# Optimization: Algorithms, Complexity \& Approximations 

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## Chapter 2

This lecture introduces smooth continuous optimization and provides the background knowledge before we delve into more specialized classes of objectives, such as convex optimization. To explore these topics, we will require several basic definitions such as gradients, Hessian matrices, Taylor Series, etc. This chapter also "scratches" the surface of properties of optimization functions: for instance, the Taylor expansion is reviewed, and types of stationary points are introduced. Several special conditions that benefit optimization, including Lipschitz and Lipschitz gradient continuity, are introduced.
The main algorithm for this chapter will be gradient descent (GD), as well as projected GD. Additionally, these notes explain convergence rates. We will see how further global assumptions lead to improved convergence guarantees.

## Lipschitz conditions | Gradient Descent

This course covers general smooth optimization, where the objective function can be pictured as a continuous curve in high dimensions. You can easily picture it: A continuous landscape parameterized by a set of unknowns, and the goal is to find the global minimum/maximum. However, other important classes of optimization problems not covered in this course follow this description, as shown in the figure 1. Some of them are typically explored as a particular topic, for example, discrete optimization and integer programming. This course is restricted only to smooth functions. The smoothness will be defined later on in the text. For now, one way to describe smoothness is by saying that we can compute gradients on these functions.


Fig. 1. Landscape of optimization

Derivatives, gradients, and Hessians. Algorithms and heuristics in optimization often involve derivatives to approach an optimal solution. Put shortly; the derivative tells you the direction (and, in some way, the magnitude) of the steepest ascent (or descent).

Definition 1. (First-order Derivative) The derivative of a univariate function $f: \mathbb{R} \rightarrow \mathbb{R}$ at a point $x$ is defined as:

$$
\frac{\partial f}{\partial x}=f^{\prime}(x)=\lim _{\epsilon \rightarrow 0} \frac{f(x+\epsilon)-f(x)}{\epsilon} .
$$

The derivative of $f$ represents the slope $f$ in a neighborhood of a point $x$. It explains how much $f$ changes within a small area when we perturb around a given point.

This suggests the second-order derivative, which is recursively defined as the derivative of the derivative and describes how rapidly the derivative changes.
Def inition 2. (Second-order Derivative) The second-order derivative of a univariate function $f: \mathbb{R} \rightarrow \mathbb{R}$ at a point $x$ is defined as:

$$
\frac{\partial^{2} f}{\partial x^{2}}=f^{\prime \prime}(x)=\lim _{\epsilon \rightarrow 0} \frac{f^{\prime}(x+\epsilon)-f^{\prime}(x)}{\epsilon}
$$

The second-order derivative represents the local curvature of $f$, i.e., how much the function's slope changes around a given point.

Some differentiation rules are:

- $(f(x) \cdot g(x))^{\prime}=f^{\prime}(x) \cdot g(x)+f(x) \cdot g^{\prime}(x)$
(Product rule)
- $\left(\frac{f(x)}{g(x)}\right)^{\prime}=\frac{f^{\prime}(x) \cdot g(x)-f(x) \cdot g^{\prime}(x)}{g^{2}(x)}$ (Quotient rule)
- $(f(x)+g(x))^{\prime}=f^{\prime}(x)+g^{\prime}(x)$
(Sum rule)
- $(f(g(x)))^{\prime}=(f \circ g)^{\prime}(x)=f^{\prime}(g(x)) \cdot g^{\prime}(x) \quad$ (Chain rule)
- $(c)^{\prime}=0$
(Derivative of a constant is zero)
- $(c \cdot f(x))^{\prime}=c \cdot f^{\prime}(x)$
(Mult. by a constant rule)
(Power rule)
The notions of derivatives have a natural generalization to higher dimensional cases. In particular, we will start by introducing the idea of a gradient.
Definition 3. (Gradient of $f$ ) The gradient of a multivariate function $f: \mathbb{R}^{p} \rightarrow \mathbb{R}$ is

$$
\nabla f(x)=\left[\begin{array}{c}
\frac{\partial f}{\partial x_{1}} \\
\vdots \\
\frac{\partial f}{\partial x_{p}}
\end{array}\right] \in \mathbb{R}^{p}
$$



Fig. 2. Graphical illustration of first-order derivative
where
$\frac{\partial f}{\partial x_{i}}$
$=\lim _{\epsilon \rightarrow 0} \frac{f\left(\ldots, x_{i-1}, x_{i}+\epsilon, x_{i+1}, \ldots\right)-f\left(\ldots, x_{i-1}, x_{i}, x_{i+1}, \ldots\right)}{\epsilon}$
$=\frac{f\left(x+\epsilon e_{i}\right)-f(x)}{\epsilon}$,
where $e_{i} \in \mathbb{R}^{p}$ denotes the basis/coordinate vector where all elements are zero, except for the one in the $i$-th position with value one.

The following definition computes the rate at which a function $f$ changes at a point $x$ in the direction of an arbitrary vector $y$. This relates linear forms of the gradient (i.e., inner product) to a one-dimensional derivative evaluated at zero.
Def inition 4. (First-order Directional Derivative) Let $f: \mathbb{R}^{p} \rightarrow \mathbb{R}$ be a differentiable function. For two points $x, y \in \mathbb{R}^{p}$ and for scalar $\gamma$, we have:

$$
\nabla_{y} f(x)=\nabla f(x)^{\top} y=\lim _{\gamma \rightarrow 0} \frac{f(x+\gamma y)-f(x)}{\gamma}
$$

$\nabla_{y} f(x)$ is called the directional derivative of $f$ at $x$ in the direction of $y$.
To verify this formula, let us first define the "helper" function:

$$
\varphi(\gamma):=f(x+\gamma y)=f(\psi(\gamma))
$$

where $\psi(\gamma):=x+\gamma y$. Computing the gradient of $\varphi(\gamma)$ with respect to $\gamma$ is equivalent to computing the gradient of $f$ along the direction $y$, for infinitesimal $\gamma$. In particular, by applying the chain rule, we obtain:

$$
\begin{aligned}
\varphi^{\prime}(\gamma) & =\sum_{i=1}^{p} \frac{\partial f(\psi(\gamma))}{\partial \psi_{i}} \cdot \nabla \psi_{i}(\gamma) \\
& =\sum_{i=1}^{p} \frac{\partial f(\psi(\gamma))}{\partial \psi_{i}} \cdot y_{i} \\
& =\langle\nabla f(\psi(\gamma)), y\rangle \\
& =\langle\nabla f(x+\gamma y), y\rangle
\end{aligned}
$$

Then, we obtain the definition of the directional derivative when we set $\gamma=0$.

The directional derivative is also often written in the notation:

$$
\nabla_{y} f(x)=y_{1} \cdot \frac{\partial f}{\partial x_{1}}+y_{2} \cdot \frac{\partial f}{\partial x_{2}}+\cdots+y_{p} \cdot \frac{\partial f}{\partial x_{p}}=\sum_{i=1}^{p} y_{i} \cdot \frac{\partial f}{\partial x_{i}}
$$

Next, we will define the derivative for a multivariate vector function.
Definition 5. (Jacobian of a function $f$ ) The Jacobian of a multivariate vector function $f: \mathbb{R}^{p} \rightarrow \mathbb{R}^{m}$ is given by:

$$
D f(x)=\left[\begin{array}{cccc}
\frac{\partial f_{1}}{\partial x_{1}} & \frac{\partial f_{1}}{\partial x_{2}} & \cdots & \frac{\partial f_{1}}{\partial x_{p}} \\
\vdots & \vdots & & \vdots \\
\frac{\partial f_{m}}{\partial x_{1}} & \frac{\partial f_{m}}{\partial x_{2}} & \cdots & \frac{\partial f_{m}}{\partial x_{p}}
\end{array}\right] \in \mathbb{R}^{m \times p}
$$

Loosely speaking, taking the Jacobian of the gradient yields the Hessian, which contains the second-order local information about $f$ :
Def inition 6. (Hessian matrix of $f$ ) The Hessian of a multivariate function $f: \mathbb{R}^{p} \rightarrow \mathbb{R}$ is

$$
\nabla^{2} f(x)=\left[\begin{array}{cccc}
\frac{\partial^{2} f}{\partial x_{1}^{2}} & \frac{\partial^{2} f}{\partial x_{1} \partial x_{2}} & \cdots & \frac{\partial^{2} f}{\partial x_{1} \partial x_{p}} \\
\vdots & \vdots & & \vdots \\
\frac{\partial^{2} f}{\partial x_{p} \partial x_{1}} & \frac{\partial^{2} f}{\partial x_{p} \partial x_{2}} & \cdots & \frac{\partial^{2} f}{\partial x_{p}^{2}}
\end{array}\right]
$$

The Hessian matrix of a continuous function is symmetric. The Hessian matrix provides information about the curvature of the function $f$. For example, given a point $x^{\star}$, when $\nabla^{2} f\left(x^{\star}\right) \succ 0$ holds, then $x^{\star}$ is (at least) a strict local minimizer of $f$. Alternatively, when $\nabla^{2} f\left(x^{\star}\right) \prec 0$, then $x^{\star}$ is a strict local maximizer of $f$. See figure 5 for a geometric interpretation of the above facts.
Similarly to gradients, we can relate quadratic forms of the Hessian matrix to one-dimensional derivatives.

Definition 7. (Second-order Directional Derivative) Let $f: \mathbb{R}^{p} \rightarrow$ $\mathbb{R}$ be a twice-differentiable function. Let $x, y \in \mathbb{R}^{p}$ and $\gamma$ a scalar. Then:
$\left\langle\nabla^{2} f(x+\gamma y) \cdot y, y\right\rangle=\lim _{\gamma \rightarrow 0} \frac{\nabla f(x+\gamma y)^{\top} y-\nabla f(x)^{\top} y}{\gamma}=\frac{\partial^{2} f(x+\gamma y)}{\partial \gamma^{2}}$.

Taylor expansion of a function $f$. Now that we know what derivates, gradients, and Hessians are, how can we use them in practice? The answer to this question will come from answering the following question: Are there any intuitive ways of approximating the behavior of a function, even locally? The answer is Yes: the Taylor expansion of the function may be used to approximate the function locally.
Definition 8. (Taylor Series) Assuming that $f$ is $n$-times differentiable, then the Taylor series of $f$ centered at $\alpha$ is:

$$
\begin{aligned}
T_{\alpha}(x) & =\sum_{k=0}^{\infty} \frac{f^{(k)}(\alpha)(x-\alpha)^{k}}{k!} \\
& =\frac{f(\alpha)}{0!}+\frac{f^{\prime}(\alpha)}{1!}(x-\alpha)+\frac{f^{\prime \prime}(\alpha)}{2!}(x-\alpha)^{2}+\cdots
\end{aligned}
$$

The $k$-th order Taylor approximation is the above series truncated at the $k^{\text {th }}$ term in the sum.

Here, $f$ is assumed to be differentiable as often as we would like. For the rest of this course, we will assume that our functions are differentiable unless stated otherwise. More often than not, we will focus on the up-to-2nd-order Taylor approximation of functions. We note that the Taylor expansion gives a reasonable (local) estimate of the function. When we keep only the first two terms, we call it a linear approximation of the function near $\alpha$, as is illustrated in figure 3.


Fig. 3. The first-order Taylor expansion provides a good estimation of the function near the point $\alpha$ but easily drifts away when we move a little bit away from it.

When we keep the first three terms, we obtain a quadratic approximation of $f$, as is illustrated in figure 4.


Fig. 5. How Hessian looks around interesting points of a two-dimensional function $f$ ( z -axis).


Fig. 4. The second-order Taylor expansion estimates a function better near point $\alpha$.

Adding more terms provides a more accurate approximation, and for a univariate function, this is attainable. However, the complexity increases significantly in high-order Taylor expansion of multivariate functions.
Definition 9. The Taylor expansion of a multivariate function $f: \mathbb{R}^{p} \rightarrow \mathbb{R}$ at point $\alpha \in \mathbb{R}^{p}$ is
$f(x) \approx f(\alpha)+\langle\nabla f(\alpha), x-\alpha\rangle+\frac{1}{2}\left\langle\nabla^{2} f(\alpha)(x-\alpha),(x-\alpha)\right\rangle+\ldots$
This is a natural generalization of the one-dimensional version. For a first-order Taylor expansion approximation, we obtain:

$$
f(x) \approx f(\alpha)+\langle\nabla f(\alpha), x-\alpha\rangle, \quad \alpha \in \mathbb{R}^{p},
$$

while for a second-order one, we obtain:
$f(x) \approx f(\alpha)+\langle\nabla f(\alpha), x-\alpha\rangle+\frac{1}{2}\left\langle\nabla^{2} f(\alpha)(x-\alpha), x-\alpha\right\rangle, \quad \alpha \in \mathbb{R}^{p}$.
For further discussions, the following fundamental theorem of calculus (part II) is useful: it will help show that the differentiation in the multivariate setting can be expressed as integrals of univariate functions. The fundamental theorem reads as follows:
Def inition 10. (Fundamental theorem of calculus, part II) Let $f$ : $[\alpha, \beta] \rightarrow \mathbb{R}$ be a continuously differentiable function. Then:

$$
\int_{\alpha}^{\beta} \frac{d}{d t} f(t) d t=f(\beta)-f(\alpha)
$$

Based on the above, Taylor's expansion implies the following:
Lemma 1. Let $f: \mathbb{R}^{p} \rightarrow \mathbb{R}$ be a differentiable function. Let two points $x, y \in \mathbb{R}^{p}$. Then:

$$
f(y)=f(x)+\langle\nabla f(x), y-x\rangle+\int_{0}^{1}(1-\gamma) \frac{\partial^{2} f(x+\gamma(y-x))}{\partial \gamma^{2}} d \gamma
$$

The above provides an idea of a local approximation of a function. This leads to the Taylor's theorem, often called the multivariate mean-value theorem. Taylor's theorem below allows the approximation of smooth functions by simple polynomials.

## Theorem 1. (Taylor's theorem)

- If $f$ is continuously differentiable, then:
$f(w)=f\left(w_{0}\right)+\left\langle\nabla f\left(t w+(1-t) w_{0}\right),\right\rangle, \quad$ for some $t \in[0,1]$.
- If $f$ is twice differentiable, then:

$$
\nabla f(w)=\nabla f\left(w_{0}\right)+\int_{0}^{1} \nabla^{2} f\left(t w+(1-t) w_{0}\right) \cdot\left(w-w_{0}\right) d t
$$

- Further, if $f$ is twice differentiable, then, for some $t \in[0,1]$ :

$$
\begin{aligned}
f(w)= & f\left(w_{0}\right)+\left\langle\nabla f\left(w_{0}\right), w-w_{0}\right\rangle \\
& +\frac{1}{2}\left\langle\nabla^{2} f\left(t w+(1-t) w_{0}\right) \cdot\left(w-w_{0}\right), w-w_{0}\right\rangle .
\end{aligned}
$$

But how are the above useful in optimization?. Consider that someone gives you the following problem $\min _{x} f(x)$ for some function $f$. Further, we are told that computing gradients $\nabla f(\cdot)$ and Hessians $\nabla^{2} f(\cdot)$ is relatively easy. Then, assuming we start from a point $x_{0}$, instead of worrying about $f$ itself, one can do the following steps:

- Compute gradient $\nabla f\left(x_{0}\right)$; name this as the $h$ vector.
- Compute Hessian $\nabla^{2} f\left(x_{0}\right)$; name this as the $H$ matrix.
- Form the second-order Taylor approximation:

$$
\begin{aligned}
& f(x) \\
& \approx f\left(x_{0}\right)+\left\langle\nabla f\left(x_{0}\right), x-x_{0}\right\rangle+\frac{1}{2}\left\langle\nabla^{2} f\left(x_{0}\right) \cdot\left(x-x_{0}\right), x-x_{0}\right\rangle \\
& =f\left(x_{0}\right)+h^{\top}\left(x-x_{0}\right)+\frac{1}{2}\left(x-x_{0}\right)^{\top} H\left(x-x_{0}\right) .
\end{aligned}
$$

Hence, instead of optimizing directly $\min _{x} f(x)$, we first compute the second-order approximation around a point $x_{0}$ :
$\min _{x}\left\{f\left(x_{0}\right)+\left\langle\nabla f\left(x_{0}\right), x-x_{0}\right\rangle+\frac{1}{2}\left\langle\nabla^{2} f\left(x_{0}\right)\left(x-x_{0}\right), x-x_{0}\right\rangle\right\}$,
which in turn is just a minimization of a quadratic function:

$$
\min _{x}\left\{h^{\top} x+\frac{1}{2} x^{\top} H x\right\}
$$

Solving quadratic problems is a type of optimization we can efficiently compute.
The above list suggests that regardless of how difficult $f$ is to optimize, one can approximate it through Taylor's expansion to get to a problem that we can solve: that of a quadratic
objective! Of course, this does not guarantee that we will get the optimum of $f$. E.g., if $x_{0}$ is far from the optimal $x^{\star}$ and the local quadratic approximation does not follow well $f$, then we have no hope of optimizing the original $f$ function. However, we can make this happen by using iterative procedures over the above motions for ever-improved $x$ points.

Optima. It is always easy to spot the minimum of a function whenever it can be drawn on paper. For a computer, though, this is a complicated problem akin to a grid search. Unfortunately, real problems are usually multidimensional, so we cannot draw the functions on paper. Furthermore, a direct search on a multidimensional grid is computationally prohibitive (the so-called "curse of dimensionality" issue). Consequently, we have yet to learn of the global shape of the function. We rely on the limited local information to search for the minimum. We want to call this agnostic optimization. See Figure 6 and its solution in Figure 7.


Fig. 6. Agnostic optimization. Given $x_{0}$ and $f\left(x_{0}\right)$ as a starting point, the landscape looks like this for a computer program: there is no clear path to move from $x_{0}$ to a point with a better objective value.


Fig. 7. However, the whole picture is unpredictable. Is it a minimum?

We use a set of notations to refer to the optima of a function. Without loss of generality, we only discuss minimization.
Definition 11. The global minimizer $x^{\star}$ of a function $f$ satisfies

$$
f\left(x^{\star}\right) \leq f(x), \quad \forall x \text { in the domain of } f .
$$

Definition 12. A local minimizer $\hat{x}$ of a function $f$ satisfies

$$
f(\hat{x}) \leq f(x), \quad \forall x \in \mathcal{N}_{\hat{x}}
$$

where $\mathcal{N}_{\hat{x}}$ defines a very small neighborhood around $\hat{x}$.
We can recognize that a solution is a local minimum by the following necessary conditions:

- 1st order optimality condition: $\nabla f(\hat{x})=0$.
- 2nd order optimality condition: $\nabla f(\hat{x})=0$ and $\nabla^{2} f(\hat{x}) \succeq$ 0.

Intuitively, the above states that $i$ ) the function is flat at the point of the minimum, and $i i$ ) the function looks like a "bowl" at this point when both conditions are satisfied. The last point relates to the notion of convexity: this will be defined later in the class.

Note that these are only necessary conditions, with $f(x)=$ $x^{3}$ as a simple counterexample at point $x=0$, which satisfies the two conditions but is not a local minimum.

## Lipschitz Conditions

Figures 6 and 7 show different points where the gradient is zero. The gradient is not unique at some points, while the function is discontinuous at other points. In such a general case, finding the global minima seems complicated unless we start making some assumptions about the objective $f$. Many of the objectives $f$ we want to optimize in practice often satisfy a form of Lipschitz continuity.
Definition 13. A function $f$ is called Lipschitz continuous, when

$$
|f(x)-f(y)| \leq M \cdot\|x-y\|_{2}, \quad \forall x, y
$$

for some constant $M>0$.
This means a function should not be too steep, where the constant $M$ controls the steepness. A Lipschitz continuous function may not have abrupt changes.


Fig. 8. Illustration of a Lipschitz continuous function, where $M$ controls the cone's width in white. In a way, Lipschitz continuity states that a function cannot abruptly change so that it will not "appear" inside the white cone in the picture above.

A similar but quite different assumption is that of Lipschitz gradient continuity, where we apply the Lipschitz condition to the gradients of the function.
Def inition 14. A function $f$ has Lipschitz continuous gradients, when

$$
\|\nabla f(x)-\nabla f(y)\|_{2} \leq L\|x-y\|_{2}, \quad \forall x, y
$$

where $L>0$ is a constant scalar. Often, such a function is also called L-smooth.

Such a condition forbids sudden changes in the gradient. Using Taylor's expansion, we can prove that

$$
f(y) \leq f(x)+\langle\nabla f(x), y-x\rangle+\frac{L}{2}\|x-y\|_{2}^{2}
$$

which means the function is upper-bounded by a quadratic function (there is also a lower quadratic bound). There are
several equivalent characterizations of Lipschitz gradient continuity to be aware of:

$$
f(y) \leq f(x)+\langle\nabla f(x), y-x\rangle+\frac{L}{2}\|x-y\|_{2}^{2}
$$

$$
\nabla^{2} f(x) \preceq L \cdot I, \text { where } I=\text { identity and }\left\|\nabla^{2} f(x)\right\|_{2} \leq L
$$

## Comparison of Lipschitz conditions:

- Lipschitz continuity implies that $f$ should not be too steep.
- Lipschitz gradient continuity implies that changes in the slope of $f$ should not happen suddenly.

Example: Linear regression. In linear regression, the objective $f(x)=\frac{1}{2}\|A x-b\|_{2}^{2}$ is not Lipschitz continuous-it gets arbitrarily steep when approaching infinity in $x$-however, it is Lipschitz gradient continuous as in:

$$
\begin{aligned}
\|\nabla f(x)-\nabla f(y)\|_{2} & =\left\|A^{\top}(A x-b)-A^{\top}(A y-b)\right\|_{2} \\
& \leq\left\|A^{\top} A\right\|_{2} \cdot\|x-y\|_{2}
\end{aligned}
$$

where $L:=\left\|A^{\top} A\right\|_{2}$, the largest singular value, serves as the parameter $L$. This also justifies the equivalent condition:

$$
\nabla^{2} f(x) \preceq L \cdot I
$$

But how can we use the Lipschitz gradient continuity in optimization?

A key product of its definition is the inequality:

$$
f(y) \leq f(\alpha)+\langle\nabla f(\alpha), y-\alpha\rangle+\frac{L}{2}\|y-\alpha\|_{2}^{2}
$$

Therefore, at a chosen point $\alpha$, we can upper bound the curve of $f$ (for any $y$ ) with a quadratic function, evaluated around $\alpha$. One can depict a one-dimensional simple example as in figure 9 .


Fig. 9. Illustration of how Lipschitz gradient continuity has algorithmic implications. We want to minimize the one-dimensional $f(y)$ (pink curve). Instead of minimizing $f$ directly-it could be a very complicated function to minimize directlywe will successively construct quadratic (upper-bound) approximations around the current putative solutions and minimize those approximations. In the figure, we are at point $f(\alpha)$; one can construct the linear local approximation of $f$ around $\alpha$ (black curve); Lipschitz gradient continuity goes further and introduces a quadratic term, "weighted" by the Lipschitz gradient continuity constant $L$ (green curve). Minimizing this quadratic approximation will provide a new point around which we can form another quadratic approximation, etc. The key observation regarding $L$ is that the larger $L$ is, the steepest the quadratic approximation around the current point is (compare green with khaki curves). The steeper these quadratic approximations are, the smaller the learning rate/step size in algorithms needs to be to guarantee provable performance; there are more details later on.

Lipschitz gradient continuity expression $f(y) \leq f(x)+$ $\langle\nabla f(x), y-x\rangle+\frac{L}{2}\|x-y\|_{2}^{2}$ can also be proved via Taylor's
expansion + other properties of Lipschitz gradient continuous functions. We know from Taylor's expansion that:

$$
f(y)=f(x)+\langle\nabla f(x), y-x\rangle+\frac{1}{2}\left\langle\nabla^{2} f(z)(y-x), y-x\right\rangle
$$

for some $z$. Knowing that for a Lipschitz gradient continuous function, we have:

$$
\nabla^{2} f(x) \preceq L \cdot I \Rightarrow\left\|\nabla^{2} f(x)\right\|_{2} \leq\|L \cdot I\|_{2} \Rightarrow\left\|\nabla^{2} f(x)\right\|_{2} \leq L .
$$

Then,

$$
\begin{aligned}
\frac{1}{2}\left\langle\nabla^{2} f(z)(y-x), y-x\right\rangle & \leq \frac{1}{2}\left\|\nabla^{2} f(z)(y-x)\right\|_{2} \cdot\|y-x\|_{2} \\
& \leq \frac{1}{2}\left\|\nabla^{2} f(z)\right\|_{2} \cdot\|y-x\|_{2}^{2} \\
& \leq \frac{L}{2}\|y-x\|_{2}^{2} .
\end{aligned}
$$

Combining this with the initial Taylor's expansion expression, we get:

$$
f(y) \leq f(x)+\langle\nabla f(x), y-x\rangle+\frac{L}{2}\|y-x\|_{2}^{2}
$$

## Gradient Descent for Lipschitz continuous gradient $f$

With Lipschitz gradient continuity, we can establish the convergence of an iterative optimization method, such as gradient descent. Gradient descent can be derived as the method of successively minimizing the quadratic approximations around the current point.

Let us elaborate a bit more before we present gradient descent as the basic algorithm for smooth optimization. Let $\min _{x \in \mathbb{R}^{p}} f(x)$ be the problem we are interested in solving. We assume that $f$ is differentiable, and we can approximate it by Taylor's expansion as:

$$
f(x+\delta)=f(x)+\langle\nabla f(x), \delta\rangle+o\left(\|\delta\|_{2}\right)
$$

Minimizing $f$ locally, a promising direction $\delta$ is such that the quantity $\langle\nabla f(x), \delta\rangle$ is as small as possible. Given that, for now, we are interested in finding a good direction (and not how far in that direction to move to), it is easy to see that a good normalized direction is:

$$
\delta=-\frac{\nabla f(x)}{\|\nabla f(x)\|_{2}}
$$

or a direction with controllable steps:

$$
\delta=-\eta \nabla f(x) .
$$

Given the above, gradient descent is defined as follows:
Definition 15. Let $f$ be a differentiable objective with gradient $\nabla f(\cdot)$. The gradient descent method optimize $f$ iteratively, as in:

$$
x_{t+1}=x_{t}-\eta_{t} \nabla f\left(x_{t}\right), \quad t=0,1, \ldots
$$

where $x_{t}$ is the current estimate, and $\eta_{t}$ is the step size or learning rate.
The idea behind gradient descent is simple: given the current point $x_{t}$, we can compute the negative gradient $-\nabla f\left(x_{t}\right)$ as the direction that $f$ has the steepest slope (locally). Following that direction, we carefully select $\eta_{t}$, the step size, to dictate how far in that direction we will move.
While gradient descent is quite simple, there are three actions needed to make it work in practice: i) how to choose step size $\eta_{t}, i i$ ) initial point $x_{0}$, and $i i i$ ) when to terminate the algorithm.
-Step size: several approaches are known, some more practical than others, including:
i) $\eta_{t}=\eta$; i.e., the step size is fixed to a value by the user and stays fixed for all the iterations;
ii) $\eta_{t}=O\left(\frac{c}{t}\right)$ or $\eta_{t}=O\left(\frac{c}{\sqrt{t}}\right)$ for a constant $c>0$; i.e., the step size keeps decreasing as we go on with the iterations. It starts aggressively (e.g., for $t=1$ it can be $c$ ), but very fast decreases;
iii) $\eta_{t}=\operatorname{argmin}_{\eta} f\left(x_{t}-\eta \nabla f\left(x_{t}\right)\right)$; i.e., find the step size that minimizes our objective along the direction of gradient descent. This approach makes sense (computationally) only for a narrow set of problems, where solving the above problem has $i$ ) a closed-form solution, and $i i$ ) it is not difficult to compute that closed-form solution. In most cases, finding the best $\eta_{t}$ is computationally prohibited to perform per iteration, and often, it requires the same effort as finding the solution to the original problem.
$i v)$ Fixed step size procedures, such as the Goldstein-Armijo rule, are out of this course's scope and often used in classical numerical analysis.
-Initial value $x_{0}$ : because we know little about the function, we usually start from points that make sense (e.g., unless the data involved in the function definition have abruptly large or small values, starting from $x_{0}=0$ makes sense for some problems) or we pick a random value. How to initialize is (almost) irrelevant for some classes of problems (e.g., convex optimization), but it is essential for a broader class of problems. E.g., in the case of neural networks, there are several different initialization techniques, such as the LeCun initialization [1], the He initialization [2] or the NTK initialization [3], mostly based on different probability distributions and proper scaling. By important, we mean that carefully selecting the starting point either leads to some theory-but in practice, several starting points lead to the same performance - or that is required to get good performance in practice.
-Termination criterion: there are various standard criteria, like "killing" the execution after $T$ iterations (irrespective of whether we converged or not), checking how much progress we make per iteration through $\left\|x_{t+1}-x_{t}\right\|_{2}$ or $f\left(x_{t+1}\right)-f\left(x_{t}\right)$, or by checking if the norm of the gradient is below a threshold, $\left\|\nabla f\left(x_{t}\right)\right\|_{2} \leq \varepsilon$.

## Performance of gradient descent under smoothness assump-

 tions.Claim 1. Assume that i) $f$ is differentiable, and ii) that $f$ has L-Lipschitz continuous gradients. Consider the gradient descent iterate: $x_{t+1}=x_{t}-\eta_{t} \nabla f\left(x_{t}\right)$. Then:

$$
f\left(x_{t+1}\right) \leq f\left(x_{t}\right)-\eta_{t}\left(1-\frac{\eta_{t} L}{2}\right) \cdot\left\|\nabla f\left(x_{t}\right)\right\|_{2}^{2} .
$$

Proof: By using the assumption of Lipschitz gradients, we have:

$$
\begin{aligned}
f\left(x_{t+1}\right) \leq & f\left(x_{t}\right)+\left\langle\nabla f\left(x_{t}\right), x_{t+1}-x_{t}\right\rangle+\frac{L}{2}\left\|x_{t+1}-x_{t}\right\|_{2}^{2} \\
= & f\left(x_{t}\right)+\left\langle\nabla f\left(x_{t}\right), x_{t}-\eta_{t} \nabla f\left(x_{t}\right)-x_{t}\right\rangle \\
& \quad+\frac{L}{2}\left\|x_{t}-\eta_{t} \nabla f\left(x_{t}\right)-x_{t}\right\|_{2}^{2} \\
= & f\left(x_{t}\right)-\eta_{t}\left\|\nabla f\left(x_{t}\right)\right\|_{2}^{2}+\frac{\eta_{t}^{2} L}{2}\left\|\nabla f\left(x_{t}\right)\right\|_{2}^{2} \\
= & f\left(x_{t}\right)-\eta_{t}\left(1-\frac{\eta_{t} L}{2}\right)\left\|\nabla f\left(x_{t}\right)\right\|_{2}^{2}
\end{aligned}
$$

The above result indicates that, $i$ ) as long as $\eta_{t}\left(1-\frac{\eta_{t} L}{2}\right)$ is positive, by performing gradient descent steps, we decrease the objective value by a non-positive quantity
$\left.-\eta_{t}\left(1-\frac{\eta_{t} L}{2}\right)\left\|\nabla f\left(x_{t}\right)\right\|_{2}^{2} ; i i\right)$ we can maximize the decrease by maximizing the quantity $\eta_{t}\left(1-\frac{\eta_{t} L}{2}\right)$.
Define $g(\eta):=\eta\left(1-\frac{\eta L}{2}\right)$. Knowing that $\eta>0$, we first require $1-\frac{\eta L}{2}>0 \Rightarrow \eta<\frac{2}{L}$. Thus, for $0<\eta<\frac{2}{L}$, we observe that the $g(\eta)$ is maximized when we require the gradient satisfies:

$$
g^{\prime}(\eta)=0 \Rightarrow 1-\eta L=0 \Rightarrow \eta=\frac{1}{L} .
$$

We will use $\eta_{t}=\eta=\frac{1}{L}$ for the rest of our theory. Observe that this step size requires the knowledge of $L$; for some objectives, this is easy to find, e.g., linear regression and logistic regression, but for others, it is not
Claim 2. Gradient descent $x_{t+1}=x_{t}-\eta_{t} \nabla f\left(x_{t}\right)$, with $\eta_{t}=$ $\eta=\frac{1}{L}$, satisfies:

$$
f\left(x_{t+1}\right) \leq f\left(x_{t}\right)-\frac{1}{2 L}\left\|\nabla f\left(x_{t}\right)\right\|_{2}^{2}
$$

Proof: This is true by substituting $\eta_{t}=\frac{1}{L}$ in the result of claim 1.

The above characterizes the drop in function values at the $t$-th iteration. The idea of convergence is based on the concept of relaxation.
Definition 16. A sequence of real numbers $\left\{\alpha_{t}\right\}_{t=0}^{\infty}$ is called a relaxation sequence if $\alpha_{t+1} \leq \alpha_{t}, t \geq 0$.
Combining all the iterations, for $T$ iterations, we have:

$$
\begin{aligned}
f\left(x_{T+1}\right) & \leq f\left(x_{T}\right)-\frac{1}{2 L}\left\|\nabla f\left(x_{T}\right)\right\|_{2}^{2} \\
f\left(x_{T}\right) & \leq f\left(x_{T-1}\right)-\frac{1}{2 L}\left\|\nabla f\left(x_{T-1}\right)\right\|_{2}^{2} \\
& \vdots \\
f\left(x_{1}\right) & \leq f\left(x_{0}\right)-\frac{1}{2 L}\left\|\nabla f\left(x_{0}\right)\right\|_{2}^{2}
\end{aligned}
$$

Summing all these inequalities, and under the observation that $f\left(x^{\star}\right) \leq f\left(x_{T+1}\right)$, we get the following claim.
Claim 3. Over $T$ iterations, gradient descent generates a sequence of points $x_{1}, x_{2}, \ldots$, such that:

$$
\frac{1}{2 L} \sum_{t=0}^{T}\left\|\nabla f\left(x_{t}\right)\right\|_{2}^{2} \leq f\left(x_{0}\right)-f\left(x^{\star}\right) .
$$

First, observe that the right-hand side is a constant quantity, as it does not depend on the number of iterations. Subsequently, the above result implies that, even if we continue running gradient descent for many iterations, the sum of gradient norms is always bounded by a constant. This indicates that the gradient norms that we eventually add over time have to be minor, which further implies convergence to a stationary point (also known as a critical point: a point that has a gradient zero, meaning that it could be a local minimum).
However, the above says nothing about the convergence rate. For that, we have the following claim.
Claim 4. Assume we run gradient descent for $T$ iterations, and we obtain $T$ gradients, $\nabla f\left(x_{t}\right)$, for $t \in\{0, \ldots, T\}$. Then,

$$
\min _{t \in\{0, \ldots, T\}}\left\|\nabla f\left(x_{t}\right)\right\|_{2} \leq \sqrt{\frac{2 L}{T+1}}\left(f\left(x_{0}\right)-f\left(x^{\star}\right)\right)^{\frac{1}{2}}=O\left(\frac{1}{\sqrt{T}}\right) .
$$

Proof: We know that

$$
(T+1) \cdot \min _{t}\left\|\nabla f\left(x_{t}\right)\right\|_{2}^{2} \leq \sum_{t=0}^{T}\left\|\nabla f\left(x_{t}\right)\right\|_{2}^{2} .
$$

Then,

$$
\begin{aligned}
\frac{T+1}{2 L} \cdot \min _{t}\left\|\nabla f\left(x_{t}\right)\right\|_{2}^{2} & \leq \frac{1}{2 L} \sum_{t=0}^{T}\left\|\nabla f\left(x_{t}\right)\right\|_{2}^{2} \leq f\left(x_{0}\right)-f\left(x^{\star}\right) \Rightarrow \\
\min _{t}\left\|\nabla f\left(x_{t}\right)\right\|_{2}^{2} & \leq \frac{2 L}{T+1} \cdot\left(f\left(x_{0}\right)-f\left(x^{\star}\right)\right) \\
\min _{t}\left\|\nabla f\left(x_{t}\right)\right\|_{2} & \leq \sqrt{\frac{2 L}{T+1}} \cdot\left(f\left(x_{0}\right)-f\left(x^{\star}\right)\right)^{\frac{1}{2}} \\
& =O\left(\frac{1}{\sqrt{T}}\right) .
\end{aligned}
$$

This is called sublinear convergence rate. To provide a perspective of what it means, focus on Figure 12. In general, this is a rather pessimistic result. However, remember that we made no assumptions other than the differentiability of $f$ and $f$ being a $L$-smooth function. We will see that making more assumptions helps improve the convergence radically.

Side note on convergence rates. There are two notations for convergence rate, one using an error level $\varepsilon$ based on our stopping criterion and the other using the number of iterations $T$. For now, we know that gradient descent has a convergence rate, with respect to the norm of the gradients, $O(1 / \sqrt{T})$. Pick a small $\varepsilon$, and assume we require $\min _{t}\left\|\nabla f\left(x_{t}\right)\right\|_{2} \leq \varepsilon$. This translates into:

$$
\begin{aligned}
\sqrt{\frac{2 L}{T+1}} \cdot\left(f\left(x_{0}\right)-f\left(x^{\star}\right)\right)^{\frac{1}{2}} & \leq \varepsilon \Rightarrow \\
T+1 & \geq \frac{2 L}{\varepsilon^{2}} \cdot\left(f\left(x_{0}\right)-f\left(x^{\star}\right)\right) \Rightarrow \\
T & \geq\left\lceil\frac{2 L}{\varepsilon^{2}} \cdot\left(f\left(x_{0}\right)-f\left(x^{\star}\right)\right)-1\right\rceil
\end{aligned}
$$

Usually, for our convergence rates to make sense, we pick a small value for $\varepsilon$, e.g. let $\varepsilon=10^{-3}$. Our result dictates that to get a solution with $\min _{t}\left\|\nabla f\left(x_{t}\right)\right\|_{2} \leq 10^{-3}$, we will need approximately $O\left(1 / \varepsilon^{2}\right)=O\left(10^{6}\right)$ iterations (hiding all other constants). This is the meaning of a sublinear convergence rate: to get $\varepsilon$ accuracy in some sense, we require $1 / \varepsilon^{2}$ iterations. This course will discuss how to achieve better than sublinear or even better-than-linear rates.

Example: Logistic regression. We already discussed the case of linear regression, where the objective $f(x)=$ $\frac{1}{2}\|A x-b\|^{2}$ has Lipschitz continuous gradients, with constant $L:=\left\|A^{\top} A\right\|_{2}$. Here, we consider another famous-and less straightforward-objective: that of logistic regression. We know that logistic regression is based on the following premise for binary classification:

Given a sample feature vector $\alpha_{i} \in \mathbb{R}^{p}$ and a binary class $y_{i} \in\{ \pm 1\}$, define the conditional probability of $y_{i}$ given $\alpha_{i}$ as :

$$
\mathbb{P}\left[y_{i} \mid \alpha_{i}, x^{\star}\right] \propto \frac{1}{1+\exp \left(-y_{i} \alpha_{i}^{\top} x^{\star}\right)} .
$$

The above generative assumption leads to the following objective:

$$
\min _{x \in \mathbb{R}^{p}}\left\{f(x)=\frac{1}{n} \sum_{i=1}^{n} \log \left(1+\exp \left(-y_{i} \alpha_{i}^{\top} x\right)\right)\right\} .
$$

Following the same recipe with linear regression, one can compute the gradient and Hessian as

$$
\begin{aligned}
\nabla f(x) & =\frac{1}{n} \sum_{i=1}^{n} \nabla\left[\log \left(1+\exp \left(-y_{i} \alpha_{i}^{\top} x\right)\right]\right. \\
& =\frac{1}{n} \sum_{i=1}^{n} \frac{1}{1+\exp \left(-y_{i} \alpha_{i}^{\top} x\right)} \cdot \nabla_{x}\left[\exp \left(-y_{i} \alpha_{i}^{\top} x\right)\right] \\
& =\frac{1}{n} \sum_{i=1}^{n} \frac{\exp \left(-y_{i} \alpha_{i}^{\top} x\right)}{1+\exp \left(-y_{i} \alpha_{i}^{\top} x\right)} \cdot \nabla_{x}\left[-y_{i} \alpha_{i}^{\top} x\right] \\
& =\frac{1}{n} \sum_{i=1}^{n} \frac{-y_{i}}{1+\exp \left(y_{i} \alpha_{i}^{\top} x\right)} \alpha_{i}^{\top}
\end{aligned}
$$

and

$$
\begin{aligned}
\nabla^{2} f(x) & =\frac{1}{n} \sum_{i=1}^{n} \frac{y_{i}}{\left(1+\exp \left(y_{i} \alpha_{i}^{\top} x\right)\right)^{2}} \cdot \nabla\left[1+\exp \left(y_{i} \alpha_{i}^{\top} x\right)\right] \cdot \alpha_{i}^{\top} \\
& =\frac{1}{n} \sum_{i=1}^{n} \frac{y_{i}^{2}}{\left(1+\exp \left(y_{i} \alpha_{i}^{\top} x\right)\right)^{2}} \cdot \exp \left(y_{i} \alpha_{i}^{\top} x\right) \cdot \alpha_{i} \alpha_{i}^{\top} \\
& =\frac{1}{n} \sum_{i=1}^{n} \underbrace{\frac{1}{\left(1+\exp \left(y_{i} \alpha_{i}^{\top} x\right)\right)^{2}}}_{\text {scalar }} \cdot \exp \left(y_{i} \alpha_{i}^{\top} x\right)
\end{aligned} \underbrace{\alpha_{i} \alpha_{i}^{\top}}_{\in \mathbb{R}^{p} \times p} .
$$

Observe that, for $\beta \in \mathbb{R}$,

$$
\frac{1}{(1+\exp (\beta))^{2}} \cdot \exp (\beta)=\frac{1}{1+\exp (\beta)} \cdot \frac{\exp (\beta)}{1+\exp (\beta)}=\frac{1}{1+\exp (\beta)} \cdot \frac{1}{1+\exp (-\beta)}
$$

Define $h(\beta)=\frac{1}{1+\exp (-\beta)}$, and observe that $h$ maps to $(0,1)$.
Also observe that $h(-\beta)=1-h(\beta)$. Then, one can check that $h(\beta) \cdot h(-\beta) \leq \frac{1}{4}$.

Going back to our Hessian derivations:

$$
\nabla^{2} f(x)=\frac{1}{n} \sum_{i=1}^{n} h\left(y_{i} \alpha_{i}^{\top} x\right) \cdot h\left(-y_{i} \alpha_{i}^{\top} x\right) \cdot \alpha_{i} \alpha_{i}^{\top} .
$$

Thus, taking spectral norm on both sides:

$$
\left\|\nabla^{2} f(x)\right\|_{2} \leq \frac{1}{4 n}\left\|\sum_{i=1}^{n} \alpha_{i} \alpha_{i}^{\top}\right\|_{2}=\frac{1}{4 n} \cdot\left\|A^{\top} A\right\|_{2}:=L .
$$

where $A$ accumulates all $\alpha_{i}$ 's as rows.
Example: $f(x)=x^{2}+3 \sin ^{2}(x)$. This is a less practical example, but it is an example that does not satisfy some of the nice properties that linear regression and logistic regression satisfy. The objective looks like:


Fig. 10. $f(x)=x^{2}+3 \sin ^{2}(x)$

Let us compute the first and second derivatives of this function:

$$
f^{\prime}(x)=2 x+6 \sin (x) \cdot \cos (x)
$$



Fig. 12. Borrowed from Wikipedia. Illustration of different convergence rates. Note that the $y$-axis is in logarithmic scale for all the plots, while the $x$-axis has a linear scale. The y-axis denotes a metric that dictates the optimum point; for example, $\left\|x_{k}-x^{\star}\right\|_{2}$. The x-axis represents the iteration count $k$. The first two plots represent linear convergence rates: it is called linear as a convention to match the linear curve in the logarithmic y-axis scale. While the second plot depicts a preferable behavior, the two plots are equivalent in the big-Oh notation. For an error level $\varepsilon$, linear convergence rate implies $O\left(\log \frac{1}{\varepsilon}\right)$. The third plot depicts a quadratic convergence rate. For an error level $\varepsilon$, linear convergence rate implies $O\left(\log \log \frac{1}{\varepsilon}\right)$. Finally, the fourth plot represents the sublinear convergence rate, much slower than the linear rate Some typical rates are: $O\left(1 / \varepsilon^{2}\right), O(1 / \varepsilon), O(1 / \sqrt{\varepsilon})$.
and

$$
f^{\prime \prime}(x)=2+6 \cos ^{2}(x)-6 \sin ^{2}(x)
$$

Plotting the Hessian function, we obtain:


Fig. 11. $f^{\prime \prime}(x)=2+6 \cos ^{2}(x)-6 \sin ^{2}(x)$

By inspection (and based on the periodicity of the Hessian function), we can bound:

$$
\left|f^{\prime \prime}(x)\right| \leq 8:=L
$$

$\qquad$
$\infty$

Deep learning, gradients, and autodiff. It makes sense to open a parenthesis here and highlight the importance of (efficient) gradient calculation in modern machine learning applications. We will focus on the case of deep learning and neural network training. (For a deeper discussion on neural networks, there are excellent sources online; e.g., one could focus on the Deep Learning Book [4]).

Deep learning has advanced the state-of-the-art in computer vision [5-7], natural language processing [8-10] and speech recognition [11, 12]. At the time of writing this chapter,
the transformer model [13] has revolutionized the NLP research, with machine translation [13], text classification [14], and image captioning [15]. Transformers are adopted for selfsupervised pre-training and transfer learning, with the proposal of BERT [14]. Some popular science successes that use transformers or a variant of transformers include the GPT3 [16], Megatron-LM [17] and T5 [18] language models, the DALLE-2 image synthesis model by OpenAI [19], and the AlphaFold2 protein-folding predictor by DeepMind [20].
Despite the impact of deep learning, the computational requirements for training such models are significant. Though pre-trained models are available online, there is a need to train such models from scratch. Training deep learning models is expensive (e.g., recent language models cost several million USD to train [21,22]). For practitioners, even moderate-scale tasks can be prohibitive in time and cost due to hyperparameter tuning, which may lead to multiple iterations of model retraining when the hyperparameter tuning process is taken into account, where it is typical to retrain models many times to achieve optimal performance.
At the core of this computational workload is the gradient calculation. To keep the discussion simple, consider the simplest version of a neural network: that of fully connected (FC) layers, or otherwise described as multi-layer perceptrons (MLPs). The mathematical description of an FC layer is as follows:

$$
z_{i+1}=\sigma\left(W \cdot z_{i}+b\right)
$$

where $z_{i} \in \mathbb{R}^{d_{i}}$ is the vector "representation" of the input at the $i$-th layer of a multilayer neural network, $W \in \mathbb{R}^{d_{i+1} \times d_{i}}$ is a trainable matrix that maps the input representation from $d_{i}$-dimensions to $d_{i+1}$-dimensions, and $b$ is a bias vector (for the rest of the discussion, we will assume that $b=0$ for simplicity). Finally, $\sigma: \mathbb{R}^{d_{i+1}} \rightarrow \mathbb{R}^{d_{i+1}}$ is a non-linear -often operating entrywise- activation function; some classic exam-

[^0]ples (either smooth or nonsmooth) are the ReLU, the sigmoid function, and the tanh function. ${ }^{1}$

To help visualize how a neural network would look like, consider the following:


Fig. 13. 2-layer MLP

We consider as input a vector $x=\left(x_{1}, x_{2}, \ldots, x_{n}\right)$ (you can think $x_{i}$ as pixels of a flattened image), and the task is to transform the input $x$ such that the decision/output of the neural network $y$ represents an "answer to a question"; e.g., the output often is a one-hot encoding where $y$ is a probability distribution over several classes and the task is that of image classification. This toy model has two hidden layers, each with trainable parameters $W_{i}$, represented as matrices of appropriate dimensions.
Let us go through this model and see how the input is transformed as we propagate "forward" from left to right:

- Given the input $x$, the first set of parameters generate the intermediate result $W_{1} \cdot x$.
- Given this intermediate result, the neural network inserts that to the first set of neurons, with a nonlinear activation function, to get $z_{1}=\sigma\left(W_{1} \cdot x\right)$. This is the output of the first layer.
- Given $z_{1}$, we now get into the second layer: we first compute the intermediate result by applying the second set of trainable parameters, $W_{2}$ : i.e., $W_{2} \cdot z_{1}=W_{2} \cdot \sigma\left(W_{1} \cdot x\right)$. Observe the compositional formulation that a neural network takes: the input is being propagated through a series of layers that transform the initial representation to extract useful features.
- Finally, this intermediate result also goes into the nonlinear activation functions of the second set of neurons to get $z_{2}=\sigma\left(W_{2} \cdot \sigma\left(W_{1} \cdot x\right)\right)$.
- (If we had more layers, this discussion would go on ...)

The goal in deep learning (and in machine learning in general) is to train these $W_{i}$ tensors (here, 2-way tensors are matrices), such that the output of the model (i.e., here, the neural network) maps to the correct "labels": i.e., given a dataset $\mathcal{D}$ of images $x$ and their corresponding correct labels $y^{\star}$, we want the outcome of our neural network, say $y=\operatorname{MLP}(x)$, to be as close as possible to the ground truth $y^{\star}$ for all the images in the data source (which further implies we know, say, how to classify the input images). ${ }^{2}$

Mathematically, a way to measure the distance of the "learned" output $y$ to the actual labels $y^{\star}$ is by using a loss function. Here, for simplicity, we will use the $\ell_{2}$-norm dis-
tance, and we will define the training problem as:

$$
\min _{W_{1}, W_{2}} \sum_{\left(x, y^{\star}\right) \in \mathcal{D}}\left\|y^{\star}-y\right\|_{2}^{2}=\sum_{\left(x, y^{\star}\right) \in \mathcal{D}}\left\|y^{\star}-\operatorname{MLP}(x)\right\|_{2}^{2}
$$

This is the optimization problem one needs to solve to train neural networks. I.e., suppose we abstract all the above (and with the abuse of notation where we use $x$ for the variables in this course). In that case, the above is no different than a regular optimization problem:

$$
\min _{x} f(x)
$$

Gradient calculation in neural networks is no different than applying the chain rule of derivatives. Automatic differentiation or autodiff is the field that provides efficient algorithms to compute any function's gradients, especially if written as a composition of differentiable building blocks. E.g., in our discussion above, the building blocks are matrix-vector multiplications (e.g., $W_{1} \cdot x$ ), nonlinear function applications (e.g., $\sigma\left(W_{1} \cdot x\right)$ ), etc. Autodiff is based on dynamic programming tools that wisely choose what quantities can be stored through the process and what the optimal sequence of operations is so that the gradient calculation is efficiently computed. This way, gradient calculations become an abstraction for practitioners (that, these days, they do not need to worry about), and this allows researchers to focus on the modeling part of neural networks, by defining other more informative/structured building blocks, based on the application: this has led to the creation of convolutional layers, residual layers, transformer layers, etc.

That being said, autodiff efficiently implements the backpropagation algorithm, calculating the gradient of a composition of functions. In particular, given the forward output of a neural network, the backpropagation algorithm measures the discrepancy of the output with respect to the ground truth $y^{\star}$ : i.e., $\left\|y^{\star}-y\right\|_{2}^{2}$. Based on this value, it is reasonable to infer the following rules:

- If the loss $\left\|y^{\star}-y\right\|_{2}^{2}$ is small, it means the neural network does not have to change much;
- If the loss $\left\|y^{\star}-y\right\|_{2}^{2}$ is significant, we need to update the trainable parameters in the direction to minimize this loss. And, this direction along the negative of the gradient! I.e., we update the parameters based on gradient-descent motions!

That being said, the following picture represents our neural network with the help of modules:


Fig. 14. Module-based representation of our toy neural network

[^1]It turns out that the backpropagation for this model is the efficient calculation and reuse of the intermediate steps, as shown in the following picture:

$$
\begin{aligned}
\nabla f_{i}(x) & =\frac{\partial \ell\left(y_{i}, \widehat{y}_{i}\right)}{\partial W} \\
& =\frac{\partial \ell\left(y_{i}, \widehat{y}_{i}\right)}{\partial \operatorname{softmax}(\cdot)} \cdot \frac{\partial \operatorname{softmax}(\cdot)}{\partial W} \\
& =\frac{\partial \ell\left(y_{i}, \widehat{y}_{i}\right)}{\partial \operatorname{softmax}(\cdot)} \cdot \frac{\partial \operatorname{softmax}(\cdot)}{\partial \varphi\left(\cdot, W_{2}\right)} \cdot \frac{\partial \varphi\left(\cdot, W_{2}\right)}{\partial W} \\
& =\cdots
\end{aligned}
$$

Fig. 15. Chain of derivatives

The idea of these calculations is that they define a graph of intermediate calculations that can be reused to complete the full gradient computation: efficiently calculating, storing, and reusing these intermediate steps is at the core of autodiff.


Fig. 16. Graph representation of autodiff

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[^0]:    ${ }^{1}$ In an actual neural network implementation, there might be additional layers or functions per layer that modify further the inputs at each layer (e.g., pooling layers, batch normalization layers, softmax that modify further the inputs at each layer (e.g., pooling layers, batch normalization layers, softmax
    layers). Still, it is out of the scope of this chapter to delve into these. At this stage, consider that a layers). Still, it is out of the scope of this chapter to delve into these. At this stage, consider that a
    deep learning model is a black box machine that transforms the input through a sequence of layers, deep learning model is a black box machine that transforms the input thr
    each of which operates differently and based on the application at hand.

[^1]:    ${ }^{2}$ This is the case of supervised learning where $\mathcal{D}$ is a dataset that has both inputs $x$ and the correct labels $y^{\star}$; these $y^{\star}$ will help the model update its parameters in order the output to be as close as possible to these $y^{\star}$

