Randomized Smoothing of All Shapes and Sizes

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Abstract

Randomized smoothing is a recently proposed defense against adversarial attacks that has achieved state-of-the-art provable robustness against ℓ_2 perturbations. Soon after, a number of works devised new randomized smoothing schemes for other metrics, such as ℓ_1 or ℓ_∞ ; however, for each geometry, substantial effort was needed to derive new robustness guarantees. This begs the question: can we find a general theory for randomized smoothing?

In this work we propose a novel framework for devising and analyzing randomized smoothing schemes, and validate its effectiveness in practice. Our theoretical contributions are as follows: (1) We show that for an appropriate notion of "optimal", the optimal smoothing distributions for any "nice" norm have level sets given by the *Wulff Crystal* of that norm. (2) We propose two novel and complementary methods for deriving provably robust radii for any smoothing distribution. Finally, (3) we show fundamental limits to current randomized smoothing techniques via the theory of Banach space cotypes. By combining (1) and (2), we significantly improve the state-of-the-art certified accuracy in ℓ_1 on standard datasets. On the other hand, using (3), we show that, without more information than label statistics under random input perturbations, randomized smoothing cannot achieve nontrivial certified accuracy against perturbations of ℓ_{∞} -norm $\Omega(1/\sqrt{d})$, when the input dimension d is large. We provide code in github.com/tonyduan/rs4a.

1. Introduction

Deep learning models are vulnerable to adversarial examples – small imperceptible perturbations to their inputs that lead to misclassification (Goodfellow et al., 2015; Szegedy et al., 2014). To solve this problem, recent works proposed heuristic defenses that are robust to specific classes of perturbations, but many would later be broken by stronger attacking algorithms (Carlini & Wagner, 2017; Athalye et al., 2018; Uesato et al., 2018). This led the community to both strengthen empirical defenses (Kurakin et al., 2016; Madry et al., 2017) as well as build *certified* defenses that provide robustness guarantees, i.e., models whose predictions are constant within a neighborhood of their inputs (Wong & Kolter, 2018; Raghunathan et al., 2018a). In particular, randomized smoothing is a recent method that has achieved state-of-the-art provable robustness (Lecuyer et al., 2018; Li et al., 2018; Cohen et al., 2019). In short, given an input, it outputs the class most likely to be returned by a base classifier, typically a neural network, under random noise perturbation of the input. This mechanism confers stability of the output against ℓ_p perturbations, even if the base classifier itself is highly non-Lipschitz. Canonically, this noise has been Gaussian, and the adversarial perturbation it protects against has been ℓ_2 (Cohen et al., 2019; Salman et al., 2019a), but some have explored other kinds of noises and adversaries as well (Lecuyer et al., 2018; Li et al., 2019; Dvijotham et al., 2019). In this paper, we seek to comprehensively understand the interaction between the choice of smoothing distribution and the perturbation norm.

- 1. We propose two new methods to compute robust certificates for additive randomized smoothing against different norms.
- 2. We show that, for $\ell_1, \ell_2, \ell_\infty$ adversaries, the optimal smoothing distributions have level sets that are their respective *Wulff Crystals* a kind of equilibrated crystal structure studied in physics since 1901 (Wulff).
- 3. Using the above advances, we obtain state-of-the-art ℓ_1 certified accuracy on CIFAR-10 and ImageNet¹.
- 4. Finally, we leverage the classical theory of Banach space *cotypes* (Wojtaszczyk, 1996) to show that current techniques for randomized smoothing cannot certify nontrivial accuracy at more than $\Omega(1/\sqrt{d}) \ell_{\infty}$ radius, if all one uses are the probabilities of labels when classifying randomly perturbed input.

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¹ We do not train the smoothed model adversarially (c.f. SmoothAdv of Salman et al. (2019a)), as we seek only to elucidate the interaction between the adversary and the noise. Nevertheless, our results are state-of-the-art even compared against previous works that do.

ImageNet	ℓ_1 Radius	0.5	1.0	1.5	2.0	2.5	3.0	3.5	4.0
	Uniform, Ours (%)	55	49	46	42	37	33	28	25
	Laplace, Teng et al. (2019) (%)	48	40	31	26	22	19	17	14
CIFAR-10	ℓ_1 Radius	0.5	1.0	1.5	2.0	2.5	3.0	3.5	4.0
	Uniform, Ours (%)	70	59	51	43	33	27	22	18
	Laplace, Teng et al. (2019) (%)	61	39	24	16	11	7	4	3

Table 1. Certified top-1 accuracies of our ℓ_1 -robust classifiers, vs previous state-of-the-art, at various radii, for ImageNet and CIFAR-10.³

2. Related Works

Defences against adversarial examples are mainly divided into *empirical* defenses and *certified* defenses.

Empirical defenses are heuristics designed to make learned models empirically robust. An example of these are *adversarial training* based defenses (Kurakin et al., 2016; Madry et al., 2017) which optimize the parameters of a model by minimizing the worst-case loss over a neighborhood around the input to these models. Such defenses may seem powerful, but has no guarantees that they are not "breakable". In fact, the majority of the empirical defenses proposed in the literature were later "broken" by stronger attacks (Carlini & Wagner, 2017; Athalye et al., 2018; Uesato et al., 2018; Athalye & Carlini, 2018). In order to mitigate the deficiency in such defenses, recent works explored certified defenses with formal robustness guarantees.

Certified defenses guarantee that for any input x, the classifier's output is constant within a small neighborhood of x. Such defenses are typically based on certification methods that are either exact or conservative. The exact methods include those based on Satisfiability Modulo Theories solvers (Katz et al., 2017; Ehlers, 2017) or mixed integer linear programming (Tjeng et al., 2019; Lomuscio & Maganti, 2017; Fischetti & Jo, 2017), which, although guaranteed to find adversarial examples if they exist, are unfortunately computationally inefficient. On the other hand, conservative methods are more computationally efficient, but might mistakenly flag a "safe" data point as vulnerable to adversarial examples (Wong & Kolter, 2018; Wang et al., 2018a;b; Raghunathan et al., 2018a;b; Wong et al., 2018; Dvijotham et al., 2018b;a; Croce et al., 2018; Salman et al., 2019b; Gehr et al., 2018; Mirman et al., 2018; Singh et al., 2018; Gowal et al., 2018; Weng et al., 2018; Zhang et al., 2018). However, none of these defences scale to practical networks. Recently, a new method called randomized smoothing has been proposed as a *probabilistically* certified defense, whose architecture-independence makes it scalable.

Randomized smoothing Randomized smoothing was first proposed as heuristic defense without any guarantees (Liu et al., 2018; Cao & Gong, 2017). Later on, Lecuyer et al. (2018) proved an robustness guarantee for randomized smoothing classifiers from a differential privacy perspective. Subsequently, Li et al. (2018) gave a stronger robustness guarantee utilizing tools from information theory. Recently, Cohen et al. (2019) provided a tight ℓ_2 robustness guarantee for randomized smoothing. Furthermore, a series of papers came out recently that developed robustness guarantees using randomized smoothing against other adversaries such as ℓ_1 -bounded (Teng et al., 2019), ℓ_∞ -bounded (Zhang* et al., 2020), ℓ_0 -bounded (Levine & Feizi, 2019a; Lee et al., 2019), and Wasserstein attacks (Levine & Feizi, 2019b). In Section 4.3, we give a more in-depth comparison on how our techniques compare to their results.

Wulff Crystal We are the first to relate to adversarial robustness the theory of *Wulff Crystals*, which has an interesting history. Just as the round soap bubble minimizes surface tension for a given volume, the Wulff Crystal minimizes certain surface energy that arises when the crystal interfaces with another material, akin to surface tension. The Russian physicist George Wulff first proposed this shape via physical arguments in 1901 (Wulff, 1901), but its energy minimization property was not proven in full generality until relatively recently, building on a century worth of work (Gibbs, 1875; Wulff, 1901; Hilton, 1903; Liebmann, 1914; von Laue; Dinghas, 1944; Burton et al., 1951; Herring; Constable, 1968; Taylor, 1975; 1978; Fonseca & Müller, 1991; Brothers & Morgan, 1994; Cerf, 2006).

No-go theorems for randomized smoothing Prior to the initial submission of this manuscript, the only other nogo theorem for randomized smoothing in the context of adversarial robustness is Zheng et al. (2020). However, they are only concerned with a non-standard notion of certified robustness that does not imply anything for the original problem. Moreover, they show that, under this different notion of robustness, if they are robust for ℓ_{∞} , then the ℓ_2 norm of the noise must be large on average. While this provides some indirect evidence for the hardness of certifying ℓ_{∞} , it does not actually address the question. Our result, on the other hand, directly rules out a large suite

³These models were trained with noise augmentation and not adversarial training. In our replication of Teng et al. (2019)'s experiments, we found that our results with noise augmentation matched their results with adversarial training.

of current techniques for deriving robust certificates for all ℓ_p norms for p > 2, for the standard notion of certified robustness.

After the initial submission of this manuscript, we became aware of two concurrent works (Kumar et al., 2020; Blum et al., 2020) that claim impossibility results for randomized smoothing. Blum et al. (2020) demonstrate that, under some mild conditions, any smoothing distribution for ℓ_p with p > 2 must have large component-wise magnitude. While this work, like Zheng et al. (2020), gives indirect evidence for the hardness of the problem, it does not directly show a limit for the *utility* of randomized smoothness. In contrast, our no-go result directly shows impossibility of the underlying robust classification problem. Kumar et al. (2020) demonstrate that certain classes of smoothing distributions cannot certify ℓ_{∞} without losing dimension-dependent factors. Our result is more general, as it rules out *any* class of smoothing distributions.

3. Randomized Smoothing

Consider a classifier f from \mathbb{R}^d to classes \mathcal{Y} and a distribution q on \mathbb{R}^d . Randomized smoothing with q is a method that constructs a new, *smoothed* classifier g from the *base* classifier f. The smoothed classifier g assigns to a query point x the class which is most likely to be returned by the base classifier f when x is perturbed by a random noise sampled from q, i.e.,

$$g(x) \stackrel{\text{def}}{=} \operatorname*{argmax}_{c \in \mathcal{Y}} q(U_c - x) \tag{1}$$

where U_c is the decision region $\{x' \in \mathbb{R}^d : f(x') = c\}, U_c - x$ denotes the translation of U_c by -x, and q(U) is the measure of U under q, i.e. $q(U) = \mathbb{P}_{\delta \sim q}(\delta \in U)$.

Robustness guarantee for smoothed classifiers For $p \in [0, 1], v \in \mathbb{R}^d$, define the *growth function*

$$\mathcal{G}_q(p,v) \stackrel{\text{def}}{=} \sup_{U \subseteq \mathbb{R}^d: q(U) = p} q(U-v), \tag{2}$$

One should think of U has the decision region of some base classifier. Thus $\mathcal{G}_q(p, v)$ gives the maximal growth of measure of a set (i.e. decision region) when q is shifted by the vector v, if we only know the initial measure p of the set.

Consider an adversary that can perturb an input additively by any vector v inside an allowed set \mathcal{B} . In the case when \mathcal{B} is the ℓ_2 ball and q is the Gaussian measure, Cohen et al. (2019) gave a simple expression for \mathcal{G}_q involving the Gaussian CDF, derived via the Neyman-Pearson lemma, which is later rederived by Salman et al. (2019a) as a nonlinear Lipschitz property. Likewise, the expression for Laplace distributions was derived by Teng et al. (2019). (See Theorem H.1 and Theorem I.1 for their expressions.) Suppose when the base classifier f classifies $x + \delta$, $\delta \sim q$, the class $c \in \mathcal{Y}$ is returned with probability $\rho = \mathbb{P}_{\delta \sim q}(f(x + \delta) = c) > 1/2$. Then the smoothed classifier g will not change its prediction under the adversary's perturbations if ⁴

$$\sup_{v \in \mathcal{B}} \mathcal{G}_q(1-\rho, v) < 1/2. \tag{3}$$

4. Methods for Deriving Robust Radii

Let q be a distribution with a density function, and we shall write $q(x), x \in \mathbb{R}^d$, for the value of the density function on x. Then, given a shift vector $v \in \mathbb{R}^d$ and a ratio $\kappa > 0$, define the *Neyman-Pearson set*

$$\mathcal{NP}_{\kappa} \stackrel{\text{def}}{=} \{ x \in \mathbb{R}^d : \kappa q(x-v) \ge q(x) \}.$$
(4)

Then the Neyman-Pearson lemma tells us that (Neyman & Pearson, 1933; Cohen et al., 2019)

$$\mathcal{G}_q(q(\mathcal{NP}_\kappa), v) = q(\mathcal{NP}_\kappa - v).$$
 (NP)

While this gives way to a simple expression for the growth function when q is Gaussian (Cohen et al., 2019), it is difficult for more general distributions as the geometry of \mathcal{NP}_{κ} becomes hard to grasp. To overcome this difficulty, we propose the *level set method* that decomposes this geometry so as to compute the growth function exactly, and the *differential method* that upper bounds the growth function derivative, loosely speaking.

4.1. The Level Set Method

For each t > 0, let U_t be the superlevel set

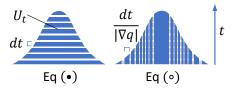
$$U_t \stackrel{\text{def}}{=} \{ x \in \mathbb{R}^d : q(x) \ge t \}$$

Then its boundary ∂U_t is the level set with q(x) = t under regularity assumptions. The integral of q's density is of course 1, but this integral can be expressed as the integral of the volumes of its superlevel sets:

$$1 = \int q(x) \, \mathrm{d}x = \int_0^\infty \operatorname{Vol}(U_t) \, \mathrm{d}t. \tag{(\bullet)}$$

If q has a differentiable density, then we may rewrite this as an integral of *level* sets (Theorem F.3):

$$1 = \int_0^\infty \int_{\partial U_t} \frac{t}{\|\nabla q(x)\|_2} \,\mathrm{d}x \,\mathrm{d}t. \tag{o}$$



⁴Many earlier works state robustness guarantees in terms of estimates of $p_A = \rho$ of the top class and p_B of the runner up class; however, their implementations are all in the form provided here, as p_B is usually taken to be $1 - p_A$.

The graphics above illustrate the two integral expressions (best viewed on screen). In this level set perspective, the Neyman-Pearson set NP_{κ} (Eq. (4)) can be written as

$$\begin{split} \mathcal{NP}_{\kappa} &= \bigcup_{t>0} \{ x : q(x) = t \text{ and } q(x-v) \ge t/\kappa \} \\ &= \bigcup_{t>0} \{ \partial U_t \cap (U_{t/\kappa} + v) \}. \end{split}$$

Then naturally, its measure is calculated by

$$q(\mathcal{NP}_{\kappa}) = \int_0^\infty \int_{\partial U_t \cap (U_{t/\kappa} + v)} \frac{t}{\|\nabla q(x)\|_2} \, \mathrm{d}x \, \mathrm{d}t. \quad (\lor)$$

Similarly, the Neyman-Pearson set can also be written from the perspective of $q(\cdot - v)$,

$$\mathcal{NP}_{\kappa} = \bigcup_{t>0} \{ x : q(x-v) = t \text{ and } q(x) \le t\kappa \}$$
$$= \bigcup_{t>0} \{ (\partial U_t + v) \setminus \mathring{U}_{t\kappa} \},$$

where \mathring{U} is the interior of the closed set U. So its measure under $q(\cdot - v)$ is

$$q(\mathcal{NP}_{\kappa} - v) = \int_{0}^{\infty} \int_{\partial U_{t} \setminus (\hat{U}_{t\kappa} - v)} \frac{t}{\|\nabla q(x)\|_{2}} \, \mathrm{d}x \, \mathrm{d}t. \quad (\wedge)$$

$$\frac{\partial U_{t} \cap (U_{t/\kappa} + v)}{U_{t/\kappa} + v} \quad \frac{\partial U_{t} \setminus (U_{t\kappa}^{\circ} - v)}{U_{t\kappa}^{\circ} - v}$$

$$Eq (V) \qquad Eq (\Lambda)$$

The graphics above illustrate the integration domains of x in Eqs. (\lor) and (\land). In general, the geometry of $\partial U_t \cap (U_{t/\kappa} + v)$ or $\partial U_t \setminus (\mathring{U}_{t\kappa} - v)$ is still difficult to handle, but in highly symmetric cases when U_t are concentric balls or cubes, Eqs. (\lor) and (\land) can be calculated efficiently.

Computing Robust Radius Eqs. (\lor) and (\land) allow us to compute the growth function by Eq. (NP). In general, this yields an *upper bound* of the robust radius

$$\sup \left\{ r : \sup_{\|v\|_p \le r} \mathcal{G}_q(1-\rho, v) < 1/2 \right\}$$

$$\leq \sup \left\{ r : \mathcal{G}_q(1-\rho, ru) < 1/2 \right\}$$

for any particular u with $||u||_p = 1$. With sufficient symmetry, e.g. with ℓ_2 adversary and distributions with spherical level sets, this upper bound becomes *tight* for well-chosen u, and we can build a lookup table of certified radii. See Algorithms 1 and 2.

Algorithm 1 Pre-Computing Robust Radius Table via Level Set Method for Spherical Distributions Againt ℓ_2 Adversary

Input: Radii $r_1 < \ldots < r_N$ Initialize $u = (1, 0, \ldots, 0) \in \mathbb{R}^d$. for i = 1 to N do Find κ s.t. $q(\mathcal{NP}_{\kappa} - r_i u) = 1/2$ (via Eq. (\wedge) or Theorem T.1) by binary search Compute $p_i \leftarrow q(\mathcal{NP}_{\kappa})$ via Eq. (\vee) or Theorem T.1 end for Output: $p_1 > \cdots > p_N$

Algorithm 2 Certification with Table				
Input: Probability of correct class ρ				
Output: Look up r_i where $p_i \ge 1 - \rho > p_{i+1}$				

4.2. The Differential Method

To derive certification (robust radius *lower bounds*) for more general distributions, we propose a *differential method*, which can be thought of as a vast generalization of the proof in Salman et al. (2019a) of the Gaussian robust radius. The idea is to compute the largest possible *infinitesimal increase in q-measure* due to an *infinitesimal adversarial perturbation*. More precisely, given a norm $\|\cdot\|$, and a smoothing measure q, we define

$$\Phi(p) \stackrel{\text{def}}{=} \sup_{\|v\|=1} \sup_{U \subseteq \mathbb{R}^d: q(U)=p} \lim_{r \searrow 0} \frac{q(U-rv)-p}{r}.$$
 (5)

Intuitively, one can then think of $1/\Phi(p)$ as the *smallest* possible perturbation in $\|\cdot\|$ needed to effect a unit of infinitesimal increase in p. Therefore,

Theorem 4.1 (Theorem G.6). *The robust radius in* $\| \cdot \|$ *is at least*

$$R \stackrel{\text{def}}{=} \int_{1-\rho}^{1/2} \frac{1}{\Phi(p)} \,\mathrm{d}p,$$

where ρ is the probability that the base classifier predicts the right label under random perturbation by q.

By exchanging differentiation and integration and applying a similar greedy reasoning as in the Neyman-Pearson lemma, $\Phi(p)$ can be derived for many distributions q and integrated symbolically to obtain expressions for R. We demonstrate the technique with a simple example below, but much of it can be automated; see Theorem G.6.

Example 4.2 (see Theorem N.1). If the smoothing distribution is $q(x) \propto \exp(-\|x\|_{\infty}/\lambda)$, then the robust radius against an ℓ_1 adversary is at least

$$R = 2d\lambda(\rho - 1/2),$$

when ρ is the probability of the correct class as in Theorem 4.1.

Randomized Smoothing of All Shapes and Sizes

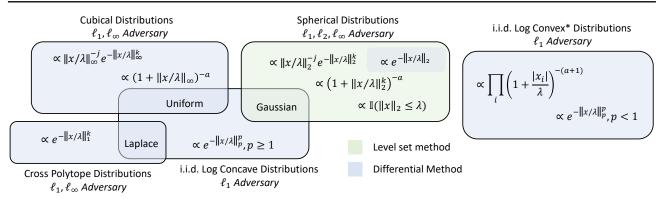


Figure 1. Smoothing distributions for which we derive robustness guarantees in this paper. Each box represents a family of distributions that obtain guarantees through similar proofs. Text beside each box indicates the name of the family and the ℓ_p adversaries against which we have guarantees. Log Convex* means log convex on the positive and negative half lines, but not necessarily on the whole line. The color indicates the basic technique used, among the two proposed techniques in this paper. We explicitly list example densities in each box. For the robust radii formulas, see Table A.1.

Proof Sketch. By linearity in λ , we WLOG assume $\lambda = 1$. By Theorem 4.1 and the monotonicity of Φ , it suffices to show that $\Phi(p) = 1/2d$ for $p \ge 1/2d$. For any fixed U with q(U) = p,

$$\lim_{r \searrow 0} \frac{q(U - rv) - p}{r} = \frac{d}{dr} \int_{U} q(x - rv) \, \mathrm{d}x \Big|_{r=0}$$
$$= \int_{U} \langle v, \nabla q(x) \rangle \, \mathrm{d}x.$$

Note $\nabla q(x) = e_x q(x)$, where $e_x = \operatorname{sgn}(x_{i^*})e_{i^*}$, e_i is the *i*th unit vector, and $i^* = \operatorname{argmax}_i |x_i|$. Additionally, the above integral is linear in v, so the supremum over $||v||_1 = 1$ is achieved on one of the vertices of the ℓ_1 ball. So we may WLOG consider only $v = \pm e_i$; furthermore, due to symmetry of $\nabla q(x)$, we can just assume $v = e_1$:

$$\Phi(p) = \sup_{U} \lim_{r \searrow 0} \frac{q(U - re_1) - p}{r} = \sup_{U} \int_{U} \langle e_1, e_x \rangle \, \mathrm{d}x,$$

where U ranges over all q(U) = p. Note $\langle e_1, e_x \rangle = 0$ if $i^* \neq 1$, and $\operatorname{sgn}(x_{i^*})$ otherwise. Thus, to maximize $\lim_{r \searrow 0} \frac{q(U-re_1)-p}{r}$ subject to the constraint that q(U) = p, we should put as much q-mass on those x with large $\langle e_1, e_x \rangle$. For $p \ge 1/2d$, we thus should occupy the entire region $\{x : \langle e_1, e_x \rangle = 1\}$, which has q-mass 1/2d, and then assign the rest of the q-mass (amounting to p - 1/2d) to the region $\{x : \langle e_1, e_x \rangle = 0\}$, which has q-mass 1 - 1/d. This shows that

$$\Phi(p) = 1/2d, \quad \forall p \in [1/2d, 1 - 1/2d]$$

as desired.

4.3. Comparison of the Two Methods and Prior Works

We summarize the distributions our methods cover in Fig. 1 and the bounds we derive in Table A.1. We highlight a few broadly applicable robustness guarantees: *Example* 4.3 (Theorem M.1). Let $\phi : \mathbb{R} \to \mathbb{R}$ be convex and even, and let CDF_{ϕ}^{-1} be the inverse CDF of the 1D random variable with density $\propto \exp(-\phi(x))$. If $q(x) \propto \prod_i e^{-\phi(x_i)}$, and ρ is the probability of the correct class, then the robust radius in ℓ_1 is

$$R = CDF_{\phi}^{-1}(\rho)$$

and this radius is *tight*. This in particular recovers the Gaussian bound of Cohen et al. (2019), Laplace bound of Teng et al. (2019), and Uniform bound of Lee et al. (2019) in the setting of ℓ_1 adversary.

Example 4.4 (Appendices N.1 and O.1). Facing an ℓ_1 adversary, cubical distributions, like that in Example 4.2, typically enjoy, via the differential method, ℓ_1 robust radii of the form

$$R = c(\rho - 1/2)$$

for some constant c depending on the distribution.

In general, the level set method always gives certificate as tight as Neyman-Pearson, while the differential method is tight only for infinitesimal perturbations, but can be shown to be tight for certain families, like in Example 4.3 above. On the other hand, the latter will often give efficiently evaluable symbolic expressions and apply to more general distributions, while the former in general will only yield a table of robust radii, and only for distributions whose level sets are sufficiently symmetric (such as a sphere or cube).

For distributions that are covered by both methods, we compare the bounds obtained and note that the differential and level set methods yield almost identical robustness certificates in high dimensions (e.g. number of pixels in CIFAR-10 or ImageNet images). See Appendix B.

Many earlier works used differential privacy or f-divergence methods to compute robust radii of smoothed models (Lecuyer et al., 2018; Li et al., 2019; Dvijotham et al.,

2019). In particular, Dvijotham et al. (2019) proposed a general f-divergence framework that subsumed all such works. Our robust radii are computed only from ρ ; Dvijotham et al. called this the "information-limited" setting, and we shall compare with their robustness guarantees of this type. While their algorithm in a certain limit becomes as good as Neyman-Pearson, in practice outside the Gaussian distribution, their robust radii are too loose. This is evident by comparing our baseline Laplace results in Table 1 with theirs, which are trained the same way. Additionally, our differential method often yields symbolic expressions for robust radii, making the certification algorithm easy to implement, verify, and run. Moreover, we derive robustness guarantees for many more (distributions, adversary) pairs (Fig. 1 and Table A.1). See Appendix E for a more detailed comparison.

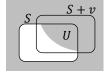
5. Wulff Crystals

A priori, it is a daunting task to understand the relationship between the adversary \mathcal{B} and the smoothing distribution q. In this section, we shall begin our investigation by looking at uniform distributions, and then end with an optimality theorem for all "reasonable" distributions.

Let q be the uniform distribution supported on a measurable set $S \subseteq \mathbb{R}^d$. WLOG, assume S has (Lebesgue) volume 1, Vol(S) = 1. Then for any $v \in \mathbb{R}^d$ and any $p \in [0, 1]$,

$$\mathcal{G}_q(p,v) = \min(1, p + \operatorname{Vol}((S+v) \setminus S))$$

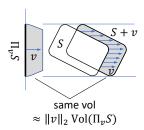
This can be seen easily by taking U in Eq. (2) to be a subset of $(S + v) \cap S$ with volume p (or any set of volume p containing $(S + v) \cap S$ if $p \ge \text{Vol}((S + v) \cap S))$ unioned with the complement of S. For example, in the figure here, U would be the gray region, if $U \cap S$ has volume p.



If S is convex, and we take v to be an infinitesimal translation, then the RHS above is infinitesimally larger than p, as follows:

$$\lim_{r \to 0} \frac{\mathcal{G}_q(p, rv) - p}{r} = \lim_{r \to 0} \frac{\operatorname{Vol}((S + rv) \setminus S)}{r}$$
$$= \|v\|_2 \operatorname{Vol}(\Pi_v S) \tag{6}$$

where $\Pi_v S$ is the projection of S along the direction $v/||v||_2$, and $\operatorname{Vol}(\Pi_v S)$ is its (d-1)dimensional Lebesgue measure. A similar formula holds when S is not convex as well (Eq. (13)). In the context of randomized smoothing, this means that the classifier q



smoothed by q is robust at x under a perturbation $\frac{\frac{1}{2}-p}{\|v\|_2 \operatorname{Vol}(\Pi_v S)} v$ when 1/2 - p is small, and p is the probability the base classifier f misclassifies $x + \delta$, $\delta \sim q$. Thus, for r small, we have

$$\sup_{v \in r\mathcal{B}} \mathcal{G}_q(p, v) \approx p + r \sup_{v \in \mathcal{B}} \|v\|_2 \operatorname{Vol}(\Pi_v S) = p + r \Phi(p),$$

with Φ as in Eq. (5). The smaller $\sup_{v \in \mathcal{B}} ||v||_2 \operatorname{Vol}(\Pi_v S)$ is, the more robust the smoothed classifier g is, for a fixed p. A natural question, then, is: among convex sets of volume 1,

which set S minimizes
$$\Phi = \sup_{v \in \mathcal{B}} ||v||_2 \operatorname{Vol}(\Pi_v S)$$
?

If \mathcal{B} is the ℓ_p ball, the reader might guess S should either be the ℓ_p ball or the ℓ_r ball with $\frac{1}{r} + \frac{1}{p} = 1$. It turns out the correct answer, at least in the case when \mathcal{B} is a highly symmetric polytope (e.g. $\ell_1, \ell_2, \ell_\infty$ balls), is a kind of *energy-minimizing* crystals studied in physics since 1901 (Wulff).

Definition 5.1. The *Wulff Crystal* (w.r.t. \mathcal{B}) is defined as the unit ball of the norm dual to $\|\cdot\|_*$, where $\|x\|_* = \mathbb{E}_{y \sim \operatorname{Vert}(\mathcal{B})} |\langle x, y \rangle|$ and y is sampled uniformly from the vertices of \mathcal{B}^{5} .

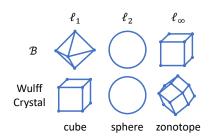
In fact, Wulff Crystals solve the more general problem without convexity constraint.

Theorem 5.2 (Theorem K.2, informal). *The Wulff Crystal* w.r.t. *B minimizes*

$$\Phi = \sup_{v \in \mathcal{B}} \lim_{r \to 0} r^{-1} \operatorname{Vol}((S + rv) \setminus S)$$

among all measurable (not necessarily convex) sets S of the same volume, when \mathcal{B} is sufficiently symmetric (e.g. $\ell_1, \ell_2, \ell_{\infty}$ balls).

When $\operatorname{Vert}(\mathcal{B})$ is a finite set, the Wulff Crystal has an elegant description as the *zonotope* of $\operatorname{Vert}(\mathcal{B})$, i.e. the Minkowski sum of the vertices of \mathcal{B} as vectors (Proposition J.4), from which we can derive the following examples. *Example* 5.3. The Wulff Crystal w.r.t. ℓ_2 ball is the ℓ_2 ball itself. The Wulff Crystal w.r.t. ℓ_1 ball is a cube (ℓ_{∞} ball). The Wulff Crystal w.r.t. ℓ_{∞} in 2 dimensions is a rhombus; in 3 dimensions, it is a rhombic dodecahedron; in higher dimension *d*, there is no simpler description of it other than the zonotope of the vectors $\{\pm 1\}^d$.



⁵Here, when \mathcal{B} is the ℓ_2 ball, $\operatorname{Vert}(\mathcal{B})$ is taken to be the entire boundary.

In fact, distributions with Wulff Crystal level sets more generally maximizs the robust radii for "hard" inputs.

Theorem 5.4 (Theorem L.2, informal). Let \mathcal{B} be sufficiently symmetric. Let q_0 be any distribution with a "reasonable"⁶ and even density function. Among all "reasonable" and even density functions q whose superlevel sets $\{x: q(x) \geq x\}$ t have the same volumes as those of q_0 , the quantity

$$\Phi(1/2) = \sup_{v \in \mathcal{B}} \sup_{q(U)=1/2} \lim_{r \searrow 0} \frac{q(U-rv) - 1/2}{r}$$

is minimized by the unique distribution q* whose superlevel sets are proportional to the Wulff Crystal w.r.t. B.

This theorem implies that distributions with Wulff Crystal level sets give the best robust radii for those hard inputs xthat a smooth classifier classifies correctly but only barely, in that the probability of the correct class $\rho = 1/2 + \epsilon$ for some small ϵ . The constraint on the volumes of superlevel sets indirectly controls the variance of the distribution. While this theorem says nothing about the robust radii for ρ away from 1/2, we find the Wulff Crystal distributions empirically to be highly effective, as we describe next in Section 6.

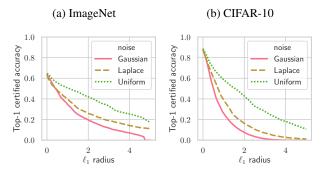
6. Experiments

We empirically study the performance of different smoothing distributions on standard image classification datasets, using the bounds derived via the level set or the differential method, and verify the predictions made by the Wulff Crystal theory. We follow the experimental procedure described by Cohen et al. (2019) and further works on randomized smoothing (Salman et al., 2019a; Li et al., 2019) using the ImageNet (Deng et al., 2009) and CIFAR-10 (Krizhevsky, 2009) datasets.

The certified accuracy at a radius ϵ is defined as the fraction of the test set for which the smoothed classifier g correctly classifies and certifies robust at an ℓ_p radius of ϵ . All results were certified with N = 100,000 samples and failure probability $\alpha = 0.001$.

For each distribution q, we train models across a range of scale parameter λ (see Table A.1), corresponding to the same range of noise variances $\sigma^2 \stackrel{\text{def}}{=} \mathbb{E}_{\delta \sim q}[\frac{1}{d} \|\delta\|_2^2]$ across different distributions. Then we calculate for each model the certified accuracies across the range of considered ϵ . Finally, in our plots, we present, for each distribution, the upper envelopes of certified accuracies attained over the range of considered σ . Further details of experimental procedures are described in Appendix D.

Figure 2. SOTA ℓ_1 Certified Accuracies. Certified ℓ_1 top-1 accuracies for ImageNet (left) and CIFAR-10 (right). For each distribution q, we train models across a range of $\sigma^2 \stackrel{\text{def}}{=} \mathbb{E}_{\delta \sim q}[\frac{1}{d} \|\delta\|_2^2]$, and at each level of ℓ_1 adversarial perturbation radius ϵ we report the best certified accuracy.



6.1. ℓ_1 Adversary

As previously mentioned, the Wulff Crystal for the ℓ_1 ball is a cube. With this motivation, we explore certified accuracies attained by distributions with cubical level sets.

- 1. Uniform, $\propto \mathbb{I}(||x||_{\infty} \leq \lambda)$
- 2. Exponential, $\propto ||x||_{\infty}^{-j} e^{-||x/\lambda||_{\infty}^{k}}$
- 3. Power law, $\propto (1 + ||x/\lambda||_{\infty})^{-a}$

We compare to previous state-of-the-art approaches using the Gaussian and Laplace distributions, as well as new noncubical distributions.

- 4. Exponential ℓ₁ (non-cubical), ∝ ||x||₁^{-j}e<sup>-||x/λ||^k₁
 5. Pareto i.i.d. (non-cubical), ∝ Π_i(1 + |x_i|/λ)^{-a}.
 </sup>

The relevant certified bounds are given in Table A.1.

We obtain state-of-the-art robust certificates for ImageNet and CIFAR-10, finding that the Uniform distribution performs best, significantly better than the Gaussian and Laplace distributions (Table 1, Fig. 2). The other distributions with cubic level sets match but do not exceed the performance of Uniform distribution, after sweeping hyperparameters. This verifies that distributions with cubical level sets are significantly better for ℓ_1 certified accuracy than those with spherical or cross-polytope level sets. Full results regarding the distributions not shown are available in Appendix C.

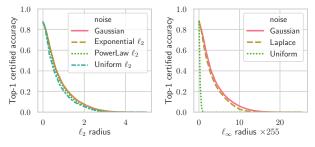
6.2. ℓ_2 Adversary

The Wulff Crystal w.r.t. the ℓ_2 ball is a sphere, so we explore distributions with spherical level sets (Table A.1):

- 1. Uniform, $\propto \mathbb{I}(\|x\|_2 \leq \lambda)$ 2. Exponential, $\propto \|x\|_2^{-j} e^{-\|x/\lambda\|_2^k}$

⁶*Reasonable* here roughly means Sobolev, i.e. has weak derivative that is integrable; for details see Theorem L.2.

Figure 3. CIFAR-10 certified accuracies for ℓ_2 (left) and ℓ_{∞} (right) adversaries. For each distribution q we train models across a range of $\sigma^2 \stackrel{\text{def}}{=} \mathbb{E}[\frac{1}{d} \|\delta\|_2^2]$, and at each level of ℓ_p adversarial perturbation radius ϵ , we pick the model that maximizes certified accuracy.



3. Power law, $\propto (1 + ||x/\lambda||_2)^{-a}$

We find these distributions perform similarly to, though do not surpass the Gaussian (Fig. 3, left).

6.3. ℓ_{∞} Adversary

The Wulff Crystal for the ℓ_{∞} ball is the zonotope of vectors $\{\pm 1\}^d$, which is a highly complex polytope hard to sample from and related to many open problems in polytope theory (Ziegler, 1995). However, we can note that it is approximated by a sphere with constant ratio (Proposition K.8), and in high dimension *d*, the sphere gets closer and closer to minimizing Φ (Theorem 5.2), but the cube and the cross polytope do not (Claim K.10). Accordingly, we find that distributions with spherical level sets outperforms those with cubical or cross polytope level sets in certifying ℓ_{∞} robustness (Fig. 3, right). In fact, in the next section we show that up to a dimension-independent factor, the Gaussian distribution is optimal for defending against ℓ_{∞} adversary if we don't use a more powerful technique than Neyman-Pearson.

7. No-Go Results for Randomized Smoothing

In this section, we demonstrate a theoretical limit for using Neyman-Pearson style bounds to derive meaningful randomized smoothing schemes for many norms, including ℓ_{∞} . In essence, we formalize an inherent tension between 1) having to have large enough noise variance to be robust and 2) having to have small enough noise variance to avoid trivializing the smoothed classifier. This tension becomes problematic against ℓ_p adversaries in high dimensions when p is large. This motivates the following definition:

Definition 7.1. Let $\|\cdot\|$ be a norm over \mathbb{R}^d , and let q be a smoothing distribution. We say that q satisfies (ε, L, ℓ) -*useful smoothing* with respect to $\|\cdot\|$ if:

- 1. (Robustness) For all x, y with $||x y|| < \varepsilon$, if $U \subseteq \mathbb{R}^d$ is any set satisfying $q(U x) \ge 1/2 + L\varepsilon$, then $q(U y) \ge 1/2$.
- 2. (Accuracy) For all x, y with $||x y|| \ge 1$, there exists $U \subseteq \mathbb{R}^d$ so that $|q(U x) q(U y)| > \ell$.

Finally, let $\kappa(q, \varepsilon, \|\cdot\|) = \min L/\ell$, minimizing over all L, ℓ such that q is (ε, L, ℓ) -useful.

We pause to interpret this definition. Recall that given a smoothing distribution q, a point $x \in \mathbb{R}^d$, and a binary base classifier $U \subseteq \mathbb{R}^d$ (identified wth its decision region), the smoothed classifier outputs $\operatorname{sgn}(\rho - 1/2)$ where $\rho = q(U - x)$ is the "confidence" of this prediction (Eq. (1)). Randomized smoothing (via Neyman-Pearson) tells us that, if ρ is large enough, then, no matter what U is, a small perturbation of x cannot decrease ρ too much to change $\operatorname{sgn}(\rho - 1/2)$ (Eq. (3)). In particular, for this guarantee to be valid, Condition (1), with some setting of parameters ε , L, ℓ , has to hold. Intuitively, if L is small, then Neyman-Pearson can certify large radius when ρ is moderately large (e.g. $\rho = 3/4$). But if q is not (ε, L, ℓ)-useful for any small L, then Neyman-Pearson cannot yield large radius unless $\rho \approx 1$ (e.g. $\rho = 0.999999$).

Condition (2) says that the resulting smoothing should not "collapse" points: in particular, if x, y are far in norm, then there should be some smoothed classifier $(q(U - \cdot))$ in Definition 7.1) that distinguishes them. If all we care about is robustness, then the optimal strategy would set q to be an arbitrarily wide distribution (say, e.g. a wide Gaussian), and the resulting smoothed classifier is roughly constant. Of course, such a smoothed classifier can never achieve good clean accuracy, so it is not useful. Condition (2) formalizes this notion. Note that it is a relatively weak assumption, for two reasons. First, for any x, y pair violating Condition (2), by linearity x', x' + (y - x) is a pair violating this condition, for all $x' \in \mathbb{R}^d$. Hence, any single violation of Condition (2) implies that there is an entire direction – namely, along the direction (y - x) — which the randomized smoothing mechanism collapses. Second, for Condition (2) to be satisfied, the U which distinguishes these two points can be completely arbitrary. Thus, if it is violated when $\ell = o(1)$, the two distributions are indistinguishable by any statistical test in high dimension, implying the impossibility of classification after smoothing.

Finally, the ratio L/ℓ measures the natural tradeoff: If $\kappa(q, \varepsilon, \|\cdot\|)$ is large, then the scheme is either not robust, or cannot attain good clean accuracy.

Randomized Smoothing as Metric Embedding A randomized smoothing scheme can be thought of as a mapping from a normed space supported on \mathbb{R}^d to the space of distributions, e.g. each point x is mapped to the distribution $q(\cdot - x)$. In this perspective, Definition 7.1 is roughly equivalent to a bi-Lipschitz condition on this mapping, where the target distributions are equipped with the total variation distance. Then the existence of a *useful* smoothing scheme is equivalent to whether $(\mathbb{R}^d, ||\cdot||)$ can be embedded *with low distortion* into the total variation space of distributions. Classical mathematics has a definitive answer to this question in the form a geometric invariant, called the *cotype*. **Definition 7.2** (see e.g. Wojtaszczyk (1996)). A normed space $T = (X, \|\cdot\|)$ is said to have *cotype* p for $2 \leq p \leq \infty$ if there exists C such that for all finite sequences $x_1, \ldots, x_n \in X$, we have $\mathbb{E}\left[\left\|\sum_{j=1}^n \sigma_j x_j\right\|\right] \leq C\left(\sum_{j=1}^n \|x_i\|^p\right)^{1/p}$, where the σ_j are independent Rademacher random variables. The smallest C which satisfies this constraint is denoted $C_p(T)$.

When the underlying space of the normed space T is \mathbb{R}^d , John's theorem (John, 1948) implies that any norm has cotype 2 with $C_2(T) \leq O(d^{1/2})$. Because C_2 lower bounds the distortion of a metric embedding of T, by the aforementioned connection with randomized smothing, C_2 also limits the usefulness of any smoothing scheme of T:

Theorem 7.3. Let $T = (\mathbb{R}^d, \|\cdot\|)$ for any norm $\|\cdot\|$ $\|$ over \mathbb{R}^d . For any distribution q, $\kappa(q, \varepsilon, \|\cdot\|) = \Omega(\min(1/\varepsilon, \varepsilon^3 C_2(T)^4)).$

In particular, it is well-known that $C_2((\mathbb{R}^d, \|\cdot\|_p)) = \Omega(d^{1/2-1/p})$. Then, setting $\varepsilon = \Omega(d^{-1/2+1/p})$, we get ⁷ **Corollary 7.4.** For any distribution q on \mathbb{R}^d and $\varepsilon = \Omega(d^{-1/2+1/p})$, $\kappa(q, \varepsilon, \|\cdot\|_p) = \Omega(\min(1/\varepsilon, d^{1/2-1/p}))$.

Moreover, we know that $\kappa(q, \varepsilon, \|\cdot\|_p) = O(d^{1/2-1/p})$ for all $p \ge 2$, by taking q to be a Gaussian. Thus, up to constants, Gaussian randomized smoothing is optimal, in the sense of Definition 7.1 (but distinguish from the optimality notion in Section 5), for all ℓ_p for $p \in [2, \infty]$.

Discussion Theorem 7.3 and Corollary 7.4 present a strong barrier to extending randomized smoothing to norms such as ℓ_{∞} . In words:

Without using more than the probability ρ of correctly classifying an input under random noise, no smoothing techniques can achieve nontrivial certified accuracy at ℓ_{∞} radius $\Omega(d^{-1/2})$.

However, we point out two ways to bypass this barrier. For one, more information about the base classifier can be collected to produce better robustness certificates. In fact, Dvijotham et al. (2019) proposed a "full-information" algorithm that computes many moments of the base classifier in a convex optimization procedure to improve certified radius, but it is 100 times slower than the "information-limited" algorithms we discuss here that use only ρ . It would be interesting to see whether this technique can be scaled up, and whether other methods can leverage more information⁸.

Another way to bypass our barrier is by going outside *translational* smoothing schemes, where any input x is made

noisy by *adding* some $\delta \sim q$, for a fixed q. This is the form covered by our no-go result. Instead of associating the distribution $q(\cdot - x)$ to every input x, one could consider more general distributions that are not just translations of some q. For example, Levine & Feizi (2019a) feeds a randomly sampled subset of pixels in the image into the base classifier to obtain certified defense against ℓ_0 perturbations. In other words, this mechanism associates to each point x the uniform distribution over the random subsets of x's pixels, and our no-go result does not apply to this scheme. However, to the best of our knowledge, no scheme for ℓ_{∞} certification has been devised that is not translational.

Finally, we formulated our no-go result in the setting of binary classification, and it is not clear whether a similarly strong barrier applies for multi-class classification. However, current techniques for certification only look at the two most likely classes, and separately reason about how much each one can change by perturbing the input. Our no-go result then straightforwardly applies to this case as well.

8. Conclusion

In this work, we have showed how far we can push randomized smoothing with different smoothing distributions against different ℓ_p adversaries, by presenting two new techniques for deriving robustness guarantees, by elucidating the geometry connecting the noise and the norm, and by empirically achieving state-of-the-art in ℓ_1 provable defense. At the same time, we have showed the limit current techniques face against ℓ_p adversaries when p > 2, especially ℓ_{∞} . Our results point out ways to bypass this barrier, by either leveraging more information about the base classifier or by designing non-translational smoothing schemes. We wish to investigate both directions in the future.

More broadly, randomized smoothing is a method for inducing stability in a mechanism while maintaining utility — precisely the bread and butter of differential privacy. We suspect our methods for deriving robustness guarantees here and for optimizing the noise distribution can be useful in that setting as well, where Laplace and Gaussian noise dominate the discussion. Whereas previous work Lecuyer et al. (2018) has applied differential privacy tools to randomized smoothing, we hope to go the other way around in the future.

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⁷In order for the robustness guarantee to be useful, ε has to be large anyway, so the reader need not be alarmed at this condition.

⁸ Lee et al. (2019) also used the decision tree structure of their base classifier to improve ℓ_0 certification, but the ℓ_0 -adversary does not fall within our framework.

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A. Table of Certified Radii

Distribution	Density	Adv.	Certified radius	Reference
iid Log Concave	$\propto e^{-\sum_i \phi(x_i)}$	ℓ_1	$\mathrm{CDF}_{\phi}^{-1}(\rho)$	Theorem M.1
iid Log Convex*	$\propto e^{-\sum_i \phi(x_i)}$	ℓ_1	$\int_{\varphi^{-1}(1-\rho)}^{\infty} \frac{1}{e^{\phi(c)-\phi(0)}-1} \mathrm{d}c \qquad \qquad \text{for } \varphi, \text{see} \to$	Theorem M.2
Exp. $\ell_p, p \ge 1$	$\propto e^{-\ rac{x}{\lambda}\ _p^p}$	ℓ_1	$\lambda \sqrt[p]{\text{GammaCDF}^{-1}(2\rho - 1; 1/p)}$	Corollary M.3
Exp. $\ell_p, p < 1$	$\propto e^{-\ rac{x}{\lambda}\ _p^p}$	ℓ_1	$\lambda \int_{\varphi^{-1}(1-\rho)}^{\infty} \frac{\mathrm{d}c}{e^{c^p}-1} \qquad \text{for } \varphi, \text{see} \to$	Theorem M.2
Gaussian	$\propto e^{-\ rac{x}{\lambda}\ _2^2/2}$	ℓ_2	$\lambda Gaussian CDF^{-1}(\rho; 0, 1)$	Theorem H.1 ^C
		ℓ_1	$\lambda \text{GaussianCDF}^{-1}(\rho; 0, 1)$	Symmetry
		ℓ_{∞}	$\lambda \text{GaussianCDF}^{-1}(\rho; 0, 1) / \sqrt{d}$	Symmetry
Laplace	$\propto e^{-\ \frac{x}{\lambda}\ _1}$	ℓ_1	$-\lambda \log(2(1-\rho))$	Theorem I.1 ^T
		ℓ_∞	$\approx \lambda \text{GaussianCDF}^{-1}(\rho; 0, 1) / \sqrt{d}$ see \rightarrow	Theorem P.2
Exp. ℓ_{∞}	$\propto e^{-\ rac{x}{\lambda}\ _{\infty}}$	ℓ_1	$2d\lambda(\rho-\frac{1}{2})$	Theorem N.1
	ll an ll	ℓ_∞	$\lambda \log \frac{1}{2(1-\rho)}$	Theorem N.4
Exp. ℓ_2	$\propto e^{-\ \frac{x}{\lambda}\ _2}$	ℓ_2	$\lambda(d-1) \operatorname{arctanh}(1-2\beta^{-1} \left(1-\rho; \frac{d-1}{2}, \frac{d-1}{2}\right)) \\ \lambda(d-1) \operatorname{arctanh}(1-2\beta^{-1} \left(1-\rho; \frac{d-1}{2}, \frac{d-1}{2}\right))$	Theorem R.1
		ℓ_1	$\lambda(d-1) \operatorname{arctanh}(1-2\beta^{-1}(1-\rho;\frac{d-1}{2},\frac{d-1}{2}))$	Symmetry
		ℓ_∞	$\frac{\lambda(d-1)}{\sqrt{d}}\operatorname{arctanh}(1-2\beta^{-1}\left(1-\rho;\frac{d-1}{2},\frac{d-1}{2}\right))$	Symmetry
Uniform ℓ_{∞}	$\propto \mathbb{I}(\ x\ _{\infty} \le \lambda)$	ℓ_1	$2\lambda(ho-rac{1}{2})$	Theorem N.3 ^L
		ℓ_∞	$2\lambda(1-\sqrt[d]{rac{3}{2}- ho})$	Theorem N.6 ^L
Uniform ℓ_2	$\propto \mathbb{I}(\ x\ _2 \le \lambda)$	ℓ_2	$ \begin{array}{l} \lambda \left(2 - 4\beta^{-1} \left(\frac{3}{4} - \frac{\rho}{2}; \frac{d+1}{2}, \frac{d+1}{2} \right) \right) \\ \lambda \left(2 - 4\beta^{-1} \left(\frac{3}{4} - \frac{\rho}{2}; \frac{d+1}{2}, \frac{d+1}{2} \right) \right) \\ \frac{\lambda}{\sqrt{d}} \left(2 - 4\beta^{-1} \left(\frac{3}{4} - \frac{\rho}{2}; \frac{d+1}{2}, \frac{d+1}{2} \right) \right) \end{array} $	Theorem S.1
		ℓ_1	$\lambda \left(2 - 4\beta^{-1} \left(\frac{3}{4} - \frac{\tilde{\rho}}{2}; \frac{d+1}{2}, \frac{d+1}{2} \right) \right)$	Symmetry
		ℓ_∞	$\frac{\lambda}{\sqrt{d}}\left(2-4\beta^{-1}\left(\frac{3}{4}-\frac{\rho}{2};\frac{d+1}{2},\frac{d+1}{2}\right)\right)$	Symmetry
General Exp. ℓ_∞	$\propto \ rac{x}{\lambda}\ _{\infty}^{-j}e^{-\ rac{x}{\lambda}\ _{\infty}^{k}}$	ℓ_1	$\frac{2d\lambda}{d-1}\Gamma\left(\frac{d-j}{k}\right)/\Gamma\left(\frac{d-1-j}{k}\right)\left(ho-\frac{1}{2}\right)$	Theorem N.2
		ℓ_∞	$\lambda \int_{1-\rho}^{1/2} \frac{1}{\Phi(p)} dp$ for Φ , see \rightarrow	Theorem N.5
General Exp. ℓ_2	$\propto \ \frac{x}{\lambda}\ _2^{-j} e^{-\ \frac{x}{\lambda}\ _2^k}$	ℓ_2	level set method	Example T.2
General Exp. ℓ_1	$\propto e^{-\ rac{x}{\lambda}\ _1^k}$	ℓ_1	$\lambda \int_{1-\rho}^{1/2} \frac{R}{\Psi(p)} dp$ for R, Ψ , see \rightarrow	Theorem P.1
		ℓ_∞	$\lambda \int_{1-\rho}^{1/2} \frac{1}{\Phi(p)} dp$ for Φ , see \rightarrow	Theorem P.3
Power Law ℓ_∞	$\propto \frac{1}{(1+\ \frac{x}{\lambda}\ _{\infty})^a}$	ℓ_1	$rac{2d\lambda}{a-d}\left(ho-rac{1}{2} ight)$	Theorem O.1
		ℓ_∞	$\frac{2\lambda}{a-d} \int_{1-\rho}^{1/2} \frac{\mathrm{d}p}{\Upsilon(\Upsilon^{-1}(2p;d,a-d);d,a+1-d)}$	Theorem O.2
Power Law ℓ_2	$\propto \frac{1}{(1+\ \frac{x}{\lambda}\ _2^k)^a}$	ℓ_2	$a-a$ $(1-\rho)$ $(1-(2p;a,a-a);a,a+1-a)$ level set method	Example T.3
Pareto (i.i.d.)	$\propto \frac{1}{\prod_{i} \left(1 + \frac{ x_i }{\lambda}\right)^{a+1}}$	ℓ_1	$\lambda \frac{2\rho - 1}{a} _2F_1\left(1, \frac{a}{a+1}, \frac{2a+1}{a+1}; (2\rho - 1)^{1+1/a}\right)$	Theorem Q.1

Table A.1. Distributions we derive bounds for and assess experimentally. Here ρ is the probability the base classifier answers correctly when input is perturbed by the smoothing noise, d is the dimensionality of the noise, CDF_{ϕ}^{-1} is the inverse CDF of the 1D random variable with density $\propto e^{-\phi(x)}$, $\beta^{-1}(\cdot; a, b)$ is the inverse Beta CDF function with shape parameters a and b, $\Upsilon(\cdot; a, b)$ (resp. $\Upsilon^{-1}(\cdot; a, b)$) is the Beta Prime (resp. inverse) CDF function with shape parameters a and b, Γ is the Gaussian hypergeometric function. Under *Reference*, superscript ^C refers to Cohen et al. (2019), superscript ^L refers to Lee et al. (2019), and superscript ^T refers to Teng et al. (2019).

Distribution q	Density	$\lambda/\sigma = \lambda/\sqrt{\mathbb{E}_{\delta \sim q} \frac{1}{d} \ \delta\ ^2}$
Exp. ℓ_p	$\propto e^{-\ rac{x}{\lambda}\ _p^p}$	$\sqrt{\frac{1}{\Gamma(3/p)/\Gamma(1/p)}}$
Gaussian	$\propto e^{-\ \frac{x}{\lambda}\ _2^2/2}$	1
Laplace	$\propto e^{-\ \frac{x}{\lambda}\ _1}$	$1/\sqrt{2}$
Exp. ℓ_∞	$\propto e^{-\ rac{x}{\lambda}\ _{\infty}}$	$\sqrt{\frac{1}{(d+1)((d-1)/3+1)}}$
Exp. ℓ_2	$\propto e^{-\ \frac{x}{\lambda}\ _2}$	$\sqrt{\frac{1}{d+1}}$
Uniform ℓ_{∞}	$\propto \mathbb{I}(\ x\ _{\infty} \le \lambda)$	$\sqrt{3}$
Uniform ℓ_2	$\propto \mathbb{I}(\ x\ _2 \le \lambda)$	$\sqrt{\frac{1}{d+2}}$
General Exp. ℓ_∞	$\propto \ \frac{x}{\lambda}\ _{\infty}^{-j} e^{-\ \frac{x}{\lambda}\ _{\infty}^{k}}$	$\sqrt{\frac{d\Gamma(\frac{d-j}{k})}{((d-1)/3+1)\Gamma(\frac{d+2-j}{k})}}$
General Exp. ℓ_2	$\propto \ \frac{x}{\lambda}\ _2^{-j} e^{-\ \frac{x}{\lambda}\ _2^k}$	$\sqrt{\frac{d\Gamma(\frac{d-j}{k})}{\Gamma(\frac{d+2-j}{k})}}$
General Exp. ℓ_1	$\propto e^{-\ \frac{x}{\lambda}\ _1^k}$	$\sqrt{rac{d(d+1)\Gamma(rac{d}{k})}{2\Gamma(rac{d+2}{k})}}$
Power Law ℓ_∞	$\propto \frac{1}{(1+\ rac{x}{\lambda}\ _{\infty})^a}$	$\sqrt{\frac{(a-d-1)(a-d-2)}{(d+1)((d-1)/3+1)}}$
Power Law ℓ_2	$\propto rac{1}{(1+\ rac{x}{\lambda}\ _2^k)^a}$	$\sqrt{\frac{(a-d-1)(a-d-2)}{d+1}}$
Pareto (i.i.d.)	$\propto rac{1}{\prod_i \left(1+rac{ x_i }{\lambda} ight)^{a+1}}$	$\sqrt{\frac{1}{2}(a-1)(a-2)}$

Table A.2. Relation between the scale parameter λ and the variance $\sigma^2 = \mathbb{E}_{\delta \sim q} \frac{1}{d} \|\delta\|^2$ of each distribution. This table is used to choose the correct λ to match σ across different distributions. All quantities can be computed easily using Lemmas U.3 and U.4. They are also tested to be numerically correct in the test suite of our code base github.com/tonyduan/rs4a.

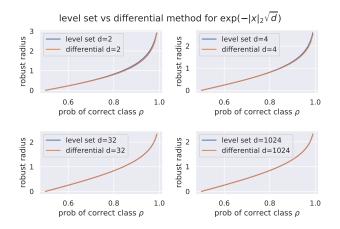


Figure B.1. **Differential Method is Tight** for practical purposes in high dimension *d*.

B. Level Set Method vs Differential Method

Here we concretely compare the robust radii obtained from the level set method and those obtained from the differential method for the distribution $\exp(-||x||_2\sqrt{d})$, for various input dimensions d (we scale the distributions this way so each coordinate has size $\Theta(1)$). For convenience, here's the robust radius from the differential method (Theorem R.1):

$$R = \frac{d-1}{\sqrt{d}} \operatorname{arctanh} \left(1 - 2\operatorname{BetaCDF}^{-1} \left(1 - \rho; \frac{d-1}{2}, \frac{d-1}{2} \right) \right).$$

The robust radii from level set method are computed as in Theorem T.1, and they are tight. As we see in Fig. B.1, the differential method is *very slightly* loose in low dimensions d = 2 and 4, but in high dimensions d = 32 or 1024, the robust radii obtained from both methods are indistinguishable.

C. Additional Experimental Results

All results in this section are described for CIFAR-10.

 ℓ_1 Adversary In addition to the Gaussian, Laplace, and Uniform distributions, we considered an Exponential distribution with cubic level sets, an Exponential distribution with ℓ_1 level sets, a power law distribution with cubic level sets, and an i.i.d Pareto distribution.

$$q_{\text{Exp}_{\infty}}(x) \propto \exp(-(\|x/\lambda\|_{\infty}^{k}))$$
$$q_{\text{Exp}_{1}}(x) \propto \exp(-(\|x/\lambda\|_{1}^{k}))$$
$$q_{\text{Power}_{\infty}}(x) \propto (1 + \|x/\lambda\|_{\infty})^{-a}$$
$$q_{\text{Pareto}}(x) \propto \prod_{i} \left(1 + \frac{|x_{i}|}{\lambda}\right)^{-(a+1)}$$

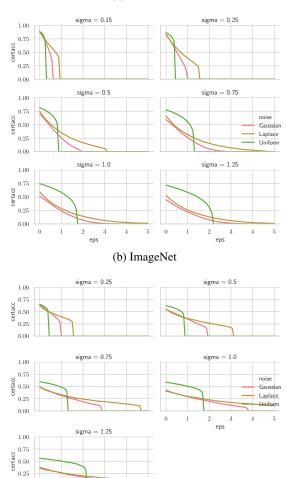


Figure C.1. Certified Accuracy per σ . Certified accuracies of ResNet-18 against an ℓ_1 adversary at each level of ϵ , across the range of σ with which models were trained (we omit $\sigma > 1.25$ for brevity). The upper envelope for each distribution is taken to be the maximum certified accuracy across values of σ .

3 4

2 eps

0.00

0

Results for these experiments are shown in Fig. C.2. The suffix of the noises in the legend denotes the value of the shape parameter k or a that was chosen (whereas we fixed shape parameter j = 0). We note that results for distributions with cubical level sets match but do not exceed that of the Uniform distribution. Meanwhile distributions without cubical level sets do not match performance of the Uniform distribution. This suggests that the tail behavior of the noise does not matter as much as the shape of level sets.

Ablation of Our ℓ_1 Improvement over Previous SOTA To understand how much of our ℓ_1 results come from improved certification vs improved training performance, we repeated our ResNet experiments with a multi-layer perceptron (MLP) and AlexNet. We find that the Uniform

(a) CIFAR-10

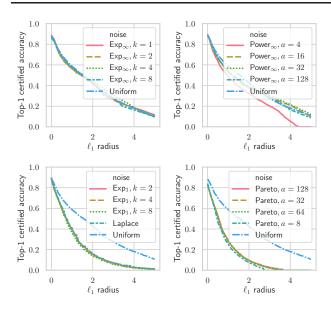


Figure C.2. More Distributions for ℓ_1 Adversary. CIFAR-10 certified top-1 accuracies of ResNet-18 against the ℓ_1 adversary, on generalized exponential law (with ℓ_{∞} and ℓ_1 level sets), power law (with ℓ_{∞} level sets), and Pareto distributions. After appropriate hyperparameter search (k or a), distributions with cubic level sets achieve performance roughly matching that of the Uniform distribution.

distribution attains a higher upper envelope of certified accuracy than Gaussian or Laplace with this model (Fig. C.3), but the improvement is less dramatic compared to Table 1 and Fig. 2. Interestingly, the clean (i.e. $\epsilon = 0$) training and testing accuracy of all three distributions are identical when fixed to the same level of σ for the fully-connected model, but for AlexNet, the Uniform noise allows much higher accuracies (Fig. C.4), and for ResNet, even more so. This training improvement leads to substantial improvement in certified accuracies (Fig. C.5).

As an additional visualization, when we plot the certified accuracy at fixed ϵ s versus the training accuracy of a ResNet on noise-augmented CIFAR-10, the Uniform distribution can be seen to significantly outperform the Gaussian and Laplace distributions at all training accuracies except those very close to 1 (Fig. C.6).

So while some of the improvement in certified accuracy in Fig. 2 is due to improved certified radius per ρ , it seems much more of it is due to the difference in how well a classifier trains when smoothed by noise.

Why Does Uniform Distribution Get Better Training Accuracy? Here we further investigate why improvement in architecture seems to amplify the advantage of uniform distribution over others, in terms of training accuracy for each level of σ . Letting $W \in \mathbb{R}^{d,d}$ denote a pre-specified rotation matrix fixed throughout training/testing, we consider:

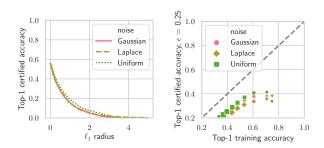


Figure C.3. **Multi-layer Perceptron.** (Left) CIFAR-10 certified top-1 accuracies for the ℓ_1 adversary, with a multi-layer perceptron. (Right) Certified accuracies at ℓ_1 perturbation $\epsilon = 0.25$ plotted against training accuracy under smoothing noise.

- 1. Smoothing with unmodified noise, rotated images:
 - $x \leftarrow Wx + \delta, \quad \delta \sim q.$
- 2. Smoothing with rotated noise, unmodified images:

$$x \leftarrow x + W\delta, \quad \delta \sim q$$

Note that certification bounds are no longer necessarily applicable, so we only compare clean training accuracy i.e. whether $\arg \max_{\mathcal{Y}} g(x) = y$. Results for ResNet are shown in Fig. C.7. We find that the difference in training performance still exists (but to a lesser degree) under alternative (1), smoothing with unmodified noise but rotated images. On the other hand, we find this difference vanishes under alternative (2), smoothing with rotated noise and unmodified images.

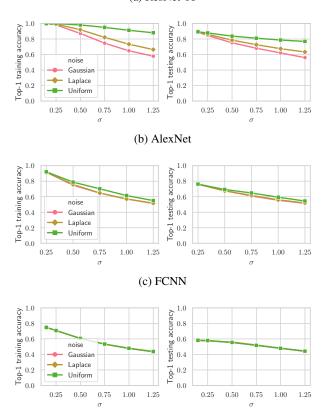
This suggests that the improvement of training accuracy under Uniform noise is due to some *synergy* of the model architecture with the data distribution and the smoothing noise. The choice of Uniform distribution induces some improvement in training accuracy but this is greatly amplified by the interaction between convolution layers and the image dataset. Thus, a good noise for randomized smoothing seems to be one that balances its robustness properties with its *compatibility* with the architecture and the data.

 ℓ_2 Adversary In addition to the Gaussian distribution, we considered an Exponential distribution with spherical level sets and a power law distribution with spherical level sets.

$$q_{\text{Exp}_2}(x) \propto (\|x\|_2/\lambda)^{-j} \exp(-(\|x\|_2^2/\lambda))$$
$$q_{\text{Power}_2}(x) \propto (1+\|x\|_2/\lambda)^{-a}$$

Results for these experiments are shown in Fig. C.8. After appropriate hyperparameter search (of j and a), performance for both distributions with spherical level sets matches that of the Gaussian.

Figure C.4. **Effect of Architecture.** Clean CIFAR-10 training (left) and testing (right) accuracies for ResNet-18, AlexNet, and a fully connected neural network, at fixed levels of $\mathbb{E}[\frac{1}{d} ||\delta||_2^2] \stackrel{\text{def}}{=} \sigma^2$. For fixed σ , there is no difference between the distributions when smoothing a fully connected network, but differences arise when the architecture improves to AlexNet and ResNet.



(a) ResNet-18



Training Methods There are three methods of training a smoothed classifier. Let f denote the base classifier, q denote the smoothing distribution, and consider an observation (x, y).

1. Noise augmentation as in Cohen et al. (2019),

$$\mathcal{L}(x,y) = -\log f(x+\delta)_y, \quad \delta \sim q$$

2. Directly training the smoothed classifier as described in Salman et al. (2019a),

$$\mathcal{L}(x,y) = -\log \mathbb{E}[f(x+\delta)]_y, \quad \delta \sim q$$

3. Adversarial training as in Salman et al. (2019a),

$$\mathcal{L}(x,y) = -\log f(\tilde{x})_y, \quad \tilde{x} \text{ found via PGD.}$$

In all experiments we trained with the first option, appropriate noise augmentation. We found that direct training was slower and did not yield superior performance in practice.

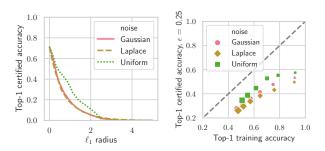


Figure C.5. AlexNet. (Left) CIFAR-10 certified top-1 accuracies for the ℓ_1 adversary, with an AlexNet architecture. (Right) Certified accuracies at $\epsilon = 0.25$, plotted against training accuracy under noise.

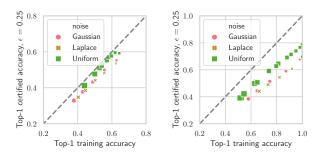


Figure C.6. Training Accuracy vs Certified Accuracy. Top-1 ℓ_1 certified accuracies for ImageNet (left) and CIFAR-10 (right) at pre-specified $\epsilon = 0.25$, controlling for fixed training accuracy. Larger sized points denote larger σ . Predictably, as σ increases, training and certified accuracy decreases. At fixed training accuracy, the Uniform distribution significantly outperforms Gaussian and Laplace.

Range of σ Recall that $\sigma^2 \stackrel{\text{def}}{=} \mathbb{E}[\frac{1}{d} ||\delta||_2^2]$. This is a fairly consistent measurement of noise level across different noise distributions, and is a natural control variate for comparing the effect (e.g. training, testing, and certified accuracies) of different noises. In addition, to obtain a good estimate of the upper envelope of certified accuracy, we need to take the pointwise maximum of the radius-vs-certified-accuracy curve (such as those in Fig. 2) for many σ s. In this work, we swept over:

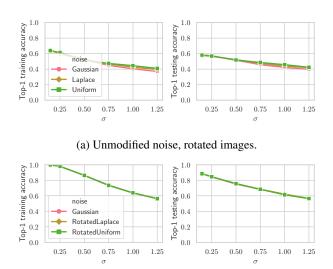
$$\sigma \in \{0.15, 0.25, 0.5, 0.75, 1.0, 1.25, 1.50\}.$$

For distributions with cubic level sets, we needed to sweep over larger σ s as well to estimate the large-radius portion of the upper envelope better:

 $\sigma \in \{1.75, 2.0, 2.25, 2.5, 2.75, 3.0, 3.25, 3.5\}.$

Table A.2 lists for each distribution the conversion constant needed to obtain λ from $\sigma = \sqrt{\mathbb{E}_{\delta \sim q} \frac{1}{d} \|\delta\|_2^2}$.

Certified Accuracy per σ In Fig. C.1 we show the certified accuracies of Gaussian, Laplace and Uniform distribu-



(b) Rotated noise, unmodified images.

Figure C.7. **Rotation Experiments.** ResNet-18 clean training/testing accuracies in the two rotation experiments.

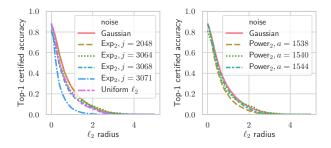


Figure C.8. Distributions with Spherical Level Sets. CIFAR-10 certified top-1 accuracies of ResNet-18 against the ℓ_2 adversary, on spherical level set exponential and power law distributions. After appropriate hyper-parameter search, performance matches that of the Gaussian distribution.

tions, for each σ , for both ImageNet and CIFAR-10. The upper envelopes reported in the main text are defined as the maximum certified accuracies over σ .

Experiment Hyperparameters For all experiments we trained with a cosine-annealed learning rate of 0.1, optimized by stochastic gradient descent with momentum of 0.9 and weight decay of 0.0001.

For ImageNet experiments we used a ResNet-50 model and trained with a batch size of 64 for 30 epochs.

For CIFAR-10 experiments we used a ResNet-18 model and trained with a batch size of 128 for 120 epochs.

Ablation studies with a fully connected neural network employed two hidden layers of 2048 and 512 nodes followed by ReLU activations, trained with a learning rate of 0.01. To compute the top categories for certification (which used N = 100,000 samples), we used 64 samples.

Our code is publicly available at: github.com/tonyduan/rs4a

E. In-Depth Comparison with Dvijotham et al. (2019)

The information-limited certification algorithm in Dvijotham et al. (2019) relaxes the optimization problem

$$\sup_{v \in \mathcal{B}} \mathcal{G}_q(p, v) = \sup_{q' \in q_{\mathcal{B}}} \sup_{U:q(U)=p} q'(U)$$
$$\leq \sup_{q' \in \mathcal{D}_F(q)} \sup_{U:q(U)=p} q'(U)$$
(7)

enlarging the set of shifted distributions $q_{\mathcal{B}} \stackrel{\text{def}}{=} \{q(\cdot - v) : v \in \mathcal{B}\}$ to the set of distributions close to q in several fdivergences $\mathcal{D}_F \stackrel{\text{def}}{=} \{q' : \mathcal{D}_f(q'||q) \le \epsilon_f, \forall f \in F\}$, for a set of functions F. Dvijotham et al. (2019) showed that when F consists of all Hockey-Stick divergences, Eq. (7) becomes tight,, but in practice this is not feasible. In fact, Dvijotham et al. (2019) admits themselves that

It turns out that the Renyi and KL divergences are computationally attractive for a broad class of smoothing measures, while the Hockey-Stick divergences are theoretically attractive as they lead to optimal certificates in the information-limited setting. However, Hockey-Stick divergences are harder to estimate in general, so we only use them for Gaussian smoothing measures.

Concretely, the looseness of their relaxation can be observed when comparing our baseline Laplace results (Table 1) with theirs.

Operationally, their algorithm proceeds as follows

1. For each distribution q and function f, manually find the f-divergence "ball" that contains $\{q(\cdot - v) : v \in B\}$, i.e. compute $\{\epsilon_f\}_{f \in F}$ such that

$$\{q(\cdot - v) : v \in \mathcal{B}\} \subseteq \{q' : \mathcal{D}_f(q' \| q) \le \epsilon_f, \forall f \in F\}.$$

2. Then they relax the original certification problem to the certification of all q' close to q in f-divergence, i.e. they solve Eq. (7) for the ϵ_f found in the previous step.

The 2nd step is a straightforward low-dimensional convex optimization problem, but the trickiness of the 1st step limits the distributions they can apply their technique to. For example, they only know how to do step 1 for $\exp(-||x||_p)$ against ℓ_p adversary, but not against ℓ_r for $r \neq p$; in contrast, our differential method computes robust radii for Laplace against ℓ_{∞} perturbation.

F. Mathematical Preliminaries

In this section, we rigorously define several mathematical notions and their properties that will recurrent throughout what follows. We will be brief here, but readers can skip this on first reading and refer back when necessary.

Note about measure and integration We will use Vol to denote measure, typically Hausdorff measure with the dimension implicit from context. When integrating over a measurable set, the underlying measure is also typically the Hausdorff measure as well.

Sobolev Functions and Regular Functions

Definition F.1. Let $\Omega \subseteq \mathbb{R}^d$ be an open set. For the purpose of this paper, we say a function $f : \Omega \to \mathbb{R}$ is *regular* if 1) f is continuous at all $x \in \Omega$ except for a possible closed set of measure 0, and if 2) f is Sobolev, $f \in W^{1,1}(\Omega)$ (see Definition F.2).

Condition 2) roughly means that f has a weak derivative ∇f such that both f and ∇f are integrable. We can replace condition 1) with the weaker condition of *precise representability* (Maly et al., 2002), but we are satisfied with the easily interpretable condition 1).

Definition F.2. For any open set $\Omega \subseteq \mathbb{R}^d$, the *Sobolev space* $W^{k,p}(\Omega)$ is defined as the functions $f \in L^p(\Omega)$ whose mixed weak derivatives (in the sense of distributions) exist up to order k and are all in $L^p(\Omega)$, i.e.

$$W^{k,p}(\Omega) \stackrel{\text{def}}{=} \{ f \in L^p(\Omega) : D^{\alpha} f \in L^p(\Omega), \forall |\alpha| \le k \}.$$

Coarea Formula and the Weak Sard's Theorem

Theorem F.3 (Coarea Formula (Maly et al., 2002)). Let $\Omega \subseteq \mathbb{R}^d$ be an open set, $g \in L^1(\Omega)$, and $f : \mathbb{R}^d \to \mathbb{R}$ regular in the sense of Definition F.1. Then

$$\int_{\Omega} g(x) \|\nabla f(x)\|_2 \, \mathrm{d}x = \int_{\mathbb{R}} \int_{f^{-1}(y)} g(x) \, \mathrm{d}x \, \mathrm{d}y$$

Here the integral in x in the RHS is over the (d-1)dimensional Hausdorff measure of $f^{-1}(y)$.

The regularity of f can be replaced by weaker conditions (see Maly et al. (2002)), but the statement here suffices for our purposes.

By setting g in Theorem F.3 to be the indicator function over the set where the gradient of f vanishes, we get

Theorem F.4 (Weak Sard). For any $f : \mathbb{R}^d \to \mathbb{R}$ regular in the sense of Definition F.1, let $Z \stackrel{\text{def}}{=} \{x \in \mathbb{R}^d : \nabla f(x) = 0\}$. Then

$$\operatorname{Vol}(Z \cap f^{-1}(y)) = 0$$

for almost every $y \in \mathbb{R}$. Here Vol denote the Hausdorff measure of dimension d - 1.

Sobolev Functions and Absolute Continuity Recall the standard definition of *absolute continuity*, which can be thought of as a more general notion of *differentiability*.

Definition F.5. A function $f : \mathbb{R} \to \mathbb{R}$ is called *absolute continuous* if there exists a Lebesgue integrable function $g : \mathbb{R} \to \mathbb{R}$ and some $a \in \mathbb{R}$, such that

$$f(x) = f(a) + \int_{a}^{x} g(t) \,\mathrm{d}t.$$

Such an f has derivative f' almost everywhere, and f' coincides with g almost everywhere.

Sobolev functions are known to be absolutely continuous on every line, and this property roughly captures all Sobolev functions.

Theorem F.6 (ACL Property of Sobolev Functions (Nikodym, 1933)). Let $\Omega \subseteq \mathbb{R}^d$ be an open set. The following statements hold.

- Let f: Ω → ℝ be Sobolev, f ∈ W^{1,p}(Ω). Then possibly after modifying f on a set of measure 0, for every u ∈ ℝ^d, the function t ↦ f(x + tu) is absolutely continuous for almost every x. Furthermore, the (classical) directional derivative D_uf is in L^p(Ω) for every u.
- Conversely, if the restriction of a function $f : \Omega \to \mathbb{R}$ on almost every line parallel to the coordinate axes is absolutely continuous, then pointwise gradient ∇f exists almost everywhere, and $f \in W^{1,p}(\Omega)$ as long as $f, \nabla f \in L^p(\Omega)$.

The ACL property of Sobolev functions yields the differentiability of the convolution of a L^∞ and a $W^{1,1}$ function.

Lemma F.7. If a function q is in $W^{1,1}(\mathbb{R}^d)$, then for every bounded measurable $F : \mathbb{R}^d \to \mathbb{R}$, the convolution F * q is differentiable, and

$$\nabla(F * q) = F * (\nabla q).$$

Proof. A function is differentiable if all of 1) its partial derivatives exist and 2) are continuous. First we show that, for any vector u, $F * (D_u q) = D_u (F * q)$. We can compute as follows.

$$F * (D_u q)(x)$$

$$= \int F(\hat{x}) D_u q(\hat{x} - x) d\hat{x}$$

$$= \frac{d}{d\tau} \int_0^\tau \int F(\hat{x}) D_u q(\hat{x} - x + tu) d\hat{x} dt \Big|_{\tau=0}$$
(8)

$$= \left. \frac{d}{d\tau} \int F(\hat{x}) \int_0^\tau D_u q(\hat{x} - x + tu) \, \mathrm{d}t \, \mathrm{d}\hat{x} \right|_{\tau=0} \tag{9}$$

$$= \left. \frac{d}{d\tau} \int F(\hat{x}) [q(\hat{x} - x + \tau u) - q(\hat{x} - x)] \,\mathrm{d}\hat{x} \right|_{\tau=0}$$
(10)

$$= D_u \int F(\hat{x})q(\hat{x} - x) \,\mathrm{d}\hat{x} = D_u(F * q)(x). \tag{11}$$

In these equations, first note that

$$\begin{split} &\int_0^\tau \int |F(\hat{x}) D_u q(\hat{x} - x + tu)| \, \mathrm{d}\hat{x} \, \mathrm{d}t \\ &\leq \tau \int |D_u q(\hat{x})| \, \mathrm{d}\hat{x} < \infty, \end{split}$$

by the ACL property of q (Theorem F.6). Thus, in Eq. (8), we introduced $\frac{d}{d\tau} \int_0^{\tau}$ innocuously by the fundamental theorem of calculus, since the inner integral is absolutely integrable in t. Then, in Eq. (9), we applied Fubini-Tonelli Theorem to swap the order of integration. In Eq. (10), we integrated out the directional derivative $D_u q$ for almost every \hat{x} where $t \mapsto q(\hat{x} - x + tu)$ is absolutely continuous. Finally, in Eq. (11), we simplified the integral by noting that $q(\hat{x} - x)$ does not depend on τ , and $q(\hat{x} - x + \tau u)$ is absolutely integrable in \hat{x} . This proves our claim that $F * (D_u q) = D_u (F * q)$ for any $u \in \mathbb{R}^d$.

Note additionally that, since $D_u q \in L^1$ (by assumption) and $F \in L^{\infty}$, their convolution $F * (D_u q)$ is bounded and continuous.

Then, taking u to be the coordinate vectors, we see the partial derivatives of F * q all exist and are continuous, proving our lemma.

G. The Differential Method

We summarize the setup of this section in the following assumption. Here we use a notion called *regularity* introduced in Definition F.1 that roughly says that a function needs to continuous almost everywhere and be "weakly" differentiable, and it and its gradient are both absolutely integrable. All concrete density functions we work with in this paper will be regular.

Assumption G.1. Let $F : \mathbb{R}^d \to [0,1]$ be a measurable function and let $G : \mathbb{R}^d \to [0,1]$ be the smoothing of F by the distribution $q(x) \propto \exp(-\psi(x))$ for some $\psi : \mathbb{R}^d \to \mathbb{R}$, such that q is regular in the sense of Definition F.1. Formally,

$$G(x) = \mathop{\mathbb{E}}_{\delta \sim q} F(x+\delta) = \int q(\delta)F(x+\delta) \,\mathrm{d}\delta$$
$$= \int q(\hat{x}-x)F(\hat{x}) \,\mathrm{d}\hat{x}.$$

Consider a norm $\|\cdot\|$ with unit ball \mathcal{B} that is a convex body. Let $\operatorname{Vert}(\mathcal{B})$ be the set of its extremal points.

Example G.2. If $\|\cdot\| = \|\cdot\|_1$ is the ℓ_1 -norm, then \mathcal{B} is what is called the cross-polytope, defined as the convex hull of the unit vectors and their negations. If $\|\cdot\| = \|\cdot\|_{\infty}$ is the ℓ_{∞} -norm, then \mathcal{B} is the cube with vertices $\{\pm 1\}^d$. If $\|\cdot\| = \|\cdot\|_2$ is the ℓ_2 -norm, then \mathcal{B} is the unit sphere, and Vert(\mathcal{B}) is its entire boundary.

Example G.3. If $\psi(x) = ||x||_2^2$, then q is the standard Gaussian distribution. If $\psi(x) = ||x||_1$, then q is the Laplace distribution.

The following definition of Φ turns out to be equivalent to Eq. (5), which will be apparent in the proof of Theorem G.6. It gives a systematic way of computing Φ .

Definition G.4. Let $q(x) \propto \exp(-\psi(x))$ be a distribution over \mathbb{R}^d as in Assumption G.1. For any vector $u \in \mathbb{R}^d$, let γ_u be the random variable $\langle u, \nabla \psi(\delta) \rangle \in \mathbb{R}$ with $\delta \sim q$. Define φ_u to be the complementary CDF of γ_u ,

$$\varphi_u(c) \stackrel{\text{def}}{=} \mathbb{P}[\gamma_u > c],$$

define the inverse complementary CDF $\varphi_u^{-1}(p)$ of γ_u to be

$$\varphi_u^{-1}(p) \stackrel{\text{def}}{=} \inf\{c : \mathbb{P}[\gamma_u > c] \le p\},\$$

For any $p \in [0, 1]$, define a new random variable $\gamma_u^{(p)}$ by

$$\gamma_{u}^{(p)} = \begin{cases} \gamma_{u}|_{(c,\infty)} & \text{with probability } \varphi_{u}(c) \\ c & \text{with probability } p - \varphi_{u}(c) \\ 0 & \text{with probability } 1 - p, \end{cases}$$

where $c \stackrel{\text{def}}{=} \varphi_u^{-1}(p)$ and $\gamma_u|_{(c,\infty)}$ is the random variable γ_u conditioned on $\gamma_u > c$. Roughly speaking, the PDF of $\gamma_u^{(p)}$ allocates probability p to the right portion of γ_u 's PDF, and puts the rest 1 - p probability on 0. One just needs to be careful when γ_u 's measure has a singular point at $\varphi_u^{-1}(p)$, which is dealt with in the middle line above.

Let \mathcal{B} be the unit ball of $\|\cdot\|$ as in Assumption G.1. Then we define $\Phi: [0,1] \to \mathbb{R}$ by

$$\Phi(p) \stackrel{\text{def}}{=} \max_{u \in \operatorname{Vert}(\mathcal{B})} \mathbb{E} \, \gamma_u^{(p)}.$$

Remark G.5. The function $p \mapsto \overline{\mathbb{E}}\gamma_u^{(p)}$ in Definition G.4 is increasing on $[0, \varphi_u(0)]$ and nonincreasing on $[\varphi_u(0), 1]$. Thus $\Phi(p)$ is also increasing on $[0, \inf_{u \in \operatorname{Vert}(\mathcal{B})} \varphi_u(0)]$.

The following theorem is the master theorem for applying the differential method. We illustrate its usage to recover the known Gaussian (Cohen et al., 2019) and Laplace (Teng et al., 2019) bounds as warmups in Appendices H and I before applying the technique at scale.

Theorem G.6 (The Differential Method). As in Assumption G.1, fix any norm $\|\cdot\|$ and let $G : \mathbb{R}^d \to [0,1]$ be the smoothing of any measurable $F : \mathbb{R}^d \to [0,1]$ by $q(x) \propto \exp(-\psi(x))$, such that q is regular in the sense of Definition F.1. Let $\Phi : [0,1] \to \mathbb{R}$ be given as in Definition G.4.

Then for any x, if G(x) < 1/2, then $G(x + \delta) < 1/2$ for any

$$\|\delta\| < \int_{G(x)}^{1/2} \frac{1}{\Phi(p)} \,\mathrm{d}p.$$
 (*)

In Theorem G.6, one should think of G(x) as the probability that the smoothed classifier assigns to *any class other than* *the* correct one. So Theorem G.6 says that, if the smoothed classifier predicts the correct class (G(x) < 1/2), then it continues to do so even when the input is perturbed by a noise with magnitude bounded by Eq. (*).

Sometimes, when φ_u is continuous for all u, for $p \in [0, 1/2]$, we can factor

$$\Phi(p) = \bar{\varphi}_u(\varphi_u^{-1}(p)), \quad \text{where} \quad \bar{\varphi}_u(c) \stackrel{\text{def}}{=} \underset{\gamma_u}{\mathbb{E}} \gamma_u \mathbb{I}(\gamma_u > c)$$

for some specific $u \in Vert(\mathcal{B})$, either due to symmetry in the vertices of \mathcal{B} (so that it doesn't matter which u it is) or because a specific u maximizes the expression for all $p \in [0, 1/2]$. Then the following lemma is very useful for simplifying the integral in Theorem G.6. It can be proved easily using change of coordinates.

Lemma G.7. Suppose $\Phi(p) = \overline{\varphi}(\varphi^{-1}(p))$ on $p \in [0, 1/2]$, where $\varphi(p)$ is differentiable and both φ and $\overline{\varphi}$ are nonincreasing. Then for any $0 \le p_0 \le 1/2$,

$$\int_{p_0}^{1/2} \frac{1}{\Phi(p)} \, \mathrm{d}p = \int_{\varphi^{-1}(1/2)}^{\varphi^{-1}(p_0)} \frac{|\varphi'(c)|}{\bar{\varphi}(c)} \, \mathrm{d}c$$

Finally, as mentioned before, the proof of Theorem G.6 will show that

Proposition G.8. The definition of Φ in Definition G.4 coincides with the definition Eq. (5) for any smoothing distribution q with regular density function supported everywhere in \mathbb{R}^d .

Proof of Theorem G.6. Consider a path $\xi_t : [0, ||\delta||] \to \mathbb{R}^d$ given by $\xi_0 = x$, $\xi_{||\delta||} = x + \delta$, and $\xi'_t = d\xi_t/dt = \delta/||\delta||$. We will show

$$dG(\xi_t)/dt \le \Phi(G(\xi_t))$$

and apply Lemma G.9 to yield the desired result.

By chain rule,

$$dG(\xi_t)/dt = \xi'_t \cdot \nabla G(\xi_t) = \frac{\delta}{\|\delta\|} \cdot \nabla G(\xi_t).$$

To upper bound this quantity, we relax

$$\frac{\delta}{\|\delta\|} \cdot \nabla G(\xi_t) \le \max_{u \in \mathcal{B}} u \cdot \nabla G(\xi_t) = \max_{u \in \operatorname{Vert}(\mathcal{B})} u \cdot \nabla G(\xi_t)$$

where \mathcal{B} is the unit ball of the norm $\|\cdot\|$, and the equality is because $u \cdot \nabla G(x)$ is linear in u, so optima are achieved on vertices. Therefore, it suffices to show that,

$$\forall u \in \operatorname{Vert}(\mathcal{B}), x \in \mathbb{R}^d, \quad u \cdot \nabla G(x) \le \Phi(G(x)).$$
 (12)

Below, we let x be any vector in \mathbb{R}^d (not just those satisfyiing $G(x) \leq 1/2$ as in the theorem statement). In general,

for any vector u and any $x \in \mathbb{R}^d$, the directional derivative $u \cdot \nabla G(x)$ of G(x) in the direction of u is given by

$$u \cdot \nabla G(x) = u \cdot \int \nabla_x q(\hat{x} - x) F(\hat{x}) \, \mathrm{d}\hat{x}$$
$$= \int \langle u, \nabla \psi(\hat{x} - x) \rangle q(\hat{x} - x) F(\hat{x}) \, \mathrm{d}\hat{x}$$

where we used Lemma F.7 and the assumption that q is regular. Then

$$\begin{split} u \cdot \nabla G(x) &= \mathop{\mathbb{E}}_{\delta \sim q} F(x+\delta) \langle u, \nabla \psi(\delta) \rangle \\ &\leq \sup_{\hat{F}: \hat{G}(x) = G(x)} \mathop{\mathbb{E}}_{\delta \sim q} \hat{F}(x+\delta) \langle u, \nabla \psi(\delta) \rangle, \end{split}$$

where we vary over all $\hat{F} : \mathbb{R}^d \to [0, 1]$ such that its smoothing \hat{G} has the same value as G at x. While at first glance, this seems like a unwieldy quantity to maximize, there's a simple intuition to find the maximizing \hat{F} :

Imagine $\hat{F}(x + \cdot)$ as some allocation of mass in \mathbb{R}^d that amounts to G(x) under the measure q. When we vary \hat{F} , we are allowed to shuffle this mass around while keeping its qmeasure equal to G(x), as long as $0 \leq \hat{F} \leq 1$. To maximize $\mathbb{E}_{\delta \sim q} \hat{F}(x+\delta) \langle u, \nabla \psi(\delta) \rangle$, we then need to allocate as much q-measure as possible toward regions where $\langle u, \nabla \psi(\cdot) \rangle$ is large.

In other words, the maximizing $\hat{F},$ which we denote as $\hat{F}^*,$ is

$$\hat{F}^*(x+\delta) = \begin{cases} 1 & \text{if } \langle u, \nabla \psi(\delta) \rangle > \varphi_u^{-1}(G(x)) \\ 0 & \text{else,} \end{cases}$$

if $\mathbb{P}[\langle u, \nabla \psi(\delta) \rangle = \varphi_u^{-1}(G(x))] = 0$, where φ_u^{-1} is the inverse complementary CDF of the random variable $\gamma_u = \langle u, \nabla \psi(\delta) \rangle$ (with randomness induced by $\delta \sim q$), as defined in Definition G.4. If there is a singular point at $\varphi_u^{-1}(G(x))$, i.e. $\mathbb{P}[\langle u, \nabla \psi(\delta) \rangle = \varphi_u^{-1}(G(x))] > 0$, then we choose a subset of $U \subseteq \{\delta : \langle u, \nabla \psi(\delta) = \varphi_u^{-1}(G(x))\}$ with *q*-measure $\mathbb{P}[\delta \in U] = G(x) - \varphi_u(\varphi_u^{-1}(G(x)))$, and define \hat{F}^* as

$$\hat{F}^*(x+\delta) = \begin{cases} 1 & \text{if } \langle u, \nabla \psi(\delta) \rangle > \varphi_u^{-1}(G(x)) \text{ or } \delta \in U \\ 0 & \text{else.} \end{cases}$$

Then

$$\mathbb{E}_{\delta \sim q} \hat{F}(x+\delta) \langle u, \nabla \psi(\delta) \rangle$$

$$\leq \mathbb{E}_{\delta \sim q} \hat{F}^*(x+\delta) \langle u, \nabla \psi(\delta) \rangle = \mathbb{E} \gamma_u^{G(x)}$$

where $\gamma_u^{(p)}$ is the random variable defined in Definition G.4.

Finally, putting everything together,

$$\begin{aligned} \max_{u \in \operatorname{Vert}(\mathcal{B})} u \cdot \nabla G(x) \\ &\leq \max_{u \in \operatorname{Vert}(\mathcal{B})} \sup_{\hat{F}: \hat{G}(x) = G(x)} \mathbb{E}_{\delta \sim q} \hat{F}(x+\delta) \langle u, \nabla \psi(\delta) \rangle \\ &= \max_{u \in \operatorname{Vert}(\mathcal{B})} \mathbb{E} \gamma_u^{G(x)} \\ &= \Phi(G(x)) \end{aligned}$$

by the definition of Φ in Definition G.4. This shows Eq. (12) and consequently the theorem as well.

Lemma G.9. Consider a function p_t differentiable in $t \in [0, \infty)$. Suppose $0 < p_0 \le 1/2$, and

$$dp_t/dt \leq \Phi(p_t)$$

for some function $\Phi : (0, \infty) \to \mathbb{R}^+$ taking only positive values, Then $p_T < 1/2$ as well for any

$$T < \int_{p_0}^{1/2} \frac{1}{\Phi(p)} \,\mathrm{d}p$$

Proof. WLOG, we can assume that $dp_t/dt > 0$ for all $t \in [0, \infty)$. Thus, p_t is increasing in t, and there exists a differentiable inverse function t(p) that expresses the time t that $p_t = p$. We then have $dt(p)/dp = \frac{1}{dp_t/dt} = \frac{1}{\Phi(p)}$, and for any $\epsilon \ge 0$,

$$t(1/2 - \epsilon) = t(1/2 - \epsilon) - t(p_0)$$

= $\int_{p_0}^{1/2 - \epsilon} \frac{dt(p)}{dp} dp = \int_{p_0}^{1/2 - \epsilon} \frac{1}{\Phi(p)} dp.$

Since this integral is continuous in ϵ , there is an $\epsilon^* > 0$ such that

$$t(1/2 - \epsilon^*) = \int_{p_0}^{1/2 - \epsilon^*} \frac{1}{\Phi(p)} \,\mathrm{d}p = T.$$

Therefore $p_T = 1/2 - \epsilon^* < 1/2$, as desired.

H. Example: Gaussian against ℓ_2 Adversary

We give a quick example of recovering the tight Gaussian bound of Cohen et al. (2019) using the differential method.

In this section, we set the norm $\|\cdot\|$ to be the ℓ_2 norm $\|x\|_2 = \sqrt{\sum_{i=1}^d x_i^2}$. Then \mathcal{B} is just the unit ball, and its "vertices" are just the points on the unit sphere. Additionally, we let q be the Gaussian measure

$$q(x) \propto \exp(-\|x\|_2^2/2)$$

so that $\psi(x) = \|x\|_2^2/2$ and $\nabla \psi(x) = x$

Below, let GaussianCDF be the CDF of the standard Gaussian in 1D.

Theorem H.1. Suppose *H* is a smoothed classifier smoothed by the Gaussian distribution

$$q(x) \propto \exp(-\|x\|_2^2/2\sigma^2),$$

such that $H(x) = (H(x)_1, \ldots, H(x)_C)$ is a vector of probabilities that H assigns to each class $1, \ldots, C$. If H correctly predicts the class y on input x, and the probability of the correct class is $\rho \stackrel{\text{def}}{=} H(x)_y > 1/2$, then H continues to predict the correct class when x is perturbed by any η with

$$\|\eta\|_2 < \sigma \text{GaussianCDF}^{-1}(\rho).$$

Proof. By linearity in σ , it suffices to show this for $\sigma = 1$. For brevity, let us denote GaussianCDF in this proof by Ψ .

We seek to apply Theorem G.6 to $G(x) = 1 - H(x)_y$, for which we need to derive random variables γ_u and $\gamma_u^{(p)}$, and most importantly, the function Φ .

For any $u \in \operatorname{Vert}(\mathcal{B})$ (i.e. any unit vector u), $\gamma_u = \langle u, \nabla \psi(\delta) \rangle = \langle u, \delta \rangle, \delta \sim q$, is a standard Gaussian random variable (in \mathbb{R}). Therefore, for $p \in [0, 1]$, the random variable $\gamma_u^{(p)}$ defined in Definition G.4 is just

$$\gamma_u^{(p)} = \begin{cases} 0 & \text{with prob. } 1-p \\ \mathcal{N}(0,1)|_{[c,\infty)} & \text{with prob. } p, \end{cases}$$

where $c \stackrel{\text{def}}{=} \Psi^{-1}(1-p)$, and $\mathcal{N}(0,1)|_{[c,\infty)}$ is a standard Gaussian z conditioned on $z \ge c$. Thus, for any $u \in \text{Vert}(\mathcal{B})$,

$$\begin{split} \Phi(p) &= \mathbb{E} \, \gamma_u^{(p)} = \mathop{\mathbb{E}}_{z \sim \mathcal{N}(0,1)} z \mathbb{I}(z \ge c) = \frac{1}{\sqrt{2\pi}} e^{-z^2/2} \bigg|_{\infty}^c \\ &= \Psi'(c) = \Psi'(\Psi^{-1}(1-p)). \end{split}$$

Then, by setting G(x) in Theorem G.6 to be $1 - H(x)_y = 1 - \rho$, we get the provably robust radius of

$$\int_{1-\rho}^{1/2} \frac{1}{\Phi(p)} dp = \int_{1-\rho}^{1/2} \frac{1}{\Psi'(\Psi^{-1}(1-p))} dp$$
$$= \int_{0}^{\Psi^{-1}(\rho)} dc = \Psi^{-1}(\rho),$$

as desired.

I. Example: Laplace against ℓ_1 Adversary

Let us give another quick example of recovering the tight Laplace bound of Teng et al. (2019) using the differential method.

In this section, we set the norm $\|\cdot\|$ to be the ℓ_1 norm $\|x\|_1 = \sum_{i=1}^d |x_i|$. Then the unit ball \mathcal{B} is the convex hull

of its vertices which are the coordinates vectors and their negations:

$$\operatorname{Vert}(\mathcal{B}) = \{\pm e_i : i \in [d]\}.$$

Consider the Laplace distribution

$$q(x) \propto \exp(-\|x\|_1) \quad \text{so that}$$

$$\psi(x) = \|x\|_1 \quad \text{and} \quad \nabla \psi(x) = (\operatorname{sgn}(x_1), \dots, \operatorname{sgn}(x_d)),$$

with $\nabla \psi(x)$ defined whenever all x_i s are nonzero.

Theorem I.1. Suppose *H* is a smoothed classifier smoothed by the Laplace distribution

$$q(x) \propto \exp(-\|x\|_1/\lambda)$$

such that $H(x) = (H(x)_1, \ldots, H(x)_C)$ is a vector of probabilities that H assigns to each class $1, \ldots, C$. If H correctly predicts the class y on input x, and the probability of the correct class is $\rho \stackrel{\text{def}}{=} H(x)_y > 1/2$, then H continues to predict the correct class when x is perturbed by any η with

$$\|\eta\|_1 < \lambda \log \frac{1}{2(1-\rho)}$$

Proof. By linearity in λ , it suffices to show this for $\lambda = 1$.

We seek to apply Theorem G.6 to $G(x) = 1 - H(x)_y$, for which we need to derive random variables γ_u and $\gamma_u^{(p)}$, and most importantly, the function Φ .

For any $u \in \operatorname{Vert}(\mathcal{B})$ (i.e. $u = \pm e_i$), $\gamma_u = \langle u, \nabla \psi(\delta) \rangle, \delta \sim q$, is a Rademacher random variable that takes values 1 and -1 with equal probability. Therefore, for $p \in [0, 1/2]$, the random variable $\gamma_u^{(p)}$ defined in Definition G.4 is

$$\gamma_u^{(p)} = \begin{cases} 1 & \text{with prob. } p \\ 0 & \text{with prob. } 1 - p \end{cases}$$

Thus, for any $u \in Vert(\mathcal{B})$,

$$\Phi(p) = \mathbb{E}\,\gamma_u^{(p)} = p.$$

Then, by setting G(x) in Theorem G.6 to be $1 - H(x)_y = 1 - \rho$, we get the provably robust radius of

$$\int_{1-\rho}^{1/2} \frac{1}{\Phi(p)} \, \mathrm{d}p = \int_{1-\rho}^{1/2} \frac{1}{p} \, \mathrm{d}p = \log \frac{1}{2(1-\rho)}.$$

J. Wulff Crystal

The following is an intuitive statement of the main isoperimetric property of Wulff Crystals. **Theorem J.1** (Isoperimetric property of Wulff Crystals (Brothers & Morgan, 1994), informal statement). Let $\|\cdot\|$ be any norm on \mathbb{R}^n . Let B be the Wulff Crystal of $\|\cdot\|$, i.e. the unit ball of the norm $\|\cdot\|_*$ dual to $\|\cdot\|_.$

$$B = \{x : \|x\|_* \le 1\}.$$

Let Ω be any measurable subset of \mathbb{R}^n with finite surface area and of the same volume as B. Then with $\mathbf{n}(\Omega, x)$ denoting the normal vector at x with respect to Ω , normalized to have ℓ_2 norm 1,

$$\int_{\partial\Omega} \|\mathbf{n}(\Omega, x)\| \, \mathrm{d}x \ge \int_{\partial B} \|\mathbf{n}(B, x)\| \, \mathrm{d}x$$

with equality holding if and only if Ω differs from a translate of B by a set of volume zero.

This statement carries across the core essence of the isoperimetry, and is a rigorous statement if Ω is restricted to have smooth boundary, but care needs to be taken to explain the concept of "surface area," "normal vector," the "boundary $\partial \Omega$," and the boundary measure on $\partial \Omega$, in the context of general, measurable Ω . We refer the interested reader to (Brothers & Morgan, 1994) for such mathematical details.

J.1. Wulff Crystals are Zonotopes

In this paper, the norm $\|\cdot\|$ referred to in Theorem J.1 will take the form of $\|x\| = \mathbb{E}_v |\langle x, v \rangle|$ where v is sampled from some distribution, as in Definition 5.1. When v is sampled uniformly from some finite set of vectors S, then the Wulff Crystal $B = \{x : \|x\|_* \le 1\}$ is proportional to the *Zonotope* of S (McMullen, 1971).

Definition J.2 (Zonotope). Given a finite collection of vectors S, the zonotope Zon(S) is defined as the Minkowski sum of the vectors of S, i.e.

$$\operatorname{Zon}(\mathcal{S}) \stackrel{\text{def}}{=} \left\{ \sum_{v \in \mathcal{S}} a_v v : a_v \in [0, 1], \forall v \in \mathcal{S} \right\}$$

The zonotope can be viewed as a linear projection of the cube $[0, 1]^S$ sending each unit vector to a vector of S.

Example J.3. If \mathcal{B} is the ℓ_1 unit ball, then $\text{Zon}(\text{Vert}(\mathcal{B}))$ is a cube. If \mathcal{B} is the ℓ_{∞} unit ball, then in 2 dimensions, $\text{Zon}(\text{Vert}(\mathcal{B}))$ is a rhombus; in 3 dimensions, it is rhombic dodecahedron; in higher dimensions, there is no simpler description of the resulting polytope.

Proposition J.4. The Wulff Crystal w.r.t. \mathcal{B} is equal to the zonotope $\frac{2}{|\operatorname{Vert}(\mathcal{B})|} \operatorname{Zon}(\operatorname{Vert}(\mathcal{B})).$

The volume of a zonotope, and thus of Wulff Crystals, can be computed easily using the following formula.

Proposition J.5. Let S be a finite set of vectors in \mathbb{R}^d . Then the d-dimensional volume of Zon(S) is given by

$$\sum_{\mathcal{T}\subseteq \mathcal{S}: |\mathcal{T}|=d} |\mathrm{Vol}(\mathrm{Zon}(\mathcal{T}))| = \sum_{\mathcal{T}\subseteq \mathcal{S}: |\mathcal{T}|=d} |\det \mathcal{T}|,$$

where det T is the determinant of the square matrix with vectors of T as columns.

K. Wulff Crystals Yield Optimal Uniform Distributions for Randomized Smoothing

In this section, we will formulate Theorem 5.2 rigorously and prove it.

Definition K.1. Let S be a finite set of vectors in \mathbb{R}^d and let G be the group of linear transformations that permute S (i.e. G is S's linear symmetry group). We say S is symmetric if G acts on S transitively, i.e. for any two elements $v, w \in S$, there is a group element $g \in G$ such that $g \cdot v = w$.

For example, the boolean cube $\{\pm 1\}^d$ is symmetric, and so is the set of coordinate vectors and their negations. The following is the main theorem of this section, stating the optimality of uniform distributions supported Wulff Crystals.

Theorem K.2. If \mathcal{B} is a full-dimensional polytope in \mathbb{R}^d symmetric around the origin, and whose vertices form a symmetric set, then the Wulff Crystal w.r.t. \mathcal{B} minimizes

$$\sup_{v \in \mathcal{B}} \lim_{r \to 0} r^{-1} \operatorname{Vol}((S + rv) \setminus S)$$

among all measurable, not necessarily convex, sets $S \subseteq \mathbb{R}^d$ of the same volume. In other words, among uniform distributions supported on measurable sets of volume 1, the one supported on the Wulff Crystal minimizes the maximal instantaneous growth $\Phi(p)$ in the measure of a set due to an instantaneous perturbation from \mathcal{B} .

Proof. By Lemma K.13, we have

$$\lim_{r \to 0} r^{-1} \operatorname{Vol}((S + rv) \setminus S) = \int_{\partial S} \Theta(\langle \mathbf{n}(x), v \rangle) \, \mathrm{d}x \quad (13)$$

where $\mathbf{n}(x)$ is the normal at x w.r.t. S, and $\Theta(x) = \max(0, x)$. Note this quantity is convex in v because Θ is convex. Then

$$\sup_{v \in \mathcal{B}} \lim_{r \to 0} r^{-1} \operatorname{Vol}((S + rv) \setminus S)$$

$$= \sup_{v \in \operatorname{Vert}(\mathcal{B})} \lim_{r \to 0} r^{-1} \operatorname{Vol}((S + rv) \setminus S)$$

$$\geq \mathbb{E} \lim_{v \sim \operatorname{Vert}(\mathcal{B})} \lim_{r \to 0} r^{-1} \operatorname{Vol}((S + rv) \setminus S) \quad (14)$$

$$= \mathbb{E} \int_{\partial S} \bigcup_{v \sim \operatorname{Vert}(\mathcal{B})} \Theta(\langle \mathbf{n}(x), v \rangle) \, \mathrm{d}x$$

$$= \int_{\partial S} \mathbb{E} \Theta(\langle \mathbf{n}(x), v \rangle) \, \mathrm{d}x$$

Since $\mathcal{B} = -\mathcal{B}$ and thus $\operatorname{Vert}(\mathcal{B}) = -\operatorname{Vert}(\mathcal{B})$,

$$\|w\| \stackrel{\text{def}}{=} \mathop{\mathbb{E}}_{v \sim \operatorname{Vert}(\mathcal{B})} \Theta(\langle w, v \rangle) = \frac{1}{2} \mathop{\mathbb{E}}_{v \sim \operatorname{Vert}(\mathcal{B})} |\langle w, v \rangle|$$

is a seminorm. This is in fact a norm because $Vert(\mathcal{B})$ spans \mathbb{R}^d , by the assumption that \mathcal{B} is full-dimensional. Then, plugging $\|\cdot\|$ into $\|\cdot\|$ in Theorem J.1, we get that the Wulff Crystal B w.r.t. \mathcal{B} minimizes

$$B = \operatorname*{argmin}_{S: \operatorname{Vol}(S) = \operatorname{Vol}(B)} \mathbb{E} \lim_{v \sim \operatorname{Vert}(\mathcal{B})} r^{-1} \operatorname{Vol}((S + rv) \setminus S).$$

Now note that the norm above is invariant under the transpose of \mathcal{B} 's symmetry group: for any linear symmetry g of $Vert(\mathcal{B})$,

$$\|g^{\top}w\| = \frac{1}{2} \mathop{\mathbb{E}}_{v} |\langle g^{\top}w, v\rangle| = \frac{1}{2} \mathop{\mathbb{E}}_{v} |\langle w, gv\rangle| = \|w\|.$$

This invariance translates to the dual norm $\|\cdot\|_*$'s invariance under the symmetry group itself. Thus the Wulff Crystal B, being the unit ball of the dual norm, is itself invariant under the symmetry group of $Vert(\mathcal{B})$. By assumption, this symmetry group acts transitively on $Vert(\mathcal{B})$, so B "looks the same" from the angle of every $v \in Vert(\mathcal{B})$, i.e.

$$\operatorname{Vol}((B+rw)\setminus B) = \operatorname{Vol}((B+rv)\setminus B)$$

for any $w, v \in Vert(\mathcal{B})$. Consequently, Eq. (14) holds with equality, and B minimizes the supremum in question as well.

K.1. Growth Calculations for Standard Shapes

Using the fact that the volume of the *d*-dimensional unit ball is $\pi^{d/2}\Gamma(d/2+1)^{-1}$, and the volume of the standard *d*-dimensional cross polytope is $2^d/d!$, as well as the identity

$$\lim_{r \to 0} r^{-1} \operatorname{Vol}((S + rv) \setminus S) = \|v\|_2 \operatorname{Vol}(\Pi_v S)$$

if S is convex, we can derive the following facts easily.

Theorem K.3. If $S \subseteq \mathbb{R}^d$ is an axis-parallel unit cube and e_1 is the first unit vector, then

$$\lim_{r \to 0} r^{-1} \operatorname{Vol}((S + re_1) \setminus S) = 1.$$

Theorem K.4. If $S \subseteq \mathbb{R}^d$ is a $(\ell_2$ -) ball of volume 1 and v is any $(\ell_2$ -)unit vector, then

$$\lim_{d\to\infty}\lim_{r\to 0}r^{-1}\mathrm{Vol}((S+rv)\setminus S)=\sqrt{e}.$$

Theorem K.5. If $S \subseteq R^d$ is the cross polytope (i.e. ℓ_1 ball) of volume 1, and e_1 is the first unit vector, then

$$\lim_{d \to \infty} \lim_{r \to 0} r^{-1} \operatorname{Vol}((S + re_1) \setminus S) = e.$$

Theorem K.6. If $S \subseteq \mathbb{R}^d$ is an axis-parallel unit cube and v = (1, ..., 1), then

$$\lim_{r \to 0} r^{-1} \operatorname{Vol}((S + rv) \setminus S) = d.$$

Theorem K.7. If $S \subseteq \mathbb{R}^d$ is the cross polytope (i.e. ℓ_1 ball) of volume 1, and v = (1, ..., 1), then

$$\lim_{d \to \infty} d^{-1/2} \lim_{r \to 0} r^{-1} \operatorname{Vol}((S + rv) \setminus S) = e\sqrt{2/\pi}.$$

Proof. It is equivalent to take S to be the standard ℓ_1 ball, and to calculate

$$\lim_{d \to \infty} \frac{\lim_{r \to 0} r^{-1} \operatorname{Vol}((S + rv) \setminus S)}{\operatorname{Vol}(S)^{\frac{d-1}{d}} \sqrt{d}}, \qquad (15)$$

and confirm it equals $e\sqrt{2/\pi}$. Note that the unit surface normals of S are $\{\pm 1\}^d/\sqrt{d}$, occurring with equal probability over the surface measure of S. Using Eq. (13), we then see that

$$\lim_{d \to \infty} \lim_{r \to 0} r^{-1} \operatorname{Vol}((S + rv) \setminus S) = W \sum_{i=0}^{\lfloor d/2 \rfloor} {d \choose i} \frac{d - 2i}{\sqrt{d}}$$

where $W = \frac{\sqrt{d}}{(d-1)!}$ is the volume of the simplex $\{x : \sum_{i} x_i = 1, x \ge 0\}$. This evaluates to

$$\frac{1}{(d-1)!} \times \begin{cases} \frac{d+2}{2} \binom{d}{\frac{d+1}{2}} & \text{if } d \text{ is even} \\ \frac{d+1}{2} \binom{d}{\frac{d+1}{2}} & \text{if } d \text{ is odd.} \end{cases}$$

Finally, since the volume of S is $2^d/d!$, we can calculate Eq. (15) directly and obtain the desired result.

K.2. Wulff Crystal of the ℓ_∞ Ball

In this section, let Z be the Wulff Crystal (Definition 5.1) w.r.t. $\mathcal{B} = \{x : \|x\|_{\infty} \leq 1\}$, i.e. Z is the unit ball of the norm dual to $\|x\|_* \stackrel{\text{def}}{=} \mathbb{E}_{v \sim \{\pm 1\}^d} |\langle x, v \rangle|$. By Proposition J.4, Z can also be described as the zonotope of $2^{-d+1}\{\pm 1\}^d$. From these descriptions, we can straightforwardly see the following properties of Z.

Proposition K.8. The vertices of Z farthest from the origin are coordinate vectors and their negations. The facets of Z closest to the origin are of the form $\{x : \pm x_i \pm x_j \le 1\}$. Therefore, with B denoting the ℓ_2 unit ball,

$$\frac{1}{\sqrt{2}}B \subseteq Z \subseteq B$$

In general, the properties of Z are elusive, and tied to many open problems in combinatorics and polytope theory (Ziegler, 1995). But we may heuristically compute $\lim_{d\to\infty} \lim_{r\to 0} r^{-1} \operatorname{Vol}((Z + rv) \setminus Z) =$ $\lim_{d\to\infty} ||v||_2 \Pi_v Z$ when $v = (1, \ldots, 1)$, as follows. (Because our computation is heuristic, we phrase the following as a claim, and not a theorem)

Claim K.9. If $S \subseteq \mathbb{R}^d$ is the Wulff Crystal w.r.t. the ℓ_{∞} unit ball, scaled to have volume 1, and v = (1, ..., 1), then

$$\lim_{d \to \infty} d^{-1/2} \lim_{r \to 0} r^{-1} \operatorname{Vol}((S + rv) \setminus S) = \sqrt{e}.$$

Derivation. Since $2^{d-1}Z$ is $\operatorname{Zon}(\{\pm 1\}^d)$, we have $\Pi_v 2^{d-1}Z = \operatorname{Zon}(\Pi_v \{\pm 1\}^d)$, the zonotope of the set of vectors $\left\{x - \frac{\langle x, v \rangle}{\|v\|_2}v : v \in \{\pm 1\}^d\right\}$. By Eq. (6), we then have

$$\lim_{r \to 0} \frac{\operatorname{Vol}((S + rv) \setminus S)}{r\sqrt{d}} = \frac{\operatorname{Vol}(\operatorname{Zon}(\Pi_v \{\pm 1\}^d))}{\operatorname{Vol}(\operatorname{Zon}(\{\pm 1\}^d))^{\frac{d-1}{d}}}.$$
 (16)

Now Lemma K.11 tells us that the Vol($\text{Zon}(\{\pm 1\}^d)$) is a multiple of the expected determinants of all $d \times d$ matrices with entries ± 1 . By a result of Nguyen et al. (2014) (Theorem K.12), the determinant of a random $d \times d$ matrix with iid ± 1 entries is distributed in high dimension d roughly as $\sqrt{(d-1)!}e^{z\sqrt{\frac{1}{2}\log d}}$ where $z \sim \mathcal{N}(0, 1)$. Thus, by Lemma K.11, we should expect (this is the first place where we argue heuristically)

$$\operatorname{Vol}(\operatorname{Zon}(\{\pm 1\}^d)) \approx \frac{1}{d!} 2^{d^2} \mathop{\mathbb{E}}_{z} \sqrt{(d-1)!} e^{z\sqrt{\frac{1}{2}\log d}} \\ = \frac{1}{d!} 2^{d^2} \sqrt{(d-1)!} d^{\frac{1}{4}}.$$
(17)

We verify this approximation to be correct numerically for moderately large d. Similarly, the uniform distribution over $\{\pm 1\}^d$ is close to a standard Gaussian when $d \gg 1$, so that $\prod_v \{\pm 1\}^d$ is close to a (d-1)-dimensional standard Gaussian. Therefore, we should expect that

$$\operatorname{Vol}(\operatorname{Zon}(\Pi_{v}\{\pm 1\}^{d})) \approx \frac{1}{(d-1)!} (2^{d}-1)^{d-1} \mathbb{E} |\det Y| \approx \frac{1}{(d-1)!} 2^{d(d-1)} \sqrt{(d-2)!} (d-1)^{\frac{1}{4}}, \quad (18)$$

where Y is a $(d-1) \times (d-1)$ Gaussian matrix, and where in Eq. (18), we used the heuristic that for large d, $|\det Y|$ is lognormal (Theorem K.12). Again, we verify these approximations numerically. Plugging in Eqs. (17) and (18) into Eq. (16) and taking the $d \to \infty$ limit yields the desired result.

Since the sphere has the same large d limit (Theorem K.4), we can say that

Claim K.10. For every $\epsilon > 0$, $S = the \ell_2$ -ball of volume 1 achieves within ϵ of

$$\min_{\operatorname{Vol}(S)=1} \sup_{\|v\|_{\infty} \le 1} d^{-1/2} \lim_{r \to 0} r^{-1} \operatorname{Vol}((S+rv) \setminus S),$$

for sufficiently large d. This is not true for $S = \text{the } \ell_{\infty}$ - or the ℓ_1 -ball.

Lemma K.11. The volume of $\operatorname{Zon}(\{\pm 1\}^d)$ is

$$\frac{1}{d!} 2^{d^2} \mathop{\mathbb{E}}_{X} |\det X|$$

where $X \in \{\pm 1\}^{d \times d}$ is a random $d \times d$ matrix whose coordinates are iid Rademacher variables (i.e. 1 or -1 with equal probability).

Proof. The above expression can be rewritten as

$$\frac{1}{d!} \sum_{X \in \{\pm 1\}^{d \times d}} |\det X|$$

Because det X = 0 if any two columns are equal, so this is equivalent to summing over all X with distinct columns.

$$\frac{1}{d!} \sum_{\substack{X \in \{\pm 1\}^{d \times d} \\ X \text{ has distinct columns}}} |\det X|$$

Finally, any given set of d distinct column vectors is represented d! times in the sum through its d! permutations, so this is equal to

$$\sum_{T \subseteq \{\pm 1\}^d : |T| = d} |\det T|,$$

which by Proposition J.5 is the volume of the zonotope in question. \Box

Theorem K.12 (Nguyen et al. (2014)). Let A_n be an $n \times n$ random matrix whose entries are independent real random variables with mean zero, variance one and with subexponential tail. Then with $\mu_n = \frac{1}{2} \log(n-1)!$ and $\sigma = \sqrt{\frac{1}{2} \log n}$

$$\sigma_n = \sqrt{\frac{1}{2}} \log n,$$

$$\sup_{x \in \mathbb{R}} \left| \mathbb{P}\left(\frac{\log(|\det A_n|) - \mu_n}{\sigma_n} \le x \right) - \mathbb{P}(\mathcal{N}(0, 1) \le x) \right|$$

$$\le \log^{-1/3 + o(1)} n.$$

In other words, det $A_n \approx \sqrt{(n-1)!}e^{z\sqrt{\frac{1}{2}\log n}}$ where $z \sim \mathcal{N}(0,1)$.

K.3. Growth Formula of a Set

Lemma K.13. Let $S \subseteq \mathbb{R}^d$ be a set of finite perimeter and $v \in \mathbb{R}^d$ be any vector. Then

$$\lim_{r \to 0} r^{-1} \operatorname{Vol}((S + rv) \setminus S) = \int_{\partial S} \Theta(\langle \mathbf{n}(x), v \rangle) \, \mathrm{d}x,$$

where $\mathbf{n}(x)$ is the normal at x w.r.t. S, and $\Theta(x) = \max(0, x)$.

Proof. Let $\partial S_v \stackrel{\text{def}}{=} \{x \in \partial S : \langle \mathbf{n}(x), v \rangle > 0\}$ be the part of *S*'s boundary whose surface normal *aligns* with *v*. For any vector *w*, let $\partial S_v + [0, w] = \{x + rw : x \in \partial S_v, r \in [0, 1]\}$ be the Minkowski sum of ∂S_v and the segment [0, w]. Then

it's clear that $\operatorname{Vol}(\partial S_v + [0, rv]) \leq r \int_{\partial S} \Theta(\langle \mathbf{n}(x), v \rangle) \, \mathrm{d}x$, and that

$$(S+rv)\setminus S\subseteq \partial S_v+[0,rv].$$

Thus,

$$\lim_{r \to 0} r^{-1} \operatorname{Vol}((S + rv) \setminus S)$$

$$\leq \lim_{r \to 0} r^{-1} r \int_{\partial S} \Theta(\langle \mathbf{n}(x), v \rangle) \, \mathrm{d}x$$

$$= \int_{\partial S} \Theta(\langle \mathbf{n}(x), v \rangle) \, \mathrm{d}x.$$

Now for the other direction, observe that the signed measure $r^{-1}(\mathbb{I}(x \in S + rv) - \mathbb{I}(x \in S))$ converges weakly to the (singular) signed measure $\langle \mathbf{n}(x), v \rangle$ supported on ∂S . Indeed, for any compactly supported C^1 function f, we have

$$\lim_{r \to 0} r^{-1} \int f(x) (\mathbb{I}(x \in S + rv) - \mathbb{I}(x \in S)) \, \mathrm{d}x$$
$$= \lim_{r \to 0} r^{-1} \int (f(x + rv) - f(x)) \mathbb{I}(x \in S) \, \mathrm{d}x$$
$$= \int_{S} D_v f(x) \, \mathrm{d}x = \int_{\partial S} \langle \mathbf{n}(x), v f(x) \rangle \, \mathrm{d}x.$$

Now, taking the supremum of the RHS over all compactly supported C^1 function $|f| \le 1$, we get

$$\begin{split} &\int_{\partial S} \Theta(\langle \mathbf{n}(x), v \rangle) \, \mathrm{d}x \\ &= \frac{1}{2} \int_{\partial S} |\langle \mathbf{n}(x), v \rangle| \, \mathrm{d}x \\ &= \frac{1}{2} \sup_{f} \int_{\partial S} \langle \mathbf{n}(x), v f(x) \rangle \, \mathrm{d}x \\ &= \frac{1}{2} \sup_{f} \lim_{r \to 0} r^{-1} \int f(x) (\mathbb{I}(x \in S + rv) - \mathbb{I}(x \in S)) \, \mathrm{d}x \\ &\leq \frac{1}{2} \liminf_{r \to 0} r^{-1} \sup_{f} \int f(x) (\mathbb{I}(x \in S + rv) - \mathbb{I}(x \in S)) \, \mathrm{d}x \\ &= \liminf_{r \to 0} r^{-1} \mathrm{Vol}((S + rv) \setminus S) \end{split}$$

as desired.

L. Optimal Smoothing Distributions Have Wulff Crystal Level Sets

Definition L.1 (Level-Equivalence). Call two distribution q_1 and q_2 *level-equivalent* if their superlevel sets have the same volumes:

$$\operatorname{Vol}\{x: q_1(x) \ge t\} = \operatorname{Vol}\{x: q_2(x) \ge t\}, \quad \forall t \in (0, \infty).$$

Theorem L.2. Let \mathcal{B} be a full-dimensional polytope in \mathbb{R}^d symmetric around the origin, and whose vertices form a

symmetric set. Let B be the Wulff Crystal w.r.t. \mathcal{B} . Let q_0 be a probability density function. Among all probability distributions q with regular (Definition F.1) and even density function that is level equivalent to q_0 , the distribution with concentric super level sets proportional to B minimize

$$\Phi(1/2) = \sup_{\|v\|=1} \sup_{q(U)=1/2} \lim_{r \searrow 0} \frac{q(U-rv) - 1/2}{r}$$

where $\|\cdot\|$ is the norm defined by \mathcal{B} .

Proof. Consider any distribution q level-equivalent to q_0 . Let U_t be its superlevel sets, so that $\partial U_t = q^{-1}(t)$.

Expanding the definition of Φ in terms of $\gamma_u^{(p)}$ (see Definition G.4), and exchanging maximization and expectation, we get

$$\Phi(1/2) = \max_{u \in \text{Vert}(\mathcal{B})} \mathbb{E} \gamma_u^{1/2} \ge \mathbb{E} u \mathbb{E} \gamma_u^{1/2}$$
$$= \mathbb{E} \int_u \max(\nabla q(x) \cdot u, 0) \, \mathrm{d}x$$
$$= \int_u \mathbb{E} \max(\nabla q(x) \cdot u, 0) \, \mathrm{d}x.$$

Since $\mathcal{B} = -\mathcal{B}$ and thus $\operatorname{Vert}(\mathcal{B}) = -\operatorname{Vert}(\mathcal{B})$,

$$|||w||| \stackrel{\text{def}}{=} \mathop{\mathbb{E}}_{u \sim \operatorname{Vert}(\mathcal{B})} \max(0, \langle w, u \rangle) = \frac{1}{2} \mathop{\mathbb{E}}_{u \sim \operatorname{Vert}(\mathcal{B})} |\langle w, u \rangle|$$

is a seminorm. This is in fact a norm because $Vert(\mathcal{B})$ spans \mathbb{R}^d , by the assumption that \mathcal{B} is full-dimensional. Thus

$$\Phi(1/2) \ge \int |||\nabla q(x)||| \,\mathrm{d}x$$

Define $g(x) \stackrel{\text{def}}{=} \frac{\||\nabla q(x)\|\|}{\|\nabla q(x)\|_2}$ if $\nabla q(x) \neq 0$, and g(x) = 0 otherwise. Then by Theorem F.3,

$$\int \||\nabla q(x)|| \, \mathrm{d}x = \int g(x) \|\nabla q(x)\|_2 \, \mathrm{d}x$$
$$= \int_0^\infty \int_{\partial U_t} g(x) \, \mathrm{d}x \, \mathrm{d}t.$$

By the Weak Sard's theorem (Theorem F.4), we may ignore the places where $\nabla q(x) = 0$, and this integral is the same as

$$\int_0^\infty \int_{\partial U_t} g(x) \, \mathrm{d}x \, \mathrm{d}t = \int_0^\infty \int_{\partial U_t} \frac{\||\nabla q(x)||}{\|\nabla q(x)\|_2} \, \mathrm{d}x \, \mathrm{d}t.$$
(19)

Now, the surface normal at x w.r.t. U_t is proportional to $\nabla q(x)$. Thus, the $(\ell_2$ -)unit normal $\mathbf{n}(x)$ at x w.r.t. U_t is given by $\frac{-\nabla q(x)}{\|\nabla q(x)\|_2}$, so $\frac{\|\nabla q(x)\|}{\|\nabla q(x)\|_2}$ is the $\||\cdot\||$ -norm of $\mathbf{n}(x)$. Therefore, the inner integral is

$$\int_{\partial U_t} \frac{\|\nabla q(x)\|}{\|\nabla q(x)\|_2} \, \mathrm{d}x = \int_{\partial U_t} \|\|\mathbf{n}(x)\|\| \, \mathrm{d}x.$$

By Theorem K.2, this is minimized for fixed $Vol(U_t)$ by $U_t \propto$ the Wulff Crystal w.r.t. \mathcal{B} . Thus, the unique distribution q^* level equivalent to q_0 and with concentric Wulff Crystal super level sets (all centered at 0) minimizes Eq. (19). But since

$$\Phi(1/2) = \max_{u} \int \max(\nabla q(x) \cdot u, 0) \, \mathrm{d}x$$

and the inner integral here is invariant in u when $q = q^*$ by the symmetry of the Wulff Crystal, as in the proof of Theorem K.2, Eq. (19) in fact holds with equality for q^* , so that q^* minimizes $\Phi(1/2)$ as well.

M. IID Distributions

In this section, we study smoothing distributions that have i.i.d. coordinates.

M.1. ℓ_1 Adversary

IID Log Concave Distributions

Theorem M.1. Let $\phi : \mathbb{R} \to \mathbb{R}$ be absolutely continuous, even, and convex such that $\exp(-\phi(x))$ is integrable. Suppose *H* is a smoothed classifier smoothed by

$$q(x) \propto \prod_{i=1}^{d} e^{-\phi(x_i)},$$

such that $H(x) = (H(x)_1, \ldots, H(x)_C)$ is a vector of probabilities that H assigns to each class $1, \ldots, C$. If H correctly predicts the class y on input x, and the probability of the correct class is $\rho \stackrel{\text{def}}{=} H(x)_y > 1/2$, then H continues to predict the correct class when x is perturbed by any η with

$$\|\eta\|_1 < \mathrm{CDF}_{\phi}^{-1}(\rho),$$

where CDF_{ϕ}^{-1} is the inverse CDF of the 1D random variable with density $\propto \exp(-\phi(x))$. This robust radius is tight.

Proof. We seek to apply Theorem G.6 to $G(x) = 1 - H(x)_y$, for which we need to derive random variables γ_u and $\gamma_u^{(p)}$, and most importantly, the function Φ .

WLOG, assume $u \in \operatorname{Vert}(\mathcal{B})$ is e_1 . Then $\gamma_u = \langle u, \nabla \log q(\delta) \rangle = \phi'(\delta_1)$, for $\delta \sim q$. Let X be the random variable whose density function is $\propto e^{-\phi(x)}$, and denote $\varphi(c) \stackrel{\text{def}}{=} \mathbb{P}[X > c]$. Thus γ_u is distributed as $\phi'(X)$. Since ϕ is convex, ϕ' is nondecreasing, so that $\gamma_u^{(p)}$ is distributed like $\phi'(X)\mathbb{I}(X > \varphi^{-1}(p))$ (using the fact that X has an atomless measure). Then with $C = \int_{-\infty}^{\infty} e^{-\phi(t)} dt$, we have

$$\Phi(p) = \mathbb{E} \gamma_u^{(p)} = C^{-1} \int_{\varphi^{-1}(p)}^{\infty} e^{-\phi(t)} \phi'(t) \, \mathrm{d}t$$
$$= C^{-1} e^{-\phi(\varphi^{-1}(p))}.$$

Then with $p_0 = 1 - \rho$, and by reparametrization the integral (Lemma G.7), the certified radius is

$$\int_{p_0}^{1/2} \frac{1}{\Phi(p)} dp = \int_{\varphi^{-1}(1/2)}^{\varphi^{-1}(p_0)} \frac{|\varphi'(c)|}{C^{-1}e^{-\phi(c)}} dc$$
$$= \int_{\varphi^{-1}(1/2)}^{\varphi^{-1}(p_0)} dc$$
$$= \varphi^{-1}(p_0) - \varphi^{-1}(1/2).$$

Simplifying this in terms of the CDF, and noting that $\varphi^{-1}(1/2) = 0$ because ϕ is even, yields the expression in the theorem statement.

This robust radius is tight, as can be seen from the case when a half-plane $\{x : x_1 \ge s\}$ is the set of inputs that the base classifier assigns the label y.

IID Log Convex* Distributions

Theorem M.2. Let $\phi : [0, \infty) \to \mathbb{R}$ be absolutely continuous, concave, and nondecreasing, such that $\exp(-\phi(|x|))$ is integrable. Suppose *H* is a smoothed classifier smoothed by

$$q(x) \propto \prod_{i=1}^{d} e^{-\phi(|x_i|)}$$

such that $H(x) = (H(x)_1, \ldots, H(x)_C)$ is a vector of probabilities that H assigns to each class $1, \ldots, C$. If H correctly predicts the class y on input x, and the probability of the correct class is $\rho \stackrel{\text{def}}{=} H(x)_y > 1/2$, then H continues to predict the correct class when x is perturbed by any η with

$$\begin{split} \|\eta\|_1 &< \int_{\varphi^{-1}(1-\rho)}^{\infty} \frac{\mathrm{d}c}{e^{\phi(c)-\phi(0)}-1} \\ &= \int_{1-\rho}^{1/2} \frac{C \,\mathrm{d}c}{e^{-\phi(0)} - e^{-\phi(\varphi^{-1}(p))}} \end{split}$$

where φ^{-1} is the inverse function of

$$\varphi(c) \stackrel{\text{def}}{=} \mathbb{P}_{z \sim q}[0 \le z_1 \le c],$$

and

$$C = \int_{-\infty}^{\infty} e^{-\phi(|t|)} \,\mathrm{d}t.$$

Proof. We seek to apply Theorem G.6 to $G(x) = 1 - H(x)_y$, for which we need to derive random variables γ_u and $\gamma_u^{(p)}$, and most importantly, the function Φ .

WLOG, assume $u \in \operatorname{Vert}(\mathcal{B})$ is e_1 . Then γ_u is the random variable $\langle u, \nabla \log q(\delta) \rangle = \phi'(|\delta_1|) \operatorname{sgn}(\delta_1)$ where $\delta \sim q$. Let $X \in \mathbb{R}$ be the random variable whose density function is $\propto e^{-\phi(|x|)}$, and so $\varphi(c) = \mathbb{P}[0 \leq X \leq c]$. Thus γ_u is distributed as $\phi'(X) \operatorname{sgn}(X)$. Since ϕ is concave, ϕ' is nonincreasing, so that for p < 1/2, $\gamma_u^{(p)}$ is distributed as $\phi'(X)\mathbb{I}(\varphi^{-1}(p) \geq X \geq 0)$ (using the fact that X has an atomless measure). Then with $C = \int_{-\infty}^{\infty} e^{-\phi(|t|)} dt$,

$$\Phi(p) = \mathbb{E} \gamma_u^{(p)} = C^{-1} \int_0^{\varphi^{-1}(p)} e^{-\phi(t)} \phi'(x) dt$$
$$= C^{-1} (e^{-\phi(0)} - e^{-\phi(\varphi^{-1}(p))}).$$

Then by change of variables $c = \varphi^{-1}(p)$ and with $p_0 = 1 - \rho$, the certified radius is

$$\int_{p_0}^{1/2} \frac{1}{\Phi(p)} dp = \int_{\varphi^{-1}(p_0)}^{\varphi^{-1}(1/2)} \frac{|\varphi'(c)|}{C^{-1}(e^{-\phi(0)} - e^{-\phi(c)})} dc$$
$$= \int_{\varphi^{-1}(p_0)}^{\varphi^{-1}(1/2)} \frac{e^{-\phi(c)}}{e^{-\phi(0)} - e^{-\phi(c)}} dc$$
$$= \int_{\varphi^{-1}(p_0)}^{\varphi^{-1}(1/2)} \frac{1}{e^{\phi(c) - \phi(0)} - 1} dc.$$

Since ϕ is even, $\varphi^{-1}(1/2) = \infty$, yielding the desired result.

Corollaries The ℓ_p based exponential distribution $\propto e^{-\|x\|_p^p}$ has each coordinate is distributed as Rademacher $\sqrt[p]{\text{Gamma}(1/p)}$. When $p \ge 1$, it satisfies Theorem M.1, so we obtain

Corollary M.3. Suppose *H* is a smoothed classifier smoothed by

$$q(x) \propto e^{-\|x/\lambda\|_p^p}, p \ge 1,$$

such that $H(x) = (H(x)_1, \ldots, H(x)_C)$ is a vector of probabilities that H assigns to each class $1, \ldots, C$. If H correctly predicts the class y on input x, and the probability of the correct class is $\rho \stackrel{\text{def}}{=} H(x)_y > 1/2$, then H continues to predict the correct class when x is perturbed by any η with

$$\|\eta\|_1 < \lambda \sqrt[p]{\text{GammaCDF}^{-1}(2\rho - 1; 1/p)},$$

where CDF_{ϕ}^{-1} is the inverse CDF of the 1D random variable with density $\propto \exp(-\phi(x))$. This robust radius is tight.

When p < 1, it satisfies Theorem M.2, so we obtain

Corollary M.4. Suppose *H* is a smoothed classifier smoothed by

$$q(x) \propto e^{-\|x/\lambda\|_p^p}, p < 1,$$

such that $H(x) = (H(x)_1, \ldots, H(x)_C)$ is a vector of probabilities that H assigns to each class $1, \ldots, C$. If H correctly predicts the class y on input x, and the probability of

the correct class is $\rho \stackrel{\text{def}}{=} H(x)_y > 1/2$, then H continues to predict the correct class when x is perturbed by any η with

$$\|\eta\|_1 < \lambda \int_{\varphi^{-1}(1-\rho)}^\infty \frac{\mathrm{d}c}{e^{c^p}-1},$$

where $\varphi^{-1}(p_0) \stackrel{\text{def}}{=} \text{GammaCDF}^{-1}(2p_0; 1/p)^{1/p}$.

The integral above can be evaluated explicitly for inverse integer p = 1/k. We show a few examples below:

$$p = 1/2: \quad R = 2\lambda(-c\log(1 - e^{-c}) + \text{polylog}(2, e^{-c}))$$

$$p = 1/3: \quad R = 3\lambda \left(-c^2\log(1 - e^{-c}) + 2\operatorname{polylog}(3, e^{-c}) + 2\operatorname{polylog}(3, e^{-c}) \right)$$

$$p = 1/4: \quad R = 4\lambda \left(-c^3\log(1 - e^{-c}) + 2\operatorname{polylog}(3, e^{-c}) \right)$$

$$+ 3c^{2} \operatorname{polylog}(2, e^{-c}) + 6c \operatorname{polylog}(3, e^{-c}) + 6 \operatorname{polylog}(4, e^{-c}) \right)$$

where $c = \text{GammaCDF}^{-1}(2(1 - \rho); 1/p)$ for each p, and polylog is the function defined as

$$\operatorname{polylog}(n, z) = \sum_{k=1}^{\infty} z^k / k^n.$$

N. ℓ_{∞} Norm-Based Exponential Law

In this section, we derive robustness guarantees for distributions of the form $q(x) \propto ||x||_{\infty}^{-j} \exp(-||x||_{\infty}^{k})$.

N.1. ℓ_1 Adversary

In this section, we set the norm $\|\cdot\|$ in Assumption G.1 to be the ℓ_1 norm $\|x\|_1 = \sum_{i=1}^d |x_i|$. Then the unit ball \mathcal{B} in Assumption G.1 is the convex hull of its vertices which are the coordinates vectors and their negations:

$$\operatorname{Vert}(\mathcal{B}) = \{ \pm e_i : i \in [d] \}.$$

Overview ℓ_{∞} norm-based distributions will in general have certified radius that is linear in $\rho - 1/2$, where ρ is the probability that the *smoothed* classifier assigns to the correct class.

We first demonstrate the differential method on $q(x) \propto \exp(-\|x\|_{\infty})$ as a warmup before stating the more general result.

Theorem N.1. Suppose *H* is a smoothed classifier smoothed by

$$q(x) \propto \exp(-\|x\|_{\infty}/\lambda),$$

such that $H(x) = (H(x)_1, \ldots, H(x)_C)$ is a vector of probabilities that H assigns to each class $1, \ldots, C$. If H correctly predicts the class y on input x, and the probability of the correct class is $\rho \stackrel{\text{def}}{=} H(x)_y > 1/2$, then H continues to predict the correct class when x is perturbed by any η with

$$\|\eta\|_{1} < \begin{cases} 2d\lambda(\rho - \frac{1}{2}) & \text{if } \rho \le 1 - \frac{1}{2d} \\ \lambda \log \frac{1}{2d(1-\rho)} + \lambda(d-1) & \text{if } \rho > 1 - \frac{1}{2d}. \end{cases}$$

Proof. By linearity in λ , it suffices to show this for $\lambda = 1$. Here, we have $q(x) \propto \exp(\psi(x))$ with

$$\psi(x) = \|x\|_{\infty}$$
 and $\nabla \psi(x) = \operatorname{sgn}(x_{i^*})e_{i^*},$

where $i^* = \operatorname{argmax}_i |x_i|$, and e_{i^*} is the i^* th coordinate vector, with $\nabla \psi(x)$ defined whenever i^* is the unique argmax.

We seek to apply Theorem G.6 to $G(x) = 1 - H(x)_y$, for which we need to derive random variables γ_u and $\gamma_u^{(p)}$, and most importantly, the function Φ .

For any
$$u \in \operatorname{Vert}(\mathcal{B})$$
 (i.e. $u = \pm e_i$), the random variable $\gamma_u = \langle u, \nabla \psi(\delta) \rangle = \langle u, \operatorname{sgn}(\delta_{i^*})e_{i^*} \rangle, \delta \sim q$, is given by

$$\gamma_u = \begin{cases} 0 & \text{with prob. } 1 - \frac{1}{d} \\ 1 & \text{with prob. } \frac{1}{2d} \\ -1 & \text{with prob. } \frac{1}{2d}. \end{cases}$$

Therefore, for $p \in [0,1/2],$ the random variable $\gamma_u^{(p)}$ defined in Definition G.4 is

$$\begin{cases} \gamma_u^{(p)} = \begin{cases} 1 & \text{with prob. } \frac{1}{2d} \\ 0 & \text{with prob. } 1 - \frac{1}{2d} \end{cases} & \text{if } p \in [\frac{1}{2d}, \frac{1}{2}], \\ \gamma_u^{(p)} = \begin{cases} 1 & \text{with prob. } p \\ 0 & \text{with prob. } 1 - p \end{cases} & \text{if } p \in [0, \frac{1}{2d}]. \end{cases}$$

Thus, for any $u \in Vert(\mathcal{B})$,

$$\Phi(p) = \mathbb{E} \gamma_u^{(p)} = \begin{cases} \frac{1}{2d} & \text{if } p \in \left[\frac{1}{2d}, \frac{1}{2}\right] \\ p & \text{if } p \in \left[0, \frac{1}{2d}\right] \end{cases}$$

Then, by setting G(x) in Theorem G.6 to be $1 - H(x)_y = 1 - \rho \stackrel{\text{def}}{=} p_0$, we get the provably robust radius of

$$\int_{p_0}^{1/2} \frac{1}{\Phi(p)} dp$$

$$= \begin{cases} \int_{p_0}^{1/2} 2d \, dp = 2d(\frac{1}{2} - p_0) & \text{if } p_0 \ge \frac{1}{2d} \\ \int_{p_0}^{1/2d} p^{-1} \, dp + \int_{1/2d}^{1/2} 2d \, dp & \text{if } p_0 \le \frac{1}{2d} \end{cases}$$

Simplifying the arithmetics yields the desired claim. $\hfill \Box$

Now we tackle the general case.

Theorem N.2. Suppose *H* is a smoothed classifier smoothed by

$$q(x) \propto (\|x\|_{\infty}/\lambda)^{-j} \exp(-(\|x\|_{\infty}/\lambda)^k), k \ge 1,$$

such that $H(x) = (H(x)_1, \ldots, H(x)_C)$ is a vector of probabilities that H assigns to each class $1, \ldots, C$. If H correctly predicts the class y on input x, and the probability of the correct class is $\rho \stackrel{\text{def}}{=} H(x)_y \in (1/2, 1 - \frac{1}{2d}]$, then Hcontinues to predict the correct class when x is perturbed by any η with

$$\|\eta\|_1 < \frac{2d\lambda}{d-1} \frac{\Gamma\left(\frac{d-j}{k}\right)}{\Gamma\left(\frac{d-1-j}{k}\right)} \left(\rho - \frac{1}{2}\right).$$

Proof. By linearity in λ , it suffices to show this for $\lambda = 1$. Here, we have

$$q(x) \propto \exp(-\|x\|_{\infty}^{k} - j \log \|x\|_{\infty}) \qquad \text{so that}$$

$$\psi(x) = \|x\|_{\infty}^{k} + j \log \|x\|_{\infty}$$

$$\nabla \psi(x) = (k\|x\|_{\infty}^{k-1} + j\|x\|_{\infty}^{-1}) \operatorname{sgn}(x_{i^{*}}) e_{i^{*}},$$

where $i^* = \operatorname{argmax}_i |x_i|$, and e_{i^*} is the i^* th coordinate vector, with $\nabla \psi(x)$ defined whenever i^* is the unique argmax.

We seek to apply Theorem G.6 to $G(x) = 1 - H(x)_y$, for which we need to derive random variables γ_u and $\gamma_u^{(p)}$, and most importantly, the function Φ .

WLOG among Vert(\mathcal{B}), let's assume $u = e_1$. Then the random variable $\gamma_u = \langle u, \nabla \psi(\delta) \rangle, \delta \sim q$, is 0 with probability $1 - \frac{1}{d}$, when $i^* \neq 1$. When $i^* = 1$ and $\operatorname{sgn}(x_{i^*}) = 1$ (which happens with probability $\frac{1}{2d}$), γ_u is $k ||x||_{k^{-1}}^{k-1} + j||x||_{\infty}^{-1}$, where $x \sim q$. By Lemma U.3, this is just the random variable $kz^{\frac{k-1}{k}} + jz^{-1}$, where $z \sim \operatorname{Gamma}(d/k)$. Likewise, with probability $\frac{1}{2d}$, γ_u is the random variable $-kz^{\frac{k-1}{k}} - jz^{-1}$. This can be summarized below.

$$\gamma_u = \begin{cases} 0 & \text{with prob. } 1 - \frac{1}{d} \\ kz^{\frac{k-1}{k}} + jz^{-1} & \text{with prob. } \frac{1}{2d} \\ -kz^{\frac{k-1}{k}} - jz^{-1} & \text{with prob. } \frac{1}{2d}, \end{cases}$$

where $z \sim \text{Gamma}(d/k)$.

Therefore, for $p \in [\frac{1}{2d}, \frac{1}{2}]$, the random variable $\gamma_u^{(p)}$ defined in Definition G.4 is

$$\gamma_u^{(p)} = \begin{cases} kz^{\frac{k-1}{k}} + jz^{-1} & \text{with prob. } \frac{1}{2d} \\ 0 & \text{with prob. } 1 - \frac{1}{2d} \end{cases}$$

where z is sampled from Gamma(d/k).

Thus, for any $u \in Vert(\mathcal{B})$, by Lemma U.4,

$$\Phi(p) = \mathbb{E} \, \gamma_u^{(p)} = \frac{1}{2d} \, \mathbb{E} \, k z^{\frac{k-1}{k}} = \frac{d-1}{2d} \frac{\Gamma\left(\frac{d-1-j}{k}\right)}{\Gamma\left(\frac{d-j}{k}\right)}$$

which does not depend on p.

Then, by setting G(x) in Theorem G.6 to be $1 - H(x)_y = 1 - \rho$, we get the provably robust radius of

$$\int_{1-\rho}^{1/2} \frac{1}{\Phi(p)} \, \mathrm{d}p = \frac{2d}{d-1} \frac{\Gamma\left(\frac{d-j}{k}\right)}{\Gamma\left(\frac{d-1-j}{k}\right)} \left(\rho - \frac{1}{2}\right)$$

as desired.

As j = 0 and $k \to \infty$, the distribution above converges to the uniform distribution, and the robust certificate converges likewise to the one computed previous for uniform distribution.

Theorem N.3 (Lee et al. (2019)). Suppose H is a smoothed classifier smoothed by the uniform distribution on the cube $[-\lambda, \lambda]^d$, such that $H(x) = (H(x)_1, \ldots, H(x)_C)$ is a vector of probabilities that H assigns to each class $1, \ldots, C$. If H correctly predicts the class y on input x, and the probability of the correct class is $\rho \stackrel{\text{def}}{=} H(x)_y > 1/2$, then H continues to predict the correct class when x is perturbed by any η with

$$\|\eta\|_1 < 2\lambda \left(\rho - \frac{1}{2}\right).$$

N.2. ℓ_{∞} Adversary

In this section, we set the norm $\|\cdot\|$ in Assumption G.1 to be the ℓ_{∞} norm $\|x\|_{\infty} = \max_{i=1}^{d} |x_i|$. Then the unit ball \mathcal{B} in Assumption G.1 is the convex hull of its vertices which are points in the Boolean cube:

$$\operatorname{Vert}(\mathcal{B}) = \{\pm 1\}^d.$$

Theorem N.4. Suppose *H* is a smoothed classifier smoothed by

$$q(x) \propto \exp(-\|x\|_{\infty}/\lambda),$$

such that $H(x) = (H(x)_1, \ldots, H(x)_C)$ is a vector of probabilities that H assigns to each class $1, \ldots, C$. If H correctly predicts the class y on input x, and the probability of the correct class is $\rho \stackrel{\text{def}}{=} H(x)_y > 1/2$, then H continues to predict the correct class when x is perturbed by any η with

$$\|\eta\|_{\infty} < \lambda \log \frac{1}{2(1-\rho)}$$

Proof. By linearity in λ , it suffices to show this for $\lambda = 1$. We seek to apply Theorem G.6 to $G(x) = 1 - H(x)_y$, for which we need to derive random variables γ_u and $\gamma_u^{(p)}$, and most importantly, the function Φ . For any $u \in \operatorname{Vert}(\mathcal{B})$, the random variable $\gamma_u = \langle u, \nabla \psi(\delta) \rangle = \langle u, \operatorname{sgn}(\delta_{i^*})e_{i^*} \rangle, \delta \sim q$, is a Rademacher random variable taking values ± 1 with equal probability. Therefore, for $p \in [0, 1/2]$, the random variable $\gamma_u^{(p)}$ defined in Definition G.4 is

$$\gamma_u^{(p)} = \begin{cases} 1 & \text{with prob. } p \\ 0 & \text{with prob. } 1-p \end{cases}$$

Thus, for any $u \in \operatorname{Vert}(\mathcal{B})$,

$$\Phi(p) = \mathbb{E}\,\gamma_u^{(p)} = p.$$

Then, by setting G(x) in Theorem G.6 to be $1 - H(x)_y = 1 - \rho$, we get the provably robust radius of

$$\int_{1-\rho}^{1/2} \frac{1}{\Phi(p)} \, \mathrm{d}p = \int_{1-\rho}^{1/2} \frac{1}{p} \, \mathrm{d}p = \log \frac{1}{2(1-\rho)}.$$

Theorem N.5. Suppose H is a smoothed classifier smoothed by

$$q(x) \propto \exp(-\|x/\lambda\|_{\infty}^{k}), k \ge 1$$

such that $H(x) = (H(x)_1, \ldots, H(x)_C)$ is a vector of probabilities that H assigns to each class $1, \ldots, C$. If H correctly predicts the class y on input x, and the probability of the correct class is $\rho \stackrel{\text{def}}{=} H(x)_y > 1/2$, then H continues to predict the correct class when x is perturbed by any η with

$$\|\eta\|_{\infty} < \lambda \int_{1-\rho}^{1/2} \frac{1}{\Phi(p)} \,\mathrm{d}p,$$
 (20)

in which

$$\Phi(p) \stackrel{\text{def}}{=} \frac{k-1}{2} \text{GammaCDF}\left(c^*; \frac{d+k-1}{k}\right) \frac{\Gamma\left(\frac{d+k-1}{k}\right)}{\Gamma\left(\frac{d}{k}\right)}$$

where $c^* = \text{GammaCDF}^{-1}\left(1-2p; \frac{d}{k}\right)$.

More generally, if the smoothing distribution is

$$q(x) \propto \|x/\lambda\|_{\infty}^{-j} \exp(-\|x/\lambda\|_{\infty}^{k}), k \ge 1, j < d-1,$$

then H is robust against ℓ_{∞} perturbation

$$\|\eta\|_{\infty} < \lambda \int_{1-\rho}^{1/2} \frac{1}{\Phi(p)} \,\mathrm{d}p,\tag{21}$$

where

$$\begin{split} \Phi(p) &\stackrel{\text{def}}{=} \frac{1}{2} \bar{\phi}(\phi^{-1}(2p)), \quad \text{where} \\ \phi(c) &\stackrel{\text{def}}{=} \mathbb{P}[\gamma > c] \\ \bar{\phi}(c) &\stackrel{\text{def}}{=} \mathbb{E} \gamma \mathbb{I}(\gamma > c) \end{split}$$

and $\gamma \stackrel{\text{def}}{=} (k-1)\xi^{\frac{k-1}{k}} + j\xi^{-\frac{1}{k}}, \xi \sim \text{Gamma}\left(\frac{d}{k} - \frac{j}{k}\right).$

Proof. By linearity in λ , it suffices to show this for $\lambda = 1$. We seek to apply Theorem G.6 to $G(x) = 1 - H(x)_y$, for which we need to derive random variables γ_u and $\gamma_u^{(p)}$, and most importantly, the function Φ .

For any $u \in Vert(\mathcal{B})$, we have

$$\begin{aligned} \gamma_u &= \langle u, -\nabla \log q(\delta) \rangle \\ &= \langle u, (k \|\delta\|_{\infty}^{k-1} + j \|\delta\|_{\infty}^{-1}) \operatorname{sgn}(\delta_{i^*}) e_{i^*} \rangle, \delta \sim q. \end{aligned}$$

Since $\|\delta\|_{\infty}$ is distributed as $\sqrt[k]{\text{Gamma}(\frac{d}{k} - \frac{j}{k})}$ and $\langle u, \operatorname{sgn}(\delta_{i^*})e_{i^*}\rangle$ is ± 1 with equal probability, γ_u is distributed as the product of random variables

$$\gamma_u = \zeta (k\xi^{\frac{k-1}{k}} + j\xi^{-\frac{1}{k}}),$$

$$\zeta \sim \text{Rademacher}, \xi \sim \text{Gamma}\left(\frac{d}{k} - \frac{j}{k}\right)$$

Let $\varphi(c) \stackrel{\text{def}}{=} \mathbb{P}[\gamma_u > c]$. Then for p < 1/2, $\phi^{-1}(2p) = \varphi^{-1}(p)$. Since γ_u has an absolutely continuous distribution, the variable $\gamma_u^{(p)} = \gamma_u|_{(\varphi^{-1}(p),\infty)}$ with probability p, and 0 otherwise. Thus

$$\begin{split} \Phi(p) &= \mathbb{E} \, \gamma_u^{(p)} = \bar{\varphi}(\varphi^{-1}(p)), \quad \text{where} \\ \bar{\varphi}(c) &= \mathbb{E} \, \gamma_u \mathbb{I}(\gamma_u > c). \end{split}$$

Note that $\bar{\varphi}(c) = \frac{1}{2}\bar{\phi}(c)$. Plugging into Theorem G.6 yields Eq. (21).

Assuming j = 0 If j = 0, $\gamma_u = \zeta k \xi^{\frac{k-1}{k}}, \zeta \sim$ Rademacher, $\xi \sim \text{Gamma}\left(\frac{d}{k}\right)$. Then for p < 1/2,

$$\begin{split} \bar{\varphi}(c) &= \frac{k}{2} \mathbb{E} \, \xi^{\frac{k-1}{k}} \mathbb{I}(\xi > c^*), \\ \text{where } c^* &= \text{GammaCDF}^{-1} \left(1 - 2p; \frac{d}{k} \right) \\ &= \frac{k}{2} \text{GammaCDF} \left(c^*; \frac{d+k-1}{k} \right) \frac{\Gamma\left(\frac{d+k-1}{k}\right)}{\Gamma\left(\frac{d}{k}\right)}. \end{split}$$

Plugging into Theorem G.6 yields Eq. (20).

Compare this with the uniform case below.

Theorem N.6 (Lee et al. (2019)). Suppose H is a smoothed classifier smoothed by the uniform distribution on the cube $[-\lambda, \lambda]^d$, such that $H(x) = (H(x)_1, \ldots, H(x)_C)$ is a vector of probabilities that H assigns to each class $1, \ldots, C$. If H correctly predicts the class y on input x, and the probability of the correct class is $\rho \stackrel{\text{def}}{=} H(x)_y > 1/2$, then H continues to predict the correct class when x is perturbed by any η with

$$\|\eta\|_{\infty} < 2\lambda \left(1 - \sqrt[d]{\frac{3}{2} - \rho}\right).$$

When $d \to \infty$, this robust radius is roughly

$$\begin{split} & 2\lambda \left(1-e^{\frac{1}{d}\log\left[1-\left(\rho-\frac{1}{2}\right)\right]}\right) \\ &\approx 2\lambda \left(1-\left(1+\frac{1}{d}\log\left[1-\left(\rho-\frac{1}{2}\right)\right]\right)\right) \\ &\approx \frac{2\lambda}{d}\left(\rho-\frac{1}{2}\right). \end{split}$$

On the other hand, when $k \to \infty$ in Eq. (20), we see that

- 1. $d/k \to 0$ while $\frac{d+k-1}{k} \to 1$
- 2. $c^* \rightarrow 0$ for any p < 1/2
- 3. GammaCDF $(c^*; \frac{d+k-1}{k}) \rightarrow$ GammaCDF(0; 1) = 1 consequently
- 4. by simple calculation $k \frac{\Gamma(\frac{d+k-1}{k})}{\Gamma(\frac{d}{k})} \to d$
- 5. so $\Phi(p) \rightarrow d/2$ for any p < 1/2.

Therefore, when $k \to \infty$, the ℓ_{∞} robust radius in Eq. (20) converges to

$$\frac{2\lambda}{d}\left(
ho-\frac{1}{2}
ight)$$

as well.

O. ℓ_{∞} Norm-Based Power Law

Now consider a power law of the ℓ_{∞} norm: For a > d,

$$q(x) \propto (1 + ||x||_{\infty})^{-a}$$
 so that
 $\psi(x) = a \log(1 + ||x||_{\infty})$
 $\nabla \psi(x) = a(1 + ||x||_{\infty})^{-1} \operatorname{sgn}(x_{i^*}) e_{i^*},$

where $i^* = \operatorname{argmax}_i |x_i|$, and e_{i^*} is the *i*^{*}th coordinate vector, with $\nabla \psi(x)$ defined whenever *i*^{*} is the unique argmax. Note that the ℓ_{∞} norm of vector sampled from *q* has distribution with CDF

$$\mathbb{P}_{\delta \sim q}[\|\delta\|_{\infty} \leq c] = \frac{\Gamma(a)}{\Gamma(a-d)\Gamma(d)} \int_{0}^{c} \frac{r^{d-1}}{(1+r)^{a}} \,\mathrm{d}r. \tag{22}$$

This is known as the *Beta prime or Beta distribution of the second kind, with shape parameters* $\alpha = d, \beta = a - d$, which we will denote by BetaPrime(d, a - d). If a > d + 1, this distribution has mean

$$\mathop{\mathbb{E}}_{\delta \sim q} \|\delta\|_{\infty} = \frac{d}{a-d-1}$$

O.1. ℓ_1 Adversary

Theorem O.1. Suppose *H* is a smoothed classifier smoothed by

$$q(x) \propto (1 + \|x\|_{\infty}/\lambda)^{-a}, a > d_{2}$$

such that $H(x) = (H(x)_1, \ldots, H(x)_C)$ is a vector of probabilities that H assigns to each class $1, \ldots, C$. If H correctly predicts the class y on input x, and the probability of the correct class is $\rho \stackrel{\text{def}}{=} H(x)_y > 1/2$, then H continues to predict the correct class when x is perturbed by any η with

$$\|\eta\|_1 < \lambda \frac{2d}{a-d} \left(\rho - \frac{1}{2}\right)$$

Proof. By linearity in λ , it suffices to show this for $\lambda = 1$. We seek to apply Theorem G.6 to $G(x) = 1 - H(x)_y$, for which we need to derive random variables γ_u and $\gamma_u^{(p)}$, and most importantly, the function Φ .

WLOG among Vert(\mathcal{B}), let's assume $u = e_1$. Then the random variable $\gamma_u = \langle u, \nabla \psi(\delta) \rangle, \delta \sim q$, is 0 with probability $1 - \frac{1}{d}$, when $i^* \neq 1$. When $i^* = 1$ and $\operatorname{sgn}(x_{i^*}) = 1$ (which happens with probability $\frac{1}{2d}$), γ_u is distributed as $a(1 + r)^{-1}$, where r has CDF Eq. (22). Likewise, with probability $\frac{1}{2d}$, γ_u is distributed as $-a(1 + r)^{-1}$.

This can be summarized below.

$$\gamma_u = \begin{cases} 0 & \text{with prob. } 1 - \frac{1}{d} \\ a(1+r)^{-1} & \text{with prob. } \frac{1}{2d} \\ -a(1+r)^{-1} & \text{with prob. } \frac{1}{2d}, \end{cases}$$

where r is a random variable with CDF Eq. (22).

Therefore, for $p \in [\frac{1}{2d}, \frac{1}{2}]$, the random variable $\gamma_u^{(p)}$ defined in Definition G.4 is

$$\gamma_u^{(p)} = \begin{cases} a(1+r)^{-1} & \text{with prob. } \frac{1}{2d} \\ 0 & \text{with prob. } 1 - \frac{1}{2d}. \end{cases}$$

Thus, for any $u \in Vert(\mathcal{B})$, by Lemma U.4,

$$\begin{split} \Phi(p) &= \mathbb{E} \, \gamma_u^{(p)} = \frac{1}{2d} \frac{\mathbb{E} \, a(1+r)^{-1}}{r} \\ &= \frac{1}{2d} \frac{a\Gamma(a)}{\Gamma(a-d)\Gamma(d)} \int_0^\infty a \frac{r^{d-1}}{(1+r)^{a+1}} \, \mathrm{d}r \\ &= \frac{1}{2d} \frac{\Gamma(a+1)}{\Gamma(a-d)\Gamma(d)} \frac{\Gamma(a+1-d)\Gamma(d)}{\Gamma(a+1)} \\ &= \frac{a-d}{2d}. \end{split}$$

which does not depend on p.

Then, by setting G(x) in Theorem G.6 to be $1 - H(x)_y = 1 - \rho$, we get the provably robust radius of

$$\int_{1-\rho}^{1/2} \frac{1}{\Phi(p)} \, \mathrm{d}p = \frac{2d}{a-d} \left(\rho - \frac{1}{2}\right)$$

as desired.

O.2. ℓ_{∞} Adversary

Theorem O.2. Suppose *H* is a smoothed classifier smoothed by

$$q(x) \propto (1 + \|x\|_{\infty}/\lambda)^{-a}, a > d,$$

such that $H(x) = (H(x)_1, \ldots, H(x)_C)$ is a vector of probabilities that H assigns to each class $1, \ldots, C$. If H correctly predicts the class y on input x, and the probability of the correct class is $\rho \stackrel{\text{def}}{=} H(x)_y > 1/2$, then H continues to predict the correct class when x is perturbed by any η with

$$\|\eta\|_1 < \frac{2\lambda}{a-d} \int_{1-\rho}^{1/2} \frac{\mathrm{d}p}{\Upsilon(\Upsilon^{-1}(2p;d,a-d);d,a+1-d)}$$

where $\Upsilon = \text{BetaPrimeCDF}$.

Proof. By linearity in λ , it suffices to show this for $\lambda = 1$.

We seek to apply Theorem G.6 to $G(x) = 1 - H(x)_y$, for which we need to derive random variables γ_u and $\gamma_u^{(p)}$, and most importantly, the function Φ .

For any $u \in Vert(\mathcal{B})$, we have

$$\gamma_u = \langle u, -\nabla \log q(\delta) \rangle$$

= $\langle u, a(1 + \|\delta\|_{\infty})^{-1} \operatorname{sgn}(\delta_{i^*}) e_{i^*} \rangle, \delta \sim q$

Since $\|\delta\|_{\infty}$ is distributed as BetaPrime(d, a - d) and $\langle u, \operatorname{sgn}(\delta_{i^*})e_{i^*} \rangle$ is ± 1 with equal probability, γ_u is distributed as the product of random variables

$$\gamma_u = \zeta a (1+\xi)^{-1},$$

 $\zeta \sim \text{Rademacher}, \xi \sim \text{BetaPrime}(d, a-d).$

Since $r \mapsto (1+r)^{-1}$ is a decreasing function on $r \in [0, \infty)$, we have, for p < 1/2,

$$\begin{split} \Phi(p) &= \mathbb{E} \, \gamma_u^{(p)} = \frac{1}{2} \mathop{\mathbb{E}}_{\xi} a(1+\xi)^{-1} \mathbb{I}(\xi < c(p)), \\ \text{where } c(p) &= \text{BetaPrimeCDF}^{-1}(2p; d, a - d). \end{split}$$

Of course, we can simplify

$$\begin{split} & \mathbb{E}(1+\xi)^{-1}\mathbb{I}(\xi < c) \\ &= \frac{\Gamma(a)}{\Gamma(a-d)\Gamma(d)} \int_0^c \frac{r^{d-1}}{(1+r)^{a+1}} \,\mathrm{d}r \\ &= \frac{\Gamma(a)}{\Gamma(a-d)\Gamma(d)} \frac{\Gamma(a+1-d)\Gamma(d)}{\Gamma(a+1)} \times \\ & \text{BetaPrimeCDF}(c;d,a+1-d) \\ &= \frac{(a-d)}{a} \text{BetaPrimeCDF}(c;d,a+1-d). \end{split}$$

Therefore,

$$\Phi(p) = \frac{a-d}{2} \operatorname{BetaPrimeCDF}(c(p); d, a+1-d).$$

Plugging into Theorem G.6 yields the desired result. \Box

P. ℓ_1 Norm-Based Exponential Law

Consider the following generalization of the Laplace distribution

$$q(x) \propto \exp(-\|x\|_1^k) \text{ so that}$$

$$\psi(x) = \|x\|_1^k$$

$$\nabla \psi(x) = k\|x\|_1^{k-1}(\operatorname{sgn}(x_1), \dots, \operatorname{sgn}(x_d)),$$

with $\nabla \psi(x)$ defined whenever all x_i s are nonzero.

P.1. ℓ_1 Adversary

Theorem P.1. Suppose *H* is a smoothed classifier smoothed by

$$q(x) \propto \exp(-(\|x\|_1/\lambda)^k),$$

such that $H(x) = (H(x)_1, \ldots, H(x)_C)$ is a vector of probabilities that H assigns to each class $1, \ldots, C$. If H correctly predicts the class y on input x, and the probability of the correct class is $\rho \stackrel{\text{def}}{=} H(x)_y > 1/2$, then H continues to predict the correct class when x is perturbed by any η with

$$\|\eta\|_1 < \lambda \int_{1-\rho}^{1/2} \frac{R}{\Psi(p)} \,\mathrm{d}p.$$

Here
$$R = \frac{2\Gamma\left(\frac{d}{k}\right)}{k\Gamma\left(\frac{d+k-1}{k}\right)}$$
, and

$$\Psi(p) \stackrel{\text{def}}{=} \begin{cases} 1 - \Upsilon\left(\Upsilon^{-1}(1-2p;\frac{d}{k});\frac{d+k-1}{k}\right) & \text{if } k \ge 1\\ \Upsilon\left(\Upsilon^{-1}(2p;\frac{d}{k});\frac{d+k-1}{k}\right) & \text{if } k \in (0,1). \end{cases}$$

where $\Upsilon = \text{GammaCDF}$.

Proof. By linearity in λ , it suffices to show this for $\lambda = 1$.

We seek to apply Theorem G.6 to $G(x) = 1 - H(x)_y$, for which we need to derive random variables γ_u and $\gamma_u^{(p)}$, and most importantly, the function Φ .

For any $u \in \operatorname{Vert}(\mathcal{B})$ (i.e. $u = \pm e_i$), $\gamma_u = \langle u, \nabla \psi(\delta) \rangle = \pm k \|\delta\|_1^{k-1}, \delta \sim q$, takes positive or negative sign with equal probability. By Lemma U.3, $\|\delta\|_1$ is the random variable $\Gamma(d/k)^{1/k}$. Therefore, γ_u is distributed as $k\operatorname{Gamma}(d/k)^{\frac{k-1}{k}}$ Rademacher(1/2).

Therefore, for $p \in [0, 1/2]$, the random variable $\gamma_u^{(p)}$ defined in Definition G.4 is

$$\gamma_u^{(p)} = \begin{cases} k z_p^{\frac{k-1}{k}} & \text{with prob. } p \\ 0 & \text{with prob. } 1-p, \end{cases}$$

where, if $k \geq 1$, z_p is sampled from $\operatorname{Gamma}(d/k)$ conditioned on $z_p > \Upsilon^{-1}(1-2p;d/k)$ (because $z^{\frac{k-1}{k}}$ is increasing in z), but if $k \in (0,1)$, then z_p is sampled from $\operatorname{Gamma}(d/k)$ conditioned on $z_p < \Upsilon^{-1}(2p;d/k)$ (because $z^{\frac{k-1}{k}}$ is decreasing in z).

Thus, for any $u \in \operatorname{Vert}(\mathcal{B})$,

$$\Phi(p) = \mathbb{E}\,\gamma_u^{(p)} = \begin{cases} \frac{k}{2}\,\mathbb{E}\,z^{\frac{k-1}{k}}\mathbb{I}(z > \Upsilon^{-1}(1-2p)) & \text{if } k \ge 1\\ \frac{k}{2}\,\mathbb{E}\,z^{\frac{k-1}{k}}\mathbb{I}(z < \Upsilon^{-1}(2p)) & \text{if } k < 1 \end{cases}$$

where $z \sim \text{Gamma}(d/k)$. This integral simplifies to $R^{-1}\Psi(p)$ (with Ψ taking different forms depending on k) by Lemma U.4.

Then, by setting G(x) in Theorem G.6 to be $1 - H(x)_y = 1 - \rho$, we get the provably robust radius of

$$\int_{1-\rho}^{1/2} \frac{1}{\Phi(p)} \, \mathrm{d}p = \int_{1-\rho}^{1/2} \frac{R}{\Psi(p)} \, \mathrm{d}p.$$

P.2. ℓ_{∞} Adversary

We first start with the Laplace distribution to highlight the basic logic behind the ℓ_{∞} radius derivation.

Theorem P.2. Suppose *H* is a smoothed classifier smoothed by the Laplace distribution

$$q(x) \propto \exp(-\|x\|_1/\lambda),$$

such that $H(x) = (H(x)_1, \ldots, H(x)_C)$ is a vector of probabilities that H assigns to each class $1, \ldots, C$. If H correctly predicts the class y on input x, and the probability of the correct class is $\rho \stackrel{\text{def}}{=} H(x)_y > 1/2$, then H continues to predict the correct class when x is perturbed by any η with

$$\|\eta\|_{\infty} < \lambda \int_{1-\rho}^{1/2} \frac{1}{\Phi(p)} \,\mathrm{d}p$$

where

$$\Phi(p) = c(p - \phi_d(c)) + d\phi_{d-1}\left(c - \frac{1}{2}\right) - d\phi_d(c),$$

in which $c = \phi_d^{-1}(p)$, and

$$\phi_d(c) \stackrel{\text{def}}{=} 2^{-d} \sum_{i=\frac{c+d}{2}+1}^d \binom{d}{i}$$
$$= 1 - \text{BinomCDF}\left(\frac{c+d}{2}; d\right)$$
for any $c \equiv d \mod 2$

$$\phi_d^{-1}(p) \stackrel{\text{def}}{=} \inf\{c : \phi_d(c) \le p\}$$

= 2BinomCDF⁻¹(1 - p) - d.

Note that when d is large,

$$\Phi(p) \approx \text{GaussianCDF}'(\text{GaussianCDF}^{-1}(p))\sqrt{d},$$

so that the bound above is roughly

$$\|\delta\|_{\infty} < \lambda \text{GaussianCDF}^{-1}(\rho)/\sqrt{d}.$$

Proof. By linearity in λ , it suffices to show this for $\lambda = 1$. We seek to apply Theorem G.6 to $G(x) = 1 - H(x)_y$, for which we need to derive random variables γ_u and $\gamma_u^{(p)}$, and most importantly, the function Φ .

WLOG, let $u \in \operatorname{Vert}(\mathcal{B})$ be $u = (1, \ldots, 1)$; arguments for other $u \in \operatorname{Vert}(\mathcal{B})$ proceeds similarly. For this u, $\gamma_u = \langle u, -\nabla \log q(\delta) \rangle, \delta \sim q$, is a sum of independent Rademacher random variable $\gamma_u = \sum_{i=1}^d R_i$, where each R_i independently takes values 1 and -1 with equal probability. Thus γ_u is distributed like $2B_d - d$, where B_d is the binomial random variable corresponding to the number of heads in d coin tosses. Then for any integer c with the same parity as d, $\phi_d(c)$ is the complementary CDF of γ_u and ϕ_d^{-1} is the corresponding inverse CDF. Then, for $p \in [0, 1/2]$, we have

$$\gamma_{u}^{(p)} = \begin{cases} \gamma_{u}|_{(c,\infty)} & \text{with probability } \phi_{d}(c) \\ c & \text{with probability } p - \phi_{d}(c) \\ 0 & \text{with probability } 1 - p, \end{cases}$$

where $c \stackrel{\text{def}}{=} \phi_d^{-1}(p)$ and $\gamma_u|_{(c,\infty)}$ is the random variable γ_u conditioned on $\gamma_u > c$. Therefore,

$$\Phi(p) = \mathbb{E} \gamma_u^{(p)} = c(p - \phi_d(c)) + 2^{-d} \sum_{i=\frac{c+d}{2}+1}^d (2i - d) \binom{d}{i}$$
$$= c(p - \phi_d(c)) + 2^{-d} \sum_{i=\frac{c+d}{2}+1}^d 2d \binom{d-1}{i-1} - d \binom{d}{i}$$
$$= c(p - \phi_d(c)) + d\phi_{d-1} \left(c - \frac{1}{2}\right) - d\phi_d(c)$$

Then, by setting G(x) in Theorem G.6 to be $1 - H(x)_y = 1 - \rho$, we get the desired robust radius.

Theorem P.3. Suppose *H* is a smoothed classifier smoothed by

$$q(x) \propto \exp(-\|x/\lambda\|_1^k), k > 1$$

such that $H(x) = (H(x)_1, \ldots, H(x)_C)$ is a vector of probabilities that H assigns to each class $1, \ldots, C$. If H correctly predicts the class y on input x, and the probability of the correct class is $\rho \stackrel{\text{def}}{=} H(x)_y > 1/2$, then H continues to predict the correct class when x is perturbed by any η with

$$\|\eta\|_\infty < \lambda \int_{1-\rho}^{1/2} \frac{1}{\Phi(p)} \,\mathrm{d} p,$$

where

$$\begin{split} \Phi(p) &= \mathbb{E} \gamma \mathbb{I}(\gamma > \varphi^{-1}(p)), \quad \text{with} \\ \gamma &= \left(\sum_{i=1}^{d} \zeta_i\right) k \xi^{\frac{k-1}{k}}, \\ \zeta_i &\sim \text{Rademacher}, \xi \sim \text{Gamma}(d/k) \\ \varphi(c) &= \mathbb{P}[\gamma > c]. \end{split}$$

Proof. By linearity in λ , it suffices to show this for $\lambda = 1$.

WLOG, let $u \in \operatorname{Vert}(\mathcal{B})$ be $u = (1, \ldots, 1)$; arguments for other $u \in \operatorname{Vert}(\mathcal{B})$ proceeds similarly. As in the proof of Theorem P.2, we find $\gamma_u = \langle u, -\nabla \log q(\delta) \rangle =$ $\langle u, k \| \delta \|_1^{k-1} \operatorname{sgn}(\delta) \rangle$, $\delta \sim q$, is distributed like γ in the theorem statement — a product of sum of Rademacher variables (coming from $\langle u, \operatorname{sgn}(\delta) \rangle$) and $k \xi^{\frac{k-1}{k}}$, $\xi \sim \Gamma(d/k)$ (coming from $k \| \delta \|_1^{k-1}$). Because k > 1, γ 's distribution is absolutely continuous (as it's a mixture of scaled versions of $\xi^{\frac{k-1}{k}}$'s distribution, which is absolutely continuous). Therefore, the random variable $\gamma_u^{(p)} = \gamma \mathbb{I}(\gamma > \varphi^{-1}(p))$. Then the theorem statement follows straightforwardly from Theorem G.6.

Q. Pareto Distribution

For $a, \lambda > 0$ and $u \in \mathbb{R}$, define the 0-centered, symmetric Pareto distribution by its PDF

Pareto
$$(x; a, \lambda) = \frac{a}{2\lambda} \left(1 + \left| \frac{x}{\lambda} \right| \right)^{-1-a}$$

= $\frac{a}{2\lambda} \exp \left[-(1+a) \log \left(1 + \left| \frac{x}{\lambda} \right| \right) \right]$

Its CDF is given by

ParetoCDF
$$(x; a, \lambda) = \begin{cases} 1 - \frac{1}{2} \left(1 + \left|\frac{x}{\lambda}\right|\right)^{-a} & \text{if } x > 0\\ \frac{1}{2} \left(1 + \left|\frac{x}{\lambda}\right|\right)^{-a} & \text{else.} \end{cases}$$

Consider smoothing distributions of the form

$$q(x) = \prod_{i=1}^{d} \operatorname{Pareto}(x_i; a, 1), \text{ so that}$$
$$\psi(x) = (1+a) \sum_{i=1}^{d} \log(1+|x_i|)$$
$$\nabla \psi(x) = \left\{ (1+a) \frac{\operatorname{sgn}(x_i)}{1+|x_i|} \right\}_{i=1}^{d},$$

with $\nabla \psi(x)$ defined when all coordinates x_i s are nonzero.

Q.1. ℓ_1 Adversary

Theorem Q.1. Suppose *H* is a smoothed classifier smoothed by

$$q(x) \propto \prod_{i=1}^{d} \operatorname{Pareto}(x_i; a, \lambda),$$

such that $H(x) = (H(x)_1, \ldots, H(x)_C)$ is a vector of probabilities that H assigns to each class $1, \ldots, C$. If H correctly predicts the class y on input x, and the probability of the correct class is $\rho \stackrel{\text{def}}{=} H(x)_y > 1/2$, then H continues to predict the correct class when x is perturbed by any η with

$$\|\eta\|_1 < \lambda \frac{2\rho - 1}{a} {}_2F_1\left(1, \frac{a}{a+1}; \frac{a}{a+1} + 1; (2\rho - 1)^{1+\frac{1}{a}}\right),$$

where $_2F_1$ is the hypergeometric function.

Proof. By linearity in λ , it suffices to show this for $\lambda = 1$. We seek to apply Theorem G.6 to $G(x) = 1 - H(x)_y$, for which we need to derive random variables γ_u and $\gamma_u^{(p)}$, and most importantly, the function Φ .

WLOG, assume $u \in \operatorname{Vert}(\mathcal{B})$ is e_1 . Then $\gamma_u = \langle u, \nabla \psi(\delta) \rangle = (1+a) \frac{\operatorname{sgn}(\delta_1)}{1+|\delta_1|}, \delta \sim q$, is distributed as $(1+a) \frac{\operatorname{sgn}(z)}{1+|z|}$ where $z \sim \operatorname{Pareto}(a,1)$. Therefore, for $p \in [0, 1/2]$, the random variable $\gamma_u^{(p)}$ defined in Definition G.4 is

$$\gamma_u^{(p)} = \begin{cases} \frac{1+a}{1+z_p} & \text{with prob. } p\\ 0 & \text{with prob. } 1-p, \end{cases}$$

where z_p is sampled from Pareto(a, 1) conditioned on the interval [0, ParetoCDF⁻¹(p + 1/2; a, 1)].

Thus, for any $u \in \operatorname{Vert}(\mathcal{B})$,

$$\Phi(p) = \mathbb{E} \gamma_u^{(p)} = \mathbb{E} \frac{1+a}{1+z} \mathbb{I}(z \in [0,c]),$$

where $z \sim \text{Pareto}(a, 1)$ and $c = \text{ParetoCDF}^{-1}(p + 1/2; a, 1)$. This can be simplified as follows:

$$\Phi(p) = \int_0^c \operatorname{Pareto}(r; a, 1) \frac{1+a}{1+r} \, \mathrm{d}r$$
$$= -\int_0^c \operatorname{Pareto}'(r; a, 1) \, \mathrm{d}r$$
$$= \operatorname{Pareto}(0; a, 1) - \operatorname{Pareto}(c; a, 1)$$
$$= \frac{a}{2} \left(1 - (1+c)^{-1-a} \right).$$

Note that for $p \in [1/2, 1]$,

ParetoCDF⁻¹
$$(p + 1/2; a, 1) = (1 - 2p)^{-1/a} - 1,$$

so $\Phi(p)$ can be further simplified:

$$\Phi(p) = \frac{a}{2} \left(1 - (1 - 2p)^{\frac{a+1}{a}} \right).$$

Then, by setting G(x) in Theorem G.6 to be $1 - H(x)_y = 1 - \rho$, we get the provably robust radius of

$$\int_{1-\rho}^{1/2} \frac{1}{\Phi(p)} dp = \int_{1-\rho}^{1/2} \frac{2}{a} \left(1 - (1-2p)^{\frac{a+1}{a}} \right)^{-1} dp$$
$$= \frac{2p-1}{a} {}_{2}F_{1} \left(1, \frac{a}{a+1}; \frac{a}{a+1} + 1; (1-2p)^{1+\frac{1}{a}} \right) \Big|_{1-\rho}^{1/2}$$
$$= \frac{2\rho-1}{a} {}_{2}F_{1} \left(1, \frac{a}{a+1}; \frac{a}{a+1} + 1; (2\rho-1)^{1+\frac{1}{a}} \right).$$

R. ℓ_2 -Norm Based Exponential Law

In this section we consider

$$\begin{split} q(x) \propto \exp(-\|x\|_2) & \text{so that} \\ \psi(x) = \|x\|_2 & \text{and} \quad \nabla \psi(x) = x/\|x\|_2, \end{split}$$

defined as long as $x \neq 0$.

R.1. ℓ_2 Adversary

Theorem R.1. Suppose *H* is a smoothed classifier smoothed by

$$q(x) \propto \exp(-\|x\|_2/\lambda),$$

such that $H(x) = (H(x)_1, \ldots, H(x)_C)$ is a vector of probabilities that H assigns to each class $1, \ldots, C$. If H correctly predicts the class y on input x, and the probability of the correct class is $\rho \stackrel{\text{def}}{=} H(x)_y > 1/2$, then H continues to predict the correct class when x is perturbed by any η with

$$\|\eta\|_{2} < \lambda(d-1)\operatorname{arctanh}\left(1-2\operatorname{BetaCDF}^{-1}\left(1-\rho;\frac{d-1}{2},\frac{d-1}{2}\right)\right).$$

Proof. By linearity in λ , it suffices to show this for $\lambda = 1$. We seek to apply Theorem G.6 to $G(x) = 1 - H(x)_y$, for which we need to derive random variables γ_u and $\gamma_u^{(p)}$, and most importantly, the function Φ .

For any $u \in \operatorname{Vert}(\mathcal{B})$ (i.e. any unit vector u), $\gamma_u = \langle u, \nabla \psi(\delta) \rangle = \langle u, \frac{\delta}{\|\delta\|_2} \rangle, \delta \sim q$, is distributed like $2\operatorname{Beta}\left(\frac{d-1}{2}, \frac{d-1}{2}\right) - 1$ by Lemma U.1. Its complementary CDF is given by

$$\begin{split} \mathbb{P}[\gamma_u > c] &= R \int_c^1 (1 - t^2)^{\frac{d-3}{2}} \,\mathrm{d}t \\ &= \mathrm{BetaCDF}\left(\frac{1 - c}{2}; \frac{d - 1}{2}, \frac{d - 1}{2}\right) \stackrel{\mathrm{def}}{=} \varphi(c), \end{split}$$

where $R \stackrel{\text{def}}{=} \frac{\Gamma(\frac{d}{2})}{\sqrt{\pi}\Gamma(\frac{d-1}{2})}$. Therefore, for $p \in [0,1]$, the random varible $\gamma_u^{(p)}$ defined in Definition G.4 is given by

$$\gamma_u^{(p)} = \begin{cases} \gamma_u|_{(\varphi^{-1}(p),\infty)} & \text{with probability } p \\ 0 & \text{with probability } 1-p, \end{cases}$$

Thus, for any $u \in \operatorname{Vert}(\mathcal{B})$,

$$\Phi(p) = \mathbb{E} \gamma_u^{(p)} = R \int_{\varphi^{-1}(p)}^1 t(1-t^2)^{\frac{d-3}{2}} dt$$
$$= \frac{R}{d-1} \left(1-\varphi^{-1}(p)^2\right)^{\frac{d-1}{2}}.$$

Then, by setting G(x) in Theorem G.6 to be $1 - H(x)_y = 1 - \rho$, we get the provably robust radius of

$$\int_{1-\rho}^{1/2} \frac{1}{\Phi(p)} dp$$

= $\frac{d-1}{R} \int_{1-\rho}^{1/2} \left(1 - \varphi^{-1}(p)^2\right)^{-\frac{d-1}{2}} dp$
= $(d-1) \int_0^{\varphi^{-1}(1-\rho)} (1 - c^2)^{-\frac{d-1}{2}} (1 - c^2)^{\frac{d-3}{2}} dc$
= $(d-1) \int_0^{\varphi^{-1}(1-\rho)} (1 - c^2)^{-1} dc$
= $(d-1) \operatorname{arctanh}(\varphi^{-1}(1-\rho)).$

Unpacking the definition of φ yields the result.

S. Uniform Distribution over a Sphere

Consider the distribution that is uniform on the ℓ_2 unit ball $\{x : ||x||_2 \le 1\}$. The *spherical cap* of height $h \le 1$ in this unit ball is the portion of the ball that is cut away by a hyperplane of distance 1 - h from the origin; see Fig. S.1. By Lemma U.2, this spherical cap has volume

$$V_d^h \stackrel{\text{def}}{=} V_d \text{BetaCDF}\left(\frac{h}{2}; \frac{d+1}{2}, \frac{d+1}{2}\right)$$

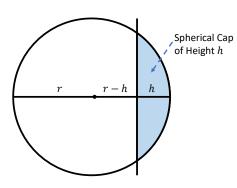


Figure S.1. Spherical Cap

where V_d is the volume of the unit sphere in \mathbb{R}^d .

Two unit radius spheres with centers ϵ apart intersects in a region that is the union of two spherical caps of height $1-\epsilon/2$. This intersection thus has volume $2V_d^{1-\epsilon/2}$, and the volume of one of the spheres outside this intersection is $V_d - 2V_d^{1-\epsilon/2} = V_d \left(1 - 2\text{BetaCDF}\left(\frac{1-\epsilon/2}{2}; \frac{d+1}{2}, \frac{d+1}{2}\right)\right)$.

Theorem S.1. Suppose *H* is a smoothed classifier smoothed by the uniform distribution *q* over a ball of radius λ centered at the origin, such that H(x) = $(H(x)_1, \ldots, H(x)_C)$ is a vector of probabilities that *H* assigns to each class $1, \ldots, C$. If *H* correctly predicts the class *y* on input *x*, and the probability of the correct class is $\rho \stackrel{\text{def}}{=} H(x)_y > 1/2$, then *H* continues to predict the correct class when *x* is perturbed by any η with

$$\|\eta\|_2 < \lambda \left(2 - 4 \text{BetaCDF}^{-1}\left(\frac{3}{4} - \frac{\rho}{2}; \frac{d+1}{2}, \frac{d+1}{2}\right)\right)$$

Proof. By linearity in λ , it suffices to show this for $\lambda = 1$.

By assumption, there is a region of probability ρ under the uniform distribution $q(x + \cdot)$ centered at x that the base classifier classifies as y. The intersection between the support of $q(x + \cdot)$ and $q(x + \delta + \cdot)$ for any $\|\delta\|_2 \leq \epsilon$ contains a region of probability at least

$$\rho - \left(1 - 2\text{BetaCDF}\left(\frac{1 - \epsilon/2}{2}; \frac{d+1}{2}, \frac{d+1}{2}\right)\right)$$

that the base classifier classifies as y. For this probability to be at least 1/2, we require

$$\begin{split} \frac{1}{2} &\leq \rho - (1 - 2 \text{BetaCDF}\left(\frac{1 - \epsilon/2}{2}; \frac{d + 1}{2}, \frac{d + 1}{2}\right))\\ \frac{3}{4} - \frac{\rho}{2} &\leq \text{BetaCDF}\left(\frac{1 - \epsilon/2}{2}; \frac{d + 1}{2}, \frac{d + 1}{2}\right)\\ 1 - \epsilon/2 &\geq 2 \text{BetaCDF}^{-1}\left(\frac{3}{4} - \frac{\rho}{2}; \frac{d + 1}{2}, \frac{d + 1}{2}\right)\\ \epsilon &\leq 2 - 4 \text{BetaCDF}^{-1}\left(\frac{3}{4} - \frac{\rho}{2}; \frac{d + 1}{2}, \frac{d + 1}{2}\right), \end{split}$$

as desired.

T. Level Set Method: General Distributions with Spherical Level Sets

T.1. ℓ_2 Adversary

Define $W_d(r, s, \epsilon)$ to be the probability a point sampled from the *surface* of a ball of radius *r* centered at the origin is outside a ball of radius *s* with center ϵ away from the origin. By Lemma U.1, we have

$$W_d(r, s, \epsilon) = \text{BetaCDF}\left(\frac{(r+\epsilon)^2 - s^2}{4\epsilon r}; \frac{d-1}{2}, \frac{d-1}{2}\right).$$
(23)

Note that W_d can be evaluated quickly using standard scipy functions.

Theorem T.1. Suppose that the density of a distribution satisfies $q(x) = \bar{q}(||x||_2)$ for some differentiable, decreasing function $\bar{q} : \mathbb{R}^{\geq 0} \to \mathbb{R}^{\geq 0}$. Then for any $\kappa > 0$ and any $v \in \mathbb{R}^d$, the growth function satisfies

$$\mathcal{G}_q(p_0, v) = p_1$$

$$p_0 = 1 - \mathop{\mathbb{E}}_{r} W_d(r, \bar{q}^{-1}(\bar{q}(r)/\kappa), ||v||_2)$$

$$p_1 = \mathop{\mathbb{E}}_{r} W_d(r, \bar{q}^{-1}(\bar{q}(r)\kappa), ||v||_2),$$

$$r \sim distribution with density \propto r^{d-1}\bar{q}(r).$$

For most \bar{q} , p_0 and p_1 can be evaluated numerically and quickly for each κ and $||v||_2$ using 1-dimensional integrals.

Proof. Let $r_t \stackrel{\text{def}}{=} \bar{q}^{-1}(t)$. Then the superlevel set $U_t = \{x : q(x) \ge t\}$ is a ball with radius r_t . Furthermore,

$$\nabla q(x) = \bar{q}'(\|x\|_2) \frac{x}{\|x\|_2}$$
$$\|\nabla q(x)\|_2^{-1} = -\bar{q}'(\|x\|_2)^{-1} = -r'_{\bar{q}(\|x\|_2)}$$

Let SA_d be the surface area of the unit sphere in \mathbb{R}^d . Then

$$q(\mathcal{NP}_{\kappa}) = -SA_d \int_0^\infty r'_t tr_t^{d-1} (1 - W_d(r_t, r_{t/\kappa}, \|v\|_2)) \,\mathrm{d}t$$
$$q(\mathcal{NP}_{\kappa} - v) = -SA_d \int_0^\infty r'_t tr_t^{d-1} W_d(r_t, r_{t\kappa}, \|v\|_2) \,\mathrm{d}t.$$

If we change coordinates from t to r, then

$$q(\mathcal{NP}_{\kappa}) = SA_d \times$$

$$\int_0^{\infty} \bar{q}(r)r^{d-1}(1 - W_d(r, \bar{q}^{-1}(\bar{q}(r)/\kappa), ||v||_2)) \, \mathrm{d}r$$

$$q(\mathcal{NP}_{\kappa} - v) = SA_d \times$$

$$\int_0^{\infty} \bar{q}(r)r^{d-1}W_d(r, \bar{q}^{-1}(\bar{q}(r)\kappa), ||v||_2) \, \mathrm{d}r.$$

Since $\bar{q}(r)r^{d-1}$ is proportional to the density of $||x||_2, x \sim q$, we can also write this as

$$q(\mathcal{NP}_{\kappa}) = \underset{r=\|x\|_{2,x\sim q}}{\mathbb{E}} (1 - W_{d}(r, \bar{q}^{-1}(\bar{q}(r)/\kappa), \|v\|_{2}))$$
$$q(\mathcal{NP}_{\kappa} - v) = \underset{r=\|x\|_{2,x\sim q}}{\mathbb{E}} W_{d}(r, \bar{q}^{-1}(\bar{q}(r)\kappa), \|v\|_{2}).$$
(24)

The distribution of r here has density $\propto r^{d-1}\bar{q}(r)$. Then setting $p_0 = q(\mathcal{NP}_{\kappa}), p_1 = q(\mathcal{NP}_{\kappa} - v)$ yields the desired result by Eq. (NP).

Example T.2. If $q(x) \propto ||x||_2^{-j} \exp(-||x||_2^k)$, then $\bar{q}(r) \propto r^{-j} \exp(-r^k)$, and the radius is distributed as $\sqrt[k]{\text{Gamma}(d/k-j/k)}$ by Lemma U.3. A table of robust radii can then be built according to Algorithm 1, and certification can be done via Algorithm 2.

Example T.3. If $q(x) \propto (1 + ||x||_2^k)^{-a}$, then $\bar{q}(r) \propto (1 + r^k)^{-a}$, and the radius is distributed as $\sqrt[k]{\text{BetaPrime}(d/k, a - d/k)}$. A table of robust radii can then be built according to Algorithm 1, and certification can be done via Algorithm 2.

U. Tools

Lemma U.1. If (x_1, \ldots, x_d) is sampled uniformly from the unit sphere $S^{d-1} \subseteq \mathbb{R}^d$, then

$$\frac{1+x_1}{2}$$
 is distributed as Beta $\left(\frac{d-1}{2}, \frac{d-1}{2}\right)$,

i.e.

$$\mathbb{P}[x_1 \ge c] = \text{BetaCDF}\left(\frac{1-c}{2}; \frac{d-1}{2}, \frac{d-1}{2}\right)$$
$$= 1 - \text{BetaCDF}\left(\frac{1+c}{2}; \frac{d-1}{2}, \frac{d-1}{2}\right).$$

Proof. From simple geometric reasoning, we get

$$\mathbb{P}[x_1 \ge c] \propto \int_c^1 (1-t^2)^{\frac{d-3}{2}} dt$$

= $\int_c^1 (1-t)^{\frac{d-3}{2}} (1+t)^{\frac{d-3}{2}} dt$
= $\int_0^{\frac{1-c}{2}} (2x)^{\frac{d-3}{2}} (2(1-x))^{\frac{d-3}{2}} dx.$

Lemma U.2. If (x_1, \ldots, x_d) is sampled uniformly from the ball $\{y : ||y||_2 \le 1\} \subseteq \mathbb{R}^d$, then

$$\frac{1+x_1}{2}$$
 is distributed as Beta $\left(\frac{d+1}{2}, \frac{d+1}{2}\right)$,

i.e.

$$\mathbb{P}[x_1 \ge c] = \text{BetaCDF}\left(\frac{1-c}{2}; \frac{d+1}{2}, \frac{d+1}{2}\right)$$
$$= 1 - \text{BetaCDF}\left(\frac{1+c}{2}; \frac{d+1}{2}, \frac{d+1}{2}\right).$$

Proof. Similar to Lemma U.1.

Lemma U.3. For any norm $\|\cdot\|$ on \mathbb{R}^d , the distribution

$$q(x) \propto ||x||^{-j} \exp(-||x||^k)$$

with j < d, can be sampled as follows:

- 1. Sample the radius $r \sim \sqrt[k]{\text{Gamma}(\frac{d}{k} \frac{j}{k})}$
- 2. Sample a point v from the unit sphere of $\|\cdot\|$
- *3. return* rv

Lemma U.4. For any $c, s \ge 0$ and r > 0,

$$\begin{split} & \underset{z \sim \operatorname{Gamma}(r)}{\mathbb{E}} z^{s} \mathbb{I}(z > c) \\ & = \frac{\Gamma(r+s)}{\Gamma(r)} (1 - \operatorname{GammaCDF}(c;r+s)) \end{split}$$

V. Proof of Theorem 7.3

In this section we prove our main impossibility result, Theorem 7.3. We will assume throughout this proof that the reader is familiar with standard notions in functional analysis. Our proof will proceed in two steps.

For all $p \in (0, 2]$ and $d' \ge 1$, we let $\ell_p^{d'}$ denote $\mathbb{R}^{d'}$ equipped with the *p*-quasinorm. First, we show that if there exists a useful smoothing scheme, this implies a low embedding distortion of our normed space into $\ell_{0,99}^{d'}$, for some d'.

Formally, let (X, d_X) and (Y, d_Y) be two metric spaces. We say an embedding $f : X \to Y$ has distortion D if there exist constants positive $\alpha < 1 < \beta$ so that

$$\alpha d_Y(f(x_1), f(x_2)) \le d_X(x_1, x_2) \le \beta d_Y(f(x_1), f(x_2))$$

for all $x_1, x_2 \in X$, where $\beta/\alpha \leq D$. We first show:

Theorem V.1. Suppose there exists an (ε, L, ℓ) -useful smoothing distribution for $\|\cdot\|$, and $L/\ell \leq 1/(16\varepsilon)$. Then, there exists d' and a linear embedding from $(\mathbb{R}^d, \|\cdot\|)$ into $\ell_{0.99}^{d'}$ with distortion at most

$$O\left(\left(\frac{L}{\ell}\right)^{1/4} \left(\frac{1}{\varepsilon}\right)^{3/4}\right) \,.$$

Next we show that any linear embedding into $\ell_{0.99}^{d'}$ will suffer distortion which is at least $C_2((\mathbb{R}^d, \|\cdot\|))$:

Theorem V.2. Suppose there exists a linear embedding from $(\mathbb{R}^d, \|\cdot\|)$ into $\ell_{0.99}^{d'}$ with distortion at most γ . Then $\gamma \geq C_2((\mathbb{R}^d, \|\cdot\|))$.

This result is essentially folklore in the metric embedding community, but we include a proof for completeness.

Combining these two theorems directly implies Theorem 7.3. The rest of this section is dedicated to proofs of these two Theorems.

V.1. Proof of Theorem V.1

Let q be an (ε, L, ℓ) -useful smoothing distribution for $\|\cdot\|$. For simplicity, throughout this proof, we will assume that q has a probability density function, denoted Q, that is, it is absolutely continuous with respect to the Lebesgue measure. It is not hard to generalize this proof to handle general probability distributions by taking the appropriate Radon-Nikodym derivatives.

First, we demonstrate that a useful smoothing scheme actually implies an embedding of the norm $\|\cdot\|$ into an infinite dimensional L_1 space, namely, the space of all distributions with distance given by total variation distance. Recall the total variation distance between two distributions q_1, q_2 , denoted $d_{\text{TV}}(q_1, q_2)$, is given by

$$d_{\mathrm{TV}}(q_1, q_2) = \sup_{U \subseteq \mathbb{R}^d} |q_1(U) - q_2(U)| = \frac{1}{2} ||Q_1 - Q_2||_1 .$$

We denote the space of probability distributions over \mathbb{R}^d by Δ_d . Note that the metric space $(\Delta_d, d_{\text{TV}})$ is an infinite dimensional L_1 space. Thus, classical results yield:

Fact V.3 (see e.g. Wojtaszczyk (1996)). For all $d \ge 1$, we have $C_2((\Delta_d, d_{\text{TV}})) = \Theta(1)$.

For all $x \in \mathbb{R}^d$, let q_x be the distribution given by $q_x(U) = q(U-x)$ for all $U \subseteq \mathbb{R}^d$, and let Q_x be the associated pdf. We will also let $F_q : \mathbb{R}^d \to \Delta_d$ be the map $F_q(x) = q_x$ which takes a point in \mathbb{R}^d and maps it to its associated distribution after smoothing.

The following lemma will be useful:

Lemma V.4. Let q be as above. Then, for all x, y, we have:

$$d_{\rm TV}(q_x, q_y) = \sup_{\substack{U \subseteq \mathbb{R}^d \\ q_x(U) \ge 1/2, q_y(U) \le 1/2}} q_x(U) - q_y(U) \ . \ (25)$$

Proof. That the left-hand side is at least the right-hand side is obvious, since the former is simply taking a supremum over a larger set. We now show the other direction. Define the following three disjoint sets:

$$U^{+} = \{ z \in \mathbb{R}^{d} : Q_{x}(z) > Q_{y}(z) \} ,$$

$$U^{-} = \{ z \in \mathbb{R}^{d} : Q_{y}(z) > Q_{x}(z) \} ,$$

$$U^{-} = \{ z \in \mathbb{R}^{d} : Q_{x}(z) = Q_{y}(z) \} .$$

Then, it is a classical fact that $q_x(U^+) - q_y(U^+) = q_y(U^-) - q_x(U^-) = d_{\text{TV}}(q_x, q_y)$. By symmetry we have that $U^- = U^+ + x - y$, and hence $q_x(U^+) = q_y(U^-)$ and similarly $q_x(U^-) = q_y(U^+)$. Therefore $q_x(U^+) + q_y(U^+) \le 1$, but on the other hand, $q_x(U^+ \cup U^-) + q_y(U^+ \cup U^-) \ge 1$. Therefore there exists some set $W \subseteq U^=$ so that $q_x(U^+ \cup W) + q_y(U^+ \cup W) = 1$. Letting $U = U^+ \cup W$, we have that $q_x(U) - q_y(U) = d_{\text{TV}}(q_x, q_y)$, but also $q_x(U) \ge 1/2$ and $q_y(U) \le 1/2$. This proves that the RHS of (25) is at least the LHS of (25), which proves the equality.

We now need another notion, introduced in Andoni et al. (2018).

Definition V.5 (Andoni et al. (2018)). A map $f : X \to Y$ between two metric spaces (X, d_X) and (Y, d_Y) is an $(s_1, s_2, \tau_1, \tau_2)$ -threshold map if it satisfies:

- If $d_X(x_1, x_2) \le s_1$, then $d_Y(f(x_1), f(x_2)) \le \tau_1$.
- If $d_X(x_1, x_2) \ge s_2$, then $d_Y(f(x_1), f(x_2)) \ge \tau_2$.

We observe the following:

Lemma V.6. Let q be a (ε, L, ℓ) -useful smoothing distribution for $\|\cdot\|$. Then F_q is a $(\varepsilon, 1, L\varepsilon, \ell)$ -threshold map between $(\mathbb{R}^d, \|\cdot\|)$ and (Δ, d_{TV}) .

Proof. We first prove that it satisfies the first condition. Let x, y be so that $||x - y|| \le \varepsilon$. Then, the robustness condition implies that

$$\sup_{\substack{U\subseteq \mathbb{R}^d\\q_x(U)\geq 1/2, q_y(U)\leq 1/2}} q_x(U) - q_y(U) \leq L\varepsilon \ .$$

By Lemma V.4 this implies that $d_{\text{TV}}(q_x, q_y) \leq L\varepsilon$.

We now prove it satisfies the second condition. But, the accuracy condition immediately implies that if x, y satisfy $||x - y|| \ge 1$, we must have $d_{\text{TV}}(q_x, q_y) \ge \ell$. This proves the claim.

We now use the following basic fact, which follows since L_1 embeds isometrically into squared- L_2 .

Fact V.7 (see e.g. Matoušek (2013)). Suppose there exists an $(s_1, s_2, \tau_1, \tau_2)$ -threshold map from $(\mathbb{R}^d, \|\cdot\|)$ to $(\Delta_d, d_{\text{TV}})$. Then there exists an $(s_1, s_2, \tau_1^{1/2}, \tau_2^{1/2})$ -threshold map from $(\mathbb{R}^d, \|\cdot\|)$ to a Hilbert space H.

This implies:

Corollary V.8. Suppose there exists an (ε, L, ℓ) -useful smoothing distribution for $\|\cdot\|$. Then there exists a $(\varepsilon, 1, L\varepsilon, \ell)$ -threshold map between $(\mathbb{R}^d, \|\cdot\|)$ and a Hilbert space H.

Combining Corollary V.8 and Theorem 4.12 in (Andoni et al., 2018), we obtain:

Corollary V.9. Suppose there exists an (ε, L, ℓ) -useful smoothing distribution for $\|\cdot\|$. Then, there exists a Hilbert space H with induced norm $\|\cdot\|_H$ and a map $h : \mathbb{R}^d \to H$ so that for all $x, y \in \mathbb{R}^d$, we have

$$\|x_1 - x_2\|^{1/2} \le \|h(x) - h(y)\|_2 \le \left(\frac{\ell}{4L\varepsilon}\right)^{1/4} \cdot \min(1,\varepsilon \|x_1 - x_2\|)$$

Finally, combining this with Theorem 5.1 in (Andoni et al., 2018), and our assumption that $L/\ell \leq 1/(16\varepsilon)$, we immediately obtain Theorem V.1.

V.2. Proof of Theorem V.2

We now turn to the proof of Theorem V.2. The only reason why this is slightly non-standard is that $\ell_p^{d'}$ for p < 1 are not norms, as they do not satisfy the triangle inequality. Despite this, we show that the standard results that the cotype constant is a lower bound on distortion of any linear embedding still holds for these semi-norms.

Theorem V.10 (Khintchine's Inequality). For any $p \in (0, \infty)$ there exist constants A_p, B_p such that for any $x_1, \ldots, x_n \in \mathbb{R}$

$$A_p\left(\sum_{i=1}^n x_i^2\right)^{1/2} \le \left(\mathbb{E}\left|\sum_{i=1}^n \sigma_i x_i\right|^p\right)^{1/p} \le B_p\left(\sum_{i=1}^n x_i^2\right)^{1/2}$$

Here $\sigma_1, \ldots, \sigma_n$ are independent Rademacher random variables. In particular, for $0 , <math>A_p = 2^{1/2-1/p}$.

The tight value for A_p is due to Haagerup.

Theorem V.11 (Cotype Estimate). For any $p \in (0, 1]$, there exists a universal constant so that $\ell_p^{d'}$ has cotype 2 with cotype constant α_p , i.e. for any $x_1 \dots, x_n \in \mathbb{R}^{d'}$

$$\left(\sum_{i=1}^{n} \|x_i\|_p^2\right)^{1/2} \le \alpha_p \left(\mathbb{E} \left\|\sum_{i=1}^{n} \sigma_i x_i\right\|_p\right)$$

where $\sigma_1, \ldots, \sigma_n$ are independent Rademacher random variables.

Proof. Let x_{ij} denote coordinate j of x_i , i.e., X is the $n \times d$ matrix whose *rows* are x_1, \ldots, x_n . By Khintchine's inequality,

$$\mathbb{E} \left\| \sum_{i=1}^{n} \sigma_{i} x_{i} \right\|_{p}^{p} = \sum_{j=1}^{d} \mathbb{E} \left| \sum_{i=1}^{n} \varepsilon_{i} x_{ij} \right|^{p}$$
$$\geq A_{p}^{p} \sum_{j=1}^{d} \left(\sum_{i=1}^{n} x_{ij}^{2} \right)^{p/2}$$

Let us now consider the case $p \leq 2$. By the triangle inequality for $\|\cdot\|_q$, where $q = 2/p \geq 1$, applied to the vectors $(|x_{1j}|^p, |x_{2j}|^p, \dots, |x_{nj}|^p), j \in [d],$

$$\sum_{j=1}^{d} \left(\sum_{i=1}^{n} x_{ij}^{2}\right)^{p/2} = \sum_{j=1}^{d} \left(\sum_{i=1}^{n} (|x_{ij}|^{p})^{2/p}\right)^{p/2}$$
$$\geq \left(\sum_{i=1}^{n} \left|\sum_{j=1}^{d} |x_{ij}|^{p}\right|^{2/p}\right)^{p/2}$$
$$= \left(\sum_{i=1}^{n} ||x_{i}||_{p}^{2}\right)^{p/2}.$$

Notice that we can take $\alpha_p = 1/A_p$. The proof follows by convexity of the function $x \mapsto x^{1/p}$ for $p \in (0, 1]$.

We now have all the tools we need to prove Theorem V.2:

Proof of Theorem V.2. Let $T : \mathbb{R}^d \to \mathbb{R}^{d'}$ be any linear map satisfying

$$\alpha \|Tx\|_p \le \|x\| \le \beta \|Tx\|_p$$

Let $C_2 = C_2((\mathbb{R}^d, \|\cdot\|))$, and let x_1, \ldots, x_n be a sequence in $(\mathbb{R}^d, \|\cdot\|)$ satisfying

$$\mathbb{E}\left[\left\|\sum_{j=1}^{n}\sigma_{j}x_{j}\right\|\right] = C_{2}\left(\sum_{j=1}^{n}\|x_{j}\|^{2}\right)^{1/2}.$$

However, we have that

$$\mathbb{E}\left[\left\|\sum_{j=1}^{n}\sigma_{i}Tx_{j}\right\|_{p}\right] \leq \beta \cdot \mathbb{E}\left[\left\|\sum_{j=1}^{n}\sigma_{j}x_{j}\right\|\right],$$

and simultaneously, we have

$$\left(\sum_{j=1}^{n} \|Tx_i\|^2\right)^{1/2} \ge \alpha \left(\sum_{j=1}^{n} \|Tx_i\|^2\right)^{1/2}$$

Combining these facts and Theorem V.11, we obtain that $\beta/\alpha \ge \alpha_2 C_2 = \Omega(C_2)$, as claimed.