

Conditional Coverage Estimation for High-Quality Prediction Intervals

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Abstract

Deep learning has been recently studied to generate high-quality prediction intervals (PIs) for uncertainty quantification in regression tasks, including recent applications in simulation metamodeling. The high-quality criterion requires PIs to be as narrow as possible, whilst maintaining a pre-specified level of data (marginal) coverage. However, most existing works for high-quality PIs lack accurate information on conditional coverage, which may cause unreliable predictions if it is significantly smaller than the marginal coverage. To address this problem, we propose an end-to-end framework which could output high-quality PIs and simultaneously provide their conditional coverage estimation. In doing so, we design a new loss function that is both easy-to-implement and theoretically justified via an exponential concentration bound. Our evaluation on real-world benchmark datasets and synthetic examples shows that our approach not only achieves competitive results on high-quality PIs in terms of average PI width, but also accurately estimates conditional coverage information that is useful in assessing model uncertainty.

Keywords: Uncertainty quantification, prediction intervals, conditional coverage, neural networks, calibration error

1 Introduction

Prediction interval (PI) is poised to play an increasingly prominent role in uncertainty quantification for regression tasks (Khosravi et al. , 2010, 2011; Galván et al. , 2017; Rosenfeld et al. , 2018; Tagasovska and Lopez-Paz , 2018, 2019; Romano et al. , 2019; Wang et al. , 2019; Kivaranovic et al. , 2020). Recently, it has also been used in simulation metamodeling Lam and Zhang (2021, 2022) in the spirit of the widely used stochastic kriging and Gaussian process methodologies Ankenman et al. (2010); Barton and Meckesheimer (2006); Staum (2009). A *high-quality* PI should be as narrow as possible,

whilst maintaining a pre-specified level of data coverage or marginal coverage (Pearce et al. , 2018; Zhu et al. , 2019). Compared with PIs obtained based on coverage-only consideration, the “high-quality” criterion is beneficial in balancing between marginal coverage probability and interval width. However, the conditional coverage given a feature, which is critical for making reliable context-based decisions, is unassessed and missing in most existing works on high-quality PIs. In the presence of heteroskedasticity and model misspecification, the marginal coverage can be very different from the conditional coverage at a given point, which affects the downstream decision-making task that relies on the

uncertainty information provided by the PI. Our main goal is to meaningfully incorporate and assess conditional coverages in high-quality PIs.

Conditional coverage estimation is challenging for two reasons. First is that the natural evaluation metric of *conditional coverage error*, an L^p distance between the estimated and ground-truth conditional coverages, is difficult to compute as it requires obtaining the conditional probability given feature x , which is arguably as challenging as the regression problem itself. Our first goal in this paper is to address this issue by developing a new metric called *calibration-based conditional coverage error* for conditional coverage estimation measurement. Our approach is inspired from the calibration notion in classification (Guo et al., 2017). The basic idea is to relax conditional coverage at any given point to being averaged over all points that bear the same estimated value. An estimator satisfying the relaxed property is regarded as well-calibrated. In regression, calibration-based conditional coverage error provides a middle ground between the enforcement of marginal coverage (lacking any conditional information) and conditional coverage (computationally intractable). Compared with conditional coverage, this middle-ground metric can be viewed as a “dimension reduction” of the conditioning variable from the original sample space to the space $[0, 1]$, so that we can easily discretize to compute the empirical metric values.

The second challenge is the discontinuity in the above metrics that hinders efficient training of PIs that are both high-quality and possess reliable conditional coverage information. To address this, we design a new loss function based on

a combination of the high-quality criterion and a *coverage assessment loss*. The latter can be flexibly added as a separate module to any neural network (NN) used to train PIs. It is based on an empirical version of a tight upper bound on the coverage error in terms of a Kullback–Leibler (KL) divergence, which can be readily employed for running gradient descent. We theoretically show how training with our proposed loss function attains this upper-bounding value via a concentration bound. We also demonstrate the empirical performance of our approach in terms of PI quality and conditional coverage assessment compared with benchmark methods.

We summarize our contributions as follows:

1. We identify the conditional coverage estimation problem as a new challenge for high-quality PIs and introduce a new evaluation metric for coverage estimation.
2. We propose an end-to-end algorithm that can simultaneously construct high-quality PIs and generate conditional coverage estimates. In addition, we provide theoretical justifications on the effectiveness of our algorithm by developing concentration bounds relating the coverage assessment loss and conditional coverage error.
3. By evaluating on benchmark datasets and synthetic examples, we empirically demonstrate that our approach achieves high performance on conditional coverage estimation without penalizing the performance on high-quality PI generation.

2 Evaluating Conditional Coverage for High-Quality PIs

Let $X \in \mathcal{X}$ and $Y \in \mathcal{Y} \subset \mathbb{R}$ be random variables denoting the input feature and label, where the pair (X, Y) follows an (unknown) ground-truth joint distribution $\pi(X, Y)$. Let $\pi(Y|X)$ be the conditional distribution of Y given X . Let $\pi(X)$ be the marginal distribution of X . We are given the training data $\mathcal{D} := \{(x_i, y_i), i = 1, 2, \dots, n\}$ where (x_i, y_i) are i.i.d. realizations of random variables (X, Y) . Here, (x_i, y_i) could denote supervised real data. It can also denote, in the simulation context, a simulation model output y_i given the design point x_i .

A PI refers to an interval $[L(x), U(x)]$ where L, U are two functions mapping from \mathcal{X} to \mathcal{Y} trained on the data \mathcal{D} . $[L(x), U(x)]$ is called a PI at *prediction level* $1 - \alpha$ ($0 \leq \alpha \leq 1$) if its *marginal coverage* is not less than $1 - \alpha$, i.e., $\mathbb{P}[Y \in [L(X), U(X)] | L, U] \geq 1 - \alpha$ where \mathbb{P} is with respect to a new test point $(X, Y) \sim \pi(X, Y)$.

We say that $[L(x), U(x)]$ is of high-quality if its marginal coverage attains a pre-specified target prediction level and has a short width on average. In particular, a best-quality PI at prediction level $1 - \alpha$ is an optimal solution to the following constrained optimization problem:

$$\begin{aligned} & \min_{L, U} \mathbb{E}[U(X) - L(X)], \\ & \text{subject to } \mathbb{P}[Y \in [L(X), U(X)] | L, U] \geq 1 - \alpha. \end{aligned} \quad (1)$$

The high-quality criterion has been widely adopted in previous work [Khosravi et al. \(2010, 2011\)](#); [Galván et al. \(2017\)](#); [Pearce et al. \(2018\)](#); [Rosenfeld et al. \(2018\)](#); [Zhu et al.](#)

[\(2019\)](#); [Chen et al. \(2021\)](#). However, this criterion alone may fail to carry important model uncertainty information at specific test points. To illustrate this statement, consider a simple example where $x \sim \text{Uniform}[0, 1]$, $y = 0$ for $x \in [0, 0.95]$ and $y|x \sim \text{Uniform}[0, 1]$ for $x \in (0.95, 1]$. Then according to Equation 1, a best-quality 95% PI is precisely $L(x) = U(x) = 0$ for all $x \in [0, 1]$. This PI has nonconstant “conditional” coverage if we condition at different input points (1 for $x \in [0, 0.95]$ and 0 for $x \in (0.95, 1]$), and can deviate significantly from the overall marginal coverage 95%. Therefore, it is sometimes important to obtain conditional coverage information, especially when the conditional coverage at some input points differs a lot from the marginal coverage. This simple example highlights the need to obtain conditional coverage information.

To mitigate the drawback of the high-quality criterion, we define:

Definition 1 (Conditional Coverage and Its Estimator). *The (ground-truth) conditional coverage associated with a PI $[L(x), U(x)]$ is defined as $A(x) := \mathbb{P}[Y \in [L(X), U(X)] | L, U, X = x]$ for any $x \in \mathcal{X}$, where \mathbb{P} is taken with respect to $\pi(Y|X = x)$. For a (conditional) coverage estimator \hat{P} , which is a measurable function from \mathcal{X} to $[0, 1]$, we define its L^p conditional coverage error (\widetilde{CE}_p) as*

$$\widetilde{CE}_p := \left\| A(X) - \hat{P}(X) \right\|_{L^p(\mathcal{X})}$$

where the L^p -norm is taken with respect to the randomness of X ($1 \leq p \leq +\infty$).

Note that evaluating \widetilde{CE}_p relies on approximating the conditional coverage $A(x)$, which can

be as challenging as the original prediction problem. To address this, we leverage the similarity of estimating $A(x)$ to generating prediction probabilities in binary classification, which motivates us to borrow the notion of calibration in classification. This idea is based on a relaxed error criterion by looking at the conditional coverage among all points that bear the same coverage estimator value, instead of conditioning at any given point. The resulting error metric then only relies on probabilities conditioned on variables in a much lower-dimensional space $[0, 1]$ than \mathcal{X} . To explain concretely, we introduce a “perfect-calibrated coverage estimator” as:

Definition 2 (Perfect Calibration). *A coverage estimator \hat{P} is called a perfect-calibrated coverage estimator associated with $[L(x), U(x)]$ if it satisfies*

$$\hat{P}(x) = \mathbb{P}[Y \in [L(X), U(X)] | L, U, \hat{P}(X) = \hat{P}(x)], \quad (2)$$

for a.e. $\hat{P}(x) \in [0, 1]$ where a.e. is with respect to the probability measure on $[0, 1]$ induced by the random variable $\hat{P}(X)$.

Equation 2 means that the “average” coverage of the PI restricted on the subset $\{x \in \mathcal{X} : \hat{P}(x) = p\}$ should be precisely p (where the “average” coverage is the average of all the conditional coverage $A(x)$ on the subset $\{x \in \mathcal{X} : \hat{P}(x) = p\}$, with respect to the distribution $\pi(X)$ conditional on the subset $\{x \in \mathcal{X} : \hat{P}(x) = p\}$). Corresponding to Definition 2, we define:

Definition 3 (Calibration-based Error). *An L^p ($1 \leq p \leq +\infty$) calibration-based conditional coverage error, or coverage error for short (CE_p),*

of a coverage estimator \hat{P} is:

$$CE_p := \left\| \mathbb{P}[Y \in [L(X), U(X)] | L, U, \hat{P}(X)] - \hat{P}(X) \right\|_{L^p(\mathcal{X})} \quad (3)$$

where the L^p -norm is taken with respect to the randomness of $\hat{P}(X)$.

In the above definition the conditional probability $\mathbb{P}[Y \in [L(X), U(X)] | L, U, \hat{P}(X)]$ is a measurable function of random variable $\hat{P}(X)$, say $\gamma(\hat{P}(X))$. By a change of variable,

$$CE_p^p := \left\| \gamma(\hat{P}(X)) - \hat{P}(X) \right\|_{L^p(\mathcal{X})}^p \quad (4)$$

$$= \int_0^1 |\gamma(t) - t|^p dF_{\hat{P}(X)}(t) \quad (5)$$

where $F_{\hat{P}(X)}$ is a probability distribution of $\hat{P}(X)$ on $[0, 1]$. Here, CE_p only requires estimating $\gamma(t)$ for $t \in [0, 1]$, which can be done easily by discretizing $[0, 1]$ for empirical calculation. To do this, we first construct a discrete version of (2) and then introduce an empirical counterpart of CE_p (3), which we refer to as L^p empirical calibration-based conditional coverage error ECE_p . The ideas behind these are natural extensions of the classification case (Guo et al. , 2017; Kull , 2019; Kumar , 2019; Nixon et al. , 2019) into PIs.

We consider the following partition Δ of $[0, 1]$. Let $[0, 1]$ be divided into M intervals $I_m = (a_{m-1}, a_m]$ ($m = 1, \dots, M$) where $0 = a_0 \leq a_1 \leq \dots \leq a_M = 1$. Let $B_m = \{i = 1, \dots, n : \hat{P}(x_i) \in I_m\}$, i.e., the set (bin) of indices i of samples whose coverage estimator $\hat{P}(x_i)$ falls into the interval I_m . Note that coverage estimations that are close to each other will fall into the same interval. The coverage probability (i.e., the proportion of successful coverage)

in B_m is defined as:

$$CP(B_m) = \frac{1}{|B_m|} \sum_{i \in B_m} \mathbf{1}_{y_i \in [L(x_i), U(x_i)]}. \quad (6)$$

The average coverage estimation in B_m is defined as:

$$AC(B_m) = \frac{1}{|B_m|} \sum_{i \in B_m} \hat{P}(x_i) \quad (7)$$

where $\hat{P}(x_i)$ is the coverage estimator for sample x_i . $CP(B_m)$ and $AC(B_m)$ approximate the left and right hand sides of (2) respectively in the interval I_m . A perfect-calibrated coverage estimator should satisfy $CP(B_m) = AC(B_m)$ for all $m \in \{1, \dots, M\}$. The diagram of $CP(B_m)$ versus $AC(B_m)$ for all $m \in \{1, \dots, M\}$ is called the reliability diagram in some literature (Guo et al., 2017).

Based on the partition Δ , we can introduce an empirical version of CE_p (which we refer to as ECE_p) as:

Definition 4. The L^p empirical calibration-based conditional coverage error (ECE_p) of a coverage estimator \hat{P} is defined as

$$ECE_p = \left\| (CP(B_{m(i)}) - AC(B_{m(i)}))_{i=1,2,\dots,n} \right\|_{l^p}. \quad (8)$$

where the l^p is the standard p -norm in \mathbb{R}^n and $B_{m(i)}$ is the bin containing sample i .

Equivalently,

$$ECE_p = \left(\sum_{m=1}^M \frac{|B_m|}{n} |CP(B_m) - AC(B_m)|^p \right)^{\frac{1}{p}}$$

for $1 \leq p < +\infty$, and

$$ECE_\infty = \max_{m=1,2,\dots,M} |CP(B_m) - AC(B_m)|.$$

A calibration-based error CE_p provides a middle ground between the enforcement of

marginal coverage and conditional coverage. The ground-truth conditional coverage is perfectly calibrated, but not vice versa. However, if we enforce the perfect calibration criterion for a coverage estimator to hold when restricted to any positive-probability measurable subset in \mathcal{X} , then the choice of the estimator will reduce uniquely to the conditional coverage. More formally,

Definition 5. A coverage estimator \hat{P} is called a perfect-calibrated coverage estimator on a measurable subset $\mathcal{S} \subset \mathcal{X}$ with $\mathbb{P}(\mathcal{S}) > 0$ associated with $[L(x), U(x)]$ if it satisfies

$$\begin{aligned} \hat{P}(x) = \\ \mathbb{P}[Y \in [L(X), U(X)] | L, U, \hat{P}(X) = \hat{P}(x), X \in \mathcal{S}] \end{aligned} \quad (9)$$

for a.e. $\hat{P}(x) \in [0, 1]$ where a.e. is with respect to the probability measure on $[0, 1]$ induced by the random variable $\hat{P}(X|_{\mathcal{S}})$. Note that the conditional probability space is standard: $(\mathcal{S}, \mathcal{F}_{\mathcal{S}} := \{A \cap \mathcal{S} : A \in \mathcal{F}\}, \mathbb{P}_{\mathcal{S}}(A \cap \mathcal{S}) := \mathbb{P}(A|_{\mathcal{S}}))$.

Lemma 1. (a) A coverage estimator is the conditional coverage if and only if it is a perfect-calibrated coverage estimator on any positive-probability measurable subset \mathcal{S} of \mathcal{X} .

(b) Suppose \hat{P} is a perfect-calibrated coverage estimator on two disjoint positive-probability measurable subsets $\mathcal{S}_1, \mathcal{S}_2$. Then \hat{P} is a perfect-calibrated coverage estimator on $\mathcal{S}_1 \cup \mathcal{S}_2$.

More details can be found in Appendix A. In the sense of Lemma 1(a), CE_p is an error metric that is a natural relaxation of \widetilde{CE}_p , and although less precise, CE_p is computationally much more tractable than \widetilde{CE}_p .

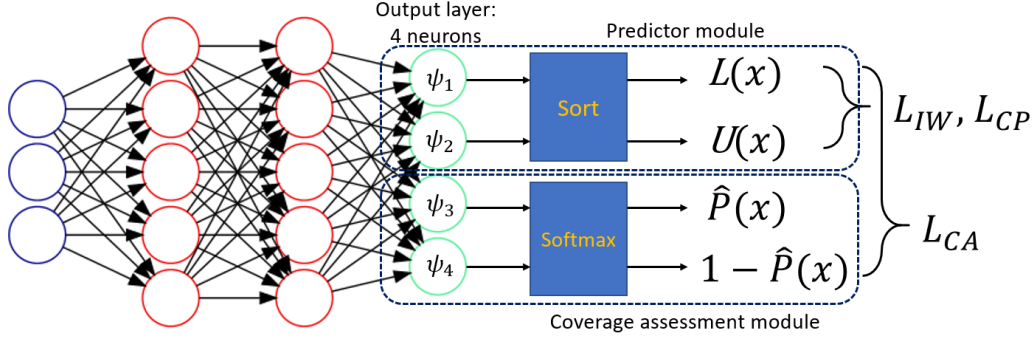


Figure 1: The framework of our proposed coverage assessment network (CaNet).

We will use ECE_1 as the primary evaluation metric to measure the quality of a coverage estimator. A high ECE_1 value of a coverage estimator indicates an unreliable coverage estimation while a small ECE_1 value indicates that the coverage estimator is close to the perfect-calibrated property. Ideally, an effective algorithm should output a coverage estimator with a small ECE_1 value.

3 Neural Network for Conditional Coverage

We propose a novel end-to-end algorithm, named coverage assessment network (CaNet), to simultaneously generate a coverage estimator along with the high-quality PI. As illustrated in Figure 1, our CaNet takes any input data x as its input layer (the leftmost layer in Figure 1), feeds x in the forward direction through the network, and obtains four output values ($\psi_1(x), \psi_2(x), \psi_3(x), \psi_4(x)$) in the output layer. Then these output values are used in two major modules in our CaNet: (1) **predictor module** and (2) **coverage assessment module** (Ca-

Module). The PI predictor module sorts the two outputs ($\psi_1(x), \psi_2(x)$) to provide the upper and lower bound of the estimated PIs. Meanwhile, the Ca-Module is added to the output layer to access the conditional coverage information of PIs from the predictor module. Our model is jointly optimized by three loss functions: coverage assessment loss L_{CA} , intervals width loss L_{IW} , and coverage probability loss L_{CP} . Benefiting from these powerful modules, the CaNet can generate and validate the coverage estimator at the same time without any requirement for further post-processing steps. In addition, we propose a novel deep ensemble technique, suitably designed for our setting, that can improve the performance of both the predictor module and the coverage assessment module.

3.1 Coverage Assessment Module

Our Ca-Module consists of two neurons ($\psi_3(x), \psi_4(x)$) fully connected to the last hidden layer to estimate the conditional coverage. After letting ($\psi_3(x), \psi_4(x)$) pass through the softmax activation function, it outputs a two-point probability distribution ($\hat{P}(x), 1 - \hat{P}(x)$)

where $\hat{P}(x)$ is viewed as the coverage estimator of the PI from the predictor module. Our Ca-Module can be easily integrated into the output layer of deep neural networks to estimate their conditional coverage.

3.2 Loss Function Design and Tuning Procedure

Our loss function is a sum of the *predictor loss* and the *coverage loss*. The predictor loss aims to narrow the prediction intervals as much as possible, while maintaining a specified marginal coverage of data. Inspired by Khosravi et al. (2010, 2011); Pearce et al. (2018); Rosenfeld et al. (2018), our predictor loss is formed by the sum of interval width (IW) loss L_{IW} and coverage probability (CP) loss L_{CP} :

$$L_{IW} = \frac{1}{n} \sum_{i=1}^n (U(x_i) - L(x_i)), \quad (10)$$

$$L_{CP} = \frac{1}{n} \sum_{i=1}^n \tilde{k}_i, \quad CP = \frac{1}{n} \sum_{i=1}^n k_i, \quad (11)$$

where k_i indicates whether each data point has been captured by the PIs: $k_i = 1$ if $L(x_i) \leq y_i \leq U(x_i)$ and $k_i = 0$ otherwise. \tilde{k}_i is a soft version of k_i , which is defined as: $\tilde{k}_i := \sigma(\lambda_3(U(x_i) - y_i)) \cdot \sigma(\lambda_3(y_i - L(x_i)))$, where $\lambda_3 \geq 0$ is a tunable parameter and $\sigma(t) := \frac{1}{1+e^{-t}}$ is the sigmoid function. Therefore, L_{CP} is a soft version of CP that can be used for gradient descent. Associated with the Ca-Module, we introduce a coverage assessment loss L_{CA} to estimate the conditional coverage:

$$L_{CA} = -\frac{1}{n} \sum_{i=1}^n \left(k_i \log(\hat{P}(x_i)) + (1 - k_i) \log(1 - \hat{P}(x_i)) \right).$$

We will show in Section 4 that the expectation of coverage assessment loss L_{CA} provides an upper bound for both the conditional coverage error (Definition 1) and calibration-based conditional coverage error (Definition 3). Hence, minimizing L_{CA} contributes to the recovery of the conditional coverage. In order to run gradient-based methods, we replace the discrete indicator $(k_i, 1 - k_i)$ in L_{CA} with its soft version $(\tilde{k}_i, 1 - \tilde{k}_i)$:

$$\tilde{L}_{CA} = -\frac{1}{n} \sum_{i=1}^n \left(\tilde{k}_i \log(\hat{P}(x_i)) + (1 - \tilde{k}_i) \log(1 - \hat{P}(x_i)) \right).$$

Our total loss function for the CaNet is defined as:

$$\text{Total Loss} = L_{IW} + \lambda_1(1 - L_{CP}) + \lambda_2 \tilde{L}_{CA} \quad (12)$$

where $\lambda_1 \geq 0$, $\lambda_2 \geq 0$ are tunable parameters. We propose an easy-to-implement yet effective tuning procedure, Algorithm 1, to pick up these parameters.

3.3 Deep Ensembles

Following previous research (Lee et al. , 2015; Lakshminarayanan , 2017; Pearce et al. , 2018; Fort et al. , 2019; Ovadia et al. , 2019; Gustafsson et al. , 2020; Pearce et al. , 2020), we apply the deep ensemble technique to provide more robust and better results. During the training period, with the same hyperparameters $\lambda_i, i = 1, 2, 3$, m networks are trained with different initializations. The prediction results from i -th network are denoted as: $([L_i(x), U_i(x)], \hat{P}_i(x), 1 -$

Algorithm 1 Tuning algorithm

Goal: Tune hyperparameters λ_1 , λ_2 , and λ_3 ;**Input:** Prediction level $1 - \alpha$, training dataset \mathcal{D} , validation dataset \mathcal{D}' ;**Procedure:** Let the coverage probability $CP_{\mathcal{D}'}$ and the average coverage estimation $AC_{\mathcal{D}'}$ on \mathcal{D}' be

$$CP_{\mathcal{D}'} = \frac{1}{|\mathcal{D}'|} \sum_{i \in \mathcal{D}'} \mathbf{1}_{y_i \in [L(x_i), U(x_i)]},$$

$$AC_{\mathcal{D}'} = \frac{1}{|\mathcal{D}'|} \sum_{i \in \mathcal{D}'} \hat{P}(x_i).$$

where $([L(x), U(x)], \hat{P}(x))$ are prediction results from the deep ensemble.

- (1) Initialize λ_i ($i = 1, 2, 3$) so that $CP_{\mathcal{D}'}$ is nontrivial, i.e., not (almost) 0 or 1.
- (2) While $CP_{\mathcal{D}'}$ is nontrivial: tune λ_2 and λ_3 so that $|CP_{\mathcal{D}'} - AC_{\mathcal{D}'}| \leq \epsilon$ (e.g., $\epsilon = 1\%$.)
- (3) Otherwise tune λ_1 such that $CP_{\mathcal{D}'}$ is nontrivial. Do step 2 again until we find λ_2 and λ_3 .
- (4) Tune λ_1 such that $CP_{\mathcal{D}'} > 1 - \alpha$ where λ_2 and λ_3 are fixed from (3).

Output: λ_1 , λ_2 , and λ_3 .

 $\hat{P}_i(x)$). Finally, the output from CaNet is:

$$\begin{aligned} \text{Lower bound} \quad \bar{L} &:= \sum_{i=1}^m \frac{1}{m} L_i, \\ \text{Upper bound} \quad \bar{U} &:= \sum_{i=1}^m \frac{1}{m} U_i, \\ \text{Coverage estimator} \quad \bar{\hat{P}} &:= \sum_{i=1}^m \frac{1}{m} \hat{P}_i. \end{aligned} \tag{13}$$

4 Theoretical Analysis

In this section, we provide theoretical insights about the coverage assessment loss L_{CA} in our approach. We show that minimizing L_{CA} is equivalent to minimizing a tight upper bound of the conditional coverage error with high probability and thus contributes to reducing the conditional coverage error of a coverage estimator. This theoretical analysis consists of two steps. First, we show that both CE_p and \widetilde{CE}_p are bounded above by the expectation of a Kullback-Leibler divergence-type risk $K_1(X)$. Second, we establish the tail bound of the generalization error $\mathbb{E}[K_1(X)]$ and its empirical coun-

terpart which is exactly L_{CA} . We will show that the generalization bound holds uniformly for any tuple of the joint output functions $(L(x), U(x), \hat{P}(x), 1 - \hat{P}(x))$ generated by our CaNet. Because multiple hypothesis classes beyond $(L(x), U(x), \hat{P}(x), 1 - \hat{P}(x))$ are involved in L_{CA} , more effort is required to build the connection among them to derive the generalization bound.

Theorem 1. *Let $A(x) := \mathbb{P}[Y \in [L(X), U(X)] | L, U, X = x]$ be the conditional coverage in Definition 1. Let $K(x) := A(x) \log \left(\frac{A(x)}{\bar{\hat{P}}(x)} \right) + (1 - A(x)) \log \left(\frac{1 - A(x)}{1 - \bar{\hat{P}}(x)} \right)$. Then*

$$CE_p \leq \widetilde{CE}_p \leq \left(\frac{1}{2} \mathbb{E}[K(X)] \right)^{\alpha_p/2}, \forall 1 \leq p \leq +\infty$$

where $\alpha_p = 1, \forall 1 \leq p \leq 2$ and $\alpha_p = \frac{2}{p}, \forall 2 \leq p \leq +\infty$. Moreover, the inequality is attainable if, e.g., $\hat{P}(x)$ equals the conditional coverage $A(x)$.

From Theorem 1, we see that minimizing $\mathbb{E}[K(x)]$ is equivalent to minimizing a tight upper

bound for the coverage error. For every x , $K(x)$ is the Kullback–Leibler divergence between the distributions represented by $(\hat{P}(x), 1 - \hat{P}(x))$ and $(A(x), 1 - A(x))$. $K(x) = K_0(x) + K_1(x)$, where $K_0(x) = A(x) \log(A(x)) + (1 - A(x)) \log(1 - A(x))$ and $K_1(x) = -A(x) \log(\hat{P}(x)) - (1 - A(x)) \log(1 - \hat{P}(x))$. Minimizing $\mathbb{E}[K(X)]$ over \hat{P} is equivalent to minimizing $\mathbb{E}[K_1(X)]$. The type of results in Theorem 1 that bounds an L^p conditional coverage error via a Kullback–Leibler-type error is new as far as we know. Next, to show L_{CA} approximates $\mathbb{E}[K_1(X)]$, we need the following assumptions:

Assumption 1. *The four classes of functions $([L(x), U(x)], \hat{P}(x), 1 - \hat{P}(x))$ output by the neural network (NN) in Figure 1 have finite VC dimensions, say they are bounded above by V_0 .*

Assumption 1 holds for a wide range of NNs (e.g., Theorem 8.14 in Anthony and Bartlett (1999), Theorem 7 in Bartlett et al. (2019)). In particular, it holds for the one we adopt in the experiments (where we use the ReLU-activated NN to construct ψ_i , $i = 1, 2, 3, 4$; see Section 5):

Theorem 2. *Suppose $\psi_i, i = 1, 2, 3, 4$ are the pre-activated output neurons of the NN in Figure 1 using the ReLU activation function. Then Assumption 1 holds. Moreover, suppose the NN has W parameters and U computation units (nodes). Then $V_0 = O(WU)$.*

Assumption 2. *There exists a constant $M < +\infty$ such that $|\log(\hat{P}(x))| \leq M$, $|\log(1 - \hat{P}(x))| \leq M$ for all x and \hat{P} .*

This is a natural assumption in practice because $\log(\hat{P}(x))$ and $\log(1 - \hat{P}(x))$ are replaced

by $\log(\hat{P}(x) + \epsilon)$ and $\log(1 - \hat{P}(x) + \epsilon)$ respectively to avoid explosion when implementing the algorithm. In particular, in our experiments in Section 5, $\epsilon = 0.1^6$ and thus $M = 14$. Let

$$\mathcal{F} = \{f(x, y) = I_{y \in [L(x), U(x)]} : \\ L, U \text{ are output by the NN}\},$$

$$\mathcal{G} = \{\hat{P}(x) : \hat{P} \text{ is output by the NN}\}.$$

Theorem 3. *Suppose Assumptions 1 and 2 hold. The training data $\mathcal{D} = \{(x_i, y_i), i = 1, 2, \dots, n\}$ where (x_i, y_i) are i.i.d. samples $\sim \pi$. Recall that the (hard) coverage estimator assessment loss is $L_{CA} = -\frac{1}{n} \sum_{i=1}^n \left(f(x_i, y_i) \log(\hat{P}(x_i)) + (1 - f(x_i, y_i)) \log(1 - \hat{P}(x_i)) \right)$. Then for any $t > 0$, we have*

$$\mathbb{P} \left(\sup_{f \in \mathcal{F}, \hat{P} \in \mathcal{G}} |L_{CA} - \mathbb{E}[K_1(X)]| \geq t \right) \leq C^* e^{-\frac{nt^2}{16M^2}}.$$

where C^* only depends on V_0 in Assumption 1.

Theorem 3 shows that the coverage assessment loss approximates $\mathbb{E}[K_1(x)]$ well with an exponential tail bound. The difficulty in analyzing Theorem 3 lies in the fact that the hypothesis classes in Assumption 1 (which are constructed by the NN) are different from the hypothesis class used in L_{CA} . To overcome this difficulty, we use the theory of VC-subgraph classes to connect the VC dimension among multiple hypothesis classes, including the class of ψ_i , the four classes of output functions, and $\mathcal{F}, \log \mathcal{G}$. Then we establish the covering number bound for the class \mathcal{F} and $\log \mathcal{G}$, and finally prove Theorem 3. To conclude, minimizing $\mathbb{E}[K_1(X)]$ over \hat{P} is equivalent to minimizing $\mathbb{E}[K(X)]$, which in turn is minimizing a tight upper bound for the coverage assessment loss. Our coverage assessment

loss empirically approximates $\mathbb{E}[K_1(X)]$ well, so that its minimization can ultimately help to reduce the conditional coverage error.

5 Experiments

Experimental Setup. We empirically verify the effectiveness of our proposed CaNet on both synthetic examples and benchmark regression datasets. These datasets have been widely used for the evaluation of methods in regression tasks (Hernández-Lobato and Adams, 2015; Gal and Ghahramani, 2016; Lakshminarayanan, 2017; Rosenfeld et al., 2018; Pearce et al., 2018; Zhu et al., 2019). In addition, we adopt the same experimental procedure in Pearce et al. (2018) for data normalization and dataset splitting. To avoid overfitting, we apply a simple network architecture with only 2 hidden layers and each hidden layer has 64 neurons. For each hidden layer, the ReLU activation function is applied to capture the non-linear features. We empirically set the ensemble number m to 5, as the smallest number leading to a stable prediction results. Please refer to Appendix D for implementation details, including those for baseline algorithms.

Evaluation Metrics. To evaluate the conditional coverage estimation of our CaNet, we examine the quality of our Ca-Module measured by the empirical coverage error ECE_1 over a partition Δ . Δ is constructed by equally dividing the width of $[0, 1]$ into M sub-intervals. The value of M depends on the size of the dataset, which is determined by the following strategy: $M = \min\{\lfloor \text{the number of data in validation}/50 \rfloor + 1, 20\}$. In addition, we also report the IW and CP , where $IW = \frac{1}{n} \sum_{i=1}^n (U(x_i) - L(x_i))$ (Equa-

tion 10) and $CP = \frac{1}{n} \sum_{i=1}^n k_i$, (Equation 11) to show the effectiveness of our predictor module under the high-quality criteria. Thanks to the deep ensemble technique (Section 3.3), our experimental results are robust and stable among multiple experimental repetitions with a standard deviation typically less than 2%. Therefore, we only report the mean value.

5.1 Conditional Coverage on Synthetic Examples

In this section, we conduct a series of experiments on synthetic examples to directly compare our prediction results with the ground-truth conditional coverage. In these examples, the conditional coverage can be analytically calculated under the known data distribution. Figure 2 compares the conditional coverage with our predicted coverage under the following settings: $x \sim \text{Uniform}[-2, 2]$ and $y|x$ is drawn from $f_i(x) = \frac{1}{3} \sin(x) + \varepsilon_i(x)$, $x \in [-2, 2]$ where

$$\varepsilon_1(x) = 0.1 \times N(0, 1),$$

$$\varepsilon_2(x) = 0.1|x| \times N(0, 1),$$

$$\varepsilon_3(x) = 0.1|x| \times t_4.$$

$N(0, 1)$ is the standard Gaussian variable and t_4 is the standard t random variable with 4 degrees of freedom. Then, the conditional coverage in Definition 1 can be analytically calculated as:

$$\begin{aligned} & \mathbb{P}[Y \in [L(X), U(X)] | L, U, X = x] \\ &= F_i(U(x) - \frac{1}{3} \sin(x)) - F_i(L(x) - \frac{1}{3} \sin(x)) \end{aligned}$$

where F_i is the cumulative distribution function of $N(0, 0.1^2)$ for $i = 1$, $N(0, (0.1x)^2)$ for $i = 2$, and $0.1|x| \times t_4$ for $i = 3$. As shown in Figure 2 (b,c), the ground-truth conditional coverages of high-quality PIs diverge among different points,

and they deviate from the marginal coverage. Thus, having access to the marginal coverage for the whole dataset is not sufficient for decision making, which highlights the need for conditional coverage. In addition, the conditional coverage estimator from our model is highly consistent with the conditional coverage on all of the synthetic examples in Figure 2. These results confirm that our CaNet can accurately estimate the conditional coverage on noisy datasets.

5.2 Performance of PIs on Benchmark Datasets

In this section, we compare the performance of PIs generated by the predictor module of our CaNet on real-world benchmark datasets with following baseline algorithms: (1) nearest-neighbors kernel conditional density estimation (NNKCDE) (Dalmasso et al., 2020), (2) quantile regression forest (QRF) (Meinshausen, 2006), (3) split conformal learning (SCL) (Lei et al., 2018) and (4) the quality-driven PI method (QD-Ens) (Pearce et al., 2018). We quote the results from Pearce et al. (2018) as a comparison since we share the same experiment setup. Table 1 reports the results of CP , IW and ECE_1 for generating PIs at 95% prediction level on benchmark datasets. We employed the criteria in Pearce et al. (2018) to evaluate the performance of PIs: the best IW is achieved by the model with the smallest IW value among those with $CP \geq 95\%$. As can be seen, our model achieves competitive results on PI generation under the high-quality criteria. With special consideration on interval width quality, it obtains the smallest average interval width (IW) while maintain-

ing high coverage probability ($CP \geq 95\%$) on all datasets. These results show that adding an additional module CA-Module does not penalize the performance of PIs generated by our CaNet.

5.3 Performance of Coverage Estimator on Benchmark Datasets

Coverage for 95% PIs. We use ECE_1 to evaluate the coverage estimation performance on real-world datasets as the conditional coverage is unknown. As shown in Table 1, the ECE_1 on all experiments are generally around or less than 1%, with better performance on larger datasets. The coverage estimators produced by CaNet have small ECE_1 values, which are very close to the perfect-calibrated coverage estimators (Definition 2). Compared with ECE_1 values obtained from the state-of-the-art algorithms in classification tasks (Guo et al., 2017; Kull, 2019), the ECE_1 values from CaNet are similar and sometimes less than their post-calibrated ECE_1 results (usually around 1% to 3%), even though the size of most regression datasets are smaller than the classification datasets. These results further demonstrate that our CaNet can accurately estimate the coverage information of 95% high-quality PIs on real-world regression tasks.

Coverage for PIs at Different Prediction Levels. We conduct multiple experiments on different PI prediction levels to show the robustness of our CaNet. By only modifying the parameter λ_1 in Equation 12, our Ca-Module could get access to different levels of coverage probability. Table 2 reports the CP , IW and ECE_1 values from the CaNet at different PI pre-

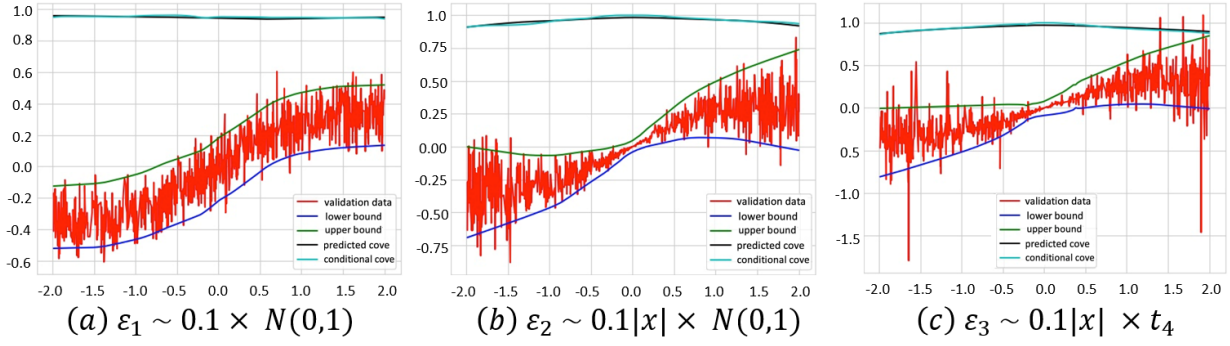


Figure 2: Prediction and conditional coverage of 95% PIs on synthetic examples. (a) $CP = 0.95$, $IW = 0.40$, $ECE_1 = 0.62\%$. (b) $CP = 0.96$, $IW = 0.40$, $ECE_1 = 0.12\%$. (c) $CP = 0.95$, $IW = 0.50$, $ECE_1 = 0.65\%$. The predicted coverage estimation from CaNet is highly consistent with the conditional coverage under different noise settings. The red curve represents the validation data. The dark blue curve represents the lower bound of the PI generated by our CaNet. The green curve represents the upper bound of the PI generated by our CaNet. The black curve represents the conditional coverage estimator generated by our CaNet. The light blue curve represents the ground-truth conditional coverage. In all figures, the conditional coverage estimator from our model (the black curve) is highly consistent with the ground-truth conditional coverage (the light blue curve). In Figure 2 (b,c), the ground-truth conditional coverages (the light blue curve) diverge among different points, highlighting the need for conditional coverage.

diction levels on three benchmark datasets. Results for more datasets can be found in Appendix D. As can be seen, all ECE_1 values in Table 2 are fairly small ($\sim 1\%$), demonstrating the stability of our proposed model. Thus, our CaNet can provide accurate coverage estimation on PIs at different prediction levels. These results demonstrate the robustness of our CaNet on real-world datasets, further suggesting its broad applicability.

6 Related Work

Finally, we discuss the connections of our paper with several lines of related work.

6.1 Prediction Interval Estimation

High-quality PIs, which can be viewed via a constrained optimization problem where the constraint concerns marginal coverage and the objective is the PI width, has been extensively studied in Khosravi et al. (2010, 2011); Galván et al. (2017); Pearce et al. (2018); Rosenfeld et al. (2018); Zhu et al. (2019); Chen et al. (2021). Such intervals are in the same spirit as the highest density intervals in statistics (Box and Tiao, 2011). While powerful, these approaches could not directly provide the conditional coverage information investigated in this work. Coverage-only criteria, on the other hand, focus solely on coverage satisfaction as the

Dataset	NNKCDE		QRF		SCL		QD-Ens		CaNet: $\lambda_2 = 0.1^5$ for all datasets				
	<i>CP</i>	<i>IW</i>	<i>CP</i>	<i>IW</i>	<i>CP</i>	<i>IW</i>	<i>CP</i>	<i>IW</i>	<i>CP</i>	<i>IW</i>	<i>ECE</i> ₁	λ_1	λ_3
Boston	0.95	1.54	0.96	2.22	0.97	1.45	0.92	1.16	0.95	1.04	1.38%	6.0	1800
Concrete	0.95	1.85	0.99	2.53	0.96	1.54	0.94	1.09	0.95	1.13	0.24%	6.5	1000
Energy	0.97	0.54	0.98	0.87	0.95	0.77	0.97	0.47	0.99	0.37	0.76%	5.0	500
Kin8nm	0.99	2.76	0.99	3.27	0.94	1.20	0.96	1.25	0.95	1.04	1.32%	3.6	300
Plant	0.95	0.88	0.97	1.05	0.95	0.88	0.95	0.86	0.95	0.84	0.38%	3.3	700
Protein	0.93	1.98	0.98	2.42	0.95	2.81	0.95	2.27	0.95	2.26	0.41%	8.3	300
Wine	0.96	2.64	0.93	3.18	0.96	3.44	0.92	2.33	0.95	2.59	0.42%	19	1100
Yacht	0.95	1.15	0.95	1.72	0.95	0.57	0.96	0.17	0.98	0.16	0.80%	1.6	500

Table 1: Evaluation metrics of different models on benchmark datasets. The *CP* values are marked in blue if they meet the 95% prediction level. The best *IW* results, marked in bold, are achieved by models with the smallest *IW* value among those that meet the 95% prediction level. Our model outperforms the baseline algorithms on high-quality PI generation. Meanwhile, it provides accurate coverage estimation on real-world datasets.

guarantee. These approaches include conformal learning (CL) and its conditional variants (Vovk et al. , 2005, 2009; Lei and Wasserman , 2014; Lei et al. , 2015, 2018; Kuchibhotla and Ramdas , 2019; Romano et al. , 2019; Barber et al. , 2019a,b). CL is desirably distribution- or model-free, and in some cases enjoys finite-sample guarantees on coverage. However, the coverage guarantees from CL are only marginal with respect to the training data (except split CL (Vovk , 2012)). Moreover, unlike high-quality PIs, they do not explicitly account for the interval width as a quality metric. Nevertheless, we note that even in some literature with coverage-only consideration, interval width has also been implicitly served as a conservative measurement of PIs, e.g., in Zhang et al. (2019); Barber et al. (2019b). We also mention conditional density estimation (Holmes et al. , 2007; Dutordoir et al. , 2018; Izbicki and Lee , 2016; Dalmaso et al. , 2020; Freeman et al. , 2017; Izbicki et al. , 2017) and closely relatedly quantile regres-

sion (Koenker and Hallock , 2001; Meinshausen , 2006) as PI construction approaches by converting from the estimated conditional quantile function. These approaches focus on the quality of conditional distribution/quantile, instead of the high-quality criterion only. In this work, we follow some recent work using deep learning to achieve state-of-the-art performance for high-quality PIs construction (Khosravi et al. , 2010; Pearce et al. , 2018; Kivaranovic et al. , 2020).

6.2 Uncertainty Measurement in Deep Learning

Standard Bayesian approaches, e.g., Bayesian linear regression (Bishop , 2006), can be applied to construct PIs based on the posterior predictive distribution. But they usually assume that data follow a certain set of parameterized distributions, which provide a different perspective from the distribution-free, frequentist view taken in this paper. On the other hand, as for NNs, the Bayesian framework also offers a

Kin8nm: $\lambda_2 = 0.1^5$, $\lambda_3 = 300$				Plant: $\lambda_2 = 0.1^5$, $\lambda_3 = 700$				Protein: $\lambda_2 = 0.1^5$, $\lambda_3 = 300$			
λ_1	CP	IW	ECE_1	λ_1	CP	IW	ECE_1	λ_1	CP	IW	ECE_1
1.9	0.83	0.73	1.54%	1.4	0.80	0.56	1.23%	4.8	0.82	1.56	0.82%
2.1	0.85	0.77	1.59%	1.5	0.84	0.62	1.36%	5.1	0.86	1.70	0.71%
2.3	0.87	0.80	1.47%	1.6	0.85	0.63	0.86%	5.4	0.88	1.79	1.25%
2.4	0.88	0.84	1.21%	1.8	0.86	0.64	1.36%	5.6	0.90	1.84	1.04%
2.6	0.90	0.87	1.50%	2.0	0.89	0.70	0.51%	6.3	0.91	1.95	0.66%
2.9	0.91	0.89	0.50%	2.4	0.91	0.75	1.15%	6.5	0.92	2.07	1.36%
3.1	0.92	0.94	0.49%	2.6	0.92	0.80	1.29%	7.0	0.93	2.11	1.23%
3.4	0.94	1.00	0.82%	3.0	0.94	0.83	0.23%	7.8	0.94	2.13	0.99%
3.6	0.95	1.04	1.32%	3.3	0.95	0.84	0.38%	8.3	0.95	2.26	0.41%
3.8	0.96	1.06	1.25%	4.0	0.96	0.89	0.80%	9.3	0.96	2.36	0.52%

Table 2: ECE_1 results of our model on benchmark datasets with different coverage probabilities. Our CaNet achieves robust performance on real-world datasets at different prediction levels.

principled approach for model uncertainty measurement by computing the posterior distribution over the NN parameters (MacKay , 1992; Neal , 2012). They focus on the parameter uncertainty of the NN instead of the coverage over a test point. In addition, exact Bayesian inference for deep Bayesian NNs is computationally intractable, making it less practical to implement. Gal and Ghahramani (2016) applied a Monte Carlo dropout method to proxy the inference. Directly generated from non-Bayesian networks, softmax response is also commonly used for uncertainty measurement on deep learning models (Bridle , 1990; Lakshminarayanan , 2017; Geifman et al. , 2018; Sensoy et al. , 2018; Ozbulak et al. , 2018). Moreover, Niculescu-Mizil and Caruana (2005) showed that NNs typically produce well-calibrated probabilities on binary classification tasks without the need for any post-hoc techniques. In this paper, we also follow this line of work and use the softmax output to access the model uncertainty information.

6.3 Simulation Metamodeling

Stochastic simulation aims to estimate a response surface of simulation outputs from complex stochastic models Barton and Meckesheimer (2006); Staum (2009). A common and versatile approach in simulation metamodeling is stochastic kriging (SK) Ankenman et al. (2008, 2010) which can handle general nonlinear input-output relations, and is also closely related to Gaussian process methodologies in Bayesian optimization Frazier (2018); Couckuyt et al. (2022). SK typically focuses on mean response surface estimation or quantile-based response measures Chen and Kim (2013); Bekki et al. (2014); Chen and Kim (2016), and can account for both epistemic and aleatory uncertainties, and as such allows to build PIs for the simulation random outputs. In addition to SK, recent work Lam and Zhang (2021, 2022) proposes to build PIs for simulation metamodeling by adapting techniques from machine learning including using the high-quality criterion and NNs, which is the approach considered in this paper. These studies mainly tar-

get at the marginal coverage of PIs while our methodology in this work aims to obtain conditional coverage information of PIs simultaneously.

7 Concluding Remarks

In this paper, we identify and investigate the conditional coverage estimation problem for high-quality PIs, which is critical for risk-based decision making in regression tasks including simulation metamodeling. To address the challenge, we propose an end-to-end algorithm with two powerful modules: the coverage assessment module and the predictor module. Benefiting from these powerful modules, our model can generate and validate the coverage estimator without any requirement for further post-processing steps. In addition, we conduct a theoretical analysis to show the effectiveness of our proposed model. Experimental results on synthetic examples and benchmark datasets further demonstrate that our model can robustly provide accurate coverage estimation while simultaneously producing a high-quality PI. Another alternative approach for conditional coverage estimation is a two-stage approach that considers the PI problem and the conditional coverage estimation problem separately in different networks. However, the two-stage approach is very complex and contains multiple steps that require two training processes for different networks. Essentially, we have to select more hyper-parameters for the two networks and also train more network weights to get the PI predictor and the conditional coverage estimator, which means more running time. Compared with the two-stage al-

gorithm, our model has several advantages: our model is an end-to-end algorithm by leveraging the common feature layer (the last hidden layer) for both the PI predictor and the conditional coverage estimator, which is very concise and easy to implement. We can also efficiently choose the parameters via our tuning procedure. Moreover, our Ca-Module can be easily integrated into other deep-learning-based algorithms to access their conditional coverage information, opening up more opportunities for broad applications. In the future, we will extend our work by conducting comparison studies with Bayesian methods.

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Appendices

This appendix presents further results and discussions and it consists of four parts. Appendix A gives more detailed properties on coverage estimator and coverage error. Appendix B contains the mathematical argument for Theorem 1. Appendix C discusses how to achieve Assumption 1 and the proof of Theorem 2 and 3. Appendix D illustrates experimental details and more experimental results.

A Further Details on Coverage Estimator and Coverage Error

A.1 Coverage Probability Types of PIs

Zhang et al. (2019) introduce the following four coverage probability types of PIs. In general, most of the coverage in PIs considered in the literature falls into one of these types.

Type I: $\mathbb{P}[Y \in [L(X), U(X)]]$ (marginal coverage);

Type II: $\mathbb{P}[Y \in [L(X), U(X)]|L, U]$ (conditional coverage given the PI);

Type III: $\mathbb{P}[Y \in [L(X), U(X)]|X = x]$ (conditional coverage given $X = x$);

Type IV: $\mathbb{P}[Y \in [L(X), U(X)]|L, U, X = x]$ (conditional coverage given the PI and $X = x$).

Note that in high-quality criteria, only Type II coverage is considered in the constraint but Type IV coverage is lacking. Since Type I and III are not considered in our paper, for simplicity Type II coverage is called the marginal coverage, and Type IV coverage is called the conditional coverage in Definition 1.

Throughout Appendix, we let $A(x) := \mathbb{P}[Y \in [L(X), U(X)]|L, U, X = x]$ denote the conditional coverage in Definition 1. Moreover, since we only concern Type II and IV coverage in this work, we make the following convention. Throughout Appendix A-B, \mathbb{P} and \mathbb{E} should be understood as probability and expectation conditional on L, U . So for $A(x)$, we could simply write $A(x) := \mathbb{P}[Y \in [L(X), U(X)]|X = x]$ and omit “conditional on L, U ”.

A.2 The Terminology “Perfect-Calibrated”

The terminology “perfect-calibrated” is borrowed from the confidence calibration in classification tasks. We first review the name of “confidence” in classification.

Confidence calibration is the problem of predicting probability estimates representative of the true correctness likelihood (Guo et al. , 2017). Intuitively, an reliable confidence should reflect the true correctness likelihood of the prediction (Kumar , 2019). For example, given 100 predictions, each with confidence of 0.8, we expect that 80 should be correctly classified (Guo et al. , 2017).

Now let h be the prediction of any models, which is a map from \mathcal{X} to \mathcal{Y} , trained on the data

\mathcal{D} . According to their definition, the “best” confidence map \hat{P} should be the true probability of correctness:

$$\hat{P}(x) = \mathbb{E}[\mathbf{1}_{h(X) \text{ is correct for } Y} | h, X = x], \quad \forall x \in \mathcal{X}. \quad (14)$$

which is a measurable function from \mathcal{X} to $[0, 1]$. In particular, we consider h as a PI in regression tasks and “correctness” of a PI is naturally defined as the success of coverage on the outcome Y . Then the right hand side of Equation (14) becomes

$$\begin{aligned} \mathbb{E}[\mathbf{1}_{h(X) \text{ is correct for } Y} | h, X = x] &= \mathbb{E}[\mathbf{1}_{Y \in [L(X), U(X)]} | L, U, X = x] \\ &= \mathbb{P}[Y \in [L(X), U(X)] | L, U, X = x]. \end{aligned}$$

which is the conditional coverage in our Definition 1.

In addition, Guo et al. (2017) introduce the perfect-calibrated confidence as follows:

$$\mathbb{P}(\hat{Y} = Y | \hat{Y}, \hat{P} = p) = p, \quad \forall p \in [0, 1]$$

where \hat{Y} is the class prediction and $\hat{Y} = Y$ means that the predicted and true class label coincide. Obviously, if $\hat{P}(X) = \mathbb{P}(\hat{Y}(X) = Y | X)$, i.e., the “best” confidence, then the above equality holds. Transferring this idea into PIs, we can naturally define the perfect-calibrated coverage estimator as

$$p = \mathbb{E}[\mathbf{1}_{h(X) \text{ is correct for } Y} | L, U, \hat{P} = p] = \mathbb{P}[Y \in [L(X), U(X)] | L, U, \hat{P} = p], \quad \forall p \in [0, 1].$$

which is the conditional coverage in our Definition 2.

A.3 Details on Coverage Estimator

A perfect-calibrated coverage estimator inherits some properties of a conditional coverage. For example, both of them have the following interpretation: If we have 1000 testing points for a PI, each with the same conditional/perfect-calibrated coverage 0.9, then approximately 900 of them are correctly covered by the PI. Note that the conditional coverage is uniquely defined, but a perfect-calibrated coverage estimator is not necessarily so. Moreover, we have the following facts: (a) The conditional coverage is always perfect-calibrated, but not vice versa. (b) A perfect-calibrated coverage estimator can be viewed as an averaged conditional coverage. (c) A perfect-calibrated coverage estimator is less “informative” than the conditional coverage.

A perfect-calibrated coverage estimator is an averaged conditional coverage. Let \hat{P}

be a (general) coverage estimator. We have

$$\begin{aligned}
& \mathbb{P}[Y \in [L(X), U(X)] | \hat{P}(X) = \hat{P}(x)] \\
&= \mathbb{P}[Y \in [L(X), U(X)] | X \in \hat{P}^{-1}(\hat{P}(x))] \\
&= \frac{\mathbb{P}[Y \in [L(X), U(X)], X \in \hat{P}^{-1}(\hat{P}(x))]}{\mathbb{P}[X \in \hat{P}^{-1}(\hat{P}(x))]} \\
&= \frac{\int_{t \in \hat{P}^{-1}(\hat{P}(x))} \mathbb{E}[\mathbf{1}_{Y \in [L(X), U(X)]} | X = t] \mathbb{P}[X \in dt]}{\mathbb{P}[X \in \hat{P}^{-1}(\hat{P}(x))]} \\
&= \frac{\int_{t \in \hat{P}^{-1}(\hat{P}(x))} A(t) \mathbb{P}[X \in dt]}{\mathbb{P}[X \in \hat{P}^{-1}(\hat{P}(x))]}
\end{aligned}$$

Suppose \hat{P} is a perfect-calibrated coverage estimator. Then we have

$$\hat{P}(x) = \frac{\int_{t \in \hat{P}^{-1}(\hat{P}(x))} A(t) \mathbb{P}[X \in dt]}{\mathbb{P}[X \in \hat{P}^{-1}(\hat{P}(x))]}$$

which implies that $\hat{P}(x)$ is a weighted average of $A(t)$ over the set $\hat{P}^{-1}(\hat{P}(x))$ with weights based on the marginal distribution of X .

A conditional coverage is perfect-calibrated. If $\hat{P}(x) = A(x)$, then $A(t) = A(x)$ for any $t \in A^{-1}(A(x))$ so

$$\frac{\int_{t \in A^{-1}(A(x))} A(t) \mathbb{P}[X \in dt]}{\mathbb{P}[X \in A^{-1}(A(x))]} = \frac{A(x) \int_{t \in A^{-1}(A(x))} \mathbb{P}[X \in dt]}{\mathbb{P}[X \in A^{-1}(A(x))]} = A(x).$$

This shows that $A(x)$ must be a perfect-calibrated coverage estimator. Another way to see this is taking conditional expectation given $A(X) = p$ in the Definition 1. Then we get

$$\begin{aligned}
p &= \mathbb{E}[A(X) | A(X) = p] \\
&= \mathbb{E}[\mathbb{P}[Y \in [L(X), U(X)] | X] | A(X) = p] \\
&= \mathbb{P}[Y \in [L(X), U(X)] | A(X) = p] \quad \text{by the tower property.}
\end{aligned}$$

A perfect-calibrated coverage estimator may be less informative and may not be the conditional coverage. Suppose we have a PI $[L(X), U(X)]$ at the exact prediction level of $1 - \alpha$, i.e., $\mathbb{P}[Y \in [L(X), U(X)]] = 1 - \alpha$. Then the constant coverage estimator

$$\hat{P}(x) = 1 - \alpha, \quad \forall x \in \mathcal{X}$$

can be viewed as an average coverage estimator over the entire space \mathcal{X} . It is a perfect-calibrated coverage estimator since by definition,

$$\mathbb{P}[Y \in [L(X), U(X)] | \hat{P}(X) = 1 - \alpha] = \mathbb{P}[Y \in [L(X), U(X)]] = 1 - \alpha = \hat{P}(X).$$

But it is not a conditional coverage in general (e.g., the second/third synthetic example in Section 5.1).

In the main paper, we have given the following extension of the definition of “perfect-calibrated” coverage estimator, allowing it to be defined on any positive-probability measurable subsets.

Definition 6 (Restated Definition 5). *A coverage estimator \hat{P} is called a perfect-calibrated coverage estimator on a measurable subset $\mathcal{S} \subset \mathcal{X}$ with $\mathbb{P}(\mathcal{S}) > 0$ associated with $[L(x), U(x)]$ if it satisfies*

$$\hat{P}(x) = \mathbb{P}[Y \in [L(X), U(X)] | L, U, \hat{P}(X) = \hat{P}(x), X \in \mathcal{S}], \quad \text{a.e. } \hat{P}(x) \in [0, 1]. \quad (15)$$

where a.e. is with respect to the probability measure on $[0, 1]$ induced by the random variable $\hat{P}(X|_{\mathcal{S}})$. Note that the conditional probability space is standard: $(\mathcal{S}, \mathcal{F}_{\mathcal{S}} := \{A \cap \mathcal{S} : A \in \mathcal{F}\}, \mathbb{P}_{\mathcal{S}}(A \cap \mathcal{S}) := \mathbb{P}(A|\mathcal{S}))$.

(As our convention in Section A.1, we will omit “conditional on L, U ” for simplicity.)

Lemma 2 (Restated Lemma 1). (a) *A coverage estimator is the conditional coverage if and only if it is a perfect-calibrated coverage estimator on any positive-probability measurable subset \mathcal{S} of \mathcal{X} .* (b) *Suppose \hat{P} is a perfect-calibrated coverage estimator on two disjoint positive-probability measurable subsets $\mathcal{S}_1, \mathcal{S}_2$. Then \hat{P} is a perfect-calibrated coverage estimator on $\mathcal{S}_1 \cup \mathcal{S}_2$.*

Proof. (a) The proof can be found in Lemma 3.

(b) We note that by law of total probability,

$$\begin{aligned} & \mathbb{P}[Y \in [L(X), U(X)] | \hat{P}(X) = p, X \in \mathcal{S}_1 \cup \mathcal{S}_2] \\ &= \mathbb{P}[Y \in [L(X), U(X)] | \hat{P}(X) = p, X \in \mathcal{S}_1] \mathbb{P}[X \in \mathcal{S}_1 | \hat{P}(X) = p, X \in \mathcal{S}_1 \cup \mathcal{S}_2] \\ & \quad + \mathbb{P}[Y \in [L(X), U(X)] | \hat{P}(X) = p, X \in \mathcal{S}_2] \mathbb{P}[X \in \mathcal{S}_2 | \hat{P}(X) = p, X \in \mathcal{S}_1 \cup \mathcal{S}_2] \\ &= p(\mathbb{P}[X \in \mathcal{S}_1 | \hat{P}(X) = p, X \in \mathcal{S}_1 \cup \mathcal{S}_2] + \mathbb{P}[X \in \mathcal{S}_2 | \hat{P}(X) = p, X \in \mathcal{S}_1 \cup \mathcal{S}_2]) \\ &= p. \end{aligned}$$

Hence \hat{P} is a perfect-calibrated coverage estimator on $\mathcal{S}_1 \cup \mathcal{S}_2$. □

Lemma 1(a) is motivated from a theoretical point of view. It provides a guidance that in order to well resemble the conditional coverage, an estimator should be perfect-calibrated on as many subsets on the feature space as possible.

A.4 Details on Coverage Error

In Section 2, we have introduced CE_p to quantify the discrepancy between a coverage estimator and a perfect-calibrated coverage estimator, and \widehat{CE}_p to quantify the discrepancy between a coverage estimator and the conditional coverage. We note that by Hölder’s inequality

$$CE_p \leq CE_q, \text{ for } 1 \leq p \leq q \leq +\infty.$$

A larger value of p corresponds to a larger CE value. Continuing Definition 5, we can further introduce the calibration-based conditional coverage error on a measurable subset as follows:

Definition 7. An L^p ($1 \leq p \leq +\infty$) calibration-based conditional coverage error, or coverage error for short, of a coverage estimator \hat{P} on a measurable subset $\mathcal{S} \subset \mathcal{X}$ with $\mathbb{P}(\mathcal{S}) > 0$ is defined as:

$$CE_p(\mathcal{S}) = \left\| \mathbb{P}[Y \in [L(X), U(X)] | L, U, \hat{P}(X), X \in \mathcal{S}] - \hat{P}(X) \right\|_{L^p(\mathcal{S})} \quad (16)$$

where L^p -norm is taken with respect to the randomness of $\hat{P}(X)$ on the conditional probability space $(\mathcal{S}, \mathcal{F}_{\mathcal{S}} := \{A \cap \mathcal{S} : A \in \mathcal{F}\}, \mathbb{P}_{\mathcal{S}}(A \cap \mathcal{S}) := \mathbb{P}(A|\mathcal{S}))$. In particular, we have $CE_p := CE_p(\mathcal{X})$.

Lemma 3. A coverage estimator \hat{P} is the conditional coverage if and only if its coverage error $CE_p(\mathcal{S}) = 0$ for any measurable subset $\mathcal{S} \subset \mathcal{X}$ with $\mathbb{P}(\mathcal{S}) > 0$. In particular, a coverage estimator is the conditional coverage if and only if it is a perfect-calibrated coverage estimator on any measurable subset \mathcal{S} of \mathcal{X} with $\mathbb{P}(\mathcal{S}) > 0$.

Proof. We first show that the conditional coverage is perfect-calibrated. Taking conditional expectation given $\{A(X) = p, X \in \mathcal{S}\}$ in the Definition 1. Then we get

$$\begin{aligned} p &= \mathbb{E}[A(X) | A(X) = p, X \in \mathcal{S}] \\ &= \mathbb{E}[\mathbb{P}[Y \in [L(X), U(X)] | X] | A(X) = p, X \in \mathcal{S}] \\ &= \mathbb{P}[Y \in [L(X), U(X)] | A(X) = p, X \in \mathcal{S}] \quad \text{by the tower property.} \end{aligned}$$

So $A(x)$ is a perfect-calibrated coverage estimator on any measurable subset \mathcal{S} with $\mathbb{P}(\mathcal{S}) > 0$. Hence $CE_p(\mathcal{S}) = 0$ for any measurable subsets $\mathcal{S} \subset \mathcal{X}$ with $\mathbb{P}(\mathcal{S}) > 0$.

On the other hand, similarly to Section A.3, we can express

$$\mathbb{P}[Y \in [L(X), U(X)] | \hat{P}(X) = \hat{P}(x), X \in \mathcal{S}] = \frac{\int_{t \in \hat{P}^{-1}(\hat{P}(x)) \cap \mathcal{S}} A(t) \mathbb{P}[X \in dt]}{\mathbb{P}[X \in \hat{P}^{-1}(\hat{P}(x)) \cap \mathcal{S}]}.$$

Suppose $\hat{P}(x)$ is not the conditional coverage, then $\mathbb{P}[\hat{P}(X) \neq A(X)] > 0$. Without loss of generality, we assume $\mathbb{P}[\hat{P}(X) > A(X)] > 0$. Let $\mathcal{S}_0 := \{x \in \mathcal{X} : \hat{P}(x) > A(x)\}$. Note that $\mathcal{S}_0 = \cup_{n=1}^{+\infty} \{x \in \mathcal{X} : \hat{P}(x) > A(x) + \frac{1}{n}\}$. Since $\mathbb{P}(\mathcal{S}_0) > 0$, there exists a n_0 such that $\mathcal{S} := \{x \in \mathcal{X} : \hat{P}(x) > A(x) + \frac{1}{n_0}\}$ and $\mathbb{P}(\mathcal{S}) > 0$.

Then for $x \in \mathcal{S}$, we have

$$\frac{\int_{t \in \hat{P}^{-1}(\hat{P}(x)) \cap \mathcal{S}} A(t) \mathbb{P}[X \in dt]}{\mathbb{P}[X \in \hat{P}^{-1}(\hat{P}(x)) \cap \mathcal{S}]} \leq \frac{\int_{t \in \hat{P}^{-1}(\hat{P}(x)) \cap \mathcal{S}} (\hat{P}(t) - \frac{1}{n_0}) \mathbb{P}[X \in dt]}{\mathbb{P}[X \in \hat{P}^{-1}(\hat{P}(x)) \cap \mathcal{S}]} = \hat{P}(x) - \frac{1}{n_0}.$$

Then we have

$$CE_p(\mathcal{S}) \geq CE_1(\mathcal{S}) \geq \frac{1}{n_0} > 0$$

so $CE_p(\mathcal{S}) > 0$, which is a contradiction. Hence $\hat{P}(x)$ is the conditional coverage. \square

B Mathematical Developments for Theorem 1

This section proves Theorem 1: we show that both coverage error and conditional coverage error are tightly bounded above by the expectation of a Kullback–Leibler divergence-type random variable $K_1(x)$. This means that minimizing $\mathbb{E}[K_1(X)]$ can recover the true conditional coverage and effectively reduce the coverage error. The proof consists of several inequalities regarding coverage errors and their relations. We first begin with the following connection between CE_p and \widetilde{CE}_p .

Theorem 4. *For any PI and its associated \hat{P} , the L^p coverage error is always less than or equal to the L^p conditional coverage error, i.e.,*

$$CE_p \leq \widetilde{CE}_p, \quad \forall 1 \leq p \leq +\infty$$

Proof. We note that the function $t \mapsto |t|^p$ is a convex function. We also note that $\sigma(\hat{P}(X)) \subset \sigma(X)$ where $\sigma(Y)$ represents the σ -field generated by a random variable Y .

$$\begin{aligned} \widetilde{CE}_p^p &= \mathbb{E} \left[\left| A(X) - \hat{P}(X) \right|^p \right] \\ &= \mathbb{E} \left[\mathbb{E} |A(x) - \hat{P}(X)|^p | \hat{P}(X) \right] \\ &\geq \mathbb{E} \left[\left| \mathbb{E} [A(x) - \hat{P}(X) | \hat{P}(X)] \right|^p \right] \quad \text{by Jensen's inequality} \\ &= \mathbb{E} \left[\left| \mathbb{E} [\mathbf{1}_{Y \in [L(X), U(X)]} | \hat{P}(X)] - \hat{P}(X) \right|^p \right] \quad \text{by the tower property} \\ &= CE_p^p \end{aligned}$$

Therefore we have

$$CE_p \leq \widetilde{CE}_p, \quad \forall 1 \leq p \leq +\infty.$$

□

Next we have the following bounds on L^p conditional coverage error:

Theorem 5. *The L^p conditional coverage error is bounded above by a power function of the L^2 conditional coverage error. Formally,*

$$\widetilde{CE}_p \leq \widetilde{CE}_2^{\alpha_p}, \quad \forall 1 \leq p \leq +\infty,$$

where $\alpha_p = 1, \forall 1 \leq p \leq 2$ and $\alpha_p = \frac{2}{p}, \forall 2 \leq p \leq +\infty$.

Proof. By Hölder's inequality,

$$\widetilde{CE}_p \leq \widetilde{CE}_2, \quad \text{if } 1 \leq p \leq 2.$$

Since

$$0 \leq |A(x) - \hat{P}(x)| \leq 1,$$

then,

$$|A(x) - \hat{P}(x)|^p \leq |A(x) - \hat{P}(x)|^2, \quad \forall p \geq 2$$

and thus

$$\widetilde{CE}_p^p = \mathbb{E}[|A(X) - \hat{P}(X)|^p] \leq \mathbb{E}[|A(X) - \hat{P}(X)|^2] \leq \widetilde{CE}_2^2, \quad \forall p \geq 2.$$

□

Next, recall that

$$\begin{aligned} K(x) &= A(x) \log \left(\frac{A(x)}{\hat{P}(x)} \right) + (1 - A(x)) \log \left(\frac{1 - A(x)}{1 - \hat{P}(x)} \right), \\ K_0(x) &= A(x) \log(A(x)) + (1 - A(x)) \log(1 - A(x)), \\ K_1(x) &= -A(x) \log(\hat{P}(x)) - (1 - A(x)) \log(1 - \hat{P}(x)). \end{aligned}$$

Theorem 6. *The L^2 conditional coverage error is bounded above by the expectation of $K(x)$. Formally,*

$$\widetilde{CE}_2^{\alpha_p} \leq \left(\frac{1}{2} \mathbb{E}[K(X)] \right)^{\alpha_p/2},$$

where α_p is defined in Theorem 5.

Proof. For any fixed x , consider two random variables with Bernoulli distributions:

$$\begin{aligned} W_1 &= \begin{cases} 1 & \text{w.p. } A(x), \\ 0 & \text{w.p. } 1 - A(x). \end{cases} \\ W_2 &= \begin{cases} 1 & \text{w.p. } \hat{P}(x), \\ 0 & \text{w.p. } 1 - \hat{P}(x). \end{cases} \end{aligned}$$

Let P_i be the distribution of W_i . It follows from Pinsker's inequality, e.g., Theorem 2.16 in (Massart, 2007), that

$$\|P_1 - P_2\|_{TV}^2 \leq \frac{1}{2} K(P_1, P_2).$$

where TV denotes the total variation distance and K denotes the KL divergence. Since P_i is the Bernoulli distribution, we can express it as

$$|A(x) - \hat{P}(x)|^2 \leq \frac{1}{2} K(x)$$

Taking expectation, we obtain

$$\mathbb{E}[|A(X) - \hat{P}(X)|^2] \leq \frac{1}{2} \mathbb{E}[K(X)].$$

Hence,

$$\widetilde{CE}_2^{\alpha_p} \leq \left(\frac{1}{2} \mathbb{E}[K(X)] \right)^{\alpha_p/2}.$$

□

Combining Theorem 4, 5 and 6, we immediately conclude that:

Theorem 7 (Restated Theorem 1).

$$CE_p \leq \widetilde{CE}_p \leq \widetilde{CE}_2^{\alpha_p} \leq \left(\frac{1}{2} \mathbb{E}[K(X)] \right)^{\alpha_p/2}, \quad \forall 1 \leq p \leq +\infty$$

where α_p is defined in Theorem 5. Moreover, all inequalities are attainable, e.g., if $\hat{P}(x)$ is the conditional coverage $A(x)$.

C Justification of Assumption 1 and Mathematical Developments for Theorem 3

In this section, we analyze the rationality of Assumption 1 and build essential ingredients for proving Theorem 2 and 3. The difficulty in analyzing Theorem 3 lies in the fact that the hypothesis classes in Assumption 1 (which are constructed by the NN) are different from the hypothesis class used in L_{CA} . To overcome this difficulty, we use the theory of VC-subgraph classes to analyze the connection between the VC dimension of the two hypothesis classes.

C.1 Review of the VC Dimension

For self-contained purpose, we first review the definition of the VC-subgraph class and VC dimension.

Definition 8. Consider an arbitrary collection $\{x_1, \dots, x_n\}$ of points in a set \mathcal{X} and a collection \mathcal{C} of subsets of \mathcal{X} . We say that \mathcal{C} shatters $\{x_1, \dots, x_n\}$ if all of 2^n possible subsets of $\{x_1, \dots, x_n\}$ can be written as $A = C \cap \{x_1, \dots, x_n\}$ for some $C \in \mathcal{C}$. The VC dimension $V(\mathcal{C})$ of the class \mathcal{C} is the smallest n for which no set of size n $\{x_1, \dots, x_n\}$ is shattered by \mathcal{C} . If \mathcal{C} shatters sets of arbitrarily large size, we set $V(\mathcal{C}) = \infty$. We say that \mathcal{C} is a VC-class if $V(\mathcal{C}) < \infty$.

In some literature, the VC dimension $V(\mathcal{C})$ of the class \mathcal{C} is alternatively defined as the largest n for which there exists a set of size n $\{x_1, \dots, x_n\}$ shattered by \mathcal{C} , i.e., it is the value in definition 8 minus 1. We can more formally define the VC dimension by the growth function as follows:

Definition 9. Define the n^{th} shatter coefficient (or growth function) of \mathcal{C} as

$$\Pi_{\mathcal{C}}(n) := \max_{x_1, \dots, x_n} |\{A : A = C \cap \{x_1, \dots, x_n\} \text{ for some } C \in \mathcal{C}\}|$$

Then

$$V(\mathcal{C}) := \inf\{n : \Pi_{\mathcal{C}}(n) < 2^n\}.$$

Definition 10. For a function $f : \mathcal{X} \rightarrow \mathbb{R}$, the subset of $\mathcal{X} \times \mathbb{R}$ given by $\{(x, t) : t < f(x)\}$ is the (open) subgraph of f . A collection \mathcal{F} of measurable real functions on the sample space \mathcal{X} is a VC-subgraph class or VC-class, if the collection of all subgraphs of functions in \mathcal{F} forms a VC-class of sets (as sets in $\mathcal{X} \times \mathbb{R}$). Let $V(\mathcal{F})$ denote the VC dimension of the set of subgraphs of \mathcal{F} .

Lemma 4. In Definition 10, the open subgraph of f , $\{(x, t) : t < f(x)\}$, can be replaced by the close subgraph $\{(x, t) : t \leq f(x)\}$, the close supergraph $\{(x, t) : t \geq f(x)\}$ or the open supergraph $\{(x, t) : t > f(x)\}$. All of them lead to the equivalent definition of the VC-class and the equal VC dimension.

Proof. This result follows from Lemma 9.33 and Lemma 9.9(iv) in Kosorok (2007). \square

For indicator functions of sets, we have the following equivalence.

Lemma 5. For any class \mathcal{C} of sets in a set \mathcal{X} , the class $\mathcal{F}_{\mathcal{C}}$ of indicator functions of sets in \mathcal{C} is a VC-class if and only if \mathcal{C} is a VC-class. Moreover, whenever at least one of \mathcal{C} or $\mathcal{F}_{\mathcal{C}}$ is VC-class, the respective VC dimensions are equal.

Proof. This is Lemma 9.8 in Kosorok (2007). Note that the sets of \mathcal{C} are in \mathcal{X} while the subgraphs of functions of $\mathcal{F}_{\mathcal{C}}$ are in $\mathcal{X} \times \mathbb{R}$. \square

C.2 Justifying Assumption 1

We first restate the assumption:

Assumption 3 (Restated Assumption 1). The four classes of functions $([L(x), U(x)], \hat{P}(x), 1 - \hat{P}(x))$ output by the neural network (NN) in Figure 1 have finite VC dimensions, say they are bounded above by V_0 .

In Figure 1, the output four neurons of the NN are denoted as $(L(x), U(x), \hat{P}(x), 1 - \hat{P}(x))$. We further let $(\psi_1(x), \psi_2(x), \psi_3(x), \psi_4(x))$ denote the pre-activated values of $(L(x), U(x), \hat{P}(x), 1 - \hat{P}(x))$. In other words,

$$L(x) = \min(\psi_1(x), \psi_2(x)),$$

$$U(x) = \max(\psi_1(x), \psi_2(x)),$$

$$\hat{P}(x) = \text{softmax}(\psi_3(x) || \psi_4(x)) = \sigma(\psi_3(x) - \psi_4(x)),$$

$$1 - \hat{P}(x) = \text{softmax}(\psi_4(x) || \psi_3(x)) = \sigma(\psi_4(x) - \psi_3(x)),$$

where σ is the sigmoid function. Let the function classes

$$\mathcal{H}_1 = \{L(x) : L \text{ is output by the the NN}\},$$

$$\mathcal{H}_2 = \{U(x) : U \text{ is output by the the NN}\},$$

$$\mathcal{G} = \{\hat{P}(x) : \hat{P} \text{ is output by the NN}\},$$

$$1 - \mathcal{G} = \{1 - \hat{P}(x) : \hat{P} \text{ is output by the NN}\}.$$

Assumption 1 holds for a wide range of NNs, in particular the one we adopt in the experiments (where we use the ReLU-activated NN to construct ψ_i , $i = 1, 2, 3, 4$; see Section 5). Our first result is to concretely show that the four NN outputs above, \mathcal{H}_1 , \mathcal{H}_2 , \mathcal{G} and $1 - \mathcal{G}$, under the ReLU setting, all have finite VC dimensions and thus satisfy Assumption 1.

Theorem 8 (Restated Theorem 2). *Suppose $\psi_i, i = 1, 2, 3, 4$ are the pre-activated output neurons of the NN in Figure 1 using the ReLU activation function. Then Assumption 1 holds. Moreover, suppose the NN has W parameters and U computation units (nodes). Then $V_0 = O(WU)$.*

Proof. First, we look at $L(x)$ and $U(x)$. Note that the class of ψ_i ($i = 1, 2$) is constructed by a NN with the ReLU activation function. Therefore by Theorem 8 in Bartlett et al. (2019) (see also 6 below),

$$V(\{\psi_1\}) = O(WU) < \infty, \quad V(\{\psi_2\}) = O(WU) < \infty.$$

By Lemma 9.9 (i) in Kosorok (2007), we have

$$V(\mathcal{H}_1) = V(\{\min(\psi_1, \psi_2)\}) \leq V(\{\psi_1\}) + V(\{\psi_2\}) - 1 = O(WU) < \infty.$$

By Lemma 9.9 (ii) in Kosorok (2007), we have

$$V(\mathcal{H}_2) = V(\{\max(\psi_1, \psi_2)\}) \leq V(\{\psi_1\}) + V(\{\psi_2\}) - 1 = O(WU) < \infty.$$

Next, we look at $\hat{P}(x)$ and $1 - \hat{P}(x)$. We add an additional neuron after the layer where ψ_3, ψ_4 stand. This neuron is defined as $\psi_5 = \psi_3 - \psi_4$ which is a linear combination of ψ_3 and ψ_4 . Note that the class of ψ_5 is constructed by a NN with the ReLU activation function and linear activation function (by adding one unit and two parameters in the original NN). Therefore by Theorem 8 in Bartlett et al. (2019),

$$V(\{\psi_5\}) = O(WU) < \infty$$

By Lemma 9.9 (viii) in Kosorok (2007), we have

$$V(\mathcal{G}) = V(\{\sigma(\psi_5)\}) \leq V(\{\psi_5\}) = O(WU) < \infty.$$

since σ is a monotone function. Again, by Lemma 9.9 (viii) in Kosorok (2007), we have

$$V(1 - \mathcal{G}) \leq V(\mathcal{G}) = O(WU) < \infty.$$

since $t \mapsto 1 - t$ is a monotone function. □

We also list some results for other activations here. From these results, and using the same argument as above, we see that Assumption 1 holds similarly for all these activations.

Lemma 6. *Suppose the class of functions is constructed by a NN with W parameters and U units with activation functions that are piecewise polynomials with at most p pieces and of degree at most d . Then it has VC dimension $O(WU \log((d+1)p))$.*

Proof. This is Theorem 8 in Bartlett et al. (2019). \square

Note that the activation functions in Lemma 6 include in particular the ReLU activation and linear activation.

Lemma 7. *Suppose the class of functions is constructed by a NN with W parameters with binary as well as linear activation function. Then it has VC dimension $O(W^2)$.*

Proof. This is Theorem 5 in Sontag (1998). \square

Lemma 8. *Suppose the class of functions is constructed by a NN with W parameters and U units with activation function that is the standard sigmoid function (except that the output unit being a linear threshold unit). Then it has VC dimension $O(W^2U^2)$.*

Proof. This is Theorem 8.13 in Anthony and Bartlett (1999). \square

C.3 Connections among Different Hypothesis Classes

To prove Theorem 3, we need to study several building blocks on the relations between different hypothesis classes. Our first observation is:

Theorem 9. *Suppose $V(\mathcal{G}) < +\infty$. Then all of the following classes have VC dimension $\leq V(\mathcal{G})$:*

$$1 - \mathcal{G} := \{1 - \hat{P}(x) : \hat{P} \text{ is output by the NN}\}.$$

$$\mathcal{G}' := \log(\mathcal{G}) := \{\log(\hat{P}(x)) : \hat{P} \text{ is output by the NN}\}.$$

$$\log(1 - \mathcal{G}) := \{\log(1 - \hat{P}(x)) : \hat{P} \text{ is output by the NN}\}.$$

Proof. The result follows from Lemma 9.9 (viii) in Kosorok (2007) since all of the transformations are monotone functions. \square

Our second observation is about

$$\mathcal{F} = \{f(x, y) = I_{y \in [L(x), U(x)]} : L, U \text{ are output by the NN}\}.$$

Note that the domain of functions in \mathcal{F} is different from the domain of functions in \mathcal{H}_i ($i = 1, 2$) as it includes the outcome space. Below we derive a result that connects the VC dimension of \mathcal{H}_i with that of \mathcal{F} .

Theorem 10. Suppose $V(\mathcal{H}_i) \leq V_0$ ($i = 1, 2$). We have that

$$V(1 - \mathcal{F}) \leq V(\mathcal{F}) \leq 10(V_0 - 1) < +\infty$$

where

$$1 - \mathcal{F} := \{1 - f(x, y) : f \in \mathcal{F}\}.$$

Proof. The first inequality follows from Lemma 9.9 (viii) in [Kosorok \(2007\)](#). We consider the following two classes:

$$\mathcal{F}_1 := \{I_{L(x) \leq t} : L \in \mathcal{H}_1, t \in \mathbb{R}\},$$

$$\mathcal{F}_2 := \{I_{U(x) \geq t} : U \in \mathcal{H}_2, t \in \mathbb{R}\}.$$

Since the functions in \mathcal{F}_1 are all indicator functions, by Lemma 5,

$$V(\mathcal{F}_1) = V(\{(x, t) : L(x) \leq t\} : L \in \mathcal{H}_1, t \in \mathbb{R}).$$

Note that the latter is the VC dimension of the close supergraphs of all functions in \mathcal{H}_1 . Then by Definition 10 and Lemma 4, we have

$$V(\{(x, t) : L(x) \leq t\} : L \in \mathcal{H}_1, t \in \mathbb{R}) = V(\mathcal{H}_1)$$

Therefore we have

$$V(\mathcal{F}_1) = V(\mathcal{H}_1) \leq V_0$$

Similarly,

$$V(\mathcal{F}_2) = V(\mathcal{H}_2) \leq V_0$$

Note that we can write

$$I_{y \in [L(x), U(x)]} = I_{L(x) \leq y} I_{U(x) \geq y}$$

By the definition of growth functions,

$$\begin{aligned} \Pi_{\mathcal{F}}(m) &:= \max_{(x_1, y_1), \dots, (x_m, y_m)} |\{(I_{y_1 \in [L(x_1), U(x_1)]}, \dots, I_{y_m \in [L(x_m), U(x_m)]}) : L \in \mathcal{H}_1, U \in \mathcal{H}_2\}| \\ &\leq \max_{(x_1, y_1), \dots, (x_m, y_m)} |\{(I_{L(x_1) \leq y_1}, \dots, I_{L(x_m) \leq y_m}) : L \in \mathcal{H}_1\}| \times \\ &\quad \max_{(x_1, y_1), \dots, (x_m, y_m)} |\{(I_{U(x_1) \geq y_1}, \dots, I_{U(x_m) \geq y_m}) : U \in \mathcal{H}_2\}| \\ &= \Pi_{\mathcal{F}_1}(m) \Pi_{\mathcal{F}_2}(m) \\ &\leq \left(\frac{em}{V_0 - 1} \right)^{2(V_0 - 1)} \end{aligned}$$

for all $m \geq V_0$ where the last inequality is due to the Sauer–Shelah lemma. Taking $m = 10(V_0 - 1)$, we obtain

$$\left(\frac{em}{V_0 - 1} \right)^{2(V_0 - 1)} = (10e)^{2(V_0 - 1)} \leq 750^{V_0 - 1} < 2^m$$

Combining the above inequality, we have

$$\Pi_{\mathcal{F}}(m) < 2^m.$$

This shows that $V(\mathcal{F}) \leq m = 10(V_0 - 1)$. □

C.4 Proof of Theorem 3

This subsection proves Theorem 3. Recall that

$$\mathcal{F} = \{I_{y \in [L(x), U(x)]} : L, U \text{ are output by the NN}\},$$

$$\mathcal{G} = \{\hat{P}(x) : \hat{P} \text{ is output by the NN}\}.$$

$$\mathcal{G}' := \log(\mathcal{G}) := \{\log(\hat{P}(x)) : \hat{P} \text{ is output by the NN}\}.$$

Let $N(\epsilon, \mathcal{F}, L^2(Q))$ denote the covering number, i.e., the minimal number of balls $\{g : \|g - h\|_{L^2(Q)} < \epsilon\}$ of radius ϵ needed to cover the set \mathcal{F} . We need the following bounds:

Lemma 9. *Suppose \mathcal{F} is a class of functions $f : \mathcal{X} \times \mathcal{Y} \rightarrow [0, 1]$ with a finite VC dimension $V(\mathcal{F})$. For every $0 < \epsilon < 1$,*

$$\sup_Q \log N(\epsilon, \mathcal{F}, L^2(Q)) \leq K_2 \left(\frac{1}{\epsilon}\right)^{\frac{1}{e}}$$

where the constant K_2 depends on $V(\mathcal{F})$ only.

Proof. It follows from Theorem 2.6.7 in [Van der Vaart and Wellner \(1996\)](#) that there exists a universal constant K such that

$$\sup_Q N(\epsilon, \mathcal{F}, L^2(Q)) \leq KV(\mathcal{F})(16e)^{V(\mathcal{F})} \left(\frac{1}{\epsilon}\right)^{V(\mathcal{F})-1}$$

for any $0 < \epsilon < 1$. Since

$$\left(\frac{1}{\epsilon}\right)^{\frac{1}{e}} \geq \max\left(\log\left(\frac{1}{\epsilon}\right), 1\right), \quad \forall 0 < \epsilon < 1,$$

we have

$$\sup_Q \log N(\epsilon, \mathcal{F}, L^2(Q)) \leq K_3 + (V(\mathcal{F}) - 1) \log\left(\frac{1}{\epsilon}\right) \leq K_2 \left(\frac{1}{\epsilon}\right)^{\frac{1}{e}}$$

where $K_3 := \log(KV(\mathcal{F})(16e)^{V(\mathcal{F})})$ and $K_2 = K_3 + V(\mathcal{F}) - 1$ only depending on $V(\mathcal{F})$. □

We remark that a similar result can also be obtained for the class $1 - \mathcal{F}$ by Theorem 10.

Lemma 10. Suppose \mathcal{G} is a class of functions $\hat{P} : \mathcal{X} \rightarrow [0, 1]$ with a finite VC dimension $V(\mathcal{G})$ and $|\log(\hat{P}(x))| \leq M$. Let $\mathcal{G}' := \{\log(\hat{P}) : \hat{P} \in \mathcal{G}\}$. Then, for every $0 < \epsilon < 1$,

$$\sup_Q \log N(\epsilon M, \mathcal{G}', L^2(Q)) \leq K_2 \left(\frac{1}{\epsilon} \right)^{\frac{1}{e}}$$

where the constant K_2 depends on $V(\mathcal{G})$ only.

Proof. First note that $\phi(t) := \log(t)$ is a monotone function. Hence $\mathcal{G}' := \{\log(\hat{P}) : \hat{P} \in \mathcal{G}\}$ is a VC-class with VC dimension $\leq V(\mathcal{G})$ by Lemma 9.9 (viii) in [Kosorok \(2007\)](#). The rest of the proof is similar to [9](#). \square

We remark that a similar result can also be obtained for the class $\log(1 - \mathcal{G})$ by Theorem [9](#).

Next we restate Assumption [2](#):

Assumption 4 (Restated Assumption [2](#)). *There exists a constant $M < +\infty$ such that $|\log(\hat{P}(x))| \leq M$, $|\log(1 - \hat{P}(x))| \leq M$ for all x and \hat{P} .*

As discussed in Section [4](#), this is a natural assumption in practice because $\log(\hat{P}(x))$ and $\log(1 - \hat{P}(x))$ are replaced by $\log(\hat{P}(x) + \epsilon)$ and $\log(1 - \hat{P}(x) + \epsilon)$ respectively to avoid explosion when implementing the algorithm. In particular, in our experiments in Section [5](#), $\epsilon = 0.1^6$ and thus $M = 14$.

We are now ready to prove Theorem [3](#):

Theorem 11 (Restated Theorem [3](#)). *Suppose Assumptions [1](#) and [2](#) hold. The training data $\mathcal{D} = \{(x_i, y_i), i = 1, 2, \dots, n\}$ where (x_i, y_i) are i.i.d. samples $\sim \pi$. Recall that the (hard) coverage estimator assessment loss is*

$$L_{CA} = -\frac{1}{n} \sum_{i=1}^n \left(f(x_i, y_i) \log(\hat{P}(x_i)) + (1 - f(x_i, y_i)) \log(1 - \hat{P}(x_i)) \right).$$

Then for any $t > 0$, we have

$$\mathbb{P} \left(\sup_{f \in \mathcal{F}, \hat{P} \in \mathcal{G}} |L_{CA} - \mathbb{E}[K_1(X)]| \geq t \right) \leq C^* e^{-\frac{nt^2}{16M^2}}$$

where C^* only depends on V_0 in Assumption [1](#).

Proof. Note that $\mathbb{E}[f(x_i, y_i)|x_i] = A(x_i)$ for any fixed L and U . Taking expectation on L_{CA} , we have

$$\begin{aligned} \mathbb{E}[L_{CA}] &= \mathbb{E}[\mathbb{E}[L_{CA}|x_1, x_2, \dots, x_n]] \\ &= \mathbb{E} \left[-\frac{1}{n} \sum_{i=1}^n \left(A(x_i) \log(\hat{P}(x_i)) + (1 - A(x_i)) \log(1 - \hat{P}(x_i)) \right) \right] \\ &= \mathbb{E}[K_1(X)]. \end{aligned}$$

We consider the first part $A(x) \log(\hat{P}(x))$. The second part can be done using the same argument. Note that by Theorem 9.15 in [Kosorok \(2007\)](#), we have

$$\sup_Q \log N(\epsilon M, \mathcal{F} \cdot \mathcal{G}', L^2(Q)) \leq \sup_Q \log N(\epsilon/2, \mathcal{F}, L^2(Q)) + \sup_Q \log N(\epsilon M/2, \mathcal{G}', L^2(Q)).$$

Consider the class $\frac{1}{2} + \frac{1}{2M} \mathcal{F} \cdot \mathcal{G}' := \{\frac{1}{2} + \frac{1}{2M} (f(x, y) \log(\hat{P}(x))) : f \in \mathcal{F}, \log(\hat{P}) \in \mathcal{G}'\}$ which consists of functions taking values in $[0, 1]$. We have

$$\begin{aligned} & \sup_Q \log N(\epsilon, \frac{1}{2} + \frac{1}{2M} \mathcal{F} \cdot \mathcal{G}', L^2(Q)) \\ &= \sup_Q \log N(2\epsilon M, \mathcal{F} \cdot \mathcal{G}', L^2(Q)) \\ &\leq \sup_Q \log N(\epsilon, \mathcal{F}, L^2(Q)) + \sup_Q \log N(\epsilon M, \mathcal{G}', L^2(Q)) \\ &\leq K_2 \left(\frac{1}{\epsilon}\right)^{1/e} \end{aligned}$$

where the last inequality follows from Lemma 9 and Lemma 10, and K_2 only depends on $V(\mathcal{F})$ and $V(\mathcal{G})$. (Recall that we have shown $V(\mathcal{G}') \leq V(\mathcal{G})$ in Lemma 10.) Moreover, by Theorems 9 and 10, we can claim that K_2 only depends on V_0 . This inequality shows that $\frac{1}{2} + \frac{1}{2M} \mathcal{F} \cdot \mathcal{G}'$ satisfies the conditions in Theorem 2.14.10 in [Van der Vaart and Wellner \(1996\)](#) and thus for every $\delta > 0$ and $t > 0$,

$$\mathbb{P} \left(\sup_{\phi \in \frac{1}{2} + \frac{1}{2M} \mathcal{F} \cdot \mathcal{G}'} \left| \frac{1}{n} \sum_{i=1}^n \phi(x_i, y_i) - \mathbb{E}[\phi(x, y)] \right| \geq t \right) \leq C e^{D(\sqrt{nt})^{U+\delta}} e^{-2nt^2}$$

where $U = \frac{\frac{1}{e}(6-\frac{1}{e})}{2+\frac{1}{e}} < 1$ and the constants C and D depend on K_2 and δ only. Let $\delta = 1 - U$. Note that

$$-2(\sqrt{nt})^2 + D(\sqrt{nt}) \leq -(\sqrt{nt})^2 + (D/2)^2.$$

Hence we have

$$\mathbb{P} \left(\sup_{\phi \in \frac{1}{2} + \frac{1}{2M} \mathcal{F} \cdot \mathcal{G}'} \left| \frac{1}{n} \sum_{i=1}^n \phi(x_i, y_i) - \mathbb{E}[\phi(x, y)] \right| \geq t \right) \leq C^* e^{-nt^2}$$

where C^* only depends on K_2 , or, only depends on V_0 .

This shows that

$$\mathbb{P} \left(\sup_{f \in \mathcal{F}, \hat{P} \in \mathcal{G}} \left| \frac{1}{n} \sum_{i=1}^n f(x_i, y_i) \log(\hat{P}(x_i)) - \mathbb{E}[A(x) \log(\hat{P}(x))] \right| \geq t \right) \leq C^* e^{-\frac{nt^2}{4M^2}}.$$

A similar result can be established for the second part since the hypothesis classes there have been studied in Theorem 9 and 10:

$$\begin{aligned} & \mathbb{P} \left(\sup_{f \in \mathcal{F}, \hat{P} \in \mathcal{G}} \left| \frac{1}{n} \sum_{i=1}^n (1 - f(x_i, y_i)) \log(1 - \hat{P}(x_i)) - \mathbb{E}[(1 - A(x)) \log(1 - \hat{P}(x))] \right| \geq t \right) \\ & \leq C^* e^{-\frac{nt^2}{4M^2}}. \end{aligned}$$

Combining the two parts and noting the following fact:

$$\begin{aligned} & \{\sup |\gamma + \beta| \geq t\} \\ & \subset \{\sup |\gamma| + \sup |\beta| \geq t\} \\ & \subset \{\sup |\gamma| \geq \frac{t}{2}\} \cup \{\sup |\beta| \geq \frac{t}{2}\}, \end{aligned}$$

we conclude that

$$\mathbb{P} \left(\sup_{f \in \mathcal{F}, \hat{P} \in \mathcal{G}} |L_{CA} - \mathbb{E}[K_1(x)]| \geq t \right) \leq C^* e^{-\frac{nt^2}{16M^2}}$$

where C^* only depends on V_0 . □

Lastly, the following corollary explicitly connects our theoretical developments to the experimental setup:

Corollary 1. *Suppose the NN is designed as the one specified in the experiments (Section 5). The training data $\mathcal{D} = \{(x_i, y_i), i = 1, 2, \dots, n\}$ where (x_i, y_i) are i.i.d. samples $\sim \pi$. Then for any $t > 0$, we have*

$$\mathbb{P} \left(\sup_{f \in \mathcal{F}, \hat{P} \in \mathcal{G}} |L_{CA} - \mathbb{E}[K_1(X)]| \geq t \right) \leq C^* e^{-\frac{nt^2}{16M^2}}.$$

where C^* only depends on V_0 in Assumption 1.

Proof. We note that Assumptions 1 and 2 hold in this case by Theorem 2 and the observation after Assumption 2. So Theorem 3 implies Corollary 1. □

D Experimental Details and More Results

This section illustrates experimental details and more experimental results from our proposed model.

D.1 Experimental Details

Table 3 gives a detailed description about the datasets we use. These open-access real-world benchmark regression datasets are widely used for the evaluation of methods in regression tasks (Hernández-Lobato and Adams, 2015; Gal and Ghahramani, 2016; Lakshminarayanan, 2017; Rosenfeld et al., 2018; Pearce et al., 2018; Zhu et al., 2019).

For synthetic datasets, 2000 i.i.d data are generated for each synthetic setting and randomly split into 1000 training data and 1000 testing data. For benchmark datasets, we first do the data normalization and then randomly split 80% data for training and 20% for testing. The choice of

80%/20% split, compared to the 90%/10% split in [Pearce et al. \(2018\)](#), is motivated from the need to increase the test size in order to get a meaningful *ECE* evaluation. The latter is due to that evaluating *ECE* requires binning, where using a larger number of bins approximates more closely *CE*, but also requires a larger test size to sustain enough statistical quality for the resulting *ECE* estimate. This delicate tradeoff motivates us to increase the share of the test set in our split. Following [Pearce et al. \(2018\)](#), our hyper-parameters are selected using the validation set from a random split. Then, they are fixed during the evaluation on other random splits.

As specified in Section 2 (Equation (8)) and Appendix A.4 (Equation (16)), *ECE* is evaluated based on dividing $[0,1]$ into M sub-intervals. The larger M is, the more precise is in using *ECE* to approximate *CE*, the latter being the ideal conditional coverage error estimator. On the other hand, a larger test set size is needed to support the use of a larger M without deteriorating the statistical quality of the *ECE*. This delicate tradeoff motivates us to increase the share of the test set in our split.

For synthetic datasets, the hyperparameters and corresponding results in Figure 2 are:

- (a) $\lambda_1 = 1.7, \lambda_2 = 10^{-5}, \lambda_3 = 1500, CP = 0.95, IW = 0.40, ECE_1 = 0.62\%$.
- (b) $\lambda_1 = 1.9, \lambda_2 = 10^{-5}, \lambda_3 = 1000, CP = 0.96, IW = 0.40, ECE_1 = 0.12\%$.
- (c) $\lambda_1 = 3.4, \lambda_2 = 10^{-5}, \lambda_3 = 1000, CP = 0.95, IW = 0.50, ECE_1 = 0.65\%$.

For benchmark datasets, the implementation details for baseline algorithms in Table 1 are:

(1) Nearest-neighbors kernel conditional density estimation (NNKCDE). The algorithm is based on Section 2.1 in [Dalmasso et al. \(2020\)](#). We use the same Python code provided by [Dalmasso et al. \(2020\)](#) with the default Gaussian kernel. Two tuning parameters, i.e., the number of nearest neighbors k and the bandwidth h of the smoothing kernel, are chosen in a principled way by minimizing the CDE loss on validation data, the same way as in [Dalmasso et al. \(2020\)](#).

(2) Quantile regression forest (QRF). The algorithm is based on [Meinshausen \(2006\)](#). We use the *RandomForestQuantileRegressor* from the package *scikit-garden* in Python.

(3) Split conformal learning (SCL). The algorithm based on Algorithm 2 in [Lei et al. \(2018\)](#). The regression algorithm inside SCL that we use is a neural network with mean square loss. The neural network has the same structure of hidden layers as in Section 5.

D.2 Additional Experimental Results

Table 4 gives experimental results for more datasets in addition to Section 5.3.

Dataset	N	d	Open-access Link
Boston: Boston Housing	506	13	kaggle.com/c/boston-housing
Concrete: Concrete Strength	1030	8	kaggle.com/aakashphadtare/concrete-data
Energy: Energy Efficiency	768	8	kaggle.com/elikplim/eergy-efficiency-dataset
Kin8nm	8192	8	openml.org/d/189
Plant: Combined Cycle Power Plant	9568	4	kaggle.com/gova26/airpressure
Protein: Protein Structure	45730	9	networkrepository.com/CASP.php
Wine: Red Wine Quality	1599	11	kaggle.com/uciml/red-wine-quality-cortez-et-al-2009
Yacht: Yacht Hydrodynamics	308	6	archive.ics.uci.edu/ml/datasets/yacht+hydrodynamics

Table 3: Full names and details of benchmarking regression datasets. N is the number of samples in the dataset and d is the dimension of the feature vector.

Dataset	λ_3	1-st experiment				2-nd experiment				3-rd experiment			
		λ_1	CP	IW	ECE_1	λ_1	CP	IW	ECE_1	λ_1	CP	IW	ECE_1
Boston	1800	3.5	0.87	0.72	1.25%	4.5	0.89	0.85	0.96%	6.0	0.95	1.04	1.38%
Concrete	1000	3.0	0.87	0.88	1.03%	5.0	0.93	1.07	1.03%	6.5	0.95	1.13	0.24%
Kin8nm	300	2.1	0.85	0.77	1.59%	2.9	0.91	0.89	0.50%	3.6	0.95	1.04	1.32%
Plant	700	1.6	0.85	0.63	0.86%	2.4	0.91	0.75	1.15%	3.3	0.95	0.85	0.38%
Protein	300	5.1	0.86	1.70	0.71%	6.3	0.91	1.95	0.66%	8.3	0.95	2.26	0.41%
Wine	1100	9.0	0.85	1.59	0.98%	15	0.91	2.14	1.28%	19	0.95	2.59	0.42%
Yacht	500	1.4	0.94	0.13	0.13%	1.5	0.96	0.14	1.51%	1.6	0.98	0.16	0.80%
Synthetic1	1500	1.0	0.89	0.34	0.73%	1.3	0.93	0.37	0.55%	1.7	0.95	0.40	0.62%
Synthetic2	1000	1.1	0.89	0.31	1.01%	1.4	0.92	0.34	1.27%	1.9	0.96	0.40	0.12%
Synthetic3	1000	1.3	0.87	0.34	0.79%	2.4	0.91	0.42	0.82%	3.4	0.95	0.50	0.65%

Table 4: Evaluation metrics of our CaNet on benchmark datasets and synthetic examples with different coverage probabilities.