1$ it follows that

$$
\begin{aligned}
& \mathcal{R}_{L, P}\left(h_{D, s}\right)-\mathcal{R}_{L, P}^{*} \\
& \leq 6\left(c_{\mathrm{MNE}} s\right)^{\beta}+4\left(\frac{8 V\left(c_{4} r s^{-d}+\tau\right)}{n}\right)^{\frac{1}{2-\theta}}+\frac{32 c_{1}\left(8^{d+1} s^{-d}+\tau\right)}{r^{\gamma} n} \\
& \leq 6\left(c_{\mathrm{MNE}} s\right)^{\beta}+4\left(\frac{8 V\left(c_{4} \tau r s^{-d}+c_{4} \tau r s^{-d}\right)}{n}\right)^{\frac{1}{2-\theta}}+\frac{32 c_{1}\left(8^{d+1} \tau s^{-d}+\tau s^{-d}\right)}{r^{\gamma} n} \\
& \leq 6\left(c_{\mathrm{MNE}} s\right)^{\beta}+4\left(\frac{c_{5} \tau r s^{-d}}{n}\right)^{\frac{1}{2-\theta}}+\frac{c_{6} \tau s^{-d}}{r^{\gamma} n} \\
& \leq 6\left(c_{\mathrm{MNE}} s\right)^{\beta}+r^{\frac{1}{2-\theta}} 4\left(\frac{c_{5} \tau}{s^{d} n}\right)^{\frac{1}{2-\theta}}+\frac{c_{6} \tau}{r^{\gamma} s^{d} n}
\end{aligned}
$$

where $c_{5}:=24 V \max \left\{12 \mathcal{H}^{d-1}(\{\eta=1 / 2\}), 1\right\}$ and $c_{6}:=64 \cdot 8^{d+1} \max \left\{c_{L C}, 2^{\gamma}\right\}$. Thus, inserting $r$, defined in (17), with the choice of $\tilde{c}_{\alpha, \gamma, d}:=$
$\left(\frac{\left(\gamma(2-\theta) c_{6}\right)^{2-\theta}}{4^{2-\theta} c_{5}}\right)^{\frac{1}{1+\gamma(2-\theta)}}$ minimizes the right-hand side and yields to

$$
\begin{aligned}
& \mathcal{R}_{L, P}\left(h_{D, s}\right)-\mathcal{R}_{L, P}^{*} \\
& \leq 6\left(c_{\mathrm{MNE}} s\right)^{\beta}+r^{\frac{1}{2-\theta}} 4\left(\frac{c_{5} \tau}{s^{d} n}\right)^{\frac{1}{2-\theta}}+\frac{c_{6} \tau}{r^{\gamma} s^{d} n} \\
& =6\left(c_{\mathrm{MNE}} s\right)^{\beta}+4\left(\tilde{c}_{\alpha, \gamma, d} c_{5}\right)^{\frac{1}{2-\theta}}\left(\frac{\tau}{s^{d} n}\right)^{\frac{2-\theta+\gamma(2-\theta)}{(1+\gamma(2-\theta))(2-\theta)}}+\frac{c_{6}}{\tilde{c}_{\alpha, \gamma, d}^{\gamma}}\left(\frac{\tau}{s^{d} n}\right)^{\frac{1+\gamma}{1+\gamma(2-\theta)}} \\
& =6\left(c_{\mathrm{MNE}} s\right)^{\beta}+\left(\frac{\tilde{c}_{\alpha, \gamma, d}^{\frac{1+\gamma(2-\theta)}{2-\theta}} 4 c_{5}^{\frac{1}{2-\theta}}+c_{6}}{\tilde{c}_{\alpha, \gamma, d}^{\gamma}}\right)\left(\frac{\tau}{s^{d} n}\right)^{\frac{1+\gamma}{1+\gamma(2-\theta)}} \\
& =6\left(c_{\mathrm{MNE}} s\right)^{\beta}+\left(\frac{\gamma(2-\theta) c_{6}+c_{6}}{\tilde{c}_{\alpha, \gamma, d}^{\gamma}}\right)\left(\frac{\tau}{s^{d} n}\right)^{\frac{1+\gamma}{1+\gamma(2-\theta)}} \\
& \leq 6\left(c_{\mathrm{MNE}} s\right)^{\beta}+\left(\frac{2 c_{6} \max \{\gamma(2-\theta), 1\}}{\tilde{c}_{\alpha, \gamma, d}^{\gamma}}\right)\left(\frac{\tau}{s^{d} n}\right)^{\frac{1+\gamma}{1+\gamma(2-\theta)}}
\end{aligned}
$$

and we find again by inserting $\theta$ that

$$
\begin{equation*}
\mathcal{R}_{L, P}\left(h_{D, s}\right)-\mathcal{R}_{L, P}^{*} \leq 6\left(c_{\mathrm{MNE}} s\right)^{\beta}+c_{\alpha, \gamma, d}\left(\frac{\tau}{s^{d} n}\right)^{\frac{(1+\gamma)(\alpha+\gamma)}{(1+\gamma)(\alpha+\gamma)+\gamma^{2}}} \tag{36}
\end{equation*}
$$

holds with probability $P^{n} \geq 1-2 e^{-\tau}$, where $c_{\alpha, \gamma, d}:=\frac{2 c_{6} \max \{\gamma(2-\theta), 1\}}{\tilde{c}_{\alpha, \gamma, d}^{\gamma}}=$ $\frac{2 c_{6} \max \left\{\frac{\gamma(\alpha+2 \gamma)}{\alpha+\gamma}, 1\right\}}{\tilde{c}_{\alpha, \gamma, d}^{\gamma}}$.

Proof of Theorem 3.5: We begin by proving that the chosen sequence $s_{n}$ satisfies assumptions (9) and (10). To this end, we define $n_{\tau, \alpha, \beta, \gamma, d}:=\left(\frac{\tilde{c}_{\alpha, \beta, \gamma, \tau, d}}{c_{1}}\right)^{\frac{1}{\zeta_{1}}}$ with $c_{1}:=\tilde{c}_{\alpha, \gamma, d}^{\frac{\kappa+\gamma^{2}}{\kappa+\gamma^{2}+d \gamma}} \tau^{\frac{\gamma}{\kappa+\gamma^{2}+d \gamma}}$, where $\tilde{c}_{\alpha, \gamma, d}$ is the constant from Theorem 3.4, and $\zeta_{1}:=\frac{\kappa\left(\kappa+\gamma^{2}+d \gamma\right)-\gamma\left(\beta\left(\kappa+\gamma^{2}\right)+d \kappa\right)}{\left(\beta\left(\kappa+\gamma^{2}\right)+d \kappa\right)\left(\kappa+\gamma^{2}+d \gamma\right)}$. We remark that $\zeta_{1} \geq 0$ since we find by $\beta \leq \gamma^{-1}(1+\gamma)(\alpha+\gamma)$ that

$$
\begin{aligned}
\kappa\left(\kappa+\gamma^{2}+d \gamma\right)-\gamma\left(\beta\left(\kappa+\gamma^{2}\right)+d \kappa\right) & =\kappa^{2}+k \gamma^{2}-\gamma \beta \kappa-\beta \gamma^{3} \\
& \geq \kappa^{2}+\kappa \gamma^{2}-\kappa^{2}-\kappa \gamma^{2} \\
& =0 .
\end{aligned}
$$

Then, for $n \geq n_{\tau, \alpha, \beta, \gamma, d}$ a simple calculation shows that the latter is equivalent to

$$
c_{1} n^{\frac{-\gamma}{\kappa+\gamma^{2}+d \gamma}} \geq \tilde{c}_{\alpha, \beta, \gamma, \tau, d} n^{-\frac{\kappa}{\beta\left(\kappa+\gamma^{2}\right)+d \kappa}}
$$

which equals assumption (9) with $s_{n}:=\tilde{c}_{\alpha, \beta, \gamma, \tau, d} n^{-\frac{\kappa}{\beta\left(\kappa+\gamma^{2}\right)+d \kappa}}$. To see that assumption (10) is satisfied we define $\tilde{n}_{\tau, \alpha, \beta, \gamma, d}:=\left(\frac{c_{2}}{\tilde{c}_{\alpha, \beta, \gamma, \tau, d}}\right)^{\frac{1}{\zeta_{2}}}$ with $c_{2}:=$
$\tau^{\frac{1}{d}}\left(\frac{\tilde{c}_{\alpha, \gamma, d}}{\min \left\{\frac{\left.\delta^{*}, 1\right\}}{3}, 1\right.}\right)^{\frac{\kappa+\gamma^{2}}{d \gamma}}$, where $\tilde{c}_{\alpha, \gamma, d}$ is the constant from Theorem $3.4, \delta^{*}$ the one from Lemma 2.1 and where $\zeta_{2}:=\frac{\beta\left(\kappa+\gamma^{2}\right)}{d\left(\beta\left(\kappa+\gamma^{2}\right)+d \kappa\right)}$. Then, a simple transformation shows again that for all $n \geq \tilde{n}_{\tau, \alpha, \beta, \gamma, d}$ we find

$$
\tilde{c}_{\alpha, \beta, \gamma, \tau, d} n^{-\frac{\kappa}{\left(\beta\left(\kappa+\gamma^{2}\right)+d \kappa\right)}} \geq c_{2} n^{-1 / d},
$$

which equals assumption (10) with $s_{n}:=\tilde{c}_{\alpha, \beta, \gamma, \tau, d} n^{-\frac{\kappa}{\left(\beta\left(\kappa+\gamma^{2}\right)+d \kappa\right)}}$.
Finally, we obtain for all $n \geq n_{0}:=\left\lceil\max \left\{n_{\tau, \alpha, \beta, \gamma, d}, \tilde{n}_{\tau, \alpha, \beta, \gamma, d}\right\}\right\rceil$ by inserting our chosen sequence $s_{n}$, satisfying (9) and (10), in (11) that

$$
\begin{aligned}
& \mathcal{R}_{L, P}\left(h_{D, s_{n}}\right)-\mathcal{R}_{L, P}^{*} \\
& \leq 6\left(c_{\mathrm{MNE}} s\right)^{\beta}+c_{\alpha, \gamma, d}\left(\frac{\tau}{s^{d} n}\right)^{\frac{\kappa}{\kappa+\gamma^{2}}}
\end{aligned}
$$

$$
\begin{aligned}
& =\left(\frac{6 c_{\mathrm{MNE}}^{\beta} \tilde{c}_{\alpha, \beta, \gamma, \tau, d}^{\frac{\beta\left(\kappa+\gamma^{2}\right)+d \kappa}{\kappa+\gamma^{2}}}+c_{\alpha, \gamma, d} \tau^{\frac{\kappa}{\kappa+\gamma^{2}}}}{\tilde{c}_{\alpha, \beta, \gamma, \tau, d}^{\kappa+\gamma^{2}}}\right) n^{-\frac{\beta \kappa}{\beta\left(\kappa+\gamma^{2}\right)+d \kappa}} \\
& =\left(\frac{\frac{d \kappa}{\beta\left(\kappa+\gamma^{2}\right)} c_{\alpha, \gamma, d} \tau^{\frac{\kappa}{\kappa+\gamma^{2}}}+c_{\alpha, \gamma, d} \tau^{\frac{\kappa}{\kappa+\gamma^{2}}}}{\tilde{c}_{\alpha, \beta, \gamma, \tau, d}^{\frac{d \kappa}{\kappa+\gamma^{2}}}}\right) n^{-\frac{\beta \kappa}{\beta\left(\kappa+\gamma^{2}\right)+d \kappa}} \\
& \leq\left(\frac{2 \max \left\{\frac{d \kappa}{\beta\left(\kappa+\gamma^{2}\right)}, 1\right\} c_{\alpha, \gamma, \delta} \tau^{\frac{\kappa}{\kappa+\gamma^{2}}}}{\tilde{c}_{\alpha, \beta, \gamma, \tau, d}^{\frac{d \kappa}{\kappa+\gamma^{2}}}}\right) n^{-\frac{\beta \kappa}{\beta\left(\kappa+\gamma^{2}\right)+d \kappa}} \\
& =c_{\alpha, \beta, \gamma, \tau, d} n^{-\frac{\beta \kappa}{\beta\left(\kappa+\gamma^{2}\right)+d \kappa}}
\end{aligned}
$$

holds with probability $P^{n} \geq 1-2 e^{-\tau}$, where $c_{\alpha, \beta, \gamma, \tau, d}:=$ $2 \max \left\{\frac{d \kappa}{\beta\left(\kappa+\gamma^{2}\right)}, 1\right\} c_{\alpha, \gamma, \delta} \tau^{\frac{\kappa}{\kappa+\gamma^{2}}} \cdot \tilde{c}_{\alpha, \beta, \gamma, \tau, d}^{-\frac{d \kappa}{\kappa+\gamma^{2}}}$.

## Appendix A: Appendix

Lemma A.1. Let $Y:=\{-1,1\}$ and $P$ be a probability measure on $X \times Y$. For $\eta(x):=P(y=1 \mid x), x \in X$ define the set $X_{1}:=\{x \in X \mid \eta(x)>1 / 2\}$. Let $L$ be the classification loss and consider for $A \subset X$ the loss $L_{A}(x, y, t):=$ $\mathbf{1}_{A}(x) L(y, t)$, where $y \in Y, t \in \mathbb{R}$. For a measurable $f: X \rightarrow \mathbb{R}$ we then have

$$
\mathcal{R}_{L_{A}, P}(f)-\mathcal{R}_{L_{A}, P}^{*}=\int_{\left(X_{1} \triangle\{f \geq 0\}\right) \cap A}|2 \eta(x)-1| d P_{X}(x)
$$

where $\triangle$ denotes the symmetric difference.

Proof of Lemma A.1: It is well known, e.g., [7, Example 3.8], that

$$
\begin{align*}
& \mathcal{R}_{L_{A}, P}(f)-\mathcal{R}_{L_{A}, P}^{*} \\
& =\int_{A}|2 \eta(x)-1| \cdot \mathbf{1}_{(-\infty, 0)}((2 \eta(x)-1) \operatorname{sign} f(x)) d P_{X}(x) \tag{37}
\end{align*}
$$

Next, for $P_{X}$-almost all $x \in A$ we have

$$
\mathbf{1}_{(-\infty, 0]}((2 \eta(x)-1) \operatorname{sign} f(x))=1 \Leftrightarrow(2 \eta(x)-1) \operatorname{sign} f(x) \leq 0
$$

The latter is true if for $x \in A$ holds that $f(x)<0$ and $\eta(x)>1 / 2$ or that $f(x) \geq 0$ and $\eta(x) \leq 1 / 2$ or that $\eta(x)=1 / 2$. However, for $\eta(x)=1 / 2$ we have $|2 \eta(x)-1|=0$ and hence this case can be ignored. Then, the latter obviously equals the set $\left(X_{1} \triangle\{f \geq 0\}\right) \cap A$ and we obtain in (37)

$$
\mathcal{R}_{L_{A}, P}(f)-\mathcal{R}_{L_{A}, P}^{*}=\int_{\left(X_{1} \triangle\{f \geq 0\}\right) \cap A}|2 \eta(x)-1| d P_{X}(x)
$$

Lemma A.2. Let $X:=[-1,1]^{d}$ and $P$ be a probability measure on $X \times\{-1,1\}$ with fixed version $\eta: X \rightarrow[0,1]$ of its posterior probability. Then, if $\eta$ is Hoeldercontinuous with exponent $\gamma$, we have that $\Delta_{\eta}$ controls the noise from below with exponent $\gamma$.
Proof of Lemma A.2: Fix w.l.o.g. an $x \in X_{1}$. Then, $\eta(x)>1 / 2$. Since $\eta$ is Hoelder-continuous with exponent $\gamma$, there exists a constant $c>0$ such that we have

$$
|2 \eta(x)-1|=2|\eta(x)-1 / 2| \leq 2\left|\eta(x)-n\left(x^{\prime}\right)\right| \leq 2 c\left(d\left(x, x^{\prime}\right)\right)^{\gamma}
$$

for all $x^{\prime} \in X_{-1}$ and hence

$$
|2 \eta(x)-1| \leq 2 c \inf _{\tilde{x} \in X_{-1}}(d(x, \tilde{x}))^{\gamma}=2 c \Delta_{\eta}^{\gamma}(x)
$$

Obviously the last inequality holds immediately for $x \in X$ with $\eta(x)=1 / 2$.

## References

[1] Audibert, J. Y. and Tsybakov, A. (2007). Fast learning rates for plug-in classifiers. Ann. Statist. 35 608-633. MR2336861
[2] Binev, P., Cohen, A., Dahmen, W. and DeVore, R. (2014). Classification algorithms using adaptive partitioning. Ann. Statist 42 2141-2163. MR3269976
[3] Breiman, L. (2001). Random Forests. Machine Learning 45 5-32.
[4] Devroye, L., Györfi, L. and Lugosi, L. (1996). A Probabilistc Theory of Pattern Recognition. Springer. MR1383093
[5] Federer, H. (1969). Geometric measure theory. Springer. MR0257325
[6] Steinwart, I. (2015). Fully adaptive density-based clustering. Ann. Statist. 43 2132-2167. MR3396981
[7] Steinwart, I. and Christmann, A. (2008). Support Vector Machines. Springer. MR2450103

Ingrid Blaschzyk
Fachbereich Mathematik, Universität Stuttgart, Pfaffenwaldring 57, 70569 Stuttgart, Germany
E-Mail: ingrid.blaschzyk@mathematik.uni-stuttgart.de
Ingo Steinwart
Fachbereich Mathematik, Universität Stuttgart, Pfaffenwaldring 57, 70569 Stuttgart, Germany
E-Mail: steinwart@mathematik.uni-stuttgart.de
WWW: http://www.isa.uni-stuttgart.de/Steinwart/

## Erschienene Preprints ab Nummer 2012-001

Komplette Liste: http://www.mathematik.uni-stuttgart.de/preprints
2016-004 Blaschzyk, I.; Steinwart, I.: Improved Classification Rates under Refined Margin Conditions
2016-003 Feistauer, M.; Roskovec, F.; Sändig, AM.: Discontinuous Galerkin Method for an Elliptic Problem with Nonlinear Newton Boundary Conditions in a Polygon
2016-002 Steinwart, I.: A Short Note on the Comparison of Interpolation Widths, Entropy Numbers, and Kolmogorov Widths
2016-001 Köster, I.: Sylow Numbers in Spectral Tables
2015-016 Hang, H.; Steinwart, I.: A Bernstein-type Inequality for Some Mixing Processes and Dynamical Systems with an Application to Learning
2015-015 Steinwart, I.: Representation of Quasi-Monotone Functionals by Families of Separating Hyperplanes
2015-014 Muhammad, F.; Steinwart, I.: An SVM-like Approach for Expectile Regression
2015-013 Nava-Yazdani, E.: Splines and geometric mean for data in geodesic spaces
2015-012 Kimmerle, W.; Köster, I.: Sylow Numbers from Character Tables and Group Rings
2015-011 Györfi, L.; Walk, H.: On the asymptotic normality of an estimate of a regression functional
2015-010 Gorodski, C, Kollross, A.: Some remarks on polar actions
2015-009 Apprich, C.; Höllig, K.; Hörner, J.; Reif, U.: Collocation with WEB-Splines
2015-008 Kabil, B.; Rodrigues, M.: Spectral Validation of the Whitham Equations for Periodic Waves of Lattice Dynamical Systems
2015-007 Kollross, A.: Hyperpolar actions on reducible symmetric spaces
2015-006 Schmid, J.; Griesemer, M.: Well-posedness of Non-autonomous Linear Evolution Equations in Uniformly Convex Spaces
2015-005 Hinrichs, A.; Markhasin, L.; Oettershagen, J.; Ullrich, T.: Optimal quasi-Monte Carlo rules on higher order digital nets for the numerical integration of multivariate periodic functions
2015-004 Kutter, M.; Rohde, C.; Sändig, A.-M.: Well-Posedness of a Two Scale Model for Liquid Phase Epitaxy with Elasticity
2015-003 Rossi, E.; Schleper, V.: Convergence of a numerical scheme for a mixed hyperbolic-parabolic system in two space dimensions
2015-002 Döring, M.; Györfi, L.; Walk, H.: Exact rate of convergence of kernel-based classification rule
2015-001 Kohler, M.; Müller, F.; Walk, H.: Estimation of a regression function corresponding to latent variables
2014-021 Neusser, J.; Rohde, C.; Schleper, V.: Relaxed Navier-Stokes-Korteweg Equations for Compressible Two-Phase Flow with Phase Transition
2014-020 Kabil, B.; Rohde, C.: Persistence of undercompressive phase boundaries for isothermal Euler equations including configurational forces and surface tension
2014-019 Bilyk, D.; Markhasin, L.: BMO and exponential Orlicz space estimates of the discrepancy function in arbitrary dimension
2014-018 Schmid, J.: Well-posedness of non-autonomous linear evolution equations for generators whose commutators are scalar
2014-017 Margolis, L.: A Sylow theorem for the integral group ring of $\operatorname{PSL}(2, q)$

2014-016 Rybak, I.; Magiera, J.; Helmig, R.; Rohde, C.: Multirate time integration for coupled saturated/unsaturated porous medium and free flow systems
2014-015 Gaspoz, F.D.; Heine, C.-J.; Siebert, K.G.: Optimal Grading of the Newest Vertex Bisection and $H^{1}$-Stability of the $L_{2}$-Projection
2014-014 Kohler, M.; Krzyżak, A.; Walk, H.: Nonparametric recursive quantile estimation
2014-013 Kohler, M.; Krzyżak, A.; Tent, R.; Walk, H.: Nonparametric quantile estimation using importance sampling
2014-012 Györfi, L.; Ottucsák, G.; Walk, H.: The growth optimal investment strategy is secure, too.
2014-011 Györfi, L.; Walk, H.: Strongly consistent detection for nonparametric hypotheses
2014-010 Köster, I.: Finite Groups with Sylow numbers $\left\{q^{x}, a, b\right\}$
2014-009 Kahnert, D.: Hausdorff Dimension of Rings
2014-008 Steinwart, I.: Measuring the Capacity of Sets of Functions in the Analysis of ERM
2014-007 Steinwart, I.: Convergence Types and Rates in Generic Karhunen-Loève Expansions with Applications to Sample Path Properties
2014-006 Steinwart, I.; Pasin, C.; Williamson, R.; Zhang, S.: Elicitation and Identification of Properties
2014-005 Schmid, J.; Griesemer, M.: Integration of Non-Autonomous Linear Evolution Equations
2014-004 Markhasin, L.: $\quad L_{2}$ - and $S_{p, q}^{r} B$-discrepancy of (order 2) digital nets
2014-003 Markhasin, L.: Discrepancy and integration in function spaces with dominating mixed smoothness
2014-002 Eberts, M.; Steinwart, I.: Optimal Learning Rates for Localized SVMs
2014-001 Giesselmann, J.: A relative entropy approach to convergence of a low order approximation to a nonlinear elasticity model with viscosity and capillarity
2013-016 Steinwart, I.: Fully Adaptive Density-Based Clustering
2013-015 Steinwart, I.: Some Remarks on the Statistical Analysis of SVMs and Related Methods
2013-014 Rohde, C.; Zeiler, C.: A Relaxation Riemann Solver for Compressible Two-Phase Flow with Phase Transition and Surface Tension
2013-013 Moroianu, A.; Semmelmann, U.: Generalized Killling spinors on Einstein manifolds
2013-012 Moroianu, A.; Semmelmann, U.: Generalized Killing Spinors on Spheres
2013-011 Kohls, K; Rösch, A.; Siebert, K.G.: Convergence of Adaptive Finite Elements for Control Constrained Optimal Control Problems
2013-010 Corli, A.; Rohde, C.; Schleper, V.: Parabolic Approximations of Diffusive-Dispersive Equations
2013-009 Nava-Yazdani, E.; Polthier, K.: De Casteljau's Algorithm on Manifolds
2013-008 Bächle, A.; Margolis, L.: Rational conjugacy of torsion units in integral group rings of non-solvable groups
2013-007 Knarr, N.; Stroppel, M.J.: Heisenberg groups over composition algebras
2013-006 Knarr, N.; Stroppel, M.J.: Heisenberg groups, semifields, and translation planes
2013-005 Eck, C.; Kutter, M.; Sändig, A.-M.; Rohde, C.: A Two Scale Model for Liquid Phase Epitaxy with Elasticity: An Iterative Procedure
2013-004 Griesemer, M.; Wellig, D.: The Strong-Coupling Polaron in Electromagnetic Fields

2013-003 Kabil, B.; Rohde, C.: The Influence of Surface Tension and Configurational Forces on the Stability of Liquid-Vapor Interfaces
2013-002 Devroye, L.; Ferrario, P.G.; Györfi, L.; Walk, H.: Strong universal consistent estimate of the minimum mean squared error
2013-001 Kohls, K.; Rösch, A.; Siebert, K.G.: A Posteriori Error Analysis of Optimal Control Problems with Control Constraints
2012-013 Diaz Ramos, J.C.; Dominguez Vázquez, M.; Kollross, A.: Polar actions on complex hyperbolic spaces
2012-012 Moroianu; A.; Semmelmann, U.: Weakly complex homogeneous spaces
2012-011 Moroianu; A.; Semmelmann, U.: Invariant four-forms and symmetric pairs
2012-010 Hamilton, M.J.D.: The closure of the symplectic cone of elliptic surfaces
2012-009 Hamilton, M.J.D.: Iterated fibre sums of algebraic Lefschetz fibrations
2012-008 Hamilton, M.J.D.: The minimal genus problem for elliptic surfaces
2012-007 Ferrario, P.: Partitioning estimation of local variance based on nearest neighbors under censoring
2012-006 Stroppel, M.: Buttons, Holes and Loops of String: Lacing the Doily
2012-005 Hantsch, F.: Existence of Minimizers in Restricted Hartree-Fock Theory
2012-004 Grundhöfer, T.; Stroppel, M.; Van Maldeghem, H.: Unitals admitting all translations
2012-003 Hamilton, M.J.D.: Representing homology classes by symplectic surfaces
2012-002 Hamilton, M.J.D.: On certain exotic 4-manifolds of Akhmedov and Park
2012-001 Jentsch, T.: Parallel submanifolds of the real 2-Grassmannian

