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# METHODS OF OPTIMAL CONTROL

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**Notations**

$(\Omega, \mathcal{F}, \mathbb{P})$	a probability space with a $\sigma$ -field $\mathcal{F}$ and probability $\mathcal{P}$
$\mathbb{F} = \{\mathcal{F}_t\}_{t \geq 0}$	a filtration in a probability space
$\mathbb{F}^X$	filtration generated by process $X$
$\mathbb{F}^{X+}$	$\cup_{s>t} \mathcal{F}_s^X$
$\bar{\mathbb{F}}$	filtration augmented by adding all the $\mathbb{P}$ -null events
$(\Omega, \bar{\mathbb{F}}, \mathbb{P})$	a filtered probability space with right-continuous augmented filtration (usual conditions) that hosts a Brownian motion adapted to the filtration $\bar{\mathbb{F}}$
$O \subseteq \mathbb{R}^d$	open set
$\text{USC}(O)$ ( $\text{LSC}(O)$ )	lower semicontinuous (upper semicontinuous) functions on $O$
$\mathcal{C}^k(O)$	$k$ times continuously differentiable functions on $O$
$\bar{\mathcal{C}}^k(\bar{O})$	functions on $\bar{O}$ $k$ times continuously differentiable on all variables over on $O$ with derivatives continuously extendable to $\bar{O}$
$\mathcal{C}^{k,l,\dots}(O)$	functions on $O$ $k$ times continuously differentiable on first variable, $l$ times on second variable, ...
$\bar{\mathcal{C}}^{k,l,\dots}(\bar{O})$	functions on $\bar{O}$ times continuously differentiable on first variable, $l$ times on second variable, ... on $O$ with derivatives continuously extendable to $\bar{O}$
$M(n, m)$	the set of all real $n$ by $m$ matrices
$A \cdot B := \text{Tr}[AB^T]$	the inner product of two matrices $A$ and $B$ in $M(n, m)$
$\ A\  := A \cdot A = \sum_{i=1}^n \sum_{j=1}^m a_{ij}^2$	$l^2$ -norm of matrices in $M(n, m)$
$\Omega$	the sample space of a random event
$\omega$	is reserved for the members of $\Omega$
$\partial_t V, \partial_x V, \partial_{xx} V$	Partial derivatives of a function $V : [0, T] \times \mathbb{R}$ once wrt $t$ , once wrt $x$ and twice wrt $x$
$\partial_t V, \nabla V, D^2 V$	Partial derivative of a function $V : [0, T] \times \mathbb{R}^2$ once wrt $t$ , gradient of $V$ wrt $x$ and Hessian of $V$ wrt $x$

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# Chapter 1

## Preliminaries

### 1.1 Optimization Versus Control

In this chapter, we provide a brief overview of the aspect in which optimization and control are different. We start the chapter with some examples.

**Example 1.** We start by a quadratic problem. Let  $\alpha : [0, T] \rightarrow \mathbb{R}$  be given.

$$\inf \left\{ \int_0^T (x_t^2 - \alpha_t x_t) dt \right\} \quad (1.1.1)$$

where the infimum is over all functions  $x : [0, T] \rightarrow \mathbb{R}$ . As  $x$  can be any mapping, one can solve the following problem for each  $x_t$  separately to obtain  $x_t^* = \alpha_t/2$ .

$$\inf_{x \in \mathbb{R}^d} \left\{ x^2 - \alpha_t x \right\} \quad (1.1.2)$$

The above problem is a dynamic optimization problem.

**Example 2.** We make the above example more complicated by specifying a dynamics for  $x_t$ :

$$\inf \left\{ \int_0^T (x_t^2 - \alpha_t x_t) dt \right\} \quad (1.1.3)$$

where the infimum is over all functions  $x : [0, T] \rightarrow \mathbb{R}$  such that for some (Borel measurable) function  $u : [0, T] \rightarrow \mathbb{R}$

$$dx_t = (-\beta x_t + u_t) dt, \quad x_0 = x \quad (1.1.4)$$

Note that the existence of the term  $u_t dt$  is necessary for the meaning of infimum. This problem is not a simple dynamic optimization problem because one can only choose  $x$  such that the dynamic equation (1.1.4) holds for some function  $u$ . **Therefore**, we can only choose  $u$  and indirectly modify  $x_t$  to minimize (1.1.3).

**Exercise 1.** Solve (1.1.4) for  $x_t$  in terms of  $u$ . Hint:  $d(e^{\beta t} x_t) = e^{\beta t} u_t dt$ .

**Exercise 2.** Show that if there exists a function  $u$  such that  $\alpha_t = 2 \int_0^t e^{\beta(s-t)} u_s ds$ , then  $u$  minimizes (1.1.3).

**Exercise 3.** If we modify an optimal control  $u$  in a countable number of points described in the exercise above, does it remain an optimal control? Does the initial value  $x_0 = x$  play a role in the problem?

**Exercise 4.** Assume that there exists no function  $u$  such that  $\alpha_t = 2 \int_0^t e^{-\beta(s-t)} u_s ds$ . What is the minimum value of (1.1.3)?

The problem (1.1.3) is a simple control problem. However, after doing the above exercises, you note that it can simply be reduced to an optimization problem. Such a solution for control problems are called myopic solutions and do not necessarily exist for more interesting control problems. Here is an example which does not allow for a myopic solution.

**Example 3.** Consider the control problem:

$$\inf \int_0^T (x_t^2 - \alpha_t x_t + u_t^2) dt \quad (1.1.5)$$

where the infimum is over all (Riemann integrable) functions  $u : [0, T] \rightarrow \mathbb{R}$  and  $x : [0, T] \rightarrow \mathbb{R}$  and  $u : [0, T] \rightarrow \mathbb{R}$  satisfy (1.1.4). Unlike the previous example, the choice of  $u$  induces a new cost, the term  $\int_0^T u_t^2 dt$ . Due to this new cost, we cannot freely choose  $u_t$  to make  $x_t = \alpha_t/2$  optimal. It is possible that  $x_t = \alpha_t/2$  is not even optimal.

In the above example, there is a tradeoff between the choice of  $u_t$  to minimize the dependence of cost on  $x_t$  and to minimize the dependence of cost on  $u_t$  itself. In such cases, the tradeoff prevents us from finding a myopic solution. Here is another example.

**Example 4.** Consider the control problem:

$$\inf \int_0^T (x_t^2 + u_t^2) dt \quad (1.1.6)$$

where the infimum is over all (Riemann integrable) functions  $u : [0, T] \rightarrow \mathbb{R}$  and  $x : [0, T] \rightarrow \mathbb{R}$  and  $u : [0, T] \rightarrow \mathbb{R}$  satisfy  $dx_t = (-\beta x_t + u_t)dt$ ,  $x_0 = x$ . On one hand, we like  $x_t$  to be zero to minimize the cost, but pushing  $x_t$  to zero requires application of  $u_t$ , which introduces another cost term.

A general control problem is described as

$$\inf_{u \in \mathcal{U}} \int_0^T C(t, x_t, u_t) dt + g(x_T) \quad (1.1.7)$$

where  $dx_t = f(x_t, u_t)dt$ . The function  $C : \mathbb{R}_+ \times \mathbb{R}^d \times \mathbb{R}^n \rightarrow \mathbb{R}$  is called the running cost and  $g : \mathbb{R}^d \rightarrow \mathbb{R}$  is called the terminal cost.

Set  $\mathcal{U}$  is a set of functions  $u : [0, T] \rightarrow \mathbb{R}^n$  called control variable, which is determined by the application or by the wellposedness of the problem. It is crucial to choose a set  $\mathcal{U}$  of control variables to fit the proper application. More over, it is also important to choose a set that makes the control problem wellposed. For instance, for the control problem

$$\inf_{u \in \mathcal{U}} \int_0^T (x_t - u_t^2) dt, \quad dx_t = (x_t - u_t) dt \quad (1.1.8)$$

if we choose  $\mathcal{U}$  to be the set of all functions  $u : [0, T] \rightarrow \mathbb{R}$ , by simply choosing  $u_t$  very large, the value of the infimum is  $-\infty$ . However, if we restrict  $\mathcal{U}$  to the set of functions  $u : [0, T] \rightarrow [-1, \infty)$  (some lower bound on the value), then

$$\inf_{u \in \mathcal{U}} \int_0^T (x_t - u_t^2) dt > -\infty, \quad dx_t = (x_t - u_t) dt. \quad (1.1.9)$$

A suitable set of controls chosen for a specific problem is referred to as the set of all admissible control. We denote this set by  $\mathcal{U}$ .

An infinite horizon control problem is accommodated by setting  $T = \infty$ . For example,

$$\inf_{u \in \mathcal{U}} \int_0^\infty e^{-t} (x_t^2 + u_t^2) dt \quad (1.1.10)$$

The following exercise is an example of infinite horizon control problem.

**Exercise 5.** Write the following problem as a generic control problems by associating the horizon  $T$ , the running cost  $C(t, x, u)$  and terminal cost  $g(x)$  in (1.1.7):

(Shortest time to exit a bounded domain) Given a bounded domain  $D \subset \mathbb{R}^d$ , find

$$\inf_u \{t \geq 0 : x_t \notin D\} \quad (1.1.11)$$

where  $dx_t = u_t dt$  with control  $|u_t| \leq 1$  and  $u_t \in \mathbb{R}^d$  and initial position  $x_0 = x \in D$ .

**Exercise 6.** In your area of study, find an optimal control problem. Then, write the cost functions and the control variable and determine a set of admissible controls.

## 1.2 Solution methods for deterministic optimal control

In this section, we propose methods for the optimal control problems. There are two groups of numerical methods, first group are based on the dynamic programming principle (DPP). DPP provides a backward recursive method to solve an optimal control problem. The second group completely avoids DPP and formulates problem as an optimization.

### 1.2.1 Dynamic programming principle (DPP)

Consider the optimal control problem (1.1.7):

$$\inf_{u \in \mathcal{U}} \int_0^T C(t, x_t, u_t) dt + g(x_T), \quad dx_t = f(x_t, u_t) dt \quad (1.2.1)$$

A key concept in DPP is the value function:

**Definition 1.2.1.** Let  $x_t = x$ . Then, the value function of (1.1.7) is defined by

$$V(t, x) := \inf_{u \in \mathcal{U}_t} \int_t^T C(s, x_s, u_s) ds + g(x_T), \quad dx_s = f(x_s, u_s) ds \quad (1.2.2)$$

where  $\mathcal{U}_t$  is the set of admissible controls restricted to time interval  $[t, T]$ .

The value function encodes the outcome of optimal control at time  $t$  at point  $x$ . The function  $V(t, \cdot) : \mathbb{R}^d \rightarrow \mathbb{R}$  returns the optimal value of the control problem for any point  $x$  at time  $t$ .

Given the value function, the DPP is explained in the following result.

**Theorem 1.2.1.** [Dynamic Programming Principle] Let  $0 \leq t < s \leq T$  and  $x_t = x$ . Then, the DPP for (1.1.7) is given by

$$V(t, x) = \inf_{u \in \mathcal{U}_{t,s}} \int_t^s C(r, x_r, u_r) dr + V(s, x_s), \quad dx_r = f(x_r, u_r) dr \quad (1.2.3)$$

where  $\mathcal{U}_{t,s}$  is the set of admissible controls restricted to time interval  $[t, s]$ .



DPP can be interpreted as the tradeoff between two intervals  $[t, s]$  and  $[s, T]$ . More precisely, to obtain the optimal value of the control problem at time  $t$ , DPP finds a balance between the running cost on  $[t, s]$ ,  $\int_t^s C(r, x_r, u_r)dr$ , and the continuation cost on  $[s, T]$ ,  $V(t, x_s)$ . Note that choosing a control  $u$  on  $[t, s]$ , determines the position of  $x_s$ . This interpretation reveals the proof of DPP.

*Proof of DPP, Theorem 1.2.3.* To show that DPP holds, we need to decompose a control  $u$  on  $[t, T]$  into two pieces,  $u_1$  on  $[t, s]$  and  $u_2$  on  $[s, T]$ . More formally,

$$u(r) = \begin{cases} u_1(r) & r \in [t, s] \\ u_2(r) & r \in [s, T] \end{cases}$$

We denote this decomposition by  $u = u_1 \oplus u_2$ . Note that  $u_1 \in \mathcal{U}_{t,s}$  and  $u_2 \in \mathcal{U}_s$ . Therefore, by (1.2.2),

$$\begin{aligned} V(t, x) &= \inf_{u \in \mathcal{U}_t} \int_t^T C(r, x_r, u_r)dr + g(x_T) \\ &= \inf_{u_1 \in \mathcal{U}_{t,s}} \inf_{u_2 \in \mathcal{U}_s} \int_t^s C(r, x_r, u_1(r))dr + \int_s^T C(r, x_r, u_2(r))dr + g(x_T) \\ &= \inf_{u_1 \in \mathcal{U}_{t,s}} \int_t^s C(r, x_r, u_1(r))dr + \inf_{u_2 \in \mathcal{U}_s} \int_s^T C(r, x_r, u_2(r))dr + g(x_T) \end{aligned}$$

Note that  $\inf_{u_2 \in \mathcal{U}_s} \int_s^T C(r, x_r, u_2(r))dr + g(x_T) = V(s, x_s)$ . Therefore,

$$V(t, x) = \inf_{u_1 \in \mathcal{U}_{t,s}} \int_t^s C(r, x_r, u_1(r))dr + V(s, x_s)$$

□

### Numerical solution based on DPP

Let  $\Delta t = \frac{T}{N}$  and set  $t_i = i\Delta t$  with  $i = 0, \dots, N$ . The DPP (1.2.3) for  $t = t_i$  and  $s = t_{i+1}$  is written by

$$V(t_i, x) = \inf_{u \in \mathcal{U}_{t_i, t_{i+1}}} \int_{t_i}^{t_{i+1}} C(r, x_r, u_r)dr + V(t_{i+1}, x_{t_{i+1}}) \quad (1.2.4)$$

We can approximate  $\int_{t_i}^{t_{i+1}} C(r, x_r, u_r)dr \approx C(t_i, x_{t_i}, u_{t_i})\Delta t$  to write

$$V(t_i, x) \approx \inf_u C(t_i, x_{t_i}, u)\Delta t + V(t_{i+1}, x_{t_{i+1}}) \quad (1.2.5)$$

This suggest to approximate the value function and optimal control for the control problem (1.1.7) in Algorithm 1. Throughout this notes,  $\hat{f}$  represents an approximation of function  $f$ . Therefore, if  $V(t_{i+1}, \cdot)$  is approximated by  $\hat{V}(t_{i+1}, \cdot)$ , then

$$\hat{V}(t_i, x) := \inf_u C(t_i, x_{t_i}, u)\Delta t + \hat{V}(t_{i+1}, x_{t_{i+1}}), \quad x_{t_{i+1}} = x + f(t_i, x, u)\Delta t \quad (1.2.6)$$

suggests the approximation for  $V(t_i, \cdot)$ . Note that unlike DPP, the infimum in the approximate DPP (1.2.6) is over a  $u \in \mathbb{R}^m$  and not the functions  $u : [t_i, t_{i+1}] \rightarrow \mathbb{R}^m$ . Therefore,  $\hat{u}^*$  given by

$$\hat{u}^* \in \operatorname{argmin}_u C(t_i, x_{t_i}, u)\Delta t + \hat{V}(t_{i+1}, x + f(t_i, x, u)\Delta t)$$

provides a constant control over the interval  $[t_i, t_{i+1}]$ .

We provide an algorithm derived from discrete DPP in Algorithm 1.

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**Algorithm 1:** Numerical DPP
 

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**Parameter:**  $T, N, f(t, x, u), C(t, x, u)$ , and  $g(x)$ ;

$\Delta t = \frac{T}{N}$

**Data:**  $\hat{V}(t_N, x) = g(x)$ ;

$x_i^j$  for  $j = 1, \dots, J$  and  $i = 0, \dots, N - 1$

1 **for**  $i \leftarrow N - 1$  **to** 0 **do**

2      $\hat{x}_{i+1}^j = x_i^j + f(t_i, x_i^j, u)\Delta t$ ;  
 3      $\tilde{V}(t_i, x_i^j) \leftarrow \inf_u C(t_i, x_i^j, u)\Delta t + \hat{V}(t_{i+1}, \hat{x}_{i+1}^j)$ ;  
 4      $\hat{V}(t_i, x)$  obtained from interpolation on  $\tilde{V}(t_i, x_i^j)$  for  $j = 1, \dots, J$ ;  
 5      $\hat{u}^*(x_i^j) \in \underset{u}{\operatorname{argmin}} C(t_i, x_i^j, u)\Delta t + \hat{V}(t_{i+1}, \hat{x}_{i+1}^j)$ ;

6 **return**  $\hat{V}$  and  $\hat{u}^*$ .

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**Exercise 7.** Why interpolation is required in Algorithm 1? Can we perform the algorithm by only knowing  $\hat{V}(t_{i+1}, x_{i+1}^j)$  for all  $j = 1, \dots, J$ ? Note the difference between  $\hat{V}(t_{i+1}, x_{i+1}^j)$  and  $\hat{V}(t_{i+1}, \hat{x}_{i+1}^j)$  with  $\hat{x}_{i+1}^j = x_i^j + f(t_i, x_i^j, u)\Delta t$ .

Next example provides a special case where we can evaluate the approximate value function from Algorithm 1 without using interpolation.

**Example 5.** Consider

$$\inf_{u \in \mathcal{U}} \int_0^T (x_t^2 + u_t^2) dt + \frac{1}{2}x_T^2 - x_T, \quad dx_t = (x_t - u_t)dt \quad (1.2.7)$$

We write the value function by (1.2.2).

$$V(t, x) := \inf_{u \in \mathcal{U}_t} \int_t^T (x_s^2 + u_s^2) ds + \frac{1}{2}x_T^2 - x_T, \quad dx_s = (x_s - u_s)ds. \quad (1.2.8)$$

We cannot find value functions using a myopic argument. We will later find a way to find the value function  $V$ . However, In the next exercise, we apply Algorithm 1 to find the approximate value function.

**Exercise 8.** In example above, write the approximate DPP from time  $t_i$  to  $t_{i+1}$ . Then, assume that  $\hat{V}(t_{i+1}, x) = a_{i+1}x^2 + b_{i+1}x + c_{i+1}$  for some known values  $a_{i+1}$ ,  $b_{i+1}$ , and  $c_{i+1}$ . Use optimization of a quadratic function to find  $\hat{V}(t_i, x)$ . Note that you need to use  $\hat{x}_{i+1} = x + (x^2 + u^2)\Delta t$ . Does  $\hat{V}(t_i, x)$  is of the form  $a_i x^2 + b_i x + c_i$ ? What is the relation between  $(a_i, b_i, c_i)$  and  $(a_{i+1}, b_{i+1}, c_{i+1})$ ?

Example (8) is a linear-quadratic optimal control problem, LQC. In LQC, the running cost is quadratic in  $x$  and  $u$ , the terminal cost is quadratic in  $x$ , and the ODE for state variable is linear in  $x$  and  $u$ .

**Remark 1.2.1.** Consider the discretized version of the control problem (1.1.7):

$$\inf \left\{ \sum_{i=0}^{N-1} C(t_i, \hat{x}_{t_i}, u_{t_i})\Delta t + g(\hat{x}_T) \mid u_{t_i} : i = 0, \dots, N - 1 \right\} \quad (1.2.9)$$

with  $\hat{x}_{t_{i+1}} = \hat{x}_{t_i} + f(t_i, \hat{x}_{t_i}, u_{t_i})\Delta t$ . The DPP for this discrete-time optimal control problem is the same as the discretized DPP (1.2.6).

### Hamilton-Jacobi equation

One of the side-products of DPP is a partial differential equation, Hamilton-Jacobi equation or HJ for short, that can be solved via different methods to provide us with the solution of the control problem. To derive the HJB, we first assume that the value function is sufficiently continuously differentiable, first order derivatives exist and are continuous. Then, we use the first Taylor polynomial with the remainder term for the value function  $V(s, x_s)$  about  $(t, x)$  as follows:

$$V(s, x_s) = V(t, x) + V_t(t, x)(s - t) + V_x(t, x)(x_s - x) + R_2$$

Note that  $x_s = x + \int_t^s f(x_r, u_r)dr$ . Let's insert the above Taylor polynomial into the DPP:

$$\begin{aligned} V(t, x) &= \inf_{u \in \mathcal{U}_{t,s}} \int_t^s C(r, x_r, u_r)dr + V(s, x_s) \\ &= \inf_{u \in \mathcal{U}_{t,s}} \int_t^s C(r, x_r, u_r)dr + V(t, x) + V_t(t, x)(s - t) + V_x(t, x) \int_t^s f(x_r, u_r)dr + R_2 \\ &= V(t, x) + V_t(t, x)(s - t) + \inf_{u \in \mathcal{U}_{t,s}} \int_t^s C(r, x_r, u_r)dr + V_x(t, x) \int_t^s f(x_r, u_r)dr + R_2 \end{aligned}$$

$V(t, x)$  on both sides can be canceled. Note that  $R_2 = o(s - t)$ , which means  $\lim_{s \rightarrow t} \frac{R_2}{s - t} = 0$ . Dividing both sides by  $s - t$  and sending  $s \rightarrow t$ , we obtain

$$0 = V_t(t, x) + \inf_u C(t, x, u) + V_x(t, x)f(t, x, u) \quad (1.2.10)$$

In the above, we used the fundamental theorem of calculus to write  $\lim_{s \rightarrow t} \frac{\int_t^s f(x_r, u_r)dr}{s - t} = f(t, x, u)$  and  $\lim_{s \rightarrow t} \frac{\int_t^s C(r, x_r, u_r)dr}{s - t} = C(t, x, u)$ . A PDE requires proper boundary condition. For HJ, the boundary condition is the terminal cost at  $T$ :

$$V(T, x) = g(x) \quad (1.2.11)$$

Let's explore the HJB inside an example.

**Example 6.** Consider the control problem in Section 1.1:

$$\inf_u \left\{ \int_0^T (x_t^2 + u_t^2) dt \right\} \quad (1.2.12)$$

where  $dx_t = (-\beta x_t + u_t)dt$ . Here,  $C(t, x, u) = x^2 + u^2$  and  $f(t, x, u) = -\beta x + u$ . The HJB is given by:

$$0 = V_t + x^2 - \beta x V_x + \inf_u u^2 + V_x u$$

Note that the infimum above is attained at  $u^* = -\frac{1}{2}V_x$  and, therefore, we can write the HJB without using infimum:

$$0 = V_t + x^2 - \beta x V_x - \frac{1}{4}V_x^2$$

Since there is no terminal cost, the boundary condition for the HJ is  $V(T, x) = 0$ .

An important question is that how can the HJ equation help us find the optimal control. For the above

example, if we solve the HJ and find the value function, then a candidate for the optimal control is  $u_t^* = -\frac{1}{2}V_x(t, x_t)$  where  $dx_t = (-\beta x_t + u_t^*)dt$ . In general, we cannot guarantee a simple way to solve the HJ. However, for certain problems there is a simple way and for other problems numerical solutions to PDEs can be used. Next, exercise shows how to solve the HJ derived for the previous example.

**Exercise 9.** We guess that the solution for the HJ

$$\begin{cases} 0 = V_t + x^2 - \beta x V_x - \frac{1}{4} V_x^2 \\ V(T, x) = 0 \end{cases}$$

takes the form  $V(t, x) = a(t)x^2 + b(t)x + c(t)$ . Insert this guess into the HJB and find ODEs for  $a(t)$ ,  $b(t)$ , and  $c(t)$ . Use terminal conditions to prescribe terminal conditions for  $a$ ,  $b$ , and  $c$ . Can you solve the system of ODEs?

### 1.2.2 Pontryagin principle

Consider the simple control problem in Section 1.1:

$$\inf_u \left\{ \int_0^T (x_t^2 - \alpha_t x_t) dt \right\} \quad (1.2.13)$$

where  $dx_t = (-\beta x_t + u_t)dt$ ,  $x_0 = x$ . One can consider the dynamics of  $x$  as a constraint. Therefore, we can formally write Lagrangian by

$$\sup_{\lambda} \inf_{u, x} \left\{ \int_0^T (x_t^2 - \alpha_t x_t) dt - \int_0^T \lambda_t (dx_t + (\beta x_t - u_t)dt) \right\} \quad (1.2.14)$$

Here, there are two important remarks. First, since the constraint is given by a differential equation for each  $t$ , the dual variable  $\lambda$  is a function of  $t$ . Second, the problem is now unconstrained, that is, the state variable  $x$  and the control variable  $u$  are now both variables in an optimization problem.

The main trick to solve this optimization problem is integration by part formula

$$\int_0^T \lambda_t dx_t = \lambda_T x_T - \lambda_0 x_0 - \int_0^T x_t d\lambda_t \quad (1.2.15)$$

to write (1.2.14) as

$$\sup_{\lambda} \inf_{u, x} \left\{ \int_0^T (x_t d\lambda_t - \lambda_t (\beta x_t - u_t) + x_t^2 - \alpha_t x_t) dt - \lambda_T x_T + \lambda_0 x_0 \right\} \quad (1.2.16)$$

Note that optimization on  $x_T$  is independent of  $x_t$  for  $t < T$ . The KKT conditions for the strong duality in the saddle point problem (1.2.16) were discovered by Lev Pontryagin in 1952. Define the Hamiltonian by

$$H(t, x, \lambda, u) := -\lambda(\beta x - u) + x^2 - \alpha_t x \quad (1.2.17)$$

Thus, (1.2.16) is written as a saddle point problem with free variables  $x_t$ ,  $u_t$ , and  $\lambda_t$ :

$$\sup_{\lambda} \inf_{u, x} \left\{ \int_0^T (x_t d\lambda_t + H(t, x_t, \lambda_t, u_t)) dt - \lambda_T x_T + \lambda_0 x_0 \right\} \quad (1.2.18)$$

$$\begin{cases} d\lambda_t^* + \partial_x H(t, x_t^*, \lambda_t^*, u_t^*) dt = d\lambda_t^* + (-\lambda_t^* \beta + 2x_t^* - \alpha_t) dt = 0 \text{ (minimize integrand wrt } x) \\ \lambda_T^* = 0 \text{ (minimize terminal wrt } x_T) \\ H(t, x_t^*, \lambda_t^*, u_t^*) \leq H(t, x_t^*, \lambda_t^*, u) \text{ for all } u \text{ (minimize integrand wrt } u) \\ dx_t^* = (-\beta x_t^* + u_t^*) dt \text{ (constraint)} \end{cases} \quad (1.2.19)$$

In the above, the first equation is obtained from taking derivative with respect to  $x_t$ , second equality corresponds to derivative with respect to  $x_T$  in  $\lambda_T x_T$ , the third line guarantees the optimality of  $u^*$ , and the last equation is the constraint of the problem which is the dynamic of the state variable in the control.

**Exercise 10.** Show that  $\lambda_t^* = 0$ ,  $x_t^* = \alpha_t/2$ , and  $u^*$  with  $\alpha_t = 2 \int_0^t e^{-\beta(s-t)} u_s^* ds$  satisfy Pontryagin principle for the above problem.

Following the steps of in the above example, the Pontryagin principle for a generic deterministic control problem, 1.1.7 is described as the following saddle point problem.

$$\sup_{\lambda} \inf_{u, x} \left\{ \int_0^T C(t, x_t, u_t) dt + g(x_T) - \int_0^T \lambda_t (dx_t + (\beta x_t - u_t) dt) \right\} \quad (1.2.20)$$

By applying (1.2.15), we obtain

$$\sup_{\lambda} \inf_{u, x} \left\{ \int_0^T \left( x_t d\lambda_t + H(t, x_t, \lambda_t, u_t) \right) dt + g(x_T) - \lambda_T x_T + \lambda_0 x_0 \right\} \quad (1.2.21)$$

where the Hamiltonian is given by

$$H(t, x, \lambda, u) := \lambda f(t, x, u) + C(t, x, u) \quad (1.2.22)$$

**Theorem 1.2.2.** [Pontryagin principle] Assume that there exists  $(x^*, u^*, \lambda^*) : [0, T] \rightarrow \mathbb{R}^d \times \mathbb{R}^n \times \mathbb{R}^d$  such that

$$\begin{cases} d\lambda_t^* + \partial_x H(t, x_t^*, \lambda_t^*, u_t^*) dt = 0 \text{ (minimize integrand wrt } x) \\ \lambda_T^* = \nabla g(x_T^*) \text{ (minimize terminal wrt } x_T) \\ H(t, x_t^*, \lambda_t^*, u_t^*) \leq H(t, x_t^*, \lambda_t^*, u) \text{ for all } u \text{ (minimize integrand wrt } u) \\ dx_t^* = f(t, x_t^*, u_t^*) dt \text{ (constraint)} \end{cases} \quad (1.2.23)$$

Then,  $u^*$  is an optimal control for (1.1.7).

The function  $\lambda_t^*$ , described by 1.2.2 is called the adjoint process. It is important to distinguish between the two ODEs in (1.2.23), namely,

$$dx_t^* = f(t, x_t^*, u_t^*) dt \text{ and } d\lambda_t^* + \partial_x H(t, x_t^*, \lambda_t^*, u_t^*) dt = 0 \quad (1.2.24)$$

The first one has an initial condition  $x_0^* = x_0$ , the initial position of the state process. However, the adjoint equation comes with a terminal condition,  $\lambda_T^* = \nabla g(x_T^*)$ . This makes solving Theorem 1.2.2 challenging. We shall see later how Pontryagin principle is applied to some specific examples such as linear quadratic linear control problem, where solving (1.2.23) is simpler to solve.

**Remark 1.2.2.** If  $u_t^*$  is an interior minimizer of  $H(t, x_t^*, \lambda_t^*, u)$ , then we can write  $\partial_u H(t, x_t^*, \lambda_t^*, u_t^*) = 0$ . However, if there are constraint on  $u$ , e.g.,  $u \geq 0$ , we shall stick to the inequality above.

We start with an example of a LQC problem.

**Example 7.** Consider the control problem with  $C(x, u) = x^2 + u^2$ ,  $g(x) = x^2 + x$ , and  $f(x, u) = x + u$ :

$$\inf_u \int_0^T (x_t^2 + u_t^2) dt + x_T^2 + x_T, \text{ subject to } dx_t = (x_t + u_t) dt \quad (1.2.25)$$

LQC problems can be solved via Pontryagin principle. More precisely, for the above example, the Hamiltonian is given by

$$H(x, \lambda, u) = x^2 + u^2 + \lambda(x + u)$$

and (1.2.23) is

$$\begin{cases} d\lambda_t^* + (2x_t^* + \lambda_t^*) dt = 0 \\ \lambda_T^* = 2x_T^* + 1 \\ H(t, x_t^*, \lambda_t^*, u_t^*) \leq H(t, x_t^*, \lambda_t^*, u) \\ dx_t^* = (x_t^* + u_t^*) dt \end{cases} \quad (1.2.26)$$

Minimizing  $H(x, \lambda, u)$  in  $u$  suggests that  $H_u(x, \lambda_t^*, u_t^*) = 2u_t^* + \lambda_t^* = 0$ , equivalently,  $u_t^* = -\frac{1}{2}\lambda_t^*$ .

**Exercise 11.** (1) Use  $u_t^* = -\frac{1}{2}\lambda_t^*$  to write the system of ODEs (1.2.26) for  $\lambda^*$  and  $x^*$  by

$$\begin{cases} d\lambda_t^* = (-\lambda_t^* - 2x_t^*) dt \\ dx_t^* = (-\frac{1}{2}\lambda_t^* + x_t^*) dt \end{cases} \quad (1.2.27)$$

(2) Find the general solution the system of ODEs in terms of  $x_0^*$  and  $\lambda_0^*$ . Hint: Use denationalization of the matrix

$$\begin{bmatrix} -1 & -2 \\ -\frac{1}{2} & 1 \end{bmatrix} = \begin{bmatrix} 2(\sqrt{2}+1) & 1 \\ 1 & 2(\sqrt{2}-1) \end{bmatrix} \begin{bmatrix} -\sqrt{2} & 0 \\ 0 & \sqrt{2} \end{bmatrix} \begin{bmatrix} \frac{2}{3}(\sqrt{2}-1) & -1 \\ -1 & \frac{2}{3}(\sqrt{2}+1) \end{bmatrix} \quad (1.2.28)$$

and write the ODEs as

$$\begin{bmatrix} d\lambda^* \\ dx^* \end{bmatrix} = \begin{bmatrix} 1 & 1 \\ 1 & -\frac{1}{2} \end{bmatrix} \begin{bmatrix} \lambda^* \\ x^* \end{bmatrix} = \begin{bmatrix} 2(\sqrt{2}+1) & 1 \\ 1 & 2(\sqrt{2}-1) \end{bmatrix} \begin{bmatrix} -\sqrt{2} & 0 \\ 0 & \sqrt{2} \end{bmatrix} \begin{bmatrix} \frac{2}{3}(\sqrt{2}-1) & -1 \\ -1 & \frac{2}{3}(\sqrt{2}+1) \end{bmatrix} \begin{bmatrix} \lambda^* \\ x^* \end{bmatrix} \quad (1.2.29)$$

(3) Consider  $x_0^*$  given. Use  $\lambda_T^* = 2x_T^* + 1$  to find  $\lambda_0^*$ , hence a special solution for the system of ODEs from Pontryagin principle as a function  $x_0^*$  and  $x_T^*$ .

### 1.3 Linear quadratic control problem

A class of control problem with closed-form solution consists of linear-quadratic optimal control problems, LQC henceforth. In LQC, the running cost is quadratic in  $x$  and  $u$ , the terminal cost is quadratic in  $x$ , and the ODE for state variable is linear in  $x$  and  $u$ :

$$\begin{aligned} C(t, x, u) &= \begin{bmatrix} x^\top & u^\top \end{bmatrix} Q_t \begin{bmatrix} x \\ u \end{bmatrix} + a_t \cdot x + b_t \cdot u, \text{ and } Q_t = \begin{bmatrix} A_t & C_t \\ C_t^\top & B_t \end{bmatrix} \\ g(x) &= x^\top A_T x + a_T \cdot x \\ f(t, x, u) &= M_t x + N_t u \end{aligned} \quad (1.3.1)$$

Here,  $A : [0, T] \rightarrow \mathbb{M}(d, d)$ ,  $B : [0, T] \rightarrow \mathbb{M}(m, m)$ ,  $C : [0, T] \rightarrow \mathbb{M}(d, m)$ ,  $a : [0, T] \rightarrow \mathbb{R}^d$ ,  $b : [0, T] \rightarrow \mathbb{R}^m$ ,  $M : [0, T] \rightarrow \mathbb{M}(d, d)$ , and  $N : [0, T] \rightarrow \mathbb{M}(d, m)$ .

For LQC problem to be well-posed, the value is not  $-\infty$ , it is sufficient that  $Q_t$  and  $A_T$  are symmetric and positive definite. In this case  $C$  and  $g$  are convex functions of  $(x, u)$  and  $x$ , respectively. See exercise below.

**Exercise 12.** In Example 7, change the running cost function to  $C(x, u) = x^2 - u^2$ . Make some effort to find a solution to this modified problem. Of course, since the cost function is not convex, your effort fails. Explain why the optimal control problem is ill-defined.

**Exercise 13.** Write the HJB equation for the LQC problem.

**Exercise 14.** Assume that the cost functions of the general LQC problem are convex. verify that there exists  $G : [0, T] \rightarrow \mathbb{M}(d, d)$ ,  $H : [0, T] \rightarrow \mathbb{R}^d$ , and  $K : [0, T] \rightarrow \mathbb{R}$  such that  $V(t, x) = x^\top G_t x + H_t \cdot x + K_t$  satisfied the HJB equation.

**Exercise 15.** Assume that the cost functions of the general LQC problem are convex. Write Pontryagin principle for the general LQC and suggest a solution method.

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## Chapter 2

# Stochastic Control

We explain the components of a stochastic control problem.

### 2.1 State process

Given a progressively measurable process<sup>1</sup>  $u_t$  taking values in a set  $\mathbf{U} \subseteq \mathbb{R}^m$ , which will be clarified later, the state process is given by the SDE

$$\begin{cases} dx_t^u = \mu(t, X_t^u, u_t)dt + \sigma(t, X_t^u, u_t)dB_t \\ X_0^u = x \in \mathbb{R}^d \end{cases} \quad (2.1.1)$$

**Remark 2.1.1.** If we assume that  $\mu(t, x, u)$  and  $\sigma(t, x, u)$  are Lipschitz continuous in  $(t, x)$  uniformly in  $u$ , i.e.,

$$\sup_{u \in \mathbf{U}} \{|\mu(t, x, u) - \mu(s, y, u)| + |\sigma(t, x, u) - \sigma(s, y, u)|\} \leq K(|t - s| + |x - y|) \quad (2.1.2)$$

Then for any progressively measurable  $u_t$ , (2.1.1) has a strong solution.

**Remark 2.1.2.** In the deterministic control,  $\sigma \equiv 0$ .

**Example 8.** Take  $u \equiv 0$  in  $dX_t^u = u_t X_t^u dt$ . Then,  $X_t^0 = X_0$ . For  $u \equiv 1$ ,  $X_t^1 = X_0 e^t$  and For  $u_t = -t$ ,  $X_t^u = X_0 e^{-t^2/2}$ . In general,  $X_t^u = X_0 e^{\int_0^t u_s ds}$ .

**Example 9.** In Black-Scholes model, the price of an asset satisfies  $ds_t = S_t(\mu dt + \sigma dB_t)$ . If the interest rate  $r = 0$  and  $u_t$  is the amount of money invested in the asset, the wealth process from a self-financing portfolio satisfies

$$dX_t^u = u_t(\mu dt + \sigma dB_t) \quad (2.1.3)$$

Then,  $X_t^u = X_0 \exp(\sigma \int_0^t u_s dB_s + (\mu - \sigma^2/2) \int_0^t u_s^2 ds)$ , where  $X_0$  is the initial wealth.

### 2.2 Objective function

The goal is to minimize (or maximize) an objective function of the following form. Set  $X_0^u = x$  for all  $u$  and define

$$J(x; u) = \mathbb{E} \left[ \int_0^T C(s, X_s^u, u_s) ds + g(X_T^u) \right]. \quad (2.2.1)$$

---

<sup>1</sup>A process is called progressively measurable if for any  $t$ , the restriction to  $[0, t] \times \Omega$  is  $\mathcal{B}[0, t] \otimes \mathcal{F}_t$ -measurable.



where  $C : [0, T] \times \mathbb{R}^d \times \mathbb{R}^m \rightarrow \mathbb{R}$  is called running cost and  $g : \mathbb{R}^d \rightarrow \mathbb{R}$  is the terminal cost. We assume that  $C(t, x, u)$  is Lipschitz in  $(t, x)$  uniformly in  $u$  and  $g$  is Lipschitz. Here, we only focus on minimization as minimization can be obtained by a negative sign.

The value of the problem is

$$V := \inf_u J(x; u). \quad (2.2.2)$$

If a process  $u^*$  exists such that  $V = J(x; u^*)$ , we call  $u^*$  and optimal control.

## 2.3 Admissibility

It is easy to model stochastic control problems from real applications. However, if not carefully done, the stochastic control problem becomes degenerate.

**Exercise 16** (St. Petersburg paradox). Consider the wealth process as described by (9) and assume that asset price  $S$  is a martingale by setting  $\mu = 0$ . In this case,  $X_t = X_0 + \sigma \int_0^t u_s dB_s$ . Further, assume  $X_0 = 0$ . Let's try to maximize expected value of wealth by applying an investment strategy  $u$ :

$$V = \sup_u \mathbb{E}[X_T^u] = \sigma \sup_u \mathbb{E} \left[ \int_0^T u_s dB_s \right]. \quad (2.3.1)$$

Is the stochastic integral above zero? Why?

Now, consider the strategy that chooses  $u_t = u > 0$ , where  $u$  is large and liquidate the investment as soon as  $X_t^u = M$ . If  $X_t^u$  never hits  $M$ , the investment yields  $X_T^u$ .

Note that for  $X_t^u = \sigma u B_t$ , probability of hitting  $M$  before  $T$  is

$$\mathbb{P}(\max_{t \leq T} X_t^u \geq M) = \mathbb{P}\left(\max_{t \leq T} B_t \geq \frac{M}{\sigma u}\right) = 2\mathbb{P}\left(B_T \geq \frac{M}{\sigma u}\right) \quad (2.3.2)$$

The last equality is coming from the Schwartz reflection principle for Brownian motion. If we send  $u \rightarrow \infty$  and  $M \rightarrow \infty$  such that  $\frac{M}{u} \rightarrow 0$ , we obtain  $\mathbb{P}(B_T \geq \frac{M}{\sigma u}) \rightarrow \frac{1}{2}$ .

On the other hand, from the strategy described above, we have

$$\mathbb{E}[X_T^u] = \mathbb{E}[M 1_{\{\max_{t \leq T} X_t^u \geq M\}}] + \mathbb{E}[X_T^u 1_{\{\max_{t \leq T} X_t^u < M\}}] = 2M\mathbb{P}\left(B_T \geq \frac{M}{\sigma u}\right) + \mathbb{E}[X_T^u 1_{\{\max_{t \leq T} X_t^u < M\}}] \quad (2.3.3)$$

Now, answer the following questions:

1. Is it true that  $\mathbb{E}[X_T^u 1_{\{\max_{t \leq T} X_t^u < M\}}] \rightarrow 0$  as  $u \rightarrow \infty$  and  $M \rightarrow \infty$  such that  $\frac{M}{u} \rightarrow 0$ ?
2. What is the limit of  $2M\mathbb{P}(B_T \geq \frac{M}{\sigma u})$ ?  $V = \infty$ ?
3. If the stochastic integral in (2.3.1) is zero, then  $V = 0$ , why did we also get  $V = \infty$ ?

The above example shows that if someone can have unlimited borrowing power, the value functions is infinite. From the practical point of view, this is not possible. If we define a set admissible controls to be all progressively measurable  $u$  such that  $c_t \geq 0$  and  $X_T^u \geq -C$  for some credit limit  $C \geq 0$ <sup>2</sup>  $\mathbb{P}$ -a.s., then the strategies in Example 16 are not admissible.

Another consideration is that the choice of control  $u$  should be made such that the state process in (2.1.1) admits a solution for which  $J(x; u)$  is defined and finite. For instance, in Example 16, there are choices for

<sup>2</sup> $C \geq 0$  is no-short-selling condition.

$\theta_t$  such that  $dX_t = \sigma\theta_t dB_t$  does not have a (strong) solution, see Tsirel'son [1976] for an example when strong solution does not exist. The set of admissible controls should be chosen such that it includes the (unknown) optimal control we are looking for. If too restricted, we would not get optimality; if too wide, we get bizaar solutions such as in St. Petersburg paradox, 16.

Based on the above discussion, we define the set of admissible controls for 16 can be

$$\mathcal{A}_C = \left\{ u : \mathbb{P}\text{-a.s.}, u \text{ is progressively measurable, SDE } dX_t^u = \sigma u_t dB_t \text{ has a "strong" solution } X_t^u, \right. \\ \left. \text{and } X_T^u \geq -C \text{ } \mathbb{P}\text{-a.s.} \right\}. \quad (2.3.4)$$

In the above,  $C \geq 0$  is a given constant. Because  $u \equiv 0$  is an optimal control for the St. Petersburg problem, we can also choose  $\mathcal{A}_0$  as the set of admissible strategies.

Next, we provide a more practical problem.

**Exercise 17** (Merton problem). *In the same setting as 16, consider*

$$V(x) = \sup_u \mathbb{E}[U(X_T^u) | X_0^u = x] \quad (2.3.5)$$

where  $U : \mathbb{R} \rightarrow \mathbb{R}$  is a differentiable strictly concave strictly increasing function such that  $U'(\infty) = 0$ , e.g.,  $U(x) = 1 - e^{-\alpha x}$ ,  $U(x) = \frac{\ln \alpha x}{\alpha}$ , or  $U(x) = \frac{x^{1-\alpha}}{1-\alpha}$  for  $\alpha \in (0, 1)$ . Does the same paradox as in St. Petersburg occurs here? Does the set of admissible controls  $\mathcal{A}_C$  solves the issue?

## 2.4 Value function

Assume that we have a generic control problem

$$\inf_{u \in \mathcal{A}} J(u), \quad J(u) := \mathbb{E} \left[ \int_0^T C(s, X_s^u, u_s) ds + g(X_T^u) \right]. \quad (2.4.1)$$

where  $X_t^u$  satisfies (2.1.1) and  $\mathcal{A}$  is a set of admissible controls, which we assume given.

One way to solve this problem is to define a value function and obtain a dynamic programming equation on the value function.

**Definition 2.4.1** (Value function). *The value function of the control problem is defined by*

$$V(t, \omega) = \text{essinf}_{u \in \mathcal{A}_{t,T}} \mathbb{E} \left[ \int_t^T C(s, X_s^u, u_s) ds + g(X_T^u) \middle| \mathcal{F}_t \right] \quad (2.4.2)$$

In the above,  $\mathcal{A}_{t,T}$  the restriction of the set of admissible controls on the time interval  $[t, T]$ .<sup>3</sup>

Note that  $V(t, \omega)$  is a  $\mathcal{F}_t$ -measurable random variable and in this form it is not useful. However, the following result gives us a lifeline. This well-celebrated result was independently found by Haussmann [1986] and El Karoui et al. [1987] by using Markov selection theorem from Krylov [1973].

<sup>3</sup>Essential infimum. Using essential infimum is necessary because the supremum of random variables is not necessarily measurable (a random variable). For a family of random variables  $\{\chi_u\}_{u \in \mathcal{A}_{t,T}}$ ,  $\text{essinf}_{u \in \mathcal{A}_{t,T}} \chi_u$  is  $\mathbb{P}$ -a.s unique and is a by the smallest random variable that is larger or equal to  $\chi_u$  for all  $u \in \mathcal{A}_{t,T}$ . In our case,  $\chi_u = \mathbb{E} \left[ \int_t^T C(s, X_s^u, u_s) ds + g(X_T^u) \middle| \mathcal{F}_t \right]$ .

**Theorem 2.4.1.** *For any well-defined (proper admissibility condition and finite value function) stochastic control problem, the value function does not change if we reduce the set of controls to **Markovian controls**, i.e.,  $u_t := \phi(t, X_t^u)$ .*

Markovian controls are also called feedback controls, especially if the problem is deterministic.

**Corollary 2.4.1.** *For the value function we have  $V(t, \omega) = V(t, X_t^u)$  where*

$$\begin{aligned} V(t, x) &= \inf_{u \in \mathcal{A}_{t,T}^f} \mathbb{E} \left[ \int_t^T C(s, X_s^u, u_s) ds + g(X_T^u) \middle| X_t^u = x \right] \\ &= \inf_{u \in \mathcal{A}_{t,T}^f} \mathbb{E} \left[ \int_t^T C(s, X_s^u, u_s) ds + g(X_T^u) \middle| X_t^u = x \right] \end{aligned} \quad (2.4.3)$$

where  $\mathcal{A}_{t,T}^f$  is the set of all Markovian controls  $u_t = \phi(t, X_t^u)$ ; in particular the following SDE has a strong solution for all  $x$

$$\begin{cases} dX_s = \mu(s, X_s, \phi(s, X_s)) ds + \sigma(s, X_s, \phi(s, X_s)) dB_s \\ X_t = x \end{cases} \quad (2.4.4)$$

**Remark 2.4.1** (Discussion of strong and weak solution). *The requirement of existence of strong solution for (2.4.4) is for the sake of simplicity. We recall that an SDE has a strong solution when for all filtered probability space hosting a Brownian motion, the SDE has a solution. An SDE has a weak solution if there exists a filtered probability space hosting a Brownian motion in which the SDE has a solution. For more discussion and example of nonexistence of strong solutions see [Tsirel'son \[1976\]](#) and the discussion in [Rogers and Williams \[2000\]](#).*

It is worth mentioning there are three type of controls:

1. Open-loop controls which do not care about the current state of the system. In our case, an open-loop control is a deterministic control. Control  $u_t = \alpha'_t + \beta \alpha_t$  in Example 1.1.7 is open-loop because it does not depend on the state of system  $x_t$ .
2. A closed-loop or feedback control is a control that depends on the path of the state process, i.e.,  $u_t = \phi(t, X_{\cdot \wedge t}) = \phi(t, X_s : s \leq t)$  or equivalently  $u_t$  is adapted to the filtration generated by  $X$ ,  $\{\mathcal{F}_t^X\}_t$ , measurable. Note that a closed-loop control is automatically adapted to the filtration generated by the Brownian motion.
3. A Markovian control is a special case of feedback control that only depends on the latest value of the state variable,  $u_t = \phi(t, X_t)$ . In this course, when we discuss the existence of optimal control, we mean Markovian optimal control unless otherwise is specified.

## 2.5 Discrete-time dynamic programming principle (DPP)

Recall from Corollary 2.4.1 that the value function satisfies

$$V(t, x) = \inf_{u \in \mathcal{A}_{t,T}^f} \mathbb{E} \left[ \int_t^T C(s, X_s^u, u_s) ds + g(X_T^u) \middle| X_t^u = x \right] \quad (2.5.1)$$

Dynamic programming principle can be better understood in discrete-time setting. So, here we spend some time to explain a stochastic control problem in discrete time. It is natural to expect that, under suitable

conditions, any continuous-time problem can be approximated by a discrete-time problem. For instance, (2.4.1) can be approximated by

$$\inf_u \mathbb{E} \left[ \sum_{t=0}^{T-1} C(t, \hat{X}_t^u, u_t) \Delta t + g(\hat{X}_T^u) \right] \quad (2.5.2)$$

where

$$\hat{X}_{t+1}^u = \hat{X}_t^u + \mu(t, \hat{X}_t^u, u_t) \Delta t + \sigma(t, \hat{X}_t^u, u_t) \Delta B_{t+1} \quad (2.5.3)$$

where  $\Delta t = T/N$  and  $\Delta B_{t+1} = B_{t+\Delta t} - B_t$ .

Without loss of generality, we can drop  $\Delta t$  and replace  $\Delta B_{t+1}$  by a standard i.i.d random variable.

$$\inf_u \mathbb{E} \left[ \sum_{t=0}^{T-1} C(t, X_t^u, u_t) + g(X_T^u) \right] \quad (2.5.4)$$

where the infimum is over all stochastic process  $u : [0, T] \times \Omega \rightarrow \mathbb{R}^m$  and

$$X_{t+1}^u = X_t^u + \mu(t, X_t^u, u_t) + \sigma(t, X_t^u, u_t) \xi_{t+1} \quad (2.5.5)$$

and  $\{\xi_t\}_{t=1}^T$  is a sequence of i.i.d. random variables with mean 0 and variance 1.

Note that we can think about the process as a sequence of random variables  $u_0, \dots, u_{T-1}$ , which can be chosen based on the discretion of the controller up to the adaptedness condition. This allows us to write the control problem as

$$\inf_{u_0} \cdots \inf_{u_{T-1}} \mathbb{E} \left[ \sum_{t=0}^{T-1} C(t, X_t^u, u_t) + g(X_T^u) \right] \quad (2.5.6)$$

The value function for this problem is written as

$$V(t, x) = \inf_{u_t} \cdots \inf_{u_{T-1}} \mathbb{E} \left[ \sum_{s=t}^{T-1} C(s, X_s^u, u_s) + g(X_T^u) \middle| X_t^u = x \right] \quad (2.5.7)$$

Since  $X_t^u = x$ ,  $u_t$  is chosen over all real numbers and, therefore,  $C(t, X_t^u, u_t)$  is deterministic. One can write

$$V(t, x) = \inf_{u_t} C(t, x, u_t) + \inf_{u_{t+1}} \cdots \inf_{u_{T-1}} \mathbb{E} \left[ \sum_{s=t+1}^{T-1} C(s, X_s^u, u_s) + g(X_T^u) \middle| X_t^u = x \right] \quad (2.5.8)$$

Tower property of conditional expectation implies that

$$\mathbb{E} \left[ \sum_{s=t+1}^{T-1} C(s, X_s^u, u_s) + g(X_T^u) \middle| X_t^u = x \right] = \mathbb{E} \left[ \mathbb{E} \left[ \sum_{s=t+1}^{T-1} C(s, X_s^u, u_s) + g(X_T^u) \middle| X_{t+1}^u \right] \middle| X_t^u = x \right] \quad (2.5.9)$$

**Exercise 18.** Show that one can take the infimum inside the first expectation, i.e.,

$$\begin{aligned} & \inf_{u_{t+1}, \dots, u_{T-1}} \mathbb{E} \left[ \mathbb{E} \left[ \sum_{s=t+1}^{T-1} C(s, X_s^u, u_s) + g(X_T^u) \middle| X_{t+1}^u \right] \middle| X_t^u = x \right] \\ &= \mathbb{E} \left[ \inf_{u_{t+1}, \dots, u_{T-1}} \mathbb{E} \left[ \sum_{s=t+1}^{T-1} C(s, X_s^u, u_s) + g(X_T^u) \middle| X_{t+1}^u \right] \middle| X_t^u = x \right] \end{aligned} \quad (2.5.10)$$

By definition of value function,

$$\inf_{u_{t+1}, \dots, u_{T-1}} \mathbb{E} \left[ \sum_{s=t+1}^{T-1} C(s, X_s^u, u_s) + g(X_T^u) \middle| X_{t+1}^u \right] = V(t+1, X_{t+1}^u) \quad (2.5.11)$$

Therefore,

$$V(t, x) = \inf_{u_t} C(t, x, u_t) + \mathbb{E}[V(t+1, X_{t+1}^u) | X_t = x] \quad (2.5.12)$$

This provides us with a recursive formula to solve a discrete-time control problem.

$$\begin{cases} V(t, x) = \inf_{u_t} C(t, x, u_t) + \mathbb{E}[V(t+1, X_{t+1}^u) | X_t = x] \\ V(T, x) = g(x) \end{cases} \quad (2.5.13)$$

## 2.6 Dynamic programming principle in continuous time

In continuous time, dynamic programming principle is more complicated. Early results used measurable selection theorems to overcome these complications. Later, quantization methods were used to avoid measurable selection theorems. For more information, see [Touzi \[2012\]](#) and references therein.

**Theorem 2.6.1.** *If the value function is continuous and have linear growth, then for any stopping time  $\tau$ , we have*

$$V(t, x) = \inf_u \mathbb{E}_{t,x} \left[ \int_t^\tau C(s, X_s^u, u_s) ds + V(\tau, X_\tau^u) \right] \quad (2.6.1)$$

## 2.7 Hamilton-Jacobi-Bellman equation

In Theorem 2.6.1, for fixed  $h, \epsilon > 0$  take  $\tau = \tau^h = \inf\{s > 0, |X_s^u - x| \geq \epsilon\} \wedge (t + h)$ . Assume that the value function is  $C^{1,2}$ , once continuously differentiable on  $t$  and twice continuously differentiable on  $x$ . By applying Itô's formula on  $V(\tau^h, X_{\tau^h}^u)$ , we obtain

$$\begin{aligned} V(\tau^h, X_{\tau^h}^u) &= V(t, x) \\ &+ \int_t^{\tau^h} \left( \partial_t V(s, X_s^u) + \frac{1}{2} D^2 V(s, X_s^u) \cdot a(s, X_s^u, u_s) + \nabla V(s, X_s^u) \cdot \mu(s, X_s^u, u_s) \right) ds \\ &+ \int_t^{\tau^h} \sigma(s, X_s^u, u_s) dB_s \end{aligned} \quad (2.7.1)$$

where  $a = \sigma^\top \sigma$ . Recall for two matrices  $A \cdot B = \text{Tr}[A^\top B]$ . Assuming  $\mathbb{E}_{t,x} \left[ \int_0^{\tau^h} \sigma(s, X_s^u, u_s) dB_s \right] = 0$ <sup>4</sup>, (2.6.1) can be written as

$$\begin{aligned} 0 &= \inf_u \mathbb{E}_{t,x} \left[ \int_t^{\tau^h} \left( C(s, X_s^u, u_s) \right. \right. \\ &\quad \left. \left. \partial_t V(s, X_s^u) + \frac{1}{2} D^2 V(s, X_s^u) \cdot a(s, X_s^u, u_s) + \nabla V(s, X_s^u) \cdot \mu(s, X_s^u, u_s) \right) ds \right] \end{aligned} \quad (2.7.2)$$

<sup>4</sup>One can always choose  $\epsilon$  in the definition of  $\tau^h$  such that the stochastic integral has zero expected value.

Note that  $\tau^h = O(h)$ . Therefore, we divide by  $h$  and send  $h \rightarrow 0$ :

$$\begin{aligned}
 0 &= \lim_{h \rightarrow 0} \inf_u \mathbb{E}_{t,x} \left[ \frac{1}{h} \int_t^{\tau^h} \left( C(s, X_s^u, u_s) \right. \right. \\
 &\quad \left. \left. \partial_t V(s, X_s^u) + \frac{1}{2} D^2 V(s, X_s^u) \cdot a(s, X_s^u, u_s) + \nabla V(s, X_s^u) \cdot \mu(s, X_s^u, u_s) \right) ds \right] \\
 &= \inf_u \lim_{h \rightarrow 0} \mathbb{E}_{t,x} \left[ \frac{1}{h} \int_t^{\tau^h} \left( C(s, X_s^u, u_s) \right. \right. \\
 &\quad \left. \left. \partial_t V(s, X_s^u) + \frac{1}{2} D^2 V(s, X_s^u) \cdot a(s, X_s^u, u_s) + \nabla V(s, X_s^u) \cdot \mu(s, X_s^u, u_s) \right) ds \right] \quad (2.7.3) \\
 &= \inf_u \mathbb{E}_{t,x} \left[ \lim_{h \rightarrow 0} \frac{1}{h} \int_t^{\tau^h} \left( C(s, X_s^u, u_s) \right. \right. \\
 &\quad \left. \left. \partial_t V(s, X_s^u) + \frac{1}{2} D^2 V(s, X_s^u) \cdot a(s, X_s^u, u_s) + \nabla V(s, X_s^u) \cdot \mu(s, X_s^u, u_s) \right) ds \right] \\
 &= \partial_t V(t, x) + \inf_u \left\{ C(t, x, u) + \frac{1}{2} D^2 V(t, x) \cdot a(t, x, u) + \nabla V(t, x) \cdot \mu(t, x, u) \right\}
 \end{aligned}$$

The HJB is given by the following PDE:

$$\begin{cases} 0 = \partial_t V(t, x) + \inf_u \left\{ C(t, x, u) + \frac{1}{2} D^2 V(t, x) \cdot a(t, x, u) + \nabla V(t, x) \cdot \mu(t, x, u) \right\} \\ V(T, x) = g(x) \end{cases} \quad (2.7.4)$$

**Remark 2.7.1.** Rigorously, the proof is more complicated. First, we need to justify  $\lim_{h \rightarrow 0} \inf_u = \inf_u \lim_{h \rightarrow 0}$ . To show this, we have to show that the expected value as a function of  $h$  is continuous uniformly on  $u$ . This can be obtained by continuity of  $C$ ,  $V$ , and first and second derivatives of  $V$ . Second, justification of  $\lim_{h \rightarrow 0} \mathbb{E} = \mathbb{E} \lim_{h \rightarrow 0}$  requires dominated convergence theorem, which allows to change the order of expected value and limit. This also requires use of mean value theorem. For more information, see [Touzi \[2012\]](#).

**Example 10.** Consider the stochastic linear quadratic problem below:

$$\inf_u \mathbb{E} \left[ \int_0^T \left( aX_t^2 + bX_t + Au_t^2 + Bu_t \right) dt + \alpha X_T^2 + \beta X_T \right] \quad (2.7.5)$$

with

$$dX_t = (cX_t + du_t)dt + \sigma dB_t \quad (2.7.6)$$

where  $a, A, \alpha > 0$  and  $b, B, c$ , and  $d$  are constants. The HJB is given by

$$\begin{cases} 0 = \partial_t V(t, x) + \frac{\sigma^2}{2} \partial_x^2 V(t, x) + cx \partial_x V(t, x) + ax^2 + bx + \inf_u \left\{ Au^2 + (B + d \partial_x V(t, x))u \right\} \\ V(T, x) = \alpha x^2 + \beta x \end{cases} \quad (2.7.7)$$

$$u = -\frac{B + d \partial_x V(t, x)}{2A}$$

$$\inf_u \{ Au^2 + (B + d \partial_x V(t, x))u \} = -\frac{(B + d \partial_x V(t, x))^2}{4A}$$

$$\begin{cases} 0 = \partial_t V(t, x) + \frac{\sigma^2}{2} \partial_x^2 V(t, x) + cx \partial_x V(t, x) + ax^2 + bx - \frac{(B + d \partial_x V(t, x))^2}{4A} \\ V(T, x) = \alpha x^2 + \beta x \end{cases} \quad (2.7.8)$$

We anticipate that  $V(t, x) = f(t)x^2 + h(t)x + k(t)$ , which is a separation of variable technique. If we plug

in  $V(t, x)$  into the HJB, we obtain

$$\begin{cases} f' = -2cf(t) + \frac{f^2(t)}{A^2} - a & f(T) = \alpha \\ h'(t) = (\frac{f(t)}{A^2} - c)h(t) + \frac{Bd}{A^2}f(t) - b & h(T) = \beta \\ k'(t) = \frac{PB^2}{4A} + \frac{bd}{2A^2}h(t) + \frac{1}{4A^2}h^2(t) + f(t) & k(T) = 0 \end{cases} \quad (2.7.9)$$

Note that the above system of ODEs is fully solvable. The first ODE is Riccati equation. We shall show in Section 2.9 that the solution of the PDE is indeed the value function of the linear quadratic optimal control problem.

**Exercise 19.** Solve modified version of Exercise 10 with modification

$$dX_t = (cX_t + du_t)dt + (eX_t + fu_t)dB_t \quad (2.7.10)$$

**Exercise 20.** Extend the stochastic linear quadratic problem in Exercise 19 to higher dimension and write the HJB.

**Example 11** (Merton optimal investment problem). Remember a self-financing portfolio with a Black-Scholes risky asset,  $dS_t = S_t(\mu dt + \sigma dB_t)$ , under zero interest rate satisfies

$$dX_t = \theta_t(\mu dt + \sigma dB_t) \quad (2.7.11)$$

where  $\theta$  is the amount of money invested in the risky asset. Merton problem is to maximize the expected utility of wealth at a time horizon  $T$ :

$$\sup_{\theta} \mathbb{E}[U(X_T^\theta)] \quad (2.7.12)$$

The HJB for this problem is given by

$$\begin{cases} 0 = \partial_t V + \sup_{\theta} \left\{ \frac{\sigma^2}{2} \partial_x^2 V + \mu \partial_x V \right\} \\ V(T, x) = U(x) \end{cases} \quad (2.7.13)$$

After simplification, we have

$$\begin{cases} 0 = \partial_t V - \frac{\mu^2 (\partial_x V)^2}{2\sigma^2 \partial_x^2 V} \\ V(T, x) = U(x) \end{cases} \quad (2.7.14)$$

We can find a closed form solution for the following cases of utility by using separation of variables,  $V(t, x) = f(t)U(x)$ :

1. Exponential utility  $U(x) = 1 - e^{-\lambda x}$ . The separation of variables is  $V(t, x) = -f(t)e^{-\lambda x}$ .
2. HARA utility  $U(x) = \frac{x^{1-\lambda}}{1-\lambda}$ , with  $\lambda \in (0, 1)$ .

**Exercise 21** (Merton optimal consumption problem). It is similar to the Merton optimal investment problem except, the investor is consuming from the account and what matters is the utility of consumption.

$$dX_t = \theta_t(\mu dt + \sigma dB_t) - c_t dt \quad (2.7.15)$$

where  $\theta$  is the amount of money invested in the risky asset and the consumption  $c$  satisfies  $c_t \geq 0$ . Merton problem is

$$\sup_{c \geq 0, \theta} \mathbb{E} \left[ \int_0^\infty e^{-\gamma t} U(c_t) dt \right], \quad \gamma > 0 \quad (2.7.16)$$

Write the HJB for this problem.

**Remark 2.7.2.**  $\gamma > 0$  in the above example represents preference of current consumption over future consumption.

### 2.7.1 Validity of HJB equation

In the previous section, we showed that if the value function is in  $C^{1,2}$ , then it satisfies (2.7.4). Therefore, HJB can be used to find the value function. If we know that the HJB has a unique  $C^{1,2}$  solution in a suitable class of function which includes the value function, we are done. However, it is not always easy to obtain existence and uniqueness results for nonlinear PDEs, which HJB equations are.

**Example 12.** Recall from Exercise 5 the fastest exit problem:

$$\inf_{|u_t| \leq 1} \int_0^\infty 1_{\{X_t^u \in O\}} dt \quad (2.7.17)$$

with  $dX_t^u = u_t dt$ . The HJB for this problem is  $0 = \inf_{|u| \leq 1} \{u \cdot \nabla V(x)\} + 1_{\{x \in O\}}$ , which is equivalent to the following boundary value problem

$$\begin{cases} 0 = -|\nabla V(x)| + 1 & x \in O \\ 0 = V(x) & x \in \partial O \end{cases} \quad (2.7.18)$$

The optimal feedback (closed-loop) exit strategy is give by  $u_t^*(x) = -\nabla V(x)$ , which is equivalent to gradient descent.

When  $O = [-1, 1] \subset \mathbb{R}$ , the equation is  $|V'(x)| = 1$ . On the other hand, this one-dimensional simplification has an obvious answer: run as fast as you can to the nearest exit point. This yields  $V(x) = 1 - |x|$ . The value function satisfies the HJB except at  $x = 0$ , which is the sole point with two optimal exit strategies.

Assume that a robot only knows how to solve HJB equations and it ignores finite number of point of irregularities of the solution. For the robot the function  $v(x)$  given below is as good of a solution as  $V(x) = 1 - |x|$ .

$$v(x) = \begin{cases} 1 - |x| & \frac{1}{2} \leq |x| \leq 1 \\ |x| & |x| \leq \frac{1}{2} \end{cases} \quad (2.7.19)$$

However, if the robot comes of with  $v$  instead of  $V$ , it suggests to move towards 0 for  $|x| \leq \frac{1}{2}$ , which is obviously incorrect.

The above example questions the validity of HJB approach. There are two ways to address the validity issue for the solution of HJB equations. One is the notion of viscosity solution, which singles out  $V(x) = 1 - |x|$ , as the unique viscosity solution of the HJB. For more information of the viscosity solution approach, see [Crandall et al. \[1992\]](#). There are simple criteria that the robot can check to see whether his approximate solution to the HJB equation converges to the unique viscosity solution. See for example, [Barles and Souganidis \[1991\]](#). However, here we closed this discussion by giving a glimpse of what it means to be a viscosity solution in the following example.

**Example 13.** Consider the exit time problem with noise:

$$\inf_{|u_t| \leq 1} \mathbb{E} \left[ \int_0^\infty 1_{\{X_t^u \in O\}} dt \right] \quad (2.7.20)$$



with  $dX_t^u = u_t dt + \epsilon dB_t$ . The HJB is given by  $0 = \frac{\epsilon^2}{2} V''(x) + \inf_{|u| \leq 1} \{u \cdot \nabla V(x)\} + 1_{\{x \in O\}}$ <sup>5</sup>, which is equivalent to the following boundary value problem

$$\begin{cases} 0 = \frac{\epsilon^2}{2} V''(x) - |\nabla V(x)| + 1 & x \in O \\ 0 = V(x) & x \in \partial O \end{cases} \quad (2.7.21)$$

Back to the simple case where  $O = [-1, 1]$ , there exists a unique bounded solution,  $V^\epsilon(x)$ , in  $C^2$ . One can find this solution in closed form and verify that as  $\epsilon \rightarrow 0$ ,  $V^\epsilon(x) \rightarrow 1 - |x|$ . This clarifies why  $v(x)$  is not a solution.

**Exercise 22.** In the previous exercise, show that as  $\epsilon \rightarrow 0$ ,  $V^\epsilon(x) \rightarrow 1 - |x|$ .

## 2.8 Viscosity solutions to the HJB equations

Before moving to the notion of viscosity solutions, we point out an important remark about partial differential equations (PDE) such as the HJBs. The solution to the simplest equations may not exist. For instance, the backward heat equation given by

$$\begin{cases} 0 = \partial_t V + \frac{\sigma^2}{2} \Delta V - rV + f(t, x) \\ V(0, x) = g(x) \end{cases} \quad (2.8.1)$$

does not have any solution unless  $f$  and  $g$  have sufficient regularity. However, the forward (regular) heat equation

$$\begin{cases} 0 = \partial_t V + \frac{\sigma^2}{2} \Delta V - rV + f(t, x) \\ V(T, x) = g(x) \end{cases} \quad (2.8.2)$$

has a classical solution regardless of the choice of  $g$  and  $f$ . The forward heat equation is a parabolic equation while the backward equation is not. Roughly speaking, in a parabolic equation,  $\partial_t V$  and  $\frac{\sigma^2}{2} \Delta V$  show up with the same sign. In the HJB equation (2.7.4), the coefficient of  $\partial_t V$  and  $\frac{\sigma^2}{2} \Delta V$ , 1 and  $\frac{\sigma^2}{2}$  respectively, have the same sign. More rigorously, a nonlinear equation of the form

$$\begin{cases} 0 = \partial_t V(t, x) + F(t, x, V(t, x), \nabla V(t, x), D^2 V(t, x)) \\ V(T, x) = g(x) \end{cases} \quad (2.8.3)$$

is parabolic if the nonlinearity,  $F(t, x, \varrho, \Pi, \Gamma)$ , is nondecreasing in the component  $\Gamma$ , which represents the Hessian matrix in the order induced by nonnegative-definite matrices. I.e., if for two matrices  $\Gamma_1$  and  $\Gamma_2$ ,  $\Gamma_2 - \Gamma_1$  is a nonnegative-definite matrix, then  $F(t, x, \varrho, \Pi, \Gamma_2) \geq F(t, x, \varrho, \Pi, \Gamma_1)$ . One can check that

$$F(t, x, \varrho, \Pi, \Gamma) := \sup_{u \in \mathbf{U}} \left\{ C(t, x, u) + \mu(t, x, u) \Pi + \frac{1}{2} \sigma^2(t, x, u) \Gamma - r(t, x, u) \varrho \right\} \quad (2.8.4)$$

is nondecreasing in  $\Gamma$ .

**Example 14.** Show that the HJB equation (2.9.16) is a parabolic equation.

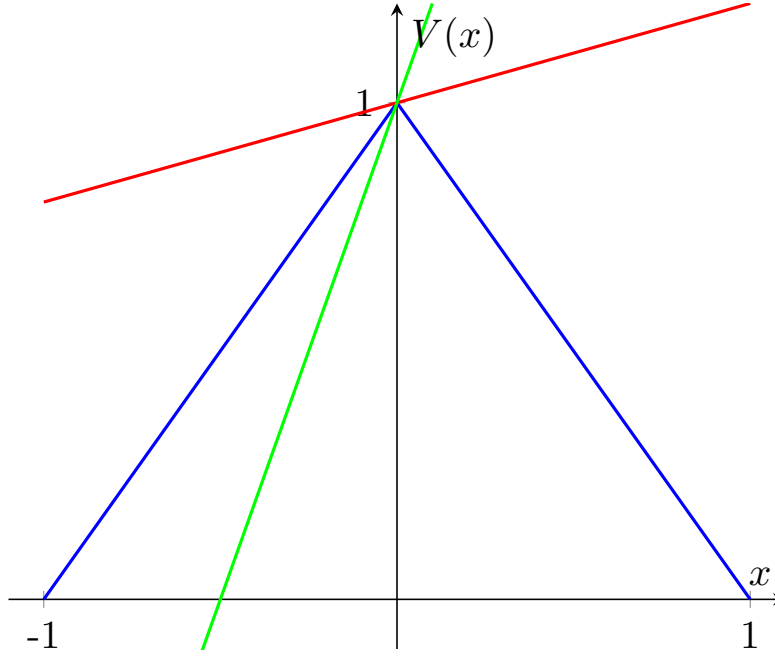
On the other hand, even an equation as simple as the heat equation may have multiple solution if we do not restrict the solution to a proper class of functions. For instance, if  $g(x)$  and  $f(t, x)$  have linear growth

<sup>5</sup>My best guess is that the term viscosity comes a fluid mechanics term for the second derivative  $\frac{\epsilon^2}{2} V''$ .

in  $x$ , then the heat equation has a unique solution in the class of function with linear growth, but not in the class of functions with super exponential growth.

A linear partial differential equation such as the heat equation has classical solution. Nonlinear equations, however, do not necessary have a classical solution. A classical solution is a solution that is continuously differentiable as much as the equation requires. For instance, if in the heat equation above,  $f(t, x)$  is a continuous function, then a solution  $V(t, x)$  is a classical solution if  $V \in C^{1,2}$ . But, as seen in Exercise 5,  $V(x) = 1 - |x|$  is the value function that we expect to satisfy the HJB equation. However, it is not differentiable as  $x = 0$ .

A suitable notion of weak solution for the HJB equations, which are nonlinear equations, is viscosity solutions. To understand the viscosity solutions, we first need to provide candidates for the weak derivatives of a function at the points of nondifferentiability. The first derivative can always be represented by the slope of a tangent plane. Locally speaking, the tangent plane is a plane that hits the graph of the function at at least one point locally and keeps the graph on one side. For example, for the function  $V(x) = 1 - |x|$  at point  $x = 0$ , the line  $y = 1 + mx$  is a tangent line by  $x = 0$  for  $m \in [-1, 1]$  (red line) but  $y = 1 + mx$  for  $m \notin [-1, 1]$  is not (green line). See Figure 2.8.1. Therefore, we have



**Figure 2.8.1:** Red line is an acceptable tangent to the graph, while the green line is not. The green line does not keep the graph on one side locally.

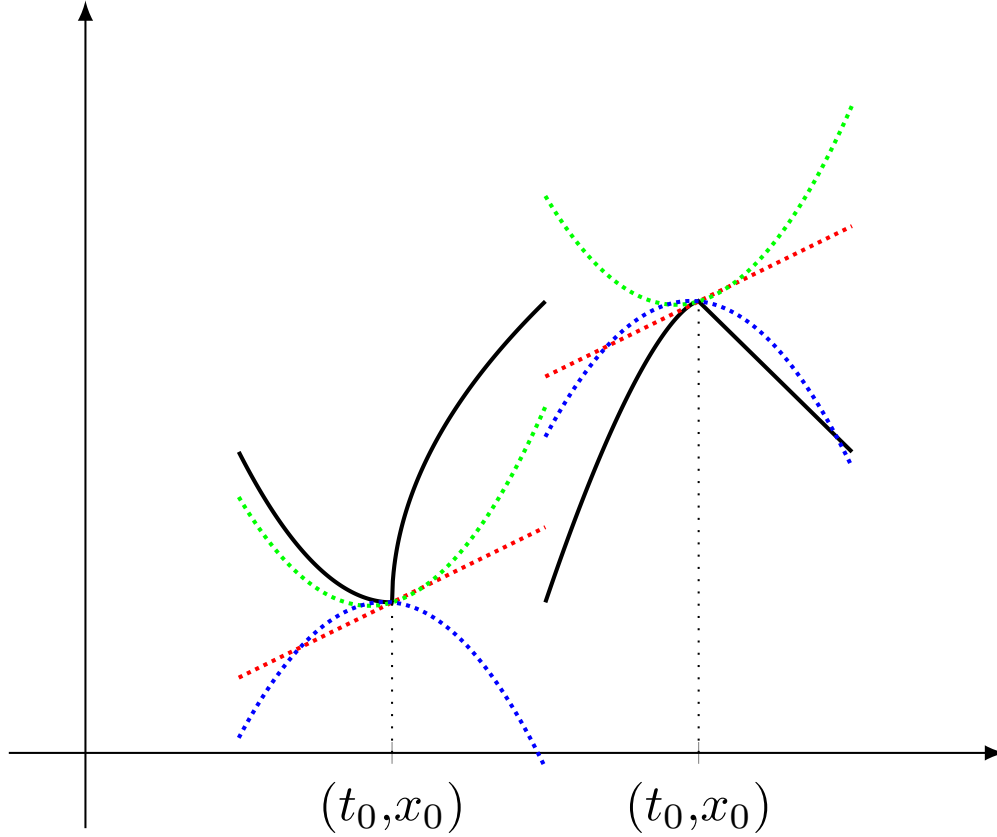
$\varphi(t, x) = a \cdot (x - x_0) + b(t - t_0) + V(t_0, x_0)$  such that  $\varphi(t, x) - f(t, x)$  has a local extrema at  $(t_0, x_0)$ .

For the second derivative tangent planes are not sufficient, because the second derivatives of the tangent planes are always zero. Therefore, we need to appeal to the quadratic functions of the form  $\varphi(t, x) = (x - x_0) \cdot A(x - x_0) + a \cdot (x - x_0) + b(t - t_0) + V(t_0, x_0)$ . Note that since we have first derivative on  $t$ , we only use first order term on  $t$ . If a function  $\varphi(t, x)$  touches  $V(t, x)$  at point  $(t_0, x_0)$  from above (resp. below), i.e.,  $(t_0, x_0)$  is a local minimum (maximum) point for  $\varphi - V$ , then we call  $(b, a, A) = (\partial_t \varphi(t_0, x_0), \nabla \varphi(t_0, x_0), D^2 \varphi(t_0, x_0))$  a superderivative (resp. sub) of  $V$  at point  $(t_0, x_0)$ . Such functions  $\varphi(t, x)$  are called test functions. See Figure 2.8.2.

A function  $\underline{V}$  (resp.  $\bar{V}$ ) is called a viscosity subsolution (resp. super) of (2.8.3) if for any superderivative

(resp. sub)  $(b, a, A)$  at point  $(t_0, x_0)$ , we have

$$\begin{cases} 0 \geq b + F(t_0, x_0, \varphi(t_0, x_0), a, A), & (\text{resp. } 0 \leq) \\ V(T, x) \leq g(x) & (\text{resp. } 0 \geq) \end{cases} \quad (2.8.5)$$



**Figure 2.8.2:** The test functions that represent the sub and superderivatives of a function at the points of nondifferentiability as well as other points that the function is differentiable.

A function is called a viscosity solution if it is a sub and a super solution.

**Remark 2.8.1.** The subsolution (resp. super) is named in accordance with submartingale (resp. super). If  $U(t, x) \in C^{1,2}$  is a subsolution (resp. super) to the equation

$$0 = \partial_t V(t, x) + C(t, x) + \mu(t, x) \partial_x V(t, x) + \frac{1}{2} \sigma^2(t, x) \partial_{xx} V(t, x), \quad (2.8.6)$$

then  $Y_t := U(t, X_t)$  is a submartingale (resp. super) martingale, where

$$dX_t = \mu(t, X_t)dt + \sigma(t, X_t)dB_t. \quad (2.8.7)$$

For the first order HJB equation, finding the direction of the inequality for the subsolution or super solutions is not particularly obvious. It can be determined while we add a singular perturbation term  $\varepsilon B_t$  to the deterministic state process and send  $\varepsilon \rightarrow 0$ . For instance, in Exercise 5, if we set the state process

$dX_t^u = u_t dt + \varepsilon dB_t$ , we obtain the HJB

$$\begin{cases} 0 = 1 - |V'(x)| + \frac{\varepsilon^2}{2} V''(x) \\ V(\pm 1) = 0 \end{cases} \quad (2.8.8)$$

In this case, the subsolution (resp. super) can satisfies

$$\begin{cases} 0 \geq 1 - |V'(x)| + \frac{\varepsilon^2}{2} V''(x) & (\text{resp. } 0 \leq) \\ V(\pm 1) \leq 0 & (\text{resp. } 0 \geq) \end{cases} \quad (2.8.9)$$

By sending  $\varepsilon \rightarrow 0$ , we obtain the following inequalities for the Eikonal equation, (2.8.8).

$$\begin{cases} 0 \geq 1 - |V'(x)| & (\text{resp. } 0 \leq) \\ V(\pm 1) \leq 0 & (\text{resp. } 0 \geq) \end{cases} \quad (2.8.10)$$

We can easily check that the function  $V(t, x) = 1 - |x|$  is a viscosity solution of the Eikonal equation, (2.8.8). For example, at all points  $x \neq 0$ , the sub and superderivatives are equal to  $-\text{sgn}(x)$ , which obviously satisfy the equation with equality. At  $x = 0$ , we only have super derivatives and the set of subderivatives is empty set. This makes the supersolution property to hold obviously by the false antecedent. The set of superderivatives contains all  $m \in [-1, 1]$ . Since  $0 \leq 1 - |m|$ , then  $V(t, x) = 1 - |x|$  is also a subsolution. All the above arguments are valid since the boundary condition  $V(t, \pm 1) = 0$  validates the sub and supersolution properties.

On the other hand, the function  $\tilde{V} = \min\{|x|, 1 - |x|\}$  only satisfies the subsolution property.

**Exercise 23.** Show that  $\tilde{V}$  is a viscosity subsolution to the Eikonal equation (2.8.10), but not a super solution. Explore the points  $x = \pm \frac{1}{2}$ .

The existence and uniqueness of the solution to the nonlinear parabolic equations in the class of functions with linear growth is studies in several papers. For instance, see Crandall et al. [1992]. For the HJB equations derived from the optimal control problems, the value function of the control problem is usually the viscosity solution of the HJB. The uniqueness is due to a technical lemma, Ishii's lemma. However, in most cases, if we manage to show that any viscosity solution is indeed a classical solution,  $C^{1,2}$ , then the verification theorem, Theorem 2.9.1 shows that any classical solution is equal to the value function, and therefore, uniqueness is obtained.

## 2.9 Verification

Instead of going through viscosity solutions, we consider the cases where the HJB equation has a  $C^{1,2}$  solution, which is a candidate for the value function. Then, we introduce a verification theorem, which approves that the solution is indeed the value function.

**Theorem 2.9.1** (Verification). *Let HJB equation, (2.7.4) has a  $C^{1,2}$  solution,  $v(t, x)$ . Then,  $v(t, x) \geq V(t, x)$ . In addition, assume that there exists  $u^*(t, x)$  such that*

$$u^*(t, x) \in \operatorname{argmin}_u \left\{ C(t, x, u) + \frac{1}{2} D^2 v(t, x) \cdot a(t, x, u) + \nabla v(t, x) \cdot \mu(t, x, u) \right\} \quad (2.9.1)$$

*and  $u^*(t, x)$  is an admissible Markovian (feedback) control, i.e., (2.4.4) has a (strong) solution,  $X_t^*$ , and  $u_t^* = u^*(t, X_t^*)$  is admissible,  $u^* \in \mathcal{A}$ . Then,  $v(t, x) = V(t, x)$ .*

*Proof.* Since  $v$  satisfies (2.7.4) in classical sense, we conclude that for any value of  $u \in U$ , we have

$$0 \geq \partial_t v(t, x) + C(t, x, u) + \mu(t, x, u) \nabla v(t, x) + \frac{1}{2} a(t, x, u) \cdot D^2 v(t, x) \quad (2.9.2)$$

Since  $v(t, x) \in C^{1,2}$ , it follows from the Itô formula that for any  $u \in \mathcal{A}_{t,T}$ ,

$$\begin{aligned} v(T, X_T^{t,x,u}) &= v(t, x) \\ &+ \int_0^T \left( \partial_t v(s, X_s^{t,x,u}) + \mu(s, X_s^{t,x,u}, u_s) \cdot \nabla v(s, X_s^{t,x,u}) + \frac{1}{2} a(s, X_s^{t,x,u}, u_s) \cdot D^2 v(s, X_s^{t,x,u}) \right) ds \\ &+ \int_t^T \sigma(s, X_s^{t,x,u}, u_s) dB_s \end{aligned} \quad (2.9.3)$$

By (2.9.2) and terminal condition  $v(T, x) = g(x)$ , we obtain

$$g(X_T^u) \geq v(t, x) - \int_0^T C(s, X_s^{t,x,u}, u_s) ds + \int_t^T \sigma(s, X_s^{t,x,u}, u_s) dB_s \quad (2.9.4)$$

If we take conditional expectation, the martingale property of the stochastic integral implies that

$$\mathbb{E}_{t,x} \left[ g(X_T^u) + \int_0^T C(s, X_s^{t,x,u}, u_s) ds \right] \geq v(t, x) \quad (2.9.5)$$

Now, by taking supremum over  $u \in \mathcal{A}_{t,T}$ , we obtain

$$V(t, x) = \sup_{u \in \mathcal{A}_{t,T}} \mathbb{E}_{t,x} \left[ g(X_T^u) + \int_0^T C(s, X_s^{t,x,u}, u_s) ds \right] \geq v(t, x). \quad (2.9.6)$$

Now, assume that for some  $u^*(t, X_t^*) \in \mathcal{A}$ ,  $Y_t^{u^*} = v(t, X_t^{u^*}) + \int_0^t C(s, X_s^{u^*}, u_s^*) ds$  is a martingale. By applying the Itô formula, we have

$$\begin{aligned} Y_t^{u^*} &= v(t, x) + \int_t^T \sigma(s, X_s^{t,x,u^*}, u_s^*) dB_s + \int_0^T \left( \partial_t v(s, X_s^{t,x,u^*}) + \mu(s, X_s^{t,x,u^*}, u_s^*) \cdot \nabla v(s, X_s^{t,x,u^*}) \right. \\ &\quad \left. + \frac{1}{2} a(s, X_s^{t,x,u^*}, u_s^*) \cdot D^2 v(s, X_s^{t,x,u^*}) + C(s, X_s^{u^*}, u_s^*) \right) ds \\ &= v(t, x) + \int_t^T \sigma(s, X_s^{t,x,u^*}, u_s^*) dB_s \end{aligned} \quad (2.9.7)$$

Therefore, the martingale property of  $Y^{u^*}$ , implies that

$$\begin{aligned} C(s, X_s^{u^*}, u_s^*) + \partial_t v(s, X_s^{t,x,u^*}) + \mu(s, X_s^{t,x,u^*}, u_s^*) \cdot \nabla v(s, X_s^{t,x,u^*}) \\ + \frac{1}{2} a(s, X_s^{t,x,u^*}, u_s^*) \cdot D^2 v(s, X_s^{t,x,u^*}) = 0, \quad \text{a.s.} \end{aligned} \quad (2.9.8)$$

Therefore, for  $u^*$ , all the inequalities in (2.9.4)-(2.9) holds as equality and we obtain  $v = V$ .  $\square$

We apply Theorem (2.9.1) to the following examples.

**Exercise 24** (Merton optimal investment problem). Remember a self-financing portfolio with a Black-Scholes risky asset,  $dS_t = S_t(\mu dt + \sigma dB_t)$ , under zero interest rate satisfies

$$dX_t = \theta_t(\mu dt + \sigma dB_t) \quad (2.9.9)$$

where  $\theta$  is the amount of money invested in the risky asset. Merton problem is

$$\sup_{\theta} \mathbb{E}[U(X_T^\theta)] \quad (2.9.10)$$

The HJB is given by

$$\begin{cases} 0 = \partial_t v + \sup_{\theta} \left\{ \frac{\theta^2 \sigma^2}{2} \partial_x^2 v + \theta \mu \partial_x v \right\} \\ v(T, x) = U(x) \end{cases} \quad (2.9.11)$$

This simplifies to

$$\begin{cases} 0 = \partial_t v - \frac{(\mu \partial_x v)^2}{\sigma^2 \partial_x^2 v} \\ v(T, x) = U(x) \end{cases} \quad (2.9.12)$$

For  $U(x) = -e^{-\lambda x}$ , with  $\lambda > 0$ , use separation of variables  $u(t, x) = -f(t)e^{-\lambda x}$  to find a closed form solution for the HJB.

**Exercise 25.** For  $U(x) = x^\lambda$ , use separation of variables  $u(t, x) = f(t)x^\lambda$  to find a closed form solution for the HJB. Is there any restriction on the value of  $\lambda$ ?

**Exercise 26.** For  $U(x) = \ln x$ , can you suggest a separation of variables?

**Exercise 27** (Merton optimal consumption problem). It is similar to the Merton optimal investment problem except, the investor is consuming from the account and what matters is the utility of consumption.

$$dX_t = \theta_t(\mu dt + \sigma dB_t) - c_t dt \quad (2.9.13)$$

where  $\theta$  is the amount of money invested in the risky asset and  $c$  is the consumption with  $c_t \geq 0$ . Merton problem is

$$\sup_{\theta, c} \mathbb{E} \left[ \int_0^\infty e^{-\gamma t} U(c_t) dt \right], \gamma > 0 \quad (2.9.14)$$

Show that the HJB is given by

$$0 = \sup_{\theta} \left\{ \frac{\theta^2 \sigma^2}{2} v'' + \theta \mu v' \right\} + \sup_{c \geq 0} \{ U(c) - cv' \} - \gamma v \quad (2.9.15)$$

or

$$0 = \frac{(\mu v')^2}{\sigma^2 v''} + \sup_{c \geq 0} \{ U(c) - cv' \} - \gamma v \quad (2.9.16)$$

For  $U(c) = \frac{c^{1-\lambda}}{1-\lambda}$ , with  $\lambda \in (0, 1)$ , we have

$$\sup_{c \geq 0} \{ U(c) - cv' \} = U(c^*) - c^* v' \quad (2.9.17)$$

with  $U'(c^*) = v'$  or  $c^*$

### 2.9.1 Martingale approach

The martingale principle for optimal control (Davis and Varaiya [1973]) is another verification result which does not require the differentiability of the value function.

**Theorem 2.9.2.** *Assume that there exists a function  $v(t, x)$  such that for all  $u \in \mathcal{A}$  the process  $\{Y_t^u\}_{t \geq 0}$*

$$Y_t^u := v(t, X_t^u) + \int_0^t C(s, X_s^u, u_s) ds \quad (2.9.18)$$

*is a super martingale and that for some  $u^* \in \mathcal{A}$ ,  $\{Y_t^{u^*}\}_{t \geq 0}$  is a martingale. Then,  $u^*$  is an optimal control and the value function is equal to  $v(t, x)$ .*

*Proof.* By the supermartingale property, we have

$$Y_t^u \leq \mathbb{E}[Y_T^u | \mathcal{F}_t] = \mathbb{E}\left[g(X_T^u) + \int_t^T C(s, X_s^u, u_s) ds \middle| \mathcal{F}_t\right] + \int_0^t C(s, X_s^u, u_s) ds \quad (2.9.19)$$

Thus,

$$v(t, x) = Y_t^u - \int_0^t C(s, X_s^u, u_s) ds \leq \mathbb{E}[Y_T^u | \mathcal{F}_t] = \mathbb{E}\left[g(X_T^u) + \int_t^T C(s, X_s^u, u_s) ds \middle| \mathcal{F}_t\right], \quad (2.9.20)$$

for all  $u \in \mathcal{A}_{t,T}$ . Given  $X_t^u = x$ , we obtain

$$v(t, x) \leq \sup_{u \in \mathcal{A}_{t,T}} \mathbb{E}\left[g(X_T^u) + \int_t^T C(s, X_s^u, u_s) ds \middle| X_t^u = x\right] = V(t, x). \quad (2.9.21)$$

If for some  $u^*$ ,  $Y^{u^*}$  is a martingale, then in all the above, the inequality turns into equality and we have  $v(t, x) = V(t, x)$ .  $\square$

The above theorem can also be regarded as a verification that a possible candidate  $V(t, x)$  is a value function of the optimal control problem. The following example shows the use of this theorem.

**Example 15.** In Example 16, let  $U(c) = \frac{c^{1-\alpha}}{1-\alpha}$  when  $c \geq 0$  and  $-\infty$  otherwise. Consider the function  $V(t, x) = e^{-\gamma(T-t)}(T - t + e^{A/\alpha t})^\alpha e^{-AT} U(x)$  with  $A = r(1 - \alpha) - \gamma + \frac{(\mu-r)(1-\alpha)}{2\sigma^2\alpha}$ .

$$Y_t^u = V(t, X_t^u) + \int_0^t e^{-\gamma s} U(c_s) ds. \quad (2.9.22)$$

By Itô formula, we have

$$\begin{aligned} dY_t^u = & \left( \partial_t V + (\theta_t(\mu - r) + rX_t^u - c_t) \partial_x V + \frac{1}{2} \sigma^2 \theta_t^2 \partial_{xx} V + e^{-\gamma t} U(c_t) \right) dt \\ & + \sigma \theta_t \partial_x V dB_t. \end{aligned} \quad (2.9.23)$$

By direct calculation, one can see that

$$\partial_t V + (\theta_t(\mu - r) + rX_t^u - c_t) \partial_x V + \frac{1}{2} \sigma^2 \theta_t^2 \partial_{xx} V + e^{-\gamma t} U(c_t) \leq 0, \quad \mathbb{P}\text{-a.s.} \quad (2.9.24)$$

for all values of  $\theta$  and  $c$ . In addition, for  $c_t^*(x) = (\partial_x V(t, x))^{-1/\alpha}$  and  $\theta_t^*(x) = -\frac{(\mu-r)\partial_x V(t, x)}{\sigma^2 \partial_{xx} V(t, x)}$ ,

$$\partial_t V + (\theta_t(\mu - r) + rX_t^u - c_t)\partial_x V + \frac{1}{2}\sigma^2\theta_t^2\partial_{xx}V + e^{-\gamma t}U(c_t) = 0, \quad \mathbb{P}\text{-a.s.} \quad (2.9.25)$$

It remains to show that  $c^*$  and  $\theta^*$  are admissible Markov controls, i.e., to show that

$$dX_t^* = \left( \theta_t(X_t^*)((\mu - r)dt + \sigma dB_t) - rX_t^*dt \right) - c_t^*(X_t^*)dt, \quad (2.9.26)$$

has a strong solution. We leave the details as an exercise.

**Remark 2.9.1.** If the control problem is with infimum instead of supremum,

$$V(t, x) = \inf_{u \in \mathcal{A}_{t,T}} \mathbb{E} \left[ \int_t^T C(s, X_s^{t,x,u}, u_s)dt + g(X_T^{t,x,u}) \right], \quad (2.9.27)$$

Theorem 2.9.2 is should be modified. More precisely,  $Y_t^u$  is a submartingale.

For deterministic cases, supermartingale (resp. submartingale) means nonincreasing (resp. nondecreasing).

Finding a candidate for a value function and an optimal control is the subject of the future sections.

## 2.10 Some nonstandard HJB equations

In this section, we provide a review of some HJB equations that come from stochastic singular control, optimal stopping time, stochastic impulse control, and switching problems. The treatment of such problems via HJB equations is similar, however, we first need to derive HJB equations.

### 2.10.1 Stochastic singular control problems

Consider the stochastic control problem below:

**Example 16.**

$$\inf_{u_t} \mathbb{E} \left[ \int_0^T \left( \frac{1}{2}X_t^2 - u_t \right) dt + X_T^2 - X_T \right] \quad (2.10.1)$$

with  $u_t \geq 0$  (takes nonnegative values) and

$$dX_t = (X_t + u_t)dt + \sigma dB_t \quad (2.10.2)$$

The HJB is formally written as

$$\begin{cases} \partial_t V + x\partial_x V + \frac{\sigma^2}{2}\partial_{xx}^2 V + \frac{1}{2}x^2 + \inf_{u \geq 0} u(\partial_x V - 1) = 0 \\ V(T, x) = x^2 - x \end{cases} \quad (2.10.3)$$

Note that

$$\inf_{u \geq 0} u(\partial_x V - 1) = \begin{cases} -\infty & \partial_x V - 1 < 0 \\ 0 & \partial_x V - 1 \geq 0 \end{cases} \quad (2.10.4)$$

Therefore, it is natural to assume that the value function satisfies  $\partial_x V - 1 \geq 0$ . In fact, we expect so find the value functions such that if  $\partial_x V > 1$ , then  $\partial_t V + x\partial_x V + \frac{\sigma^2}{2}\partial_{xx}^2 V + \frac{1}{2}x^2 = 0$  and when  $\partial_x V = 1$ , the we only have  $\partial_t V + x\partial_x V + \frac{\sigma^2}{2}\partial_{xx}^2 V + \frac{1}{2}x^2 \leq 0$ . More precisely, inequality  $\partial_x V - 1 \geq 0$  divides  $[0, T] \times \mathbb{R}$  into two regions:



- $\mathbf{N} = \{(t, x) : \partial_x V(t, x) > 1\}$
- $\mathbf{R} = \{(t, x) : \partial_x V(t, x) = 1\}$

Inside  $\mathbf{N}$ , we expect that  $\partial_t V + x\partial_x V + \frac{\sigma^2}{2}\partial_x^2 V + \frac{1}{2}x^2 = 0$  holds. This equations has a solution of the form  $v(t, x) = a(t)x^2 + b(t)x + c(t)$ .

We denote the boundary of  $\mathbf{N}$  by  $\partial\mathbf{N}$  and we define  $R(t)$  such that  $(t, R(t)) \in \partial\mathbf{N}$ , assuming that  $R(t)$  is uniquely determined. In addition, we expect to see that the value function and its first derivative are the same in the both sides of  $\partial\mathbf{N}$ , particularly,  $\partial_x V(t, R(t)) = 1$ . Therefore, if  $v$  is the solution to

$$\begin{cases} \partial_t v + x\partial_x v + \frac{\sigma^2}{2}\partial_x^2 v + \frac{1}{2}x^2 = 0 \\ v(T, x) = x^2 - x \end{cases} \quad (2.10.5)$$

then, we anticipate to write

$$V(t, x) = \begin{cases} v(t, x) & x \geq R(t) \\ x - R(t) + v(t, R(t)) & x < R(t) \end{cases} \quad (2.10.6)$$

Solving (2.10.5), we obtain that  $v(t, x) = \frac{5e^{2(T-t)}-1}{4}x^2 - e^{T-t}x - \frac{\sigma^2}{4}\left(T-t - \frac{5e^{2(T-t)}-1}{2}\right)$  and, therefore,  $R(t) = 2\frac{e^{T-t}+1}{5e^{2(T-t)}-1}$ .

**Exercise 28.**

$$\inf_{u_t \geq 0} \mathbb{E} \left[ \int_0^\tau e^{-rt} (\mu - u_t) dt \right] \quad (2.10.7)$$

where  $\tau = \inf\{s \geq 0 : X_t = 0\}$  with

$$dX_t = (\gamma X_t - u_t)dt + \sigma dB_t, \quad X_0 = x \geq 0 \quad (2.10.8)$$

Write the HJB and obtain the nonlinear term by evaluating the infimum.

Singular control problems are easy to detect; the running cost and drift of the SDE are both linear in the control. To see this, consider the control problem

$$\inf_u \mathbb{E} \left[ \int_0^T (C(s, X_s^u) + au_s) ds + g(X_T^u) \right] \quad (2.10.9)$$

where  $u_t$  takes values in a closed convex cone  $\mathcal{C}$  and

$$\begin{cases} dX_t^u = (\mu(t, X_t^u) + Au_t)dt + \sigma(t, X_t^u)dB_t \\ X_0^u = x \in \mathbb{R}^d \end{cases} \quad (2.10.10)$$

The heuristic derivation of the HJB yields

$$\begin{cases} 0 = \partial_t V(t, x) + \frac{1}{2}a(t, x, u) \cdot D^2 V(t, x) + \nabla V(t, x) \cdot \mu(t, x) + C(t, x) + \inf_u \{au + Au \cdot \nabla V(t, x)\} \\ V(T, x) = g(x) \end{cases} \quad (2.10.11)$$

$a = \sigma^\top \sigma$ . The infimum above is given by

$$\inf_{\hat{u} \in \mathcal{C}, |\hat{u}|=1, \lambda \geq 0} \lambda \hat{u} \cdot (a + A \nabla V(t, x)) = \inf_{\lambda \geq 0} \lambda \inf_{|\hat{u}|=1} \hat{u} \cdot (a + A \nabla V(t, x)) = \inf_{\lambda \geq 0} \lambda \mathcal{H}(\nabla V(t, x)) \quad (2.10.12)$$

where  $\mathcal{H}(\nabla V(t, x)) := \inf_{|\hat{u}|=1} \hat{u} \cdot (a + A\nabla V(t, x))$ . If  $\mathcal{H}(\nabla V(t, x)) < 0$ , then the infimum is  $-\infty$  and the problem becomes degenerate. To avoid degeneracy, we must have  $\mathcal{H}(\nabla V(t, x)) \geq 0$ , in which case infimum is attained in  $\lambda = 0$  and we have

$$\begin{cases} 0 = \partial_t V(t, x) + \frac{1}{2}a(t, x, u) \cdot D^2 V(t, x) + \nabla V(t, x) \cdot \mu(t, x) + C(t, x) \\ V(T, x) = g(x) \end{cases} \quad (2.10.13)$$

For simplicity of notation, we set

$$\mathcal{L}V(t, x) := \partial_t V(t, x) + \frac{1}{2}a(t, x, u) \cdot D^2 V(t, x) + \nabla V(t, x) \cdot \mu(t, x) + C(t, x) \quad (2.10.14)$$

The interpretation of the HJB requires some *variational inequalities*:<sup>6</sup>

1. for all  $(t, x)$ ,  $-\mathcal{L}V(t, x) \geq 0$  and  $\mathcal{H}\nabla V(t, x) \geq 0$ .
2. If  $\mathcal{H}(\nabla V(t, x)) > 0$ , then  $-\mathcal{L}V(t, x) = 0$ .

The description of the optimal control in singular control problems need the notion of local time of SDEs. For instance, in Exercise 28, the optimal control is the local time of the process,  $dX_t = \gamma X_t dt + \sigma dB_t$  at the point  $\hat{x}$ , existence of which is obtained through solving the HJB and going through verification step.

### 2.10.2 Optimal stopping problem

Let  $X$  be given by the SDE

$$dX_t = \mu(t, X_t)dt + \sigma(t, X_t)dB_t \quad (2.10.15)$$

in a filtered probability space and consider the following problem

$$\sup \mathbb{E} \left[ \int_0^{\tau \wedge T} e^{-rt} C(t, X_t) dt + g(X_{\tau \wedge T}) \right] \quad (2.10.16)$$

where the supremum is over all stopping time  $\tau$  adapted to the filtration. These class of problems, which are not stochastic control problems, called optimal stopping problems. To solve optimal stopping problems, we can write the following DPP:

$$V(t, x) = \mathbb{E}_{t,x} \left[ \int_0^\eta e^{-r(s-t)} C(t, X_s) ds + e^{-r(\eta-t)} g(X_\eta) \right] \quad (2.10.17)$$

for any stopping time  $\eta$  with values in  $[t, T]$ . The DPP above leads to the following variational HJB equation:

1.  $V(t, x) \geq g(x)$  and  $-\mathcal{L}V(t, x) \geq 0$
2. if  $V(t, x) > g(x)$ , then  $-\mathcal{L}V(t, x) = 0$

where  $\mathcal{L}$  is given in (2.10.14).

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<sup>6</sup>For more details on why we write the variational inequality this way, see [Fleming and Soner, 2006, Chapter 11] and the notion of sub and super solutions.

## 2.11 Backward stochastic differential equations

The backward stochastic differential equations (BSDE) is an alternative way to tackle the optimal control problems. One advantage of the BSDEs is that they can cover the nonMarkovian optimal control problems. The derivation of the BSDEs from an optimal control problem follows from the Itô martingale representation theorem.

**Theorem 2.11.1.** *In a probability space which hosts a Brownian motion  $B$ , let  $\mathbb{F}^B := \{\mathcal{F}_t^B\}_t$  be the right-continuous augmented filtration generated by a Brownian motion  $B$ . Let  $X$  be a  $\mathcal{F}_T^B$ -measurable random variable with finite expectation. Then, there exists a  $\mathbb{F}^B$ -progressively measurable process  $Z_t(\omega)$  such that*

$$X = \mathbb{E}[X] + \int_0^T Z_t dB_t \quad (2.11.1)$$

*In addition, if  $M$  is a continuous martingale with respect to  $\mathbb{F}$ , then, there exists a  $\mathbb{F}^B$ -progressively measurable process  $\phi(\omega)$  such that*

$$M_t = M_0 + \int_0^t Z_s dB_s, \quad \text{for } t \geq 0. \quad (2.11.2)$$

In the above theorem,  $Z$  can interpreted as the sensitivity of the martingale  $M$  with respect to the Brownian noise  $B$ .

Consider for a given stochastic process  $C_t(\omega)$

$$Y_t := \mathbb{E} \left[ \int_t^T L_s(\omega) ds + \xi \middle| \mathcal{F}_t^B \right]. \quad (2.11.3)$$

We simply assume that  $r_s(\omega)$  and  $L_s(\omega)$  are  $\mathbb{F}^B$ -progressively measurable processes and  $\xi$  is a  $\mathcal{F}_T^X$ -measurable and integrable random variable. For the martingale define by

$$M_t := \mathbb{E} \left[ \int_0^T L_s(\omega) ds + \xi \middle| \mathcal{F}_t^B \right], \quad (2.11.4)$$

and by the Itô martingale representation theorem, we have

$$M_T = \int_0^T L_s(\omega) ds + \xi = M_t + \int_t^T Z_s dB_s, \quad (2.11.5)$$

for some  $\mathbb{F}^B$ -progressively measurable process  $Z_t(\omega)$ . Since  $Y_T = \xi$ , we have

$$M_T = \int_0^T L_s(\omega) ds + Y_T = M_t + \int_t^T Z_s dB_s, \quad (2.11.6)$$

On the other hand, since  $\int_0^t L_s(\omega) ds$  is  $\mathcal{F}_t^B$ -measurable, we can write

$$M_t = \int_0^t L_s(\omega) ds + \mathbb{E} \left[ \int_t^T L_s(\omega) ds + \xi \middle| \mathcal{F}_t^B \right] = \int_0^t L_s(\omega) ds + Y_t. \quad (2.11.7)$$

Therefore,

$$\int_0^T L_s(\omega) ds + Y_T = \int_0^t L_s(\omega) ds + Y_t + \int_t^T Z_s dB_s, \quad (2.11.8)$$

or,

$$Y_t = Y_T + \int_t^T L_s(\omega) ds - \int_t^T Z_s dB_s. \quad (2.11.9)$$

The BSDE (2.11.9) can be written formally by

$$\begin{cases} dY_t = -L(t, \omega)dt + Z_t dB_t \\ Y_T = g(X_T) \end{cases} \quad (2.11.10)$$

Note that we could have written the above forwardly, i.e.,  $Y_t = Y_0 + \int_0^t L_s(\omega) ds - \int_0^t Z_s dB_s$ . However, this is not very useful, because  $Y_0$  is not known. See example below.

**Example 17.** Recall that the solution to the linear equation

$$\begin{cases} 0 = \partial_t v(t, x) + C(t, x) + [\mu \cdot \nabla v](t, x) + \frac{1}{2}[\sigma^\top \sigma \cdot D^2 v](t, x) \\ v(T, x) = g(x) \end{cases} \quad (2.11.11)$$

is given by the Feynmann-Kac formula:

$$V(t, x) = \mathbb{E} \left[ \int_t^T C(s, X_s) ds + g(X_T) \middle| X_t = x \right], \quad (2.11.12)$$

where

$$dX_t = \mu(t, X_t)dt + \sigma(t, X_t)dB_t. \quad (2.11.13)$$

If we define  $Y_t = V(t, X_t)$ , then (2.11.9) is given by

$$Y_t = g(X_T) + \int_t^T C(s, X_s) ds - \int_t^T Z_s dB_s. \quad (2.11.14)$$

In the above,  $g$  and  $L$  are known. Given that we can find  $Z$ ,  $Y_t$  is known. However, if we write the equation forward

$$Y_t = Y_0 + \int_0^t C(s, X_s) ds - \int_0^t Z_s dB_s. \quad (2.11.15)$$

Here,  $Y_0 = V(0, X_0)$  is not known.

The following theorem shows that a BSDE can generalize the Feynmann-Kac formula to nonlinear equations.

**Theorem 2.11.2.** Assume that the following semilinear<sup>7</sup> PDE has a solution  $V(t, x) \in C^{1,2}$ .

$$\begin{cases} 0 = \partial_t v(t, x) + [\mu \cdot \nabla v](t, x) + \frac{1}{2}[\sigma^\top \sigma \cdot D^2 V](t, x) + C(t, x, v(t, x), \nabla v(t, x)) \\ V(T, x) = g(x). \end{cases} \quad (2.11.16)$$

Then,  $Y_t = V(t, X_t)$  and  $Z_t = \sigma(t, X_t) \nabla V(t, X_t)$  satisfy the BSDE

$$\begin{cases} dY_t = -C(t, X_t, Y_t, \sigma^{-1}(t, X_t)Z_t)dt + Z_t dB_t \\ Y_T = g(X_T) \end{cases} \quad (2.11.17)$$

<sup>7</sup>The an equation with linear second order term and possibly nonlinear first order term is called semilinear.

where

$$dX_t = \mu(t, X_t)dt + \sigma(t, X_t)dB_t. \quad (2.11.18)$$

*Proof.* By applying Itô lemma on  $Y_t = V(t, X_t)$ , we obtain

$$\begin{aligned} dY_t &= (\partial_t V(t, X_t) + \mu(t, X_t)\nabla V(t, X_t) + \frac{1}{2}\sigma^2(t, X_t)\partial_{xx}V(t, X_t))dt + \sigma(t, X_t)\nabla V(t, X_t)dB_t \\ &= -C(t, X_t, V(t, X_t), \nabla V(t, X_t))dt + \sigma(t, X_t)\nabla V(t, X_t)dB_t \\ &= -C(t, X_t, Y_t, \sigma^{-1}(t, X_t)Z_t)dt + Z_tdB_t \end{aligned} \quad (2.11.19)$$

□

In the above, the second equality is obtained from the PDE and the third equality is from the definition of  $Y$  and  $Z$ . In addition, by the terminal condition we have,  $Y_T = V(T, X_T) = g(X_T)$ . In general, the use of BSDEs over the HJB is preferable when the regularity of the function  $V(t, x)$  is not established. Because by the existence theorem for the BSDEs, we already know that  $Z_t$  exists as stochastic process even if  $\sigma(t, X_t)\nabla V(t, X_t)$  does not make sense in cases when  $\nabla V$  does not exist as specific points.

Motivated by Theorem 2.11.2, we can define a general BSDE as

$$\begin{cases} dY_t = -C(t, Y_t, Z_t, \omega)dt + Z_sdB_s \\ Y_T = \xi \end{cases}, \quad (2.11.20)$$

The argument  $\omega$  inside  $L$  represents a general dependence on the randomness. In particular, it can represent solution  $X_t$  of a possible path-dependent SDE

$$dX_t = \mu(t, X_t)dt + \sigma(t, X_t)dB_t, \quad (2.11.21)$$

a possible control process  $u_t$ , or both  $X_t$  and  $u_t$ . We discuss such dependencies in further details in Section 2.11.3.

**Remark 2.11.1.** In (2.11.20), we blended  $\sigma^{-1}$  inside the Lagrangian  $L(s, X_s, Y_s, \sigma^{-1}Z_s, \omega)$  into it and simply write  $L(s, X_s, Y_s, Z_s, \omega)$ .

**Theorem 2.11.3** (Existence and uniqueness theorem for BSDEs). *Assume that*

- i)  $\mu(t, x)$  and  $\sigma(t, x)$  are Lipschitz in  $t$  and  $x$
- ii) for all  $(t, x)$  with  $0 < \lambda|x|^2 \leq x^\top \sigma^\top \sigma x$  for all  $x$  and for some  $\lambda$  (in particular  $\sigma$  is invertible)
- iii)  $L(t, x, \varrho, \Pi, \omega)$  is progressively measurable in  $\omega$  and Lipschitz in other variables  $(t, x, \varrho, \Pi)$  and is decreasing in  $\varrho$  a.s.

Then, for any  $\mathcal{F}_T$ -measurable square-integrable random variable  $\xi$ , there exists a couple  $(Y_t, Z_t)$  such that (2.11.20) holds, i.e.,

$$Y_t = \xi + \int_t^T L(s, X_s, Y_s, Z_s, \omega)ds - Z_sdB_s, \quad (2.11.22)$$

and there exists a constant  $C$  that only depends on the Lipschitz constant of  $\mu$ ,  $\sigma$ , and  $L$  and the constant  $\lambda$  that

$$\mathbb{E} \left[ \sup_{t \in [0, T]} |Y_t|^2 + \int_0^T Z_s^2 ds \right] \leq C (\mathbb{E}[\xi^2]) \quad (2.11.23)$$

### 2.11.1 Linear BSDEs

Linear BSDEs take the form

$$dY_t = -(\alpha_t Y_t + \beta_t Z_t + L_t)dt + Z_t dB_t, \quad (2.11.24)$$

where  $\alpha$ ,  $\beta$ , and  $L$  are arbitrary progressively measurable processes. One can write a closed-form solution for a linear BSDE. To do so, note that

$$dY_t = -(\alpha_t Y_t + L_t)dt + Z_t dB_t^\beta, \quad (2.11.25)$$

where by the Girsanov theorem,  $dB_t^\beta := \beta_t dt + dB_t$  is a Brownian motion under the probability  $\mathbb{P}^\beta$  given in

$$\frac{d\mathbb{P}^\beta}{d\mathbb{P}} \Big|_{\mathcal{F}_t} := \exp \left( \int_0^t \beta_s dB_s - \frac{1}{2} \int_0^t \beta_s^2 ds \right) \quad (2.11.26)$$

If we define  $\tilde{Y}_t := e^{\int_0^t \alpha_s ds} Y_t$ , then

$$d\tilde{Y}_t = e^{\int_0^t \alpha_s ds} (L_t dt + Z_t dB_t^\beta). \quad (2.11.27)$$

In other words,

$$\tilde{Y}_t = \tilde{Y}_T - \int_t^T e^{\int_t^s \alpha_z dz} (L_s dt + Z_s dB_s^\beta), \quad (2.11.28)$$

or,

$$Y_t = e^{\int_t^T \alpha_s ds} \tilde{Y}_T - \int_t^T e^{\int_t^s \alpha_z dz} (L_s dt + Z_s dB_s^\beta) = e^{\int_t^T \alpha_s ds} \xi - \int_t^T e^{\int_t^s \alpha_z dz} (L_s dt + Z_s dB_s^\beta). \quad (2.11.29)$$

After taking conditional expectation with respect to  $\mathbb{P}^\beta$ , we obtain

$$Y_t = \mathbb{E}^\beta \left[ e^{\int_t^T \alpha_s ds} \xi - \int_t^T e^{\int_t^s \alpha_z dz} L_s dt \Big| \mathcal{F}_t \right]. \quad (2.11.30)$$

If we define

$$\Gamma_t = e^{\int_t^T \alpha_s ds} \frac{d\mathbb{P}^\beta}{d\mathbb{P}} \Big|_{\mathcal{F}_t} = \exp \left( \int_0^t \beta_s dB_s + \int_0^t (\alpha_s - \frac{1}{2} \beta_s^2) ds \right), \quad (2.11.31)$$

by changing the measure back to  $\mathbb{P}$

$$Y_t = \Gamma_t^{-1} \mathbb{E} \left[ \Gamma_T \xi - \int_t^T \Gamma_s L_s dt \Big| \mathcal{F}_t \right]. \quad (2.11.32)$$

### 2.11.2 Comparison principle

Consider

$$\begin{cases} dY_t^{(i)} = -L^{(i)}(t, Y_t^{(i)}, Z_t^{(i)}, \omega)dt + Z_s^{(i)} dB_s \\ Y_T^{(i)} = \xi^{(i)} \end{cases}, \quad (2.11.33)$$

for  $i = 1, 2$ . The following theorem provides a sufficient condition for comparing  $Y^{(1)}$  and  $Y^{(2)}$ , and hence, it is called comparison principle.

**Theorem 2.11.4.** Assume that for  $i = 1, 2$ ,  $(Y^{(i)}, Z^{(i)})$  is the solution for (2.11.33) and  $L(t, y, z, \omega)$  is

Lipschitz in  $(y, z)$  uniformly in  $\omega$  and  $t$ . Further assume that  $\xi^{(1)} \geq \xi^{(2)}$  and, for each value  $(t, y, z)$ ,  $L^{(1)}(t, y, z, \omega) \leq L^{(2)}(t, y, z, \omega)$ , a.s. Then,  $Y^{(1)} \geq Y^{(2)}$ , a.s.

Proof. □

### 2.11.3 Maximum principle

For the purpose of control theory, we assume that  $L$  takes the form  $L(t, x, \varrho, \Pi, u_t)$ , where  $u_t$  is a progressively measurable process that represents a control at time  $t$ . In this case, we denote the solution to the BSDE by  $(Y^u, Z^u)$ , where

$$\begin{cases} Y_t^u = \xi + \int_t^T F(s, X_s, Y_s^u, Z_s^u, u_s) ds - \int_t^T Z_s^u dB_s \\ Y_T = \xi \in \mathcal{F}_T^X \\ dX_t = \mu(t, X_t)dt + \sigma(t, X_t)dB_t \end{cases} \quad (2.11.34)$$

and  $F(s, X_s^u, Y_s^u, Z_s^u, u_s) = L(s, X_s, Y_s, \sigma^{-1}(s, X_s^u, u_s)Z_s, u_s)$ . Then, a control problem can be written as

$$Y_t = \text{esssup}_{u \in \mathcal{A}_{t,T}} Y_t^u. \quad (2.11.35)$$

Note that the set of admissible controls are defined specific to a particular problem. In addition, if the terminal condition  $\xi = g(X_t^u)$ , the problem is a Markovian control problem that was studied in previous sections. Otherwise, it is not Markovian, and therefore, the value cannot be written as a function  $V(t, x)$ . In fact, value function takes a more general form of

$$Y_t = \text{esssup}_{u \in \mathcal{A}_{t,T}} \mathbb{E} \left[ \int_t^T F(s, X_s^u, Y_s^u, Z_s^u, u_s) ds + \xi \middle| \mathcal{F}_t^X \right]. \quad (2.11.36)$$

In the above, the Lagrangian  $F$  not only depends on the state of the system  $X_s^u$  and control  $u_s$ , but also depends on the history of the value  $Y_s^u$  and the sensitivity of the value with respect to the Brownian noise  $Z_t^u$ .

To solve such optimal control problems, we require the following comparison principle.

**Theorem 2.11.5** (Maximum principle for BSDEs). *The process*

$$Y_t = \text{esssup}_{u \in \mathcal{A}_{t,T}} \mathbb{E} \left[ \int_t^T F(s, X_s^u, Y_s^u, Z_s^u, u_s) ds + \xi \middle| \mathcal{F}_t^X \right] = \text{esssup}_{u \in \mathcal{A}_{t,T}} Y_t^u \quad (2.11.37)$$

there exists a  $Z$  such that  $(Y, Z)$  satisfies the BSDE

$$\begin{cases} Y_t = \xi + \int_t^T F^*(s, X_s, Y_s, Z_s) ds - \int_t^T Z_s dB_s \\ Y_T = \xi \in \mathcal{F}_T^X \\ dX_t = \mu(t, X_t)dt + \sigma(t, X_t)dB_t \end{cases} \quad (2.11.38)$$

where

$$F^*(s, x, y, z) := \sup_{u \in U} F(s, x, y, z, u), \quad \text{and} \quad u^*(t, x, y, z) := \arg\max_{u \in U} F(t, x, y, z, u) \quad (2.11.39)$$

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## Chapter 3

# Numerical evaluation of stochastic control problems

### 3.1 DPP based approximation

Recall from (2.5.13), the value function and the optimal control in discrete stochastic control problem (2.5.4) can be evaluated through.

$$\begin{cases} V(t, x) = \inf_u C(t, x, u) + \mathbb{E}[V(t+1, X_{t+1}^u) | X_t = x] \\ V(T, x) = g(x) \\ u_t^*(x) \in A(x) := \operatorname{argmax}_u C(t, x, u) + \mathbb{E}[V(t+1, X_{t+1}^u) | X_t = x] \end{cases} \quad (3.1.1)$$

where

$$X_{t+1}^u = x + \mu(t, x, u) + \sigma(t, x, u)\xi_{t+1} \quad (3.1.2)$$

and  $\{\xi_t\}_{t=1}^T$  is a sequence of i.i.d. random variables with mean 0 and variance 1. In the above, we need to evaluate (1) conditional expectation  $\mathbb{E}[\cdot | X_t = x]$  and (2) the infimum over control  $u$ . Given  $V(t+1, \cdot)$ , these calculations can be done separately. However, the separate evaluation creates inefficiency in the calculations. Therefore, we propose the one-shot approximation of  $\inf_u C(t, x, u) + \mathbb{E}[V(t+1, X_{t+1}^u) | X_t = x]$  through the following methods.

#### 3.1.1 Evaluation of infimum

First note that

$$\begin{aligned} C(t, x, u(x)) + \mathbb{E}[V(t+1, X_{t+1}^u) | X_t = x] \\ = \mathbb{E}[C(t, X, u(X)) + V(t+1, X + \mu(t, X, u(X)) + \sigma(t, X, u(X))\xi_{t+1}) | X = x] \end{aligned} \quad (3.1.3)$$

Denoting the right-hand side above by  $\Phi(x, u)$ , then we seek a function  $\hat{u}^*(x)$  such that

$$\hat{u}^*(\cdot) \in A(X) := \operatorname{argmin}_{u(\cdot)} \mathbb{E}[\Phi(X, u(X))] \quad (3.1.4)$$

The following Lemma guarantees that  $\hat{u}^*$  solves (3.1.1).

**Lemma 3.1.1.** *Assume that  $\hat{u}^* \in A(X)$  defined above. Then,  $\hat{u}^*(x) \in A(x)$  for all  $x$  in the set of values of  $X$ .*



*Proof.* For a complete proof, we need a measurable selection theorem and some other conditions. For simplicity, we only provide the sketch of the proof. Assume that  $u^*(x) \notin A(x)$  with a positive probability on the set of values of  $X$ . We denote the set by  $B$ . Then,  $x \in B$ , we define  $\tilde{u}^*(x)$  such that

$$\tilde{u}^*(x) \in \operatorname{argmax}_u C(t, x, u) + \mathbb{E}[V(t+1, X_{t+1}^u) | X_t = x]$$

and  $\tilde{u}^*(x) = \hat{u}^*(x)$ . Therefore,

$$\mathbb{E}[\Phi(X, \tilde{u}^*(X))] < \mathbb{E}[\Phi(X, \hat{u}^*(X))]$$

which contradicts the definition of  $\hat{u}^*$ .  $\square$

From the nonparametric point of view, we can approximate  $\hat{u}^*(\cdot)$  in (3.1.4) by a nonparametric model  $u(\cdot; \theta)$  via the following optimization

$$\theta^* \in \operatorname{argmin}_{\theta} \mathbb{E}[\Phi(X, u(X; \theta))] \quad (3.1.5)$$

For instance,  $u(\cdot; \theta)$  can be a neural network with parameter  $\theta$ .

## 3.2 Multi-step evaluation

While (3.1.1) is a classic way to evaluate optimal control problems, it is not very efficient. There are two reasons for the lack of efficiency. First, we have to run a loop over the number of time steps. Second, after evaluation of each nonparametric, we need to evaluate the value function. The total number of parameters is the number of parameters in each step times the number of steps, which for nonparametric models is massive and potentially needs a lot of memory. Therefore, we propose the following optimization in place of (3.1.1).

$$\inf_{u(\cdot, \cdot)} \mathbb{E} \left[ \sum_{t=0}^{T-1} C(t, X_t^u, u(t, X_t^u)) + g(X_T^u) \right] \quad (3.2.1)$$

where

$$\begin{cases} X_{t+1}^u = X_t^u + \mu(t, X_t^u, u(t, X_t^u)) + \sigma(t, X_t^u, u(t, X_t^u))\xi_{t+1} \\ X_0 \text{ is a random variable.} \end{cases} \quad (3.2.2)$$

The justification for the above problem is the same as Lemma 3.1.1 for one-step DPP method. Note that we can use a nonparametric model,  $u(t, x; \theta)$ , to approximate the minimization problem by

$$\theta^* \in \operatorname{argmin}_{\theta} \mathbb{E} \left[ \sum_{t=0}^{T-1} C(t, X_t^u, u(t, X_t^u; \theta)) + g(X_T^u) \right] \quad (3.2.3)$$

with

$$\begin{cases} X_{t+1}^u = X_t^u + \mu(t, X_t^u, u(t, X_t^u; \theta)) + \sigma(t, X_t^u, u(t, X_t^u; \theta))\xi_{t+1} \\ X_0 \text{ is a random variable.} \end{cases} \quad (3.2.4)$$

When an approximate optimal strategy,  $u(t, x; \theta^*)$ , is found, then one can find the value function through evaluating

$$\mathbb{E}_{t,x} \left[ \sum_{s=t}^{T-1} C(s, X_s^u, u(s, X_s^u; \theta^*)) + g(X_T^u) \right] \quad (3.2.5)$$

with

$$\begin{cases} X_{s+1}^u = X_s^u + \mu(s, X_s^u, u(s, X_s^u; \theta^*)) + \sigma(s, X_s^u, u(s, X_s^u; \theta^*))\xi_{s+1} \\ X_t = x \end{cases} \quad (3.2.6)$$

Of course, the evaluation of the conditional expectation above is a different problem.

### 3.3 Numerical methods based on BSDEs

Recall from Section 2.11 that the solution  $V(t, x)$  to the semilinear equation

$$\begin{cases} 0 = \partial_t v(t, x) + [\mu \cdot \nabla v](t, x) + \frac{1}{2}[\sigma^\top \sigma \cdot D^2 V](t, x) + C(t, x, v(t, x), \nabla v(t, x)) \\ V(T, x) = g(x). \end{cases} \quad (3.3.1)$$

is related to the BSDE

$$\begin{cases} dY_t = -C(t, X_t, Y_t, \sigma^{-1}(t, X_t)Z_t)dt + Z_t dB_t \\ Y_T = g(X_T) \end{cases} \quad (3.3.2)$$

by

$$Y_t = V(t, X_t) \text{ and } Z_t = [\sigma \nabla V](t, X_t) \quad (3.3.3)$$

where

$$dX_t = \mu(t, X_t)dt + \sigma(t, X_t)dB_t. \quad (3.3.4)$$

One way to interpret the solution to the PDE is to find functions  $\mathbf{Y}(t, x)$  and  $\mathbf{Z}(t, x)$  such that

$$\mathbf{Y}(t, X_t) = g(X_T) - \int_t^T \left( C(s, X_s, \mathbf{Y}(s, X_s), \mathbf{Z}(s, X_s))ds + \mathbf{Z}(s, X_s)dB_s \right) \quad (3.3.5)$$

Functions  $\mathbf{Y}$  and  $\mathbf{Z}$  can be approximated by neural networks  $\hat{\mathbf{Y}}(t, x, \alpha)$  and  $\hat{\mathbf{Z}}(t, x; \beta)$  and the equations (3.3.4) and (3.3.5) can be approximated discretely by

$$\begin{cases} \hat{\mathbf{Y}}(t_n, \hat{X}_{t_n}) = \hat{\mathbf{Y}}(t_{n+1}, \hat{X}_{t_{n+1}}) - C(t_n, \hat{X}_{t_n}, \hat{\mathbf{Y}}(t_n, \hat{X}_{t_n}), \hat{\mathbf{Z}}(t_n, \hat{X}_{t_n}))\Delta t + \hat{\mathbf{Z}}(t_n, \hat{X}_{t_n})\Delta B_{t_{n+1}} \\ \hat{X}_{t_{n+1}} = \hat{X}_{t_n} + \mu(t_n, \hat{X}_{t_n})\Delta t + \sigma(t_n, \hat{X}_{t_n})\Delta B_{t_{n+1}}. \end{cases} \quad (3.3.6)$$

Specifically, we shall determine  $\hat{\mathbf{Y}}$  and  $\hat{\mathbf{Z}}$  such that  $g(\hat{X}_T)$  is as close to  $\hat{\mathbf{Y}}(T, \hat{X}_T)$  as possible. In other words, we can find the approximate solution  $\hat{\mathbf{Y}}(t, x, \alpha^*)$  and  $\hat{\mathbf{Z}}(t, x; \beta^*)$  with

$$(\alpha^*, \beta^*) \in \underset{\alpha, \beta}{\operatorname{argmin}} \mathbb{E} \left[ (\hat{\mathbf{Y}}(T, \hat{X}_T) - g(\hat{X}_T))^2 \right] \quad (3.3.7)$$

## .1 Results from optimization

In this sections, we quickly review two of the most influential results from optimization, namely the Lagrange multiplier and gradient descent, which later be used in the context of control problems.

### .1.1 Lagrange multiplier

If we add a constraint to a simple minimization problem such as  $\min_x f(x)$ , the Lagrange multiplier method is the way to proceed. In a nutshell, the Lagrange multiplier method turns a constrained optimization problem into a saddle point problem without constraints by adding more variables.

Consider the constrained problem below:

$$\inf_x f(x) \quad \text{subject to} \quad g(x) = 0 \quad (.1.1)$$

Define the Lagrangian:

$$L(x, \lambda) := f(x) - \lambda \cdot g(x) \quad (.1.2)$$

Then, under proper conditions, the following saddle point problem yields the solution to (.1.1).

$$\sup_{\lambda} \inf_x L(x, \lambda) \quad (.1.3)$$

The function  $H(\lambda) = \inf_x L(x, \lambda)$  is called Hamiltonian and the Lagrange multiplier  $\lambda$  is called the dual variable. The unconstrained problem

$$\sup_{\lambda} H(\lambda) \quad (.1.4)$$

is called the dual problem for the primal problem (.1.1).

The key to the success of the Lagrange multiplier method is the strong duality.

$$\sup_{\lambda} \inf_x L(x, \lambda) = \inf_x \sup_{\lambda} L(x, \lambda) \quad (.1.5)$$

With strong duality, if for some  $x$ ,  $g(x) \neq (0)$ , then

$$\sup_{\lambda} L(x, \lambda) = \sup_{\lambda} \{f(x) - \lambda \cdot g(x)\} = \infty \quad (.1.6)$$

Therefore,  $\inf_x \sup_{\lambda} L(x, \lambda)$  is restricted to  $x$  with  $g(x) = 0$  and

$$\inf_{x: g(x)=0} \sup_{\lambda} L(x, \lambda) = \inf_{x: g(x)=0} f(x) \quad (.1.7)$$

It is not always easy to check if strong duality holds. However, Karush-Kuhn-Tucker conditions (KKT) provide some necessary and sufficient conditions for strong duality.

**Theorem .1.1.** Assume the differentiability of  $f$  and  $g$ . If  $x^*$  solves (.1.1) and  $\lambda^*$  solves (.1.4) such that

$$\begin{cases} \nabla f(x^*) - \lambda^* \cdot \nabla g(x^*) = 0 \\ g(x^*) = 0 \end{cases} \quad (.1.8)$$

then, strong duality holds and  $(x^*, \lambda^*)$  is a saddle point for  $\sup_{\lambda} \inf_x L(x, \lambda)$ .

Conversely, if strong duality holds, then any saddle point  $(x^*, \lambda^*)$  satisfies (.1.8). In particular,  $x^*$  solve (.1.1).

A geometric interpretation of  $\lambda^*$  in KKT conditions is explained in See Figure .1.1.

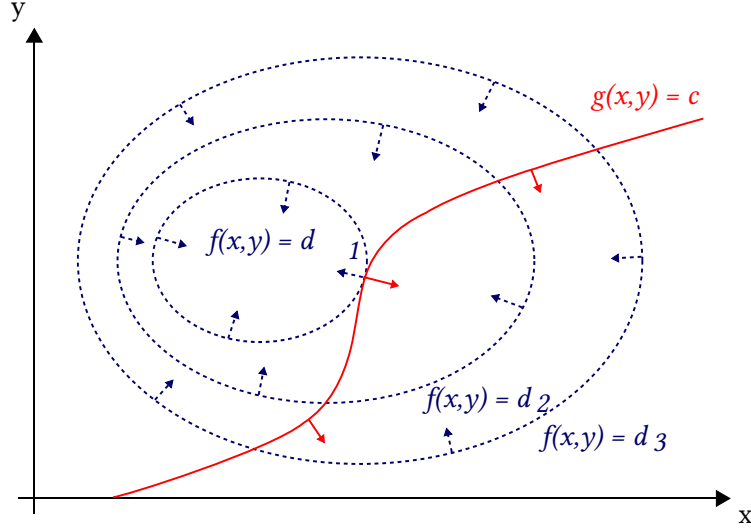
When the constraint is given by some inequalities and equalities, i.e.,

$$\inf_x f(x) \quad \text{subject to} \quad g(x) = 0, \quad h(x) \geq 0 \quad (.1.9)$$

the range of dual variable in Hamiltonian changes:

$$L(x, \lambda, \mu) := f(x) - \lambda \cdot g(x) - \mu \cdot h(x) \quad (.1.10)$$

$$\sup_{\mu \geq 0} \sup_{\lambda} \inf_x L(x, \lambda, \mu) \quad (.1.11)$$



**Figure .1.1:** KKT conditions for  $g(x) = 0$ . Source WIKIPEDIA

The reason for such modification can formally be seen after switching  $\inf_x$  and  $\sup_{\mu \geq 0}$  in the above

$$\sup_{\mu \geq 0} \inf_x L(x, \lambda, \mu) = \inf_x \sup_{\mu \geq 0} L(x, \lambda, \mu) \quad (.1.12)$$

If  $x_1$  is such that  $h(x_1) < 0$ , then

$$\sup_{\mu \geq 0} L(x_1, \lambda, \mu) = f(x_1) - \lambda g(x_1) - h(x_1) \sup_{\mu \geq 0} \mu = \infty$$

Therefore,  $\inf_x \sup_{\mu \geq 0} L(x, \lambda, \mu)$  is not attained at  $x_1$ . Thus, any saddle point  $(x^*, \lambda^*, \mu^*)$  with  $\mu^* > 0$  must satisfy  $h(x^*) \geq 0$ . More general KKT conditions guarantee the strong duality in this case:

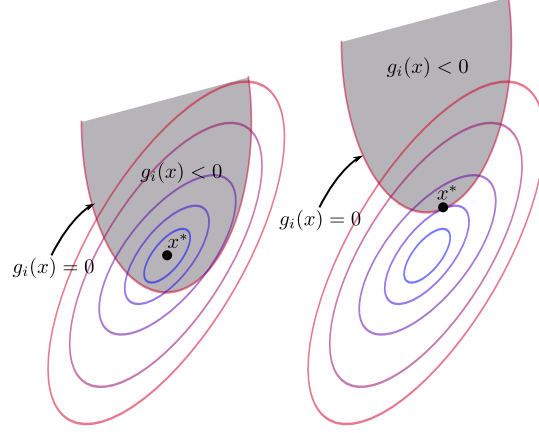
$$\begin{cases} \nabla f(x^*) - \lambda^* \nabla g(x^*) - \mu^* \nabla h(x^*) = 0 \\ g(x^*) = 0 \\ h(x^*) \geq 0 \\ \mu^* \geq 0 \\ \mu^* \cdot h(x^*) = 0 \end{cases} \quad (.1.13)$$

The last equality emphasizes that either  $h(x^*) > 0$  holds, in which case  $\mu^* = 0$ , or  $h(x^*) = 0$ , in which case  $\mu^*$  is irrelevant. See Figure .1.2.

### Lagrange multiplier and constrained dynamic optimization

Let's first put the Lagrange multiplier in the context of an optimization problem:

$$\inf \left\{ \int_0^T (x_t^2 - \alpha_t x_t) dt \right\} \quad (.1.14)$$



**Figure .1.2:** KKT conditions for  $h(x) \geq 0$ . The figure uses  $g$  for  $h$ . Source WIKIPEDIA

subject to  $x_T = 0$ . The Lagrangian is

$$\int_0^T (x_t^2 - \alpha_t x_t) dt + \lambda x_T \quad (.1.15)$$

and KKT condition suggests that  $x_t^* = 2\alpha_t$  for  $t < T$  and  $x_T^* = 0$  solves the problem. However, this is not an interesting problem. One can simply argue that changing  $x$  at  $T$  does not change the value of the integral, and therefore it is not really a constraint. Even if we impose a more restricted constraint such as  $x_t \geq 0$  on all  $t$ , the myopic solution is simply  $x_t^* = \max\{\alpha_t, 0\}/2$ .

However, for the constraint  $\int_0^T x_t dt = 0$ , the optimization problem becomes more interesting.

$$\sup_{\lambda} \inf_x \int_0^T (x_t^2 - \alpha_t x_t + \lambda x_t) dt \quad (.1.16)$$

KKT condition becomes  $x_t^* = (\alpha_t - \lambda^*)/2$  and  $\int_0^T x_t^* dt = 0$ . This implies that  $\lambda^* = \frac{1}{T} \int_0^T \alpha_t dt$ , and therefore  $x_t^* = (\alpha_t - \frac{1}{T} \int_0^T \alpha_t dt)/2$ .

Adding restrictions to a control problem is slightly more subtle due to the dynamics of  $x$ . We postpone the study of such problems to the future endeavors. However, a simple control problem can be described as a constrained optimization and can be solved via the Lagrange multiplier. We discuss this approach in the next section.

## .1.2 Gradient descent

Let  $f : \mathbb{R}^d \rightarrow \mathbb{R}$  be a differentiable function. We know that if  $x_*$  is a local minimum of  $f$ , then  $\nabla f(x_*) = 0$ . Therefore, solving  $\nabla f(x) = 0$  yields critical points of the function  $f$  including the local and global minimums, if there is any. In very high-dimensional problems such as those in deep learning,  $f$  is a highly nonlinear function. The set of critical points of  $f$  can be very irregular and hard to find via solving  $\nabla f(x) = 0$ , [Petersen et al. \[2021\]](#). In such situation, finding a global minimum requires some luck. However, a local minimum can be found via the method of gradient descent, GD henceforth. GD is based on a simple fact. At any point  $x$ , direction of maximum descent (ascent) of  $f$  is parallel to  $-\nabla f(x)$  ( $\nabla f(x)$ ). Therefore, if we move from point  $x$  toward  $-\nabla f(x)$ , the value of  $f$  decreases and we are closer to a local minimum than we were at  $x$ . In other words, for some small  $\alpha > 0$   $x - \alpha \nabla f(x)$  is closer to a local minimum than  $x$  was. The most important parameter in GD is  $\alpha$ , learning rate. Ideally, when we are further from a local minimum,

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**Algorithm 3:** Gradient descent algorithm

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**Data:**  $x$  randomly chosen from the domain.

**Parameter:** Learning rate  $\{\alpha_n\}_n$ , number of iterations  $N$ , and tolerance  $\epsilon > 0$ .

1 **while**  $|\nabla f(x_n)| > \epsilon$  &  $n \leq N$  **do**

2      $x \leftarrow x - \alpha_n \nabla f(x);$

3      $n \leftarrow n + 1$

4 **return**  $x$

---

we like to move faster, therefore, we choose a larger learning rate. When we are very close to the minimum, we like the learning rate to be smaller. A general rule of thumb suggests  $\alpha_n$  should be such that  $\alpha_n \rightarrow 0$  and  $\sum_n \alpha_n = \infty$ . Decrease of  $\alpha_n$  to 0, allows DG to slow down when it gets closer to the minimum. If  $\alpha_n$  is constant, it can potentially jump over the minimum. The justification for  $\sum_n \alpha_n = \infty$  can be rigorously explained. The following theorem shows why this condition holds under some strong assumption. The theorem shows the continuous gradient descent to the global minimum of a convex function. In practice, GD is a discrete algorithm.

**Theorem .1.2.** Let  $\int_0^\infty \alpha_t dt = \infty$  and  $f$  be a convex function with a global minimum at  $x_*$  such that for all  $x$

$$(x - x_*) \cdot \nabla f(x) \geq \lambda |x - x_*|^2$$

Then, the solution of the ODE given by  $dx_t = -\alpha_t \nabla f(x_t) dt$  converges to  $x_*$ .

*Proof.* We use chain rule to evaluate

$$d|x_t - x_*|^2 = 2(x_t - x_*) \cdot dx_t = -2\alpha(x_t - x_*) \cdot \nabla f(x_t) dt$$

By the assumption of the theorem, we have

$$d|x_t - x_*|^2 = -2\alpha_t(x_t - x_*) \cdot \nabla f(x_t) dt \leq -2\lambda\alpha_t|x_t - x_*|^2 dt$$

or

$$d|x_t - x_*|^2 + 2\lambda\alpha_t|x_t - x_*|^2 dt \leq 0$$

By multiplying the above by  $\exp(2\lambda \int_0^t \alpha_s ds)$  and integrating, we obtain

$$\exp(2\lambda \int_0^t \alpha_s ds) |x_t - x_*|^2 \leq |x_0 - x_*|^2$$

or

$$|x_t - x_*|^2 \leq |x_0 - x_*|^2 \exp(-2\lambda \int_0^t \alpha_s ds)$$

The right-hand above goes to zero as  $t \rightarrow \infty$ . □

Plain GD is not working for many applications and several versions of GD are introduced overcome the issues arising in applications. For example, if dimension is very large,  $d > 1e6$ , evaluation of gradient of a function is very time consuming. In such cases, stochastic gradient descent, SGD henceforth, is used. SGD randomly chooses small number of directions from  $\{1, \dots, d\}$  e.g. 10 out of  $1e6$  and gradient is only evaluated in those directions. To stabilize SGD, the introduced momentum into gradient descent. A popular SGD algorithm is ADAMS algorithm.



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