Syllabus part I: advanced planning Artificial Intelligence II Dr Sean Holden New things to be looked at include some more advanced material on *planning* algorithms: Computer Laboratory, Room FC06 • Heuristics and GraphPlan: incorporating heuristics into partial-order plan-Telephone extension 63725 ning, planning graphs, the GraphPlan algorithm. [1 lecture] Email: sbh11@cl.cam.ac.uk • *Planning using propositional logic:* representing planning problems using www.cl.cam.ac.uk/~sbh11/ propositional logic, and generating plans using satisfiability solvers. [1 lecturel • Planning using constraint satisfaction: representing planning problems so that they can be solved using constraint satisfaction solvers. [1 lecture] There is no warranty attached to the stated lecture timings. Copyright © Sean Holden 2002-12. 2 1 Syllabus part II: uncertainty in AI Syllabus part III: uncertainty and time We then delve into some more modern material which takes account of uncer-We then look at how uncertain reasoning and learning can take place when time is tainty: to be taken into account: • Uncertainty and Bayesian networks: review of probability as applied to AI, • Markov processes: transition and sensor models. Bayesian networks, inference in Bayesian networks using both exact and ap-• Inference in temporal models: filtering, prediction, smoothing and finding the proximate techniques, other ways of dealing with uncertainty. [4 lectures] most likely explanation. • Utility and decision-making: maximising expected utility, decision networks, • Hidden Markov models. [2 lectures] the value of information. [1 lecture] Please read the supplementary notes on probability handout.

Syllabus part IV: learning

Finally, we apply probability to *supervised learning* to obtain [1 lecture] more sophisticated models of learning.

- Bayes theorem as applied to supervised learning. [1 lecture]
- The maximum likelihood and maximum a posteriori hypotheses. [1 lecture]
- Applying the Bayesian approach to neural networks. [3 lectures]

We finish the course by taking a brief look at reinforcement learning.

- How can we learn from rewards and punishments?
- The *Q*-learning algorithm. [1 lecture]

Reinforcement learning can be thought of as combining many of the elements covered in this course and in AI I, and thus provides a natural place to stop.

Books

Once again, the main single text book for the course is:

• Artificial Intelligence: A Modern Approach. Stuart Russell and Peter Norvig, Prentice Hall.

There is an accompanying web site at

aima.cs.berkeley.edu

Either the second or third edition should be fine, but avoid the first edition as it does not fit this course so well.

Chapter numbers given in these notes refer to the third edition.

<u>Books</u>

5

For some of the new material on neural networks you might also like to take a look at:

• *Pattern Recognition and Machine Learning*. Christopher M. Bishop. Springer, 2006.

For some of the new material on reinforcement learning you might like to consult:

• Machine Learning. Tom Mitchell. McGraw Hill, 1997.

For further material on planning try:

• Automated Planning: Theory and Practice. Malik Ghallab, Dana Nau and Paolo Traverso. Morgan Kaufmann, 2004.

7

Dire Warning

6

DIRE WARNING

This course contains quite a lot of:

- 1. Probability
- 2. Matrix algebra
- 3. Calculus

As I am an evil and vindictive person who likes to be unkind to kittens I will assume that you know everything on these subjects that was covered in earlier courses.

If you don't it is *essential* that you re-visit your old notes and make sure that you're at home with that material.

YOU HAVE BEEN WARNED

How's your maths?

To see if you're up to speed on the maths, have a go at the following:

Evaluate the integral

$$\int_{-\infty}^{\infty} \exp(-x^2) \, dx$$

Hint: this is a pretty standard result. Square the integral and change to polar coordinates.

How's your maths?

Following on from that, here's something a bit more challenging.

Evaluate the integral

$$\int_{-\infty}^{\infty} \cdots \int_{-\infty}^{\infty} \exp\left(-\frac{1}{2}\left(\mathbf{x}^{T}\boldsymbol{\Sigma}\mathbf{x} + \mathbf{x}^{T}\boldsymbol{\alpha} + \beta\right)\right) dx_{1} \cdots dx_{n}$$

where Σ is a symmetric $n \times n$ matrix with real elements, $\pmb{\alpha} \in \mathbb{R}^n$, $\beta \in \mathbb{R}$ and

$$\mathbf{x}^T = \left[\begin{array}{ccc} x_1 & x_2 & \cdots & x_n \end{array} \right] \in \mathbb{R}$$

(This second one is a bit tricky. I'll show you the answer later...)

Planning II

9

We now examine:

- The way in which *basic heuristics* might be defined for use in planning problems.
- The construction of *planning graphs* and their use in obtaining more sensible heuristics.
- Planning graphs as the basis of the GraphPlan algorithm.
- Planning using propositional logic.
- Planning using constraint satisfaction.

Reading: Russell and Norvig, relevant sections of chapter 11.

A quick review

10

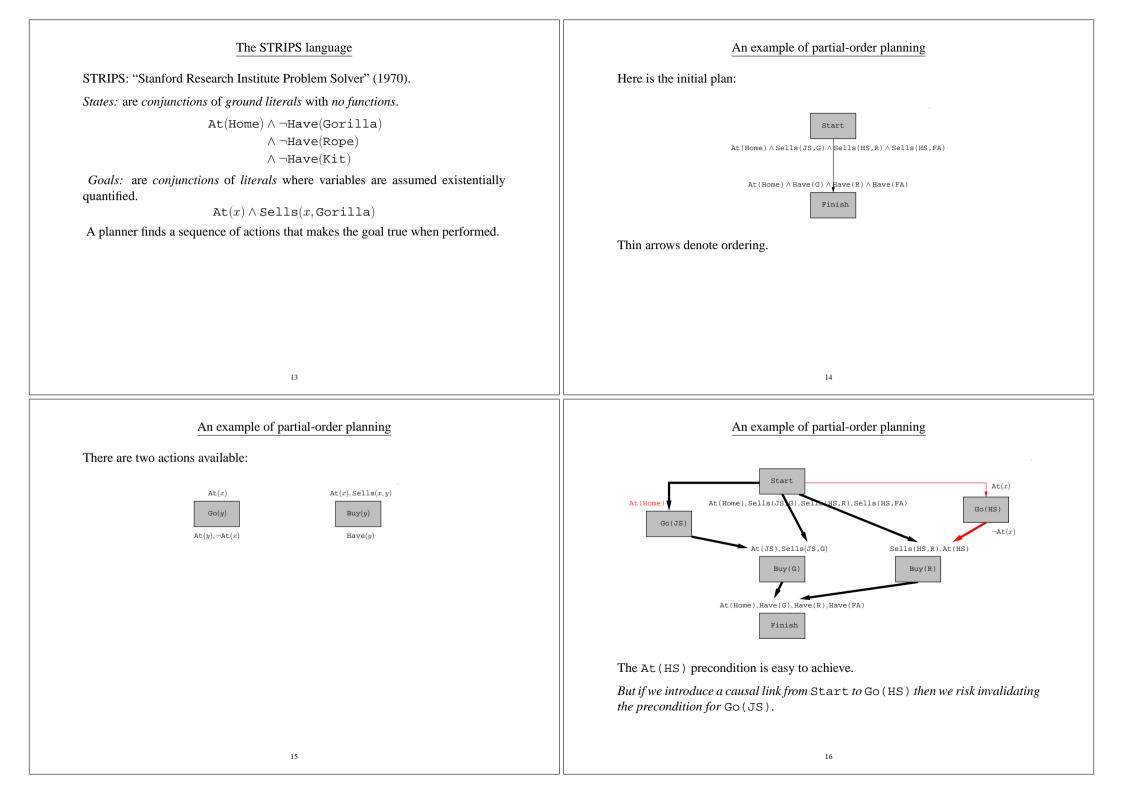
We used the following simple example problem.

The intrepid little scamps in the *Cambridge University Roof-Climbing Society* wish to attach an inflatable gorilla to the spire of a famous College. To do this they need to leave home and obtain:

- An *inflatable gorilla*: these can be purchased from all good joke shops.
- Some *rope*: available from a hardware store.
- A *first-aid kit*: also available from a hardware store.

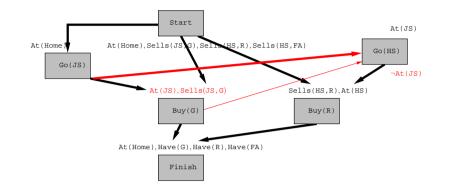
They need to return home after they've finished their shopping.

How do they go about planning their jolly escapade?



An example of partial-order planning

The planner could backtrack and try to achieve the ${\tt At}(x)$ precondition using the existing Go(JS) step.



This involves a threat, but one that can be fixed using promotion.

17

Using heuristics in planning

We can quickly suggest some possibilities.

For example

h = number of unsatisfied preconditions

or

h =number of unsatisfied preconditions - number satisfied by the start state

These can lead to underestimates or overestimates:

• Underestimates if actions can affect one another in undesirable ways.

• Overestimates if actions achieve many preconditions.

Using heuristics in planning

We found in looking at search problems that *heuristics* were a helpful thing to have.

Note that now:

- There is no simple representation of a *state*.
- Consequently it is harder to measure the distance to a *goal*.

Defining heuristics for planning is therefore more difficult than it was for search problems.

Using heuristics in planning

18

We can go a little further by learning from *Constraint Satisfaction Problems* and adopting the *most constrained variable* heuristic:

• Prefer the precondition satisfiable in the smallest number of ways.

This can be computationally demanding but two special cases are helpful:

- Choose preconditions for which no action will satisfy them.
- Choose preconditions that can only be satisfied in one way.

Planning graphs

Planning graphs can be used:

- To compute more sensible heuristics.
- To generate entire plans.

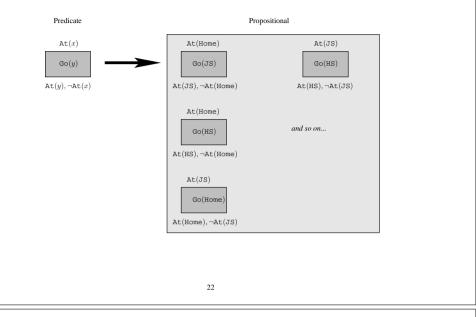
Also, planning graphs are easy to construct.

They apply only when it is possible to work entirely using *propositional* representations of plans.

Luckily, STRIPS can always be propositionalized...



For example: the triumphant return of the gorilla-purchasing roof-climbers...



Planning graphs

21

A planning graph is constructed in levels:

- Level 0 corresponds to the *start state*.
- At each level we keep *approximate* track of all things that *could* be true at the corresponding time.
- At each level we keep *approximate* track of what actions *could* be applicable at the corresponding time.

The approximation is due to the fact that not all conflicts between actions are tracked. *So*:

- The graph can *underestimate* how long it might take for a particular proposition to appear, and therefore ...
- ... a heuristic can be extracted.

Planning graphs: a simple example

Our intrepid student adventurers will of course need to inflate their *gorilla* before attaching it to a *distinguished roof*. It has to be purchased before it can be inflated.

Start state: Empty.

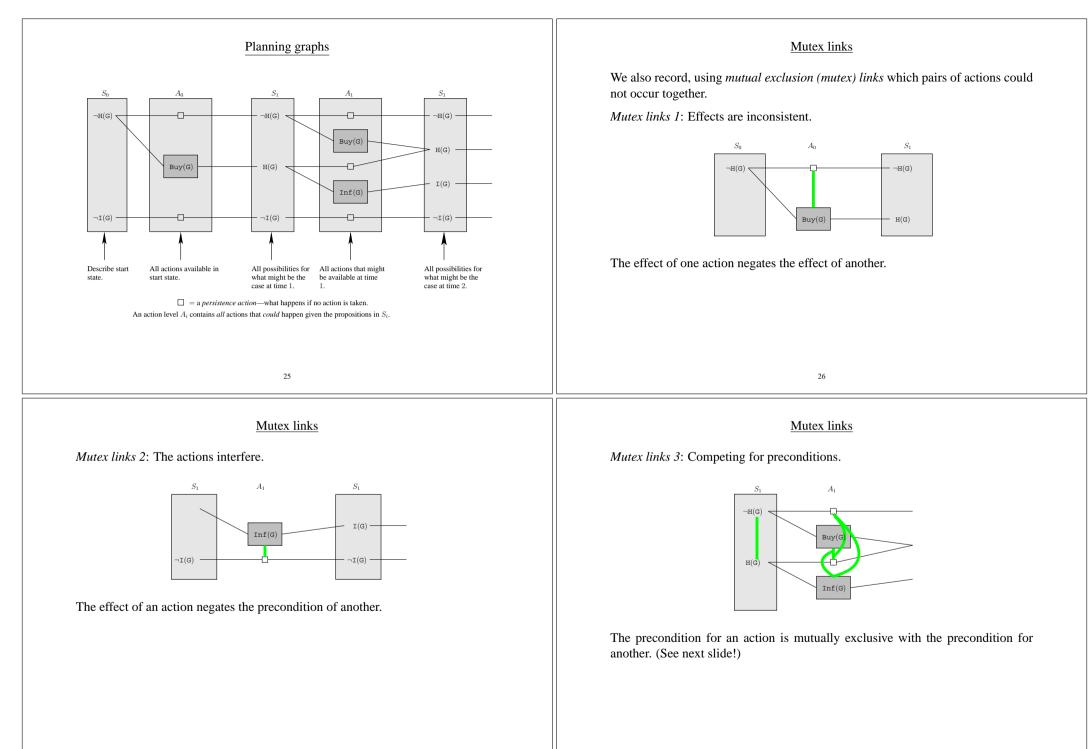
We assume that anything not mentioned in a state is false. So the state is actually

¬Have(Gorilla) and ¬Inflated(Gorilla)

Actions:

¬Have(Gorilla)	Have(Gorilla)
Buy(Gorilla)	Inflate(Gorilla)
Have(Gorilla)	Inflated(Gorilla)

Goal: Have(Gorilla) and Inflated(Gorilla).



Mutex links

A state level S_i contains *all* propositions that *could* be true, given the possible preceding actions.

 S_1

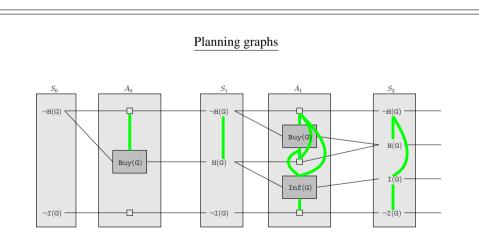
-H(G)

H(G)

29

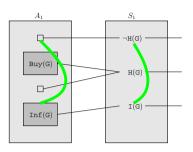
We also use mutex links to record pairs that can not be true simultaneously:

Possibility 1: pair consists of a proposition and its negation.



Mutex links

Possibility 2: all pairs of actions that could achieve the pair of propositions are mutex.



The construction of a planning graph is continued until two identical levels are obtained.

Obtaining heuristics from a planning graph

30

To estimate the cost of reaching a single proposition:

- Any proposition not appearing in the final level has *infinite cost* and *can never be reached*.
- The *level cost* of a proposition is the level at which it first appears *but* this may be inaccurate as several actions can apply at each level and this cost does not count the *number of actions*. (It is however *admissible*.)
- A *serial planning graph* includes mutex links between all pairs of actions except persistence actions.

Level cost in serial planning graphs can be quite a good measurement.

Obtaining heuristics from a planning graph

How about estimating the cost to achieve a *collection* of propositions?

- *Max-level*: use the maximum level in the graph of any proposition in the set. Admissible but can be inaccurate.
- *Level-sum*: use the sum of the levels of the propositions. Inadmissible but sometimes quite accurate if goals tend to be decomposable.
- *Set-level*: use the level at which *all* propositions appear with none being mutex. Can be accurate if goals tend *not* to be decomposable.

Other points about planning graphs

A planning graph guarantees that:

1. If a proposition appears at some level, there may be a way of achieving it.

2. If a proposition does not appear, it can not be achieved.

The first point here is a loose guarantee because only *pairs* of items are linked by mutex links.

Looking at larger collections can strengthen the guarantee, but in practice the gains are outweighed by the increased computation.

Graphplan

33

The *GraphPlan* algorithm goes beyond using the planning graph as a source of heuristics.

```
Start at level 0;
while(true) {
    if (all goal propositions appear in the current level
        AND no pair has a mutex link) {
        attempt to extract a plan;
        if (a solution is obtained)
            return the solution;
        else if (graph indicates there is no solution)
```

```
else
```

return fail;

```
expand the graph to the next level;
```

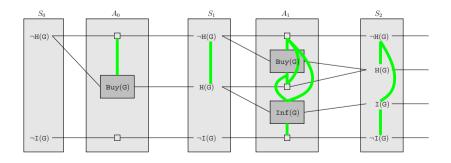
```
}
```

We *extract a plan* directly from the planning graph. Termination can be proved but will not be covered here.

Graphplan in action

34

Here, at levels S_0 and S_1 we do not have both H(G) and I(G) available with no mutex links, and so we expand first to S_1 and then to S_2 .



At S_2 we try to extract a solution (plan).

Extracting a plan from the graph

Extraction of a plan can be formalised as a *search problem*.

States contain a level, and a collection of unsatisfied goal propositions.

Start state: the current final level of the graph, along with the relevant goal propositions.

Goal: a state at level S_0 containing the initial propositions.

Extracting a plan from the graph

Actions: For a state S with level S_i , a valid action is to select any set X of actions in A_{i-1} such that:

1. no pair has a mutex link;

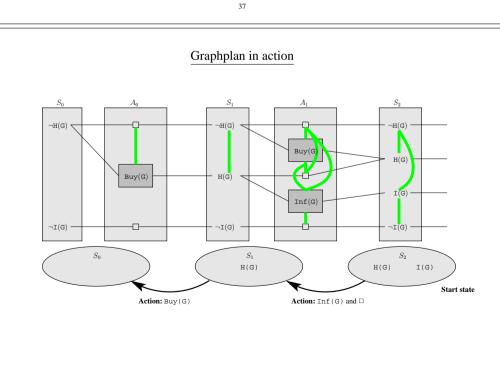
2. no pair of their preconditions has a mutex link;

3. the effects of the actions in X achieve the propositions in S.

The effect of such an action is a state having level S_{i-1} , and containing the preconditions for the actions in X.

Each action has a cost of 1.

38



Heuristics for plan extraction

We can of course also apply *heuristics* to this part of the process.

For example, when dealing with a set of propositions:

- Choose the proposition having maximum level cost first.
- For that proposition, attempt to achieve it using the action for which the *maximum/sum level cost of its preconditions is minimum*.

Planning III: planning using propositional logic	Propositional logic for planning
Last year we saw that plans might be extracted from a knowledge base via <i>theorem</i> proving, using first order logic (FOL) and situation calculus.	We attempt to construct this sentence such that:
BUT: this might be computationally infeasible for realistic problems.	• If M is a model of the sentence then M assigns \top to a proposition if and only if it is in the plan.
Sophisticated techniques are available for testing <i>satisfiability</i> in <i>propositional logic</i> , and these have also been applied to planning.	 Any assignment denoting an incorrect plan will not be a model as the goal description will not be ⊤.
The basic idea is to attempt to find a model of a sentence having the form	• The sentence is unsatisfiable if no plan exists.
description of start state	
41	42
Propositional logic for planning	Propositional logic for planning
Start state: $S = At^{0}(a, spire) \land At^{0}(b, ground)$ $\land \neg At^{0}(a, ground) \land \neg At^{0}(b, spire)$	Goal: $G = At^{i}(a, ground) \land At^{i}(b, spire)$ $\land \neg At^{i}(a, spire) \land \neg At^{i}(b, ground)$ Actions: can be introduced using the equivalent of successor-state axioms
+ The two climbers want to swap places	$\begin{array}{c} At^{1}(a,ground) \leftrightarrow \\ (At^{0}(a,ground) \wedge \neg(At^{0}(a,ground) \wedge Move^{0}(a,ground,spire))) \\ \vee (At^{0}(a,spire) \wedge Move^{0}(a,spire,ground)) \end{array} $ (1) Denote by A the collection of all such axioms.
Remember that an expression such as ${\tt At}^0({\tt a}, {\tt spire})$ is a proposition. The superscripted number now denotes time.	
43	44

Propositional logic for planning

We will now find that $S \land A \land G$ has a model in which $Move^0(a, pire, ground)$ and $Move^0(b, ground, spire)$ are \top while all remaining actions are \bot .

In more realistic planning problems we will clearly not know in advance at what time the goal might expect to be achieved.

We therefore:

- Loop through possible final times T.
- Generate a goal for time T and actions up to time T.
- Try to find a model and extract a plan.

• Until a plan is obtained or we hit some maximum time.

Propositional logic for planning

45

Life becomes more complicated still if a third location is added: hospital.

 $Move^{0}(a, spire, ground) \land Move^{0}(a, spire, hospital)$

is perfectly valid and so we need to specify that he can't move to two places simultaneously

 $\neg(\texttt{Move}^{i}(\texttt{a},\texttt{spire},\texttt{ground}) \land \texttt{Move}^{i}(\texttt{a},\texttt{spire},\texttt{hospital})) \\ \neg(\texttt{Move}^{i}(\texttt{a},\texttt{ground},\texttt{spire}) \land \texttt{Move}^{i}(\texttt{a},\texttt{ground},\texttt{hospital}))$

and so on.

These are action-exclusion axioms.

Unfortunately they will tend to produce *totally-ordered* rather than *partially-ordered* plans.

Propositional logic for planning

Unfortunately there is a problem—we may, if considerable care is not applied, also be able to obtain less sensible plans.

In the current example

 $\begin{array}{l} {\tt Move}^0({\tt b}, {\tt ground}, {\tt spire}) = \top \\ {\tt Move}^0({\tt a}, {\tt spire}, {\tt ground}) = \top \\ \hline \\ {\tt Move}^0({\tt a}, {\tt ground}, {\tt spire}) \end{array} = \top \end{array}$

is a model, because the successor-state axiom (1) does not in fact preclude the application of $Move^{0}(a, ground, spire)$.

We need a *precondition axiom*

```
Move^{i}(a, ground, spire) \rightarrow At^{i}(a, ground)
```

and so on.

46

Propositional logic for planning

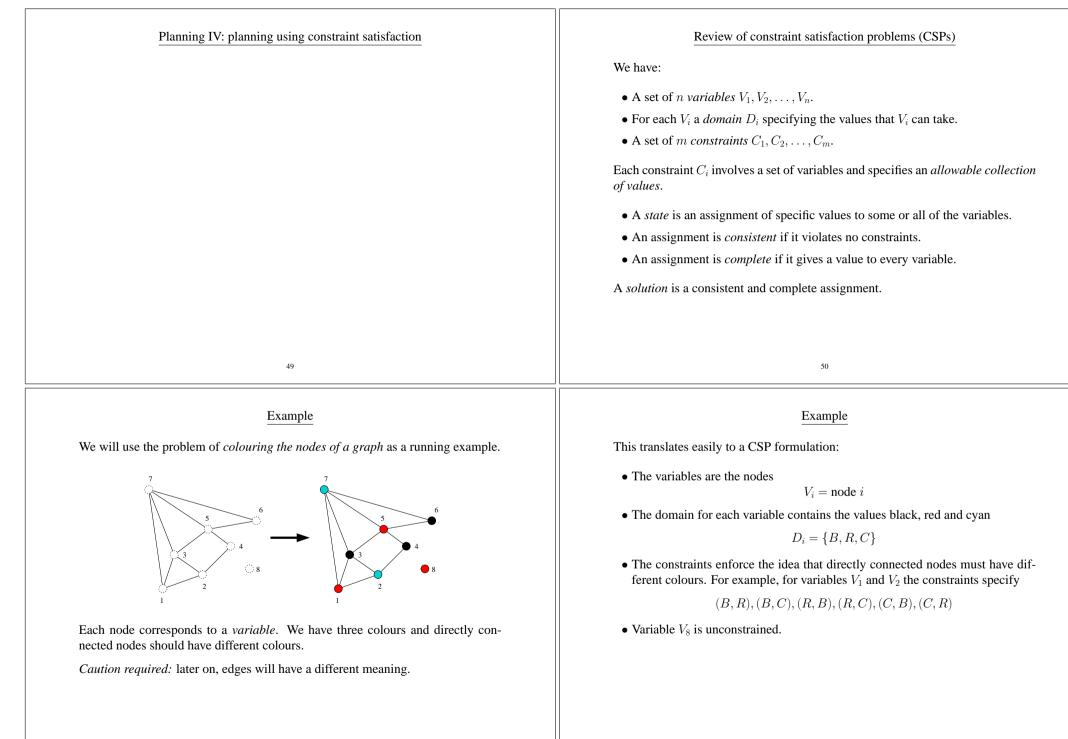
Alternatively:

- 1. Prevent actions occurring together if one negates the effect or precondition of the other.
- 2. Or, specify that something can't be in two places simultaneously

 $\forall x, i, \texttt{l1}, \texttt{l2} \quad \texttt{l1} \neq \texttt{l2} \rightarrow \neg(\texttt{At}^i(x, \texttt{l1}) \land \texttt{At}^i(x, \texttt{l2}))$

This is an example of a *state constraint*.

Clearly this process can become very complex, but there are techniques to help deal with this.



Different kinds of CSP

This is an example of the simplest kind of CSP: it is *discrete* with *finite domains*. We will concentrate on these.

We will also concentrate on *binary constraints*; that is, constraints between *pairs of variables*.

- Constraints on single variables—*unary constraints*—can be handled by adjusting the variable's domain. For example, if we don't want V_i to be *red*, then we just remove that possibility from D_i .
- *Higher-order constraints* applying to three or more variables can certainly be considered, but...
- ...when dealing with finite domains they can always be converted to sets of binary constraints by introducing extra *auxiliary variables*.

How does that work?

The state-variable representation

Another planning language: the state-variable representation.

Things of interest such as people, places, objects etc are divided into domains:

- $D_1 = \{\texttt{climber1}, \texttt{climber2}\}$
- $D_2 = \{\texttt{home, jokeShop, hardwareStore, pavement, spire, hospital}\}$
- $D_3 = \{ \texttt{rope}, \texttt{inflatableGorilla} \}$

Part of the specification of a planning problem involves stating which domain a particular item is in. For example

 $D_1(\texttt{climber1})$

and so on.

Relations and functions have arguments chosen from unions of these domains.

$$above(x,y) \subseteq \mathcal{D}_1^{above} \times \mathcal{D}_2^{abov}$$

is a relation. The $\mathcal{D}_i^{\text{above}}$ are unions of one or more D_i .

54

The state-variable representation

Note:

- For properties such as a *location* a function might be considerably more suitable than a relation.
- For locations, everything has to be *somewhere* and it can only be in *one place at a time*.

So a function is perfect and immediately solves some of the problems seen earlier.

The state-variable representation

53

The relation above is in fact a *rigid relation (RR)*, as it is unchanging: it does not depend upon *state*. (Remember *fluents* in situation calculus?)

Similarly, we have *functions*

$$\operatorname{at}(x_1,s): \mathcal{D}_1^{\operatorname{at}} \times S \to \mathcal{D}^{\operatorname{at}}.$$

Here, at(x, s) is a *state-variable*. The domain \mathcal{D}_1^{at} and range \mathcal{D}^{at} are unions of one or more D_i . In general these can have multiple parameters

$$sv(x_1,\ldots,x_n,s): \mathcal{D}_1^{sv}\times\cdots\times\mathcal{D}_n^{sv}\times S\to\mathcal{D}^{sv}$$

A state-variable denotes assertions such as

$$at(gorilla, s) = jokeShop$$

where s denotes a *state* and the set S of all states will be defined later.

The state variable allows things such as locations to change—again, much like *fluents* in the situation calculus.

Variables appearing in relations and functions are considered to be typed.

The state-variable representation

Actions as usual, have a name, a set of preconditions and a set of effects.

- Names are unique, and followed by a list of variables involved in the action.
- *Preconditions* are expressions involving state variables and relations.
- *Effects* are assignments to state variables.

For example:

$\mathtt{at}(x,s) = l$
${\tt sells}(l,y)$
${\tt has}(y,s) = l$
$\verb+has(y,s) = x$

The state-variable representation

Goals are sets of expressions involving state variables.

For example:

Goal:at(climber, s) = homehas(rope, s) = climberat(gorilla, s) = spire

From now on we will generally suppress the state s when writing state variables.

The state-variable representation

57

We can essentially regard a *state* as just a statement of what values the state variables take at a given time.

Formally:

• For each state variable sv we can consider all ground instances such as sv(climber, rope)—with arguments that are *consistent* with the *rigid relations*.

Define X to be the set of all such ground instances.

• A state s is then just a set

$$s = \{(v = c) | v \in X\}$$

where c is in the range of v.

This allows us to define the *effect of an action*.

A planning problem also needs a *start state* s_0 , which can be defined in this way.

The state-variable representation

58

Considering all the ground actions consistent with the rigid relations:

• An action is *applicable in s* if all expressions v=c appearing in the set of preconditions also appear in *s*.

Finally, there is a function γ that maps a state and an action to a new state

$$\gamma(s,a) = s'$$

Specifically, we have

$$\gamma(s,a) = \{(v=c) | v \in X\}$$

where either c is specified in an effect of a, or otherwise v = c is a member of s. Note: the definition of γ implicitly solves the *frame problem*.

The state-variable representation

A solution to a planning problem is a sequence (a_0, a_1, \ldots, a_n) of actions such that...

- a_0 is applicable in s_0 and for each i, a_i is applicable in $s_i = \gamma(s_{i-1}, a_{i-1})$.
- For each goal g we have

 $g \in \gamma(s_n, a_n).$

What we need now is a method for *transforming* a problem described in this language into a CSP.

We'll once again do this for a fixed upper limit T on the number of steps in the plan.

61

Converting to a CSP

Step 2: encode *ground state variables* as *CSP variables*, with a complete copy of all the state variables *for each time step*.

So, for each t where $0 \le t \le T$ we have a CSP variable

 $\mathbf{sv}_i^t(c_1,\ldots,c_n)$

with domain \mathcal{D}^{sv_i} . (That is, the *domain* of the CSP variable is the *range* of the state variable.)

Example: at some point in searching for a plan we might attempt to find the solution to the corresponding CSP involving

 $location^9(climber1) = hospital.$

Converting to a CSP

Step 1: encode actions as CSP variables.

For each time step t where $0 \le t \le T - 1$, the CSP has a variable

 \texttt{action}^t

with domain

 $D^{\operatorname{action}^t} = \{ a | a \text{ is the ground instance of an action} \} \cup \{ \operatorname{none} \}$

Example: at some point in searching for a plan we might attempt to find the solution to the corresponding CSP involving

 $\texttt{action}^5 = \texttt{attach}(\texttt{inflatableGorilla},\texttt{spire})$

WARNING: be careful in what follows to distinguish between *state variables, actions etc* in the planning problem and *variables* in the CSP.

62

Converting to a CSP

Step 3: encode the preconditions for actions in the planning problem as constraints in the CSP problem.

For each time step t and for each ground action $a(c_1, \ldots, c_n)$ with arguments *consistent with the rigid relations in its preconditions*:

For a precondition of the form $sv_i = v$ include constraint pairs

$$(\texttt{action}^t = \texttt{a}(c_1, \dots, c_n), \\ \texttt{sv}_i^t = v)$$

Example: consider the action buy(x, y, l) introduced above, and having the preconditions at(x) = l, sells(l, y) and has(y) = l.

Assume sells(y, l) is only true for

l = jokeShop

and

y = inflatableGorilla

(it's a very strange town) so we only consider these values for l and y. Then for each time step t we have the constraints...

Converting to a CSP

 $\begin{tabular}{|c|c|c|c|c|} \hline action^t = buy(climber1, inflatableGorilla, jokeShop) \\ paired with \\ at^t(climber1) = jokeShop \\ \hline action^t = buy(climber1, inflatableGorilla, jokeShop) \\ paired with \\ has^t(inflatableGorilla) = jokeShop \\ \hline action^t = buy(climber2, inflatableGorilla, jokeShop) \\ paired with \\ at^t(climber2) = jokeShop \\ \hline action^t = buy(climber2, inflatableGorilla, jokeShop) \\ paired with \\ has^t(inflatableGorilla) = jokeShop \\ \hline adt (inflatableGorilla) = jokeShop \\ \hline and so on... \\ \hline ext{terms}$

65

Converting to a CSP

Step 5: encode the frame axioms as constraints in the CSP problem.

An action must not change things not appearing in its effects. So:

For:

- 1. Each time step t.
- 2. Each ground action $a(c_1, \ldots, c_n)$ with arguments *consistent with the rigid relations in its preconditions*.
- 3. Each sv_i that does not appear in the effects of a, and each $v \in \mathcal{D}^{sv_i}$

include in the CSP the ternary constraint

$$(\texttt{action}^t = \texttt{a}(c_1, \dots, c_n)$$

 $\texttt{sv}_i^t = v,$
 $\texttt{sv}_i^{t+1} = v)$

Converting to a CSP

Step 4: encode the effects of actions in the planning problem as constraints in the CSP problem.

For each time step t and for each ground action $a(c_1, \ldots, c_n)$ with arguments *consistent with the rigid relations in its preconditions*:

For an effect of the form $sv_i = v$ include constraint pairs

$$(\texttt{action}^t = \texttt{a}(c_1, \dots, c_n)$$

 $\texttt{sv}_i^{t+1} = v)$

Example: continuing with the previous example, we will include constraints

66

Finding a plan

Finally, having encoded a planning problem into a CSP, we solve the CSP.

The scheme has the following property:

A solution to the planning problem with at most T steps exists if and only if there is a a solution to the corresponding CSP.

Assume the CSP has a solution.

Then we can extract a plan simply by looking at the values assigned to the $action^t$ variables in the solution of the CSP.

It is also the case that:

There is a solution to the planning problem with at most T steps if and only if there is a solution to the corresponding CSP from which the solution can be extracted in this way.

For a proof see:

Automated Planning: Theory and Practice

Malik Ghallab, Dana Nau and Paolo Traverso. Morgan Kaufmann 2004.

Uncertainty I: Probability as Degree of Belief Uncertainty in AI We now examine: The (predominantly logic-based) methods covered so far have assorted shortcomings: • How *probability theory* might be used to represent and reason with knowledge when we are uncertain about the world. • Limited *epistemological commitment*—true/false/unknown. • How *inference* in the presence of uncertainty can in principle be performed • Actions are possible when sufficient knowledge is available... using only basic results along with the *full joint probability distribution*. • ...but this is not generally the case. • How this approach *fails* in practice. • In practice there is a need to cope with *uncertainty*. • How the notions of *independence* and *conditional independence* may be used For example in the Wumpus World: to solve this problem. • We can not make observations further afield than the current locality. Reading: Russell and Norvig, chapter 13. • Consequently inferences regarding pit/wumpus location *etc* will not usually be possible. 69 70 Uncertainty in AI Logic (as seen so far) has major shortcomings A couple of more subtle problems have also presented themselves: An example: $\forall x \text{ symptom}(x, \text{toothache}) \rightarrow \text{problem}(x, \text{cavity})$ • The *Qualification Problem:* it is not generally possible to guarantee that an action will succeed—only that it will succeed if many other preconditions This is plainly incorrect. Toothaches can be caused by things other than cavities. do/don't hold. $\forall x \text{ symptom}(x, \text{toothache}) \rightarrow \text{problem}(x, \text{cavity}) \lor$ • Rational action depends on the likelihood of achieving different goals, and $problem(x, abscess) \lor$ their relative desirability. $problem(x, gum-disease) \lor$. . . BUT: • It is *impossible to complete* the list. • There's no clear way to take account of the relative likelihoods of different causes.

72

Logic (as seen so far) has major shortcomings	Logic (as seen so far) has major shortcomings
If we try to make a <i>causal rule</i> $\forall x \text{ problem}(x, abscess) \rightarrow symptom(x, toothache)$ it's still wrong—abscesses do not always cause pain. We need further information in addition to problem(x, abscess) and it's still not possible to do this correctly.	 FOL can fail for essentially three reasons: 1. <i>Laziness:</i> it is not feasible to assemble a set of rules that is sufficiently exhaustive. If we could, it would not be feasible to apply them. 2. <i>Theoretical ignorance:</i> insufficient knowledge <i>exists</i> to allow us to write the rules. 3. <i>Practical ignorance:</i> even if the rules have been obtained there may be insufficient information to apply them. Instead of thinking in terms of the <i>truth</i> or <i>falsity</i> of a statement we want to deal with an agent's <i>degree of belief</i> in the statement. <i>Probability theory</i> is the perfect tool for application here. <i>Probability theory</i> allows us to <i>summarise</i> the uncertainty due to laziness and ignorance.
73	74
An important distinction There is a fundamental difference between <i>probability theory</i> and <i>fuzzy logic</i> : • When dealing with probability theory, statements remain <i>in fact</i> either <i>true</i> or <i>false</i> . • A probability denotes an agent's <i>degree of belief</i> one way or another. • Fuzzy logic deals with <i>degree of truth</i> . In practice the use of probability theory has proved spectacularly successful.	Belief and evidence An agent's beliefs will depend on what it has <i>perceived</i> : probabilities are based on <i>evidence</i> and may be altered by the acquisition of new evidence: • Prior (unconditional) probability denotes a degree of belief in the absence of evidence. • Posterior (conditional) probability denotes a degree of belief after evidence is perceived. As we shall see Bayes' theorem is the fundamental concept that allows us to update one to obtain the other.

Making rational decisions under uncertainty

When using *logic*, we concentrated on finding an action sequence guaranteed to achieve a goal, and then executing it.

When dealing with *uncertainty* we need to define *preferences* among states of the world and take into account the *probability* of reaching those states.

Utility theory is used to assign preferences.

Decision theory combines probability theory and utility theory.

A rational agent should act in order to maximise expected utility.

Probability

We want to assign degrees of belief to propositions about the world.

We will need:

- *Random variables* with associated *domains*—typically Boolean, discrete, or continuous.
- All the usual concepts—events, atomic events, sets etc.
- Probability distributions and densities.
- Probability axioms (Kolmogorov).
- Conditional probability and Bayes' theorem.

So if you've forgotten this stuff now is a good time to re-read it.

Probability

77

The standard axioms are:

• Range

 $0 \le \Pr(x) \le 1$

• Always true propositions

```
\Pr(\texttt{always true proposition}) = 1
```

• Always false propositions

 $\Pr(\text{always false proposition}) = 0$

• Union

$$\Pr(x \lor y) = \Pr(x) + \Pr(y) - \Pr(x \land y)$$

Origins of probabilities I

78

Historically speaking, probabilities have been regarded in a number of different ways:

- Frequentist: probabilities come from measurements.
- *Objectivist:* probabilities are actual "properties of the universe" which frequentist measurements seek to uncover.

An excellent example: quantum phenomena.

A bad example: coin flipping—the uncertainty is due to our uncertainty about the initial conditions of the coin.

• *Subjectivist:* probabilities are an agent's degrees of belief. This means the agent is allowed to make up the numbers!

Origins of probabilities II Origins of probabilities III The reference class problem: even frequentist probabilities are subjective. The principle of indifference (Laplace). *Example:* Say a doctor takes a frequentist approach to diagnosis. She examines • Give equal probability to all propositions that are syntactically symmetric with a large number of people to establish the prior probability of whether or not they respect to the available evidence. have heart disease. • Refinements of this idea led to the attempted development by Carnap and oth-To be accurate she tries to measure "similar people". (She knows for example that ers of *inductive logic*. gender might be important.) • The aim was to obtain the correct probability of any proposition from an arbi-Taken to an extreme, *all* people are *different* and there is therefore no *reference* trary set of observations. class. It is currently thought that no unique inductive logic exists. Any inductive logic depends on prior beliefs and the effect of these beliefs is overcome by evidence. 81 82 Prior probability Notation A prior probability denotes the probability (degree of belief) assigned to a propo-

For example

Pr(Cavity = true) = 0.05

denotes the degree of belief that a random person has a cavity *before we make any actual observation of that person.*

To keep things compact, we will use

sition in the absence of any other evidence.

 $\Pr(\texttt{Cavity})$

to denote the entire probability distribution of the random variable Cavity.

Instead of

 $\Pr(\text{Cavity} = \text{true}) = 0.05$

Pr(Cavity = false) = 0.95

write

 $\Pr(\text{Cavity}) = (0.05, 0.95)$

A similar convention will apply for joint distributions. For example, if Decay can take the values severe, moderate or low then

 $\Pr(\texttt{Cavity},\texttt{Decay})$

is a 2 by 3 table of numbers.

	severe	moderate	low
true	0.26	0.1	0.01
false	0.01	0.02	0.6

Similarly

Pr(true, Decay)

denotes 3 numbers *etc*.

The full joint probability distribution

The *full joint probability distribution* is the joint distribution of *all* random variables that describe the state of the world.

This can be used to answer any query.

(But of course life's not really that simple!)

Conditional probability

We use the conditional probability

$\Pr(x|y)$

to denote the probability that a proposition x holds given that *all the evidence we* have so far is contained in proposition y.

From basic probability theory

$$\Pr(x|y) = \frac{\Pr(x \land y)}{\Pr(y)}$$

Conditional probability is not analogous to logical implication.

- Pr(x|y) = 0.1 does *not* mean that if y is true then Pr(x) = 0.1.
- Pr(x) is a *prior probability*.
- The notation $\Pr(x|y)$ is for use when y is the *entire evidence*.
- $\Pr(x|y \wedge z)$ might be very different.

Using the full joint distribution to perform inference

85

We can regard the full joint distribution as a knowledge base.

We want to use it to obtain answers to questions.

		CP	-	¬CP
	HBP	¬HBP	HBP	¬HBP
HD	0.09	0.05	0.07	0.01
¬HD	0.02	0.08	0.03	0.65

We'll use this medical diagnosis problem as a running example.

• HD = Heart disease

 \bullet CP = Chest pain

• HBP = High blood pressure

Using the full joint distribution to perform inference

86

The process is nothing more than the application of basic results:

• Sum atomic events:

$$\begin{aligned} \Pr(\texttt{HD} \lor \texttt{CP}) &= \Pr(\texttt{HD} \land \texttt{CP} \land \texttt{HBP}) \\ &+ \Pr(\texttt{HD} \land \texttt{CP} \land \neg \texttt{HBP}) \\ &+ \Pr(\texttt{HD} \land \neg \texttt{CP} \land \texttt{HBP}) \\ &+ \Pr(\texttt{HD} \land \neg \texttt{CP} \land \neg \texttt{HBP}) \\ &+ \Pr(\neg \texttt{HD} \land \texttt{CP} \land \neg \texttt{HBP}) \\ &+ \Pr(\neg \texttt{HD} \land \texttt{CP} \land \neg \texttt{HBP}) \\ &= 0.09 + 0.05 + 0.07 + 0.01 + 0.02 + 0.08 \\ &= 0.32 \end{aligned}$$

\bullet Marginalisation: if A and B are sets of variables then

$$\Pr(A) = \sum_{b} \Pr(A \land b) = \sum_{b} \Pr(A|b) \Pr(b)$$

Using the full joint distribution to perform inference

Usually we will want to compute the *conditional probability* of *some variable(s)* given *some evidence*.

For example

$$\Pr(\text{HD}|\text{HBP}) = \frac{\Pr(\text{HD} \land \text{HBP})}{\Pr(\text{HBP})} = \frac{0.09 + 0.07}{0.09 + 0.07 + 0.02 + 0.03} = 0.76$$

and

$$\Pr(\neg \text{HD}|\text{HBP}) = \frac{\Pr(\neg \text{HD} \land \text{HBP})}{\Pr(\text{HBP})} = \frac{0.02 + 0.03}{0.09 + 0.07 + 0.02 + 0.03} = 0.24$$

Using the full joint distribution to perform inference

89

The general inference procedure is as follows:

$$\Pr(Q|e) = \frac{1}{Z} \Pr(Q \land e) = \frac{1}{Z} \sum_{u} \Pr(Q, e, u)$$

where

- Q is the query variable.
- \bullet *e* is the evidence.
- \bullet *u* are the unobserved variables.
- 1/Z normalises the distribution.

Using the full joint distribution to perform inference

The process can be simplified slightly by noting that

$$\alpha = \frac{1}{\Pr(\texttt{HBP})}$$

is a constant and can be regarded as a *normaliser* making relevant probabilities sum to 1.

So a short cut is to avoid computing it as above. Instead:

 $\Pr(\text{HD}|\text{HBP}) = \alpha \Pr(\text{HD} \land \text{HBP}) = (0.09 + 0.07)\alpha$

$$\Pr(\neg \texttt{HD}|\texttt{HBP}) = \alpha \Pr(\neg \texttt{HD} \land \texttt{HBP}) = (0.02 + 0.03)\alpha$$

and we need

 $\Pr(\text{HD}|\text{HBP}) + \Pr(\neg \text{HD}|\text{HBP}) = 1$

 $\alpha = \frac{1}{0.09 + 0.07 + 0.02 + 0.03}$

so

Using the full joint distribution to perform inference

90

Simple eh?

Well, no...

- For n Boolean variables the table has 2^n entries.
- Storage and processing time are both $O(2^n)$.
- You need to establish 2^n numbers to work with.

In reality we might well have n > 1000, and of course it's *even worse* if *variables* are non-Boolean.

How can we get around this?

Exploiting independence

If I toss a coin and roll a dice, the full joint distribution of outcomes requires $2 \times 6 = 12$ numbers to be specified.

	1	2	3	4	5	6
head	0.014	0.028	0.042	0.057	0.071	0.086
tail	0.033	0.067	0.1	0.133	0.167	0.2

Here $\Pr(\texttt{Coin}=\texttt{head})=0.3$ and the dice has probability i/21 for the ith outcome.

BUT: if we assume the outcomes are independent then

 $\Pr(\texttt{Coin},\texttt{Dice}) = \Pr(\texttt{Coin})\Pr(\texttt{Dice})$

93

Exploiting independence

In general you need to identify such independence through knowledge of the prob-

Where $\Pr(\texttt{Coin})$ has two numbers and $\Pr(\texttt{Dice})$ has six.

So instead of 12 numbers we only need 8.

• It generally does not work as clearly as this.

• The independent subsets themselves can be big.

lem.

BUT:

Exploiting independence

Similarly, say instead of just considering HD, HBP and CP we also consider the outcome of the *Oxford versus Cambridge tiddlywinks competition* TC:

$$Pr(TC = Oxford) = 0.2$$

$$Pr(TC = Cambridge) = 0.7$$

$$\Pr(\text{TC} = \text{Draw}) = 0.1$$

Now

 $\Pr(\text{hd}, \text{hbp}, \text{cp}, \text{tc}) = \Pr(\text{tc}|\text{hd}, \text{hbp}, \text{hd}) \Pr(\text{hd}, \text{hbp}, \text{hd})$

Assuming that the patient is not an *extraordinarily keen fan of tiddlywinks*, their cardiac health has nothing to do with the outcome, so

 $\Pr(TC|HD, HBP, HD) = \Pr(TC)$

and $2 \times 2 \times 2 \times 3 = 24$ numbers has been reduced to 3 + 8 = 11.

94

Bayes theorem

From first principles

 $\Pr(x, y) = \Pr(x|y) \Pr(y)$

 $\Pr(x, y) = \Pr(y|x) \Pr(x)$

so

 $\Pr(x|y) = \frac{\Pr(y|x)\Pr(x)}{\Pr(y)}$

The most important equation in modern AI?

When evidence e is involved this can be written

$$\Pr(Q|R, e) = \frac{\Pr(R|Q, e) \Pr(Q|e)}{\Pr(R|e)}$$

Bayes theorem

Taking another simple medical diagnosis example: *does a patient with a fever have malaria?* A doctor might know that

$$\Pr(\texttt{fever}|\texttt{malaria}) = 0.99$$

$$\Pr(\texttt{malaria}) = \frac{1}{10000}$$

 $\Pr(\texttt{fever}) = \frac{1}{20}$

Consequently we can try to obtain $\Pr(\texttt{malaria}|\texttt{fever})$ by direct application of Bayes theorem

$$\Pr(\texttt{malaria}|\texttt{fever}) = \frac{0.99 \times 0.0001}{0.05} = 0.00198$$

or using the alternative technique

 $\Pr(\text{malaria}|\text{fever}) = \alpha \Pr(\text{fever}|\text{malaria}) \Pr(\text{malaria})$

if the relevant further quantity $\Pr(\texttt{fever}|\neg\texttt{malaria})$ is known.

97

Conditional independence

What happens if we have multiple pieces of evidence?

We have seen that to compute

Pr(HD|CP, HBP)

directly might well run into problems.

We could try using Bayes theorem to obtain

 $\Pr(HD|CP, HBP) = \alpha \Pr(CP, HBP|HD) \Pr(HD)$

However while HD is probably manageable, a quantity such as $\Pr(\mathtt{CP},\mathtt{HBP}|\mathtt{HD})$ might well still be problematic especially in more realistic cases.

Bayes theorem

- Sometimes the first possibility is easier, sometimes not.
- Causal knowledge such as

Pr(fever|malaria)

might well be available when diagnostic knowledge such as

Pr(malaria|fever)

is not.

- \bullet Say the incidence of malaria, modelled by $\Pr(\texttt{Malaria}),$ suddenly changes. Bayes theorem tells us what to do.
- The quantity

Pr(fever|malaria)

would not be affected by such a change. *Causal knowledge* can be more robust.

98

Conditional independence

However although in this case we might not be able to exploit independence directly we *can* say that

Pr(CP, HBP|HD) = Pr(CP|HD) Pr(HBP|HD)

which simplifies matters.

Conditional independence:

- $\Pr(A, B|C) = \Pr(A|C) \Pr(B|C).$
- If we know that C is the case then A and B are independent.

Although CP and HBP are *not* independent, they do not directly influence one another *in a patient known to have heart disease*.

This is much nicer!

$$\Pr(\text{HD}|\text{CP},\text{HBP}) = \alpha \Pr(\text{CP}|\text{HD}) \Pr(\text{HBP}|\text{HD}) \Pr(\text{HD})$$

Naive Bayes

Conditional independence is often assumed even when it does not hold.

Naive Bayes:

 $\Pr(A, B_1, B_2, \dots, B_n) = \Pr(A) \prod_{i=1}^n \Pr(B_i | A)$

Also known as Idiot's Bayes.

Despite this, it is often surprisingly effective.

Uncertainty II - Bayesian Networks

Having seen that in principle, if not in practice, the full joint distribution alone can be used to perform any inference of interest, we now examine a practical technique.

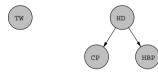
- We introduce the *Bayesian Network (BN)* as a compact representation of the full joint distribution.
- We examine the way in which a BN can be *constructed*.
- We examine the *semantics* of BNs.
- We look briefly at how *inference* can be performed.

Reading: Russell and Norvig, chapter 14.

Bayesian networks

101

Also called probabilistic/belief/causal networks or knowledge maps.



- Each node is a *random variable (RV)*.
- Each node N_i has a distribution

 $\Pr(N_i | \texttt{parents}(N_i))$

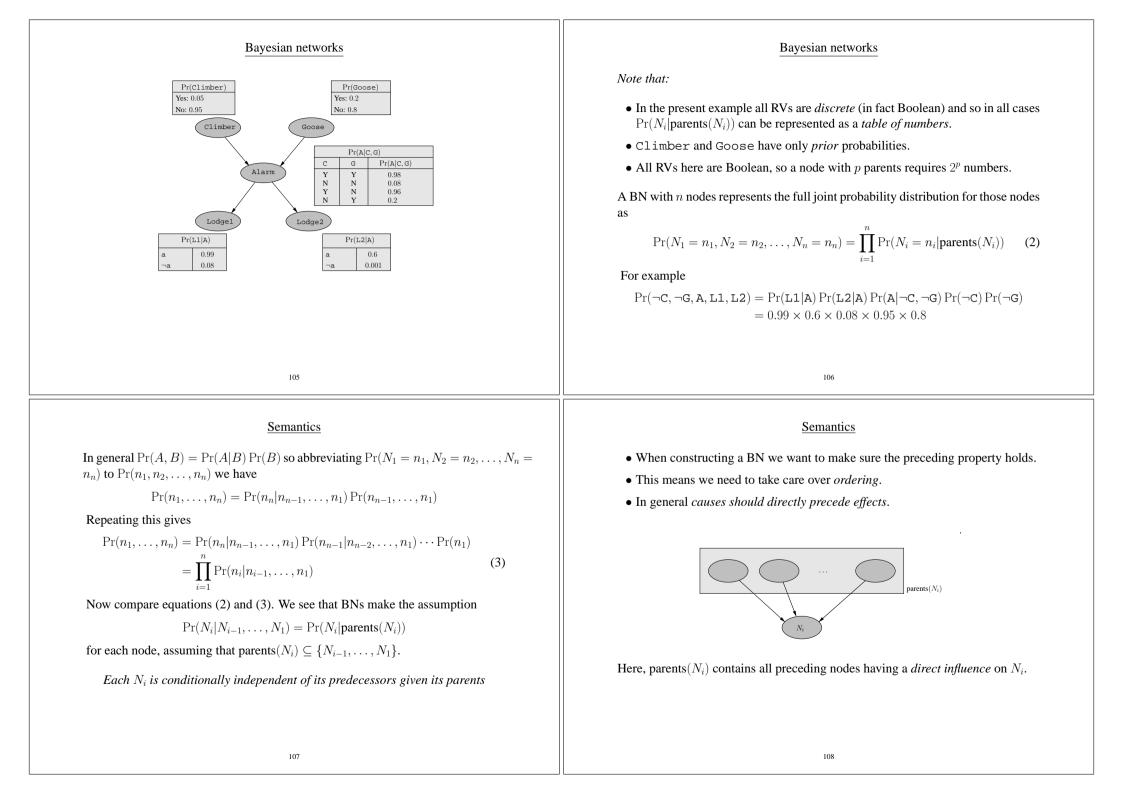
- A Bayesian network is a *directed acyclic graph*.
- Roughly speaking, an arrow from N to M means N directly affects M.

Bayesian networks

102

After a *regrettable incident* involving an *inflatable gorilla*, a famous College has decided to install an alarm for the detection of roof climbers.

- The alarm is very good at detecting climbers.
- Unfortunately, it is also sometimes triggered when one of the *extremely fat geese* that lives in the College lands on the roof.
- One porter's lodge is near the alarm, and inhabited by a chap with *excellent hearing* and a *pathological hatred* of roof climbers: he *always* reports an alarm. His hearing is so good that he sometimes thinks he hears an alarm, *even when there isn't one.*
- Another porter's lodge is a good distance away and inhabited by an *old chap* with *dodgy hearing* who likes to listen to his collection of *DEATH METAL* with the sound turned up.



Semantics

Deviation from this rule can have major effects on the complexity of the network. *That's bad!* We want to keep the network simple:

- If each node has at most p parents and there are n Boolean nodes, we need to specify at most $n2^p$ numbers...
- ...whereas the full joint distribution requires us to specify 2^n numbers.

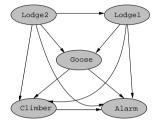
So: there is a trade-off attached to the inclusion of *tenuous* although *strictly-speaking correct* edges.

<u>Semantics</u>

As a rule, we should include the *most basic causes* first, then *the things they influence directly etc.*

What happens if you get this wrong?

Example: add nodes in the order L2,L1,G,C,A.



109

Semantics

In this example:

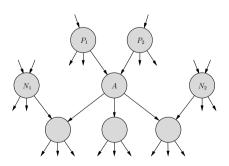
- Increased connectivity.
- Many of the probabilities here will be quite unnatural and hard to specify.

Once again: causal knowledge is preferred to diagnostic knowledge.

Semantics

110

As an alternative we can say directly what conditional independence assumptions a graph should be interpreted as expressing. There are two common ways of doing this.



Any node A is conditionally independent of the N_i —its non-descendants—given the P_i —its parents.

	7
Semantics	More complex nodes
$ \begin{array}{c} & \begin{array}{c} & \begin{array}{c} & \end{array} \\ & \end{array} \\ & \begin{array}{c} & \end{array} \\ & \begin{array}{c} & \end{array} \\ & \end{array} \\ & \begin{array}{c} & \end{array} \\ & \begin{array}{c} & \end{array} \\ & \end{array} \\ & \begin{array}{c} & \end{array} \\ & \begin{array}{c} & \end{array} \\ & \end{array} \\ & \begin{array}{c} & \end{array} \\ & \end{array} \\ & \begin{array}{c} & \end{array} \\ & \end{array} \\ & \begin{array}{c} & \end{array} \\ & \begin{array}{c} & \end{array} \\ & \begin{array}{c} & \end{array} \\ & \end{array} \\ & \begin{array}{c} & \end{array} \\ & \end{array} \\ & \begin{array}{c} & \end{array} \\ & \begin{array}{c} & \end{array} \\ & \end{array} \\ \\ & \end{array} \\ \\ & \end{array} \\ & \begin{array}{c} & \end{array} \\ & \end{array} \\ & \end{array} \\ \\ \\ & \end{array} \\ \\ & \end{array} \\ \\ \\ & \end{array} \\ \\ & \end{array} \\ \\ \\ \\$	 How do we represent Pr(N_i parents(N_i)) when nodes can denote <i>general discrete and/or continuous RVs</i>? BNs containing both kinds of RV are called <i>hybrid BNs</i>. Naive <i>discretisation</i> of continuous RVs tends to result in both a reduction in accuracy and large tables. O(2^p) might still be large enough to be unwieldy. We can instead attempt to use <i>standard and well-understood</i> distributions, such as the <i>Gaussian</i>. This will typically require only a small number of parameters to be specified.
113	114
More complex nodes	More complex nodes
Example: functional relationships are easy to deal with.	<i>Example:</i> a continuous RV with one continuous and one discrete parent.
$N_i = f(parents(N_i))$	Pr(Speed of car Throttle position, Tuned engine)
$\int 1 \text{ if } p = f(\text{parents}(N))$	where SC and TP are continuous and TE is Boolean.
$\Pr(N_i = n_i parents(N_i)) = \begin{cases} 1 & \text{if } n_i = f(parents(N_i)) \\ 0 & \text{otherwise} \end{cases}$	• For a specific setting of ET = true it might be the case that SC increases with TP, but that some uncertainty is involved
	$\Pr(\texttt{SC} \texttt{TP},\texttt{et}) = N(g_{\texttt{et}}\texttt{TP} + c_{\texttt{et}},\sigma_{\texttt{et}}^2)$
	• For an un-tuned engine we might have a similar relationship with a different behaviour $\Pr(\texttt{SC} \texttt{TP}, \neg\texttt{et}) = N(g_{\neg\texttt{et}}\texttt{TP} + c_{\neg\texttt{et}}, \sigma_{\neg\texttt{et}}^2)$
	There is a set of parameters $\{g, c, \sigma\}$ for each possible value of the discrete RV.
115	116

More complex nodes

Example: a discrete RV with a continuous parent

Pr(Go roofclimbing|Size of fine)

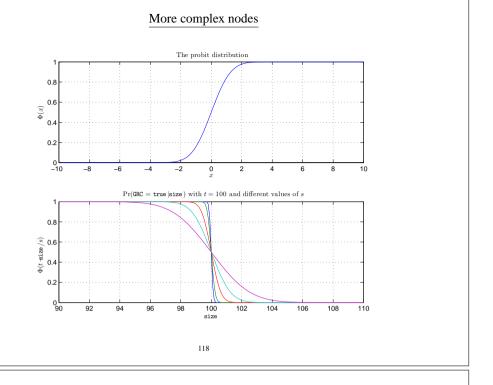
We could for example use the *probit distribution*

$$\Pr(\texttt{Go roofclimbing} = \texttt{true}|\texttt{size}) = \Phi\left(\frac{t - \texttt{size}}{s}\right)$$

where

$$\Phi(x) = \int_{-\infty}^{x} N(y) dy$$

and N(x) is the Gaussian distribution with zero mean and variance 1.



Basic inference

We saw earlier that the full joint distribution can be used to perform *all inference tasks*:

$$\Pr(Q|e) = \frac{1}{Z} \Pr(Q \land e) = \frac{1}{Z} \sum_{u} \Pr(Q, e, u)$$

where

- Q is the query variable
- $\bullet \ e$ is the evidence
- \bullet *u* are the unobserved variables
- 1/Z normalises the distribution.

More complex nodes

117

Alternatively, for this example we could use the logit distribution

$$\Pr(\text{Go roofclimbing} = \text{true}|\text{size}) = \frac{1}{1 + e^{(-2(t-\text{size})/s)}}$$

1

which has a similar shape.

- Tails are longer for the logit distribution.
- The logit distribution tends to be easier to use...
- ...but the probit distribution is often more accurate.

Basic inference

As the BN fully describes the full joint distribution

$$\Pr(Q, u, e) = \prod_{i=1}^{n} \Pr(N_i | \operatorname{parents}(N_i))$$

It can be used to perform inference in the obvious way

$$\Pr(Q|e) = \frac{1}{Z} \sum_{u} \prod_{i=1}^{n} \Pr(N_i | \text{parents}(N_i))$$

but as we'll see this is in practice problematic.

- More sophisticated algorithms aim to achieve this more efficiently.
- For complex BNs we resort to *approximation techniques*.

Other approaches to uncertainty: Default reasoning

One criticism made of probability is that it is *numerical* whereas human argument seems fundamentally different in nature:

- On the one hand this seems quite defensible. I certainly am not aware of doing *logical thought* through direct *manipulation of probabilities*, but...
- ... on the other hand, neither am I aware of *solving differential equations* in order to *walk*!

Default reasoning:

- Does not maintain degrees of belief.
- Allows something to be believed until a reason is found not to.

Other approaches to uncertainty: rule-based systems

122

Problems that can arise:

1. Say I have the causal rule

Heart disease $\xrightarrow{0.95}$ Chest pain

and the diagnostic rule

Chest pain $\xrightarrow{0.7}$ Heart disease

Without taking very great care to keep track of the reasoning process, these can form a *loop*.

2. If in addition I have

Chest pain $\xrightarrow{0.6}$ Recent physical exertion

then it is quite possible to form the conclusion that with some degree of certainty *heart disease is explained by exertion*, which may well be incorrect.

121

Other approaches to uncertainty: rule-based systems

Rule-based systems have some desirable properties:

- *Locality*: if we establish the evidence X and we have a rule $X \to Y$ then Y can be concluded regardless of any other rules.
- *Detachment*: once any Y has been established it can then be assumed. (It's justification is irrelevant.)
- *Truth-functionality*: truth of a complex formula is a function of the truth of its components.

These are not in general shared by probabilistic systems. What happens if:

- We try to attach measures of belief to rules and propositions.
- We try to make a truth-functional system by, for example, making belief in $X \wedge Y$ a function of beliefs in X and Y?

Other approaches to uncertainty: rule-based systems

In addition, we might argue that because heart disease is an explanation for chest pain the belief in physical exertion should *decrease*.

In general when such systems have been successful it has been through very careful control in setting up the rules.

Other approaches to uncertainty: Dempster-Shafer theory

Dempster-Shafer theory attempts to distinguish between *uncertainty* and *igno-rance*.

Whereas the probabilistic approach looks at the *probability* of X, we instead look at the *probability* that the *available evidence supports* X.

This is denoted by the *belief function* Bel(X).

Example: given a coin but no information as to whether it is fair I have no reason to think one outcome should be preferred to another

 $\mathbf{Bel}(\mathbf{outcome} = \mathbf{head}) = \mathbf{Bel}(\mathbf{outcome} = \mathbf{tail}) = 0$

These beliefs can be updated when new evidence is available. If an expert tells us there is n percent certainty that it's a fair coin then

$$Bel(outcome = head) = Bel(outcome = tail) = \frac{n}{100} \times \frac{1}{2}$$

We may still have a gap in that

```
Bel(outcome = head) + Bel(outcome = tail) \neq 1.
```

Dempster-Shafer theory provides a coherent system for dealing with belief functions.

126

Uncertainty III: exact inference in Bayesian networks

We now examine:

- The basic equation for inference in Bayesian networks, the latter being hard to achieve if approached in the obvious way.
- The way in which matters can be improved a little by a small modification to the way in which the calculation is done.
- The way in which much better improvements might be possible using a still more informed approach, although not in all cases.

Reading: Russell and Norvig, chapter 14, section 14.4.

125

Other approaches to uncertainty: Dempster-Shafer theory

Problems:

- The Bayesian approach deals more effectively with the quantification of how *belief changes* when *new evidence is available*.
- The Bayesian approach has a better connection to the concept of *utility*, whereas the latter is not well-understood for use in conjunction with Dempster-Shafer theory.

Performing exact inference

Performing exact inference

We know that in principle any query Q can be answered by the calculation

$$\Pr(Q|e) = \frac{1}{Z} \sum_{u} \Pr(Q, e, u)$$

where Q denotes the query, e denotes the evidence, u denotes unobserved variables and 1/Z normalises the distribution.

The naive implementation of this approach yields the *Enumerate-Joint-Ask* algorithm, which unfortunately requires $O(2^n)$ time and space for n Boolean random variables (RVs).

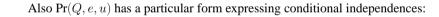
In what follows we will make use of some abbreviations.

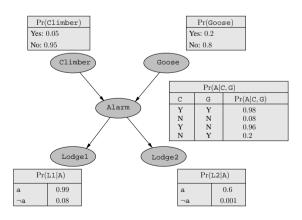
- $\bullet \ C$ denotes Climber
- $\bullet\;G$ denotes <code>Goose</code>
- A denotes Alarm
- L1 denotes Lodge1
- L2 denotes Lodge2

Instead of writing out $\Pr(C = \top)$, $\Pr(C = \bot)$ *etc* we will write $\Pr(c)$, $\Pr(\neg c)$ and so on.

Performing exact inference

129





$\Pr(C,G,A,L1,L2) = \Pr(C)\Pr(G)\Pr(A|C,G)\Pr(L1|A)\Pr(L2|A)$

Performing exact inference

130

Consider the computation of the query $\Pr(C|l1,l2)$

We have

$$\Pr(C|l1, l2) = \frac{1}{Z} \sum_{A} \sum_{G} \Pr(C) \Pr(G) \Pr(A|C, G) \Pr(l1|A) \Pr(l2|A)$$

Here there are 5 multiplications for each set of values that appears for summation, and there are 4 such values.

In general this gives time complexity $O(n2^n)$ for n Boolean RVs.

Looking more closely we see that

$$\Pr(C|l1, l2) = \frac{1}{Z} \sum_{A} \sum_{G} \Pr(C) \Pr(G) \Pr(A|C, G) \Pr(l1|A) \Pr(l2|A)$$

$$= \frac{1}{Z} \Pr(C) \sum_{A} \Pr(l1|A) \Pr(l2|A) \sum_{G} \Pr(G) \Pr(A|C, G)$$

$$= \frac{1}{Z} \Pr(C) \sum_{G} \Pr(G) \sum_{A} \Pr(A|C, G) \Pr(l1|A) \Pr(l2|A)$$
(4)

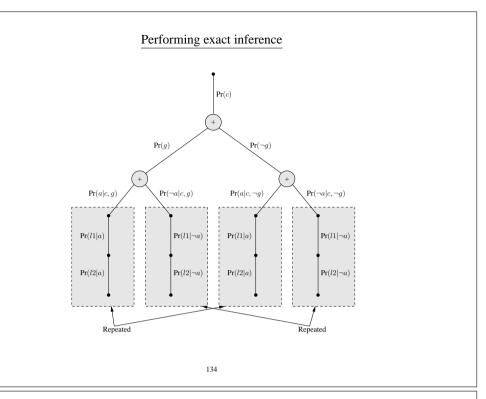
So for example...

Performing exact inference

$$\begin{aligned} \mathbf{Pr}(c|l1,l2) = & \frac{1}{Z} \mathbf{Pr}(c) \left(\mathbf{Pr}(g) \left\{ \begin{array}{c} \mathbf{Pr}(a|c,g) \mathbf{Pr}(l1|a) \mathbf{Pr}(l2|a) \\ + \mathbf{Pr}(\neg a|c,g) \mathbf{Pr}(l1|\neg a) \mathbf{Pr}(l2|\neg a) \end{array} \right\} \\ & + \mathbf{Pr}(\neg g) \left\{ \begin{array}{c} \mathbf{Pr}(a|c,\neg g) \mathbf{Pr}(l1|a) \mathbf{Pr}(l2|a) \\ + \mathbf{Pr}(\neg a|c,\neg g) \mathbf{Pr}(l1|\neg a) \mathbf{Pr}(l2|\neg a) \end{array} \right\} \end{aligned} \end{aligned}$$

with a similar calculation for $Pr(\neg c|l1, l2)$.

Basically straightforward, *BUT* optimisations can be made.



Optimisation 1: Enumeration-Ask

133

The *enumeration-ask* algorithm improves matters to $O(2^n)$ time and O(n) space by performing the computation *depth-first*.

However matters can be improved further by avoiding the *duplication of computations* that clearly appears in the example tree.

Optimisation 2: variable elimination

Looking again at the fundamental equation (4)

$$\frac{1}{Z} \underbrace{\Pr(C)}_{C} \sum_{G} \underbrace{\Pr(G)}_{G} \sum_{A} \underbrace{\Pr(A|C,G)\Pr(l1|A)\Pr(l2|A)}_{A} \underbrace{\Pr(A|C,G)\Pr(l1|A)\Pr(l2|A)}_{L1}$$

where C, G, A, L1, L2 denote the relevant factors.

The basic idea is to evaluate (4) from right to left (or in terms of the tree, bottom up) *storing results* as we progress and *re-using them* when necessary.

Pr(l1|A) depends on the value of A. We store it as a table $F_{L1}(A)$. Similarly for Pr(l2|A).

$$\mathbf{F}_{L1}(A) = \begin{pmatrix} 0.99\\ 0.08 \end{pmatrix} \mathbf{F}_{L2}(A) = \begin{pmatrix} 0.6\\ 0.001 \end{pmatrix}$$

as $\Pr(l1|a) = 0.99$, $\Pr(l1|\neg a) = 0.08$ and so on.

$\underline{Optimisation \ 2: variable \ elimination}$ Similarly for $\Pr(A C, G)$, which is dependent on A, C and G $\mathbf{F}_{A}(A, C, G) = \begin{bmatrix} \overline{A} & \overline{C} & \overline{G} & \overline{\mathbf{F}}_{A}(A, C, G) \\ \overline{T} & \overline{T} & \overline{T} & 0.98 \\ \overline{T} & \overline{T} & 1 & 0.96 \\ \overline{T} & \pm & \overline{T} & 0.2 \\ \overline{T} & \pm & \pm & 0.08 \\ \pm & \overline{T} & \overline{T} & 0.02 \\ \pm & \overline{T} & \pm & 0.04 \\ \pm & \pm & \overline{T} & 0.8 \\ \pm & \pm & \pm & 0.92 \end{bmatrix}$ Can we write $\Pr(A C, G)\Pr(l1 A)\Pr(l2 A) \qquad (5)$ as $\mathbf{F}_{A}(A, C, G)\mathbf{F}_{L1}(A)\mathbf{F}_{L2}(A) \qquad (6)$ in a reasonable way?	$\underbrace{\text{Optimisation 2: variable elimination}}_{\text{The answer is "yes" provided multiplication of factors is defined correctly. Looking at (4) \frac{1}{Z} \Pr(C) \sum_{G} \Pr(G) \sum_{A} \Pr(A C, G) \Pr(l1 A) \Pr(l2 A)note that the values of the product (5) in the summation depend on the values of C and G external to it, and the values of A themselves. So (6) should be a table collecting values for (5) where correspondences between RVs are maintained. This leads to a definition for multiplication of factors best given by example.$
¹³⁷ Optimisation 2: variable elimination	138 Optimisation 2: variable elimination
$\mathbf{F}(A,B)\mathbf{F}(B,C) = \mathbf{F}(A,B,C)$ where $\frac{A \ B \ \mathbf{F}(A,B) \ B \ C \ \mathbf{F}(B,C) \ A \ B \ C \ \mathbf{F}(A,B,C)}{\top \ T \ 0.3 \ \top \ T \ 0.1 \ T \ T \ T \ 0.3 \times 0.1}$ $\frac{T \ \bot \ 0.9 \ T \ \bot \ 0.8 \ T \ T \ \bot \ 0.3 \times 0.8}{\bot \ T \ 0.4 \ \bot \ T \ 0.8 \ T \ \bot \ T \ 0.9 \times 0.8}$ $\frac{L \ T \ 0.4 \ L \ T \ 0.8 \ T \ \bot \ 1 \ 0.9 \times 0.8}{\bot \ T \ T \ 0.4 \times 0.1}$ $\frac{L \ T \ 1 \ 0.4 \times 0.1}{\bot \ T \ 0.1 \times 0.8}$	This process gives us $\mathbf{F}_{A}(A, C, G)\mathbf{F}_{L1}(A)\mathbf{F}_{L2}(A) = \begin{bmatrix} A & C & G \\ \hline \top & \top & \top & 0.98 \times 0.99 \times 0.6 \\ \hline \top & \top & \downarrow & 0.96 \times 0.99 \times 0.6 \\ \hline \top & \bot & \top & 0.2 \times 0.99 \times 0.6 \\ \hline \top & \bot & \bot & 0.08 \times 0.99 \times 0.6 \\ \bot & \top & \top & 0.02 \times 0.08 \times 0.001 \\ \bot & \top & \bot & 0.04 \times 0.08 \times 0.001 \\ \bot & \bot & \downarrow & 0.92 \times 0.08 \times 0.001 \\ \downarrow & \bot & \bot & 0.92 \times 0.08 \times 0.001 \\ \end{bmatrix}$
139	140

Optimisation 2: variable elimination

How about

$$\mathbf{F}_{\overline{A},L1,L2}(C,G) = \sum_{A} \mathbf{F}_{A}(A,C,G) \mathbf{F}_{L1}(A) \mathbf{F}_{L2}(A)$$

To denote the fact that A has been summed out we place a bar over it in the notation.

$$\sum_{A} \mathbf{F}_{A}(A, C, G) \mathbf{F}_{L1}(A) \mathbf{F}_{L2}(A) = \mathbf{F}_{A}(a, C, G) \mathbf{F}_{L1}(a) \mathbf{F}_{L2}(a) + \mathbf{F}_{A}(\neg a, C, G) \mathbf{F}_{L1}(\neg a) \mathbf{F}_{L2}(\neg a)$$

where

$$\mathbf{F}_{A}(a,C,G) = \frac{\begin{vmatrix} C & G \\ \hline \top & \top & 0.98 \\ \hline \top & \bot & 0.96 \\ \bot & \top & 0.2 \\ \bot & \bot & 0.08 \end{vmatrix}} \mathbf{F}_{L1}(a) = 0.99 \ \mathbf{F}_{L2}(a) = 0.6$$

and similarly for
$$\mathbf{F}_A(\neg a, C, G)$$
, $\mathbf{F}_{L1}(\neg a)$ and $\mathbf{F}_{L2}(\neg a)$

141

Optimisation 2: variable elimination

Now, say for example we have $\neg c$, g. Then doing the calculation explicitly would give

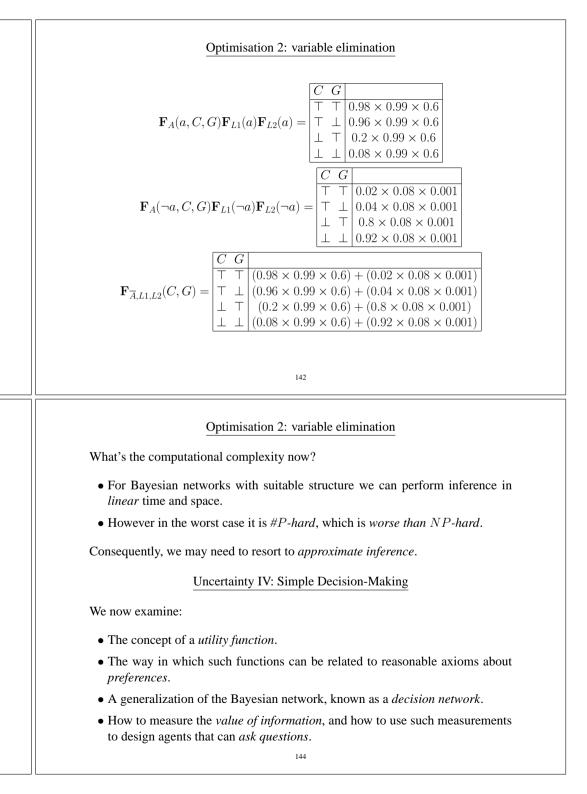
$$\sum_{A} \Pr(A|\neg c, g) \Pr(l1|A)) \Pr(l2|A) = \Pr(a|\neg c, g) \Pr(l1|a) \Pr(l2|a) + \Pr(\neg a|\neg c, g) \Pr(l1|\neg a) \Pr(l2|\neg a) = (0.2 \times 0.99 \times 0.6) + (0.8 \times 0.08 \times 0.001)$$

which matches!

Continuing in this manner form

$$\begin{split} \mathbf{F}_{G,\overline{A},L1,L2}(C,G) &= \mathbf{F}_G(G)\mathbf{F}_{\overline{A},L1,L2}(C,G)\\ \text{sum out }G \text{ to obtain } \mathbf{F}_{\overline{G},\overline{A},L1,L2}(C) &= \sum_G \mathbf{F}_G(G)\mathbf{F}_{\overline{A},L1,L2}(C,G), \text{ form}\\ \mathbf{F}_{C,\overline{G},\overline{A},L1,L2} &= \mathbf{F}_C(C)\mathbf{F}_{\overline{G},\overline{A},L1,L2}(C) \end{split}$$

and normalise.



Reading: Russell and Norvig, chapter 16.		Simple decision-making
		We now look at choosing an action by maximising <i>expected utility</i> .
		A utility function $U(s)$ measures the desirability of a state.
		If we can express a probability distribution for the states resulting from alternative actions, then we can act in order to maximise expected utility.
		For an action a , let $\text{Result}(a) = \{s_1, \ldots, s_n\}$ be a set of states that might be the result of performing action a . Then the expected utility of a is
		$\mathtt{EU}(a E) = \sum_{s \in \mathtt{Result}(a)} \Pr(s a, E) U(s)$
		Note that this applies to <i>individual actions</i> . Sequences of actions will not be covered in this course.
	145	146
	Simple decision-making: all of AI?	Utility in more detail
Much as this looks like a complete and highly attractive method	as this looks like a complete and highly attractive method for an agent to	Overall, we now want to express <i>preferences</i> between different things.
decide	how to act, it hides a great deal of complexity:	Let's use the following notation:
	hay be hard to compute $U(s)$. You generally don't know how good a state antil you know where it might lead on to: planning etc	X > Y : X is preferred to $YX = Y$: we are indifferent regarding X and Y
2. Kn	owing what state you're currently in involves most of AI!	$X \ge Y : X$ is preferred, or we're indifferent
3. Dealing with $Pr(s a, E)$ involves <i>Bayesian networks</i> .		X, Y and so on are <i>lotteries</i> . A lottery has the form
		$X = [p_1, O_1 p_2, O_2 \cdots p_n, O_n]$
		where O_i are the outcomes of the lottery and p_i their respective probabilities. Outcomes can be <i>other lotteries</i> or actual states.

Axioms for utility theory

Given we are dealing with preferences it seems that there are some clear properties that such things should exhibit:

Transitivity: if X > Y and Y > Z then X > Z.

Orderability: either X > Y or Y > X or X = Y.

Continuity: if X > Y > Z then there is a probability p such that

[p, X|(1-p), Z] = Y

Substitutability: if X = Y then

$$[p, X|(1-p), L] = [p, Y|(1-p), L]$$

Axioms for utility theory

Monotonicity: if X > Y then for probabilities p_1 and p_2 , $p_1 \ge p_2$ if and only if

$$[p_1, X|(1-p_1), Y] \ge [p_2, X|(1-p_2), Y]$$

Decomposability:

 $[p_1, X|(1-p_1), [p_2, Y|(1-p_2), Z]] = [p_1, X|(1-p_1)p_2, Y|(1-p_1)(1-p_2), Z]$

If an agent's preferences conform to the utility theory axioms—and note that we are *only* considering preferences, not numbers—then it is possible to define a utility function U(s) for states such that:

1. $U(s_1) > U(s_2) \longleftrightarrow s_1 > s_2$ 2. $U(s_1) = U(s_2) \longleftrightarrow s_1 = s_2$ 3. $U([p_1, s_1 | p_2, s_2 | \cdots | p_n, s_n]) = \sum_{i=1}^n p_i U(s_i).$

We therefore have a justification for the suggested approach.

Designing utility functions

150

The expected monetary value (EMV) of this lottery is

 $(0.5 \times 1,000,000) + (0.5 \times 0) = 500,000$

whereas the EMV of the initial amount is 100,000.

BUT: most of us would probably refuse to take the coin flip.

The story is not quite as simple as this though: our attitude probably depends on how much money we have to start with. If I have M pounds to start with then I am in fact choosing between expected utility of

U(M + 100, 000)

and expected utility of

 $(0.5 \times U(M)) + (0.5 \times U(M+1,000,000))$

If M is 50,000,000 my attitude is much different to if it is 10,000.

Designing utility functions

149

There is complete freedom in how a utility function is defined, but clearly it will pay to define them carefully.

Example: the utility of money (for most people) exhibits a *monotonic preference*. That is, we *prefer* to have *more of it*.

But we need to talk about preferences between lotteries.

Say you've won 100,000 pounds in a quiz and you're offered a coin flip:

• For heads: you win a total of 1,000,000 pounds.

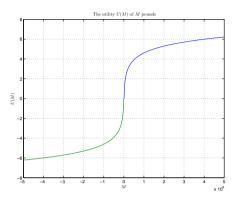
• For tails: you walk away with nothing!

Would you take the offer?

Designing utility functions

Decision networks

In fact, research shows that the utility of M pounds is for most people almost exactly proportional to $\log M$ for M > 0...



^{...} and follows a similar shape for M < 0.

153

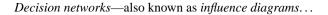
Decision networks

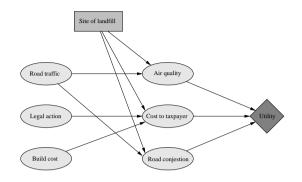
A decision network has three types of node:

Chance nodes: are denoted by ovals. These are random variables (RVs) represented by a distribution conditional on their parents, as in Bayesian networks. Parents can be other chance nodes or a decision node.

Decision nodes: are denoted by squares. They describe possible outcomes of the decision of interest. Here we deal only with *single* decisions: multiple decisions require alternative techniques.

Utility nodes: are denoted by diamonds. They describe the utility function relevant to the problem, as a function of the values of the node's parents.



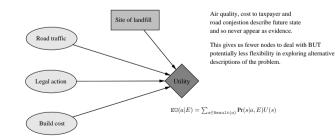


... allow us to work *actions* and *utilities* into the formalism of *Bayesian networks*. A decision network has three types of node...

154

Decision networks

Sometimes such diagrams are simplified by leaving out the RVs describing the new state and converting current state and decision directly to utility:



This is an *action-utility table*. The utility no longer depends on a state but is the expected utility for a given action.

Evaluation of decision networks The value of information Once a *specific* action is selected for a decision node it acts like a chance node for We have been assuming that a decision is to be made with *all evidence available* which a specific value is being used as evidence. *beforehand*. This is unlikely to be the case. Knowing what questions one should ask is a central, and important part of making 1. Set the current state chance nodes to their evidence values. decisions. *Example*: 2. For each potential action • Doctors do not diagnose by first obtaining results for all possible tests on their • Fix the decision node. patients. • Compute the probabilities for the utility node's parents. • They ask questions to decide what tests to do. • Compute the expected utility. • They are informed in formulating which tests to perform by probabilities of 3. Return the action that maximised EU(a|E). test outcomes, and by the manner in which knowing an outcome might improve treatment. • Tests can have associated costs. 157 158 The value of perfect information The value of perfect information Information value theory provides a formal way in which we can reason about This leads to the definition of the *value of perfect information* (VPI) what further information to gather using sensing actions. $\mathrm{VPI}_E(E') = \left\{ \sum \mathrm{Pr}(E' = e' | E) \mathrm{EU}(\mathrm{action}' | E, E' = e') \right\} - \mathrm{EU}(\mathrm{action} | E)$ Say we have evidence E, so $\texttt{EU}(\texttt{action}|E) = \max_{a} \sum_{s \in \texttt{Result}(a)} \Pr(s|a, E) U(s)$ VPI has the following properties: • $\operatorname{VPI}_E(E') \ge 0$ denotes how valuable the best action based on E must be. • It is not necessarily additive, that is, it is possible that How valuable would it be to learn about a further piece of evidence? $\operatorname{VPI}_E(E', E'') \neq \operatorname{VPI}_E(E') + \operatorname{VPI}_E(E'')$ If we examined another RV E' and found that E' = e' then the best action might

• It is independent of ordering

$$\begin{split} \mathtt{VPl}_E(E',E'') &= \mathtt{VPl}_E(E') + \mathtt{VPl}_{E,E'}(E'') \\ &= \mathtt{VPl}_E(E'') + \mathtt{VPl}_{E,E''}(E') \end{split}$$

 $\mathrm{EU}(\mathrm{action}'|E,E') = \max_{a} \sum_{s \in \mathrm{Result}(a)} \mathrm{Pr}(s|a,E,E') U(s)$

BUT: because E' is a RV, and in advance of testing we don't know its value, we

need to average over its possible values using our current knowledge.

be altered as we'd be computing

Agents that can gather information

In constructing an agent with the ability to ask questions, we would hope that it would:

- Use a good order in which to ask the questions.
- Avoid asking irrelevant questions.
- Trade off the *cost* of obtaining information against the *value* of that information.
- Choose a good time to *stop* asking questions.

We now have the means with which to approach such a design.

161

- The Viterbi algorithm for computing the most likely sequence.
- A slightly simplified system within this framework called a *hidden Markov model* (HMM), and the way in which some inference tasks can be simplified in the HMM case.

Reading: Russell and Norvig, chapter 15.

Agents that can gather information

Assuming we can associate a cost C(E') with obtaining the knowledge that E' = e' an agent can act as follows:

- Given a decision network and current percept.
- Find the piece of evidence E' maximising $VPI_E(E') C(E')$.
- If $VPI_E(E') C(E')$ is positive then find the value of E', else take the action indicated by the decision network.

This is known as a *myopic* agent as it requests a single piece of evidence at once.

Uncertainty V: probabilistic reasoning through time

We now examine:

- How an agent might operate by keeping track of the state of its environment in an uncertain world, and how alterations in world state and uncertainty in observing the world can be modelled using probability distributions.
- How inferences can be performed regarding the current state, past state and future states.

162

Probabilistic reasoning through time

A fundamental idea throughout the AI courses has been that an agent should keep track of the *state of the environment*:

- The environment's state *changes over time*.
- The knowledge of *how the state changes* may be *uncertain*.
- The agent's *perception* of the state of the environment *may be uncertain*.

For all the usual reasons related to *uncertainty*, we need to move beyond logic, situation calculus *etc*.

States and evidence

We model the (unobservable) state of the environment as follows:

• We use a *sequence*

 (S_0, S_1, S_2, \ldots)

of sets of random variables (RVs).

• Each S_t is a *set* of RVs

 $S_t = \{S_t^{(1)}, \dots, S_t^{(n)}\}$

denoting the state of the environment at time t, where t = 0, 1, 2, ...

Think of the state as changing over time.

$$S_0 \to S_1 \to S_2 \to \cdots$$

States and evidence

At each time t there is also an *observable* set

$$E_t = \{E_t^{(1)}, \dots, E_t^{(m)}\}$$

of random variables denoting the *evidence that an agent obtains about the state* at time t.

As usual capitals denote RVs and lower case denotes actual values. So actual values for the assorted RVs are denoted

 $S_t = \{s_t^{(1)}, \dots, s_t^{(n)}\} = s_t$ $E_t = \{e_t^{(1)}, \dots, e_t^{(m)}\} = e_t$

166

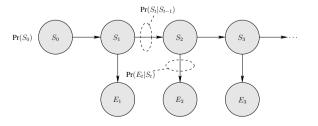
Stationary and Markov processes

Assumption 3: We assume that evidence only depends on the current state

$$\Pr(E_t|S_{0:t}, E_{1:t-1}) = \Pr(E_t|S_t)$$
(8)

Then

 $\Pr(E_t|S_t)$ is called the *sensor model*.



 $Pr(S_0)$ is the *prior probability* of the starting state. We need this as there has to be some way of getting the process started.

As $t\ {\rm can}$ in principle increase without bound we now need some simplifying assumptions.

Stationary and Markov processes

165

Assumption 1: We deal with *stationary processes*: probability distributions do not change over time.

Assumption 2: We deal with Markov processes

$$\Pr(S_t|S_{0:t-1}) = \Pr(S_t|S_{t-1})$$
(7)

where $S_{0:t-1} = (S_0, S_1, \dots, S_{t-1})$.

(Strictly speaking this is a *first order Markov Process*, and we'll only consider these.)

 $Pr(S_t|S_{t-1})$ is called the *transition model*.

The full joint distribution

Given:

- 1. The prior $Pr(S_0)$.
- 2. The transition model $Pr(S_t|S_{t-1})$.
- 3. The sensor model $Pr(E_t|S_t)$.

 $\Pr(S_0, S_1, \dots, S_t, E_1, E_2, \dots, E_t) = \Pr(S_0) \prod_{i=1}^t \Pr(S_i | S_{i-1}) \Pr(E_i | S_i)$

This follows from basic probability theory as for example

 $\begin{aligned} \Pr(S_0, S_1, S_2, E_1, E_2) &= \Pr(E_2 | S_{0:2}, E_1) \Pr(S_2 | S_{0:1}, E_1) \Pr(E_1 | S_{0:1}) \Pr(S_1 | S_0) \Pr(S_0) \\ &= \Pr(E_2 | S_2) \Pr(S_2 | S_1) \Pr(E_1 | S_1) \Pr(S_1 | S_0) \Pr(S_0) \end{aligned}$

169

Example: two biased coins

We'll use the following numbers:

- The prior $\Pr(S_0 = \texttt{coin1}) = 0.5$.
- The transition model

 $\Pr(S_t = \texttt{coin1}|S_{t-1} = \texttt{coin1}) = \Pr(S_t = \texttt{coin2}|S_{t-1} = \texttt{coin2}) = 0.8$ $\Pr(S_t = \texttt{coin1}|S_{t-1} = \texttt{coin2}) = \Pr(S_t = \texttt{coin2}|S_{t-1} = \texttt{coin1}) = 0.2$

• The sensor model

 $\Pr(E_t = \texttt{head}|S_t = \texttt{coin1}) = 0.1$ $\Pr(E_t = \texttt{head}|S_t = \texttt{coin2}) = 0.9$

Example: two biased coins

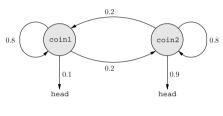
Here's a simple example with only two states and two observations.

I have two biased coins.

I flip one and tell you the outcome.

I then either *stay* with the same coin, or *swap* them.

This continues, producing a succession of outcomes:



Example: two biased coins

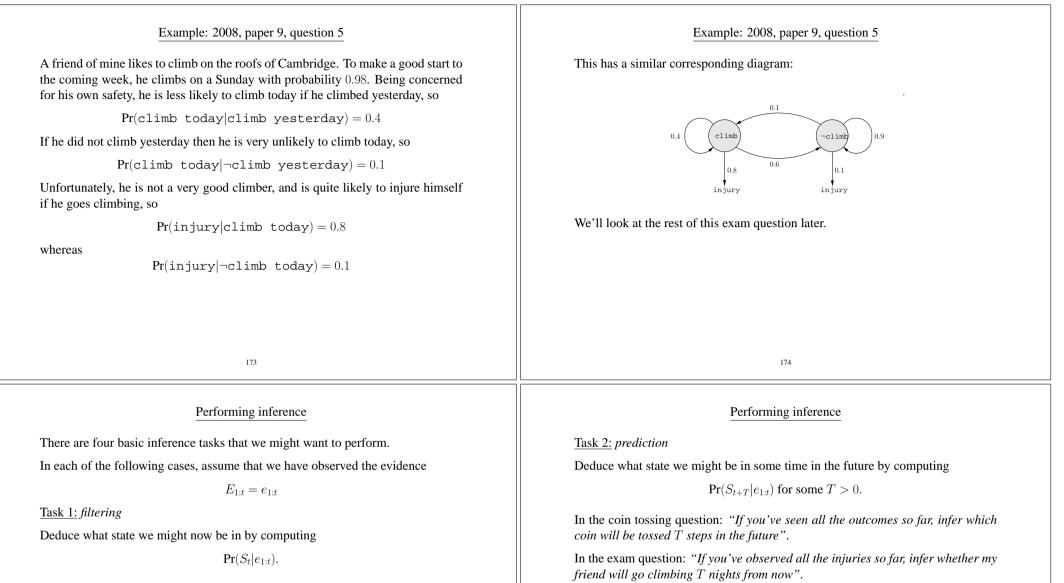
As expected, we tend to see runs of a single coin, and might expect to be able to

guess which is being used as one favours heads and the other tails.

170

This is straightforward to simulate.

Here's an example of what happens:



In the coin tossing question: "If you've seen all the outcomes so far, infer which coin was used last".

In the exam question: "If you observed all the injuries so far, infer whether my friend climbed today".

175

Performing inference

Task 3: Smoothing

Deduce what state we might have been in at some point in the past by computing

 $\Pr(S_t | e_{1:T}) \text{ for } 0 \le t < T.$

In the coin tossing question: "If you've seen all the outcomes so far, infer which coin was tossed at time t in the past".

In the exam question: "If you've observed all the injuries so far, infer whether my friend climbed on night t in the past".

Performing inference

Task 4: Find the most likely explanation

Deduce the most likely sequence of states so far by computing

 $\operatorname*{argmax}_{s_{1:t}} \Pr(s_{1:t}|e_{1:t})$

In the coin tossing question: "If you've seen all the outcomes so far, infer the most probable sequence of coins used".

In the exam question: "If you've observed all the injuries so far, infer the most probable collection of nights on which my friend climbed".

177

Filtering

We want to compute $\Pr(S_t|e_{1:t}).$ This is often called the forward message and denoted

 $f_{1:t} = \Pr(S_t | e_{1:t})$

for reasons that are about to become clear.

Remember that S_t is an RV and so $f_{1:t}$ is a *probability distribution* containing a probability for each possible value of S_t .

It turns out that this can be done in a simple manner with a *recursive estimation*. Obtain the result at time t + 1:

1. using the result from time t and...

2. ...incorporating new evidence e_{t+1} .

 $f_{1:t+1} = g(e_{t+1}, f_{1:t})$

for a suitable function g that we'll now derive.

Filtering

178

Step 1:

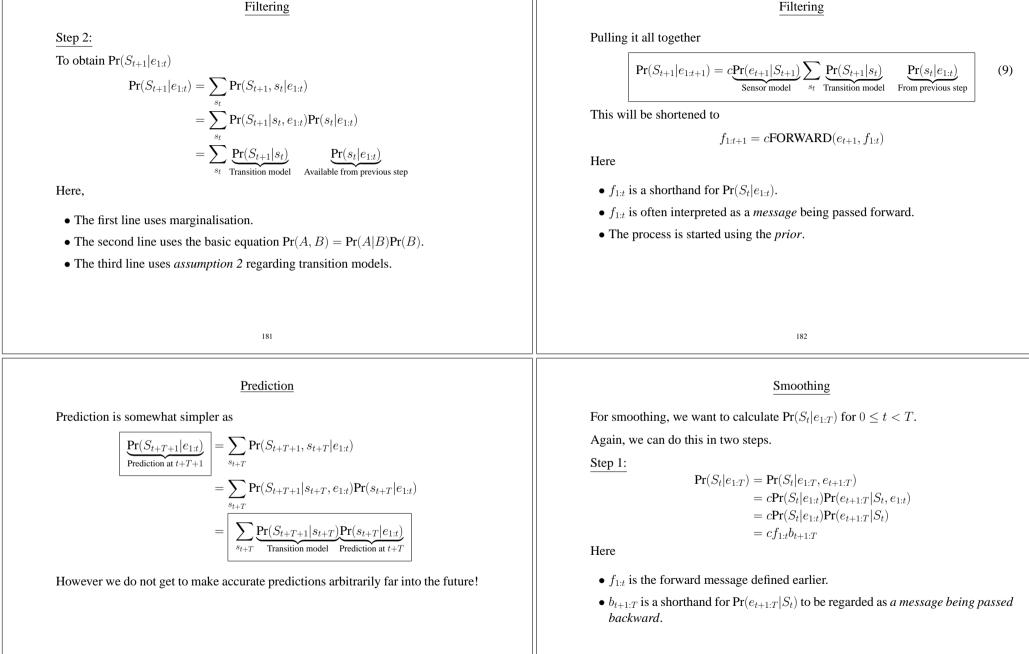
Project the current state distribution forward

$$Pr(S_{t+1}|e_{1:t+1}) = Pr(S_{t+1}|e_{1:t}, e_{t+1}) = cPr(e_{t+1}|S_{t+1}, e_{1:t})Pr(S_{t+1}|e_{1:t}) = cPr(e_{t+1}|S_{t+1})Pr(S_{t+1}|e_{1:t}) Sensor model Needs more work$$

where as usual c is a constant that normalises the distribution. Here,

- The first line does nothing but split $e_{1:t+1}$ into e_{t+1} and $e_{1:t}$.
- The second line is an application of Bayes' theorem.
- The third line uses *assumption 3* regarding sensor models.

Filtering



$$\frac{\text{Supping}}{\text{Supping}} = \frac{P(\zeta_{1,1},1;\beta_{1,1}^{2}) - \sum_{\substack{k \in 1 \\ k \in 1$$

	7
Computing the most likely sequence: the Viterbi algorithm	Computing the most likely sequence: the Viterbi algorithm
We therefore have $ \begin{aligned} & \qquad \qquad$	 There is however a way to visualise it that leads to a dynamic programming algorithm called the <i>Viterbi algorithm</i>. Step 1: Simplify the notation. Assume there are n states s₁,, s_n and m possible observations e₁,, e_m at any given time. Denote Pr(S_t = s_j S_{t-1} = s_i) by p_{i,j}(t). Denote Pr(e_t S_t = s_i) by q_i(t). It's important to remember in what follows that the <i>observations are known</i> but that we're <i>maximising over all possible state sequences</i>.
189	190
Computing the most likely sequence: the Viterbi algorithm	Computing the most likely sequence: the Viterbi algorithm
The equation we're interested in is now of the form	Step 2: Make a grid: columns denote time and rows denote state.
$P = \prod_{t=1}^{T} p_{i,j}(t)q_i(t)$ (The prior $\Pr(S_0)$ has been dropped out for the sake of clarity, but is easy to put back in in what follows.)	$1 2 3 \cdots k k+1 \cdots t$ $s_1 \bullet \bullet \bullet \bullet$ $s_2 \bullet \bullet \bullet \bullet$
The equation P will be referred to in what follows.	
It is in fact a <i>function</i> of <i>any given sequence of states</i> .	
	s_{n-1} \bullet \bullet \bullet \bullet \bullet \bullet \bullet \bullet \bullet
191	192

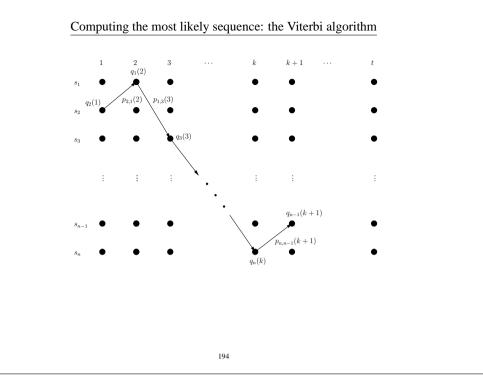
Computing the most likely sequence: the Viterbi algorithm

Step 3: Label the nodes:

- Say at time t the actual observation was e_t . Then label the node for s_i in column t with the value $q_i(t)$.
- Any sequence of states through time is now a path through the grid. So for any transition from s_i at time t 1 to s_j at time t label the transition with the value $p_{i,j}(t)$.

In the following diagrams we can often just write $p_{i,j}$ and q_i because the time is clear from the diagram.

So for instance...



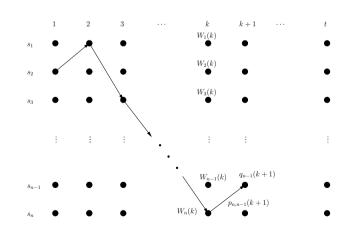
Computing the most likely sequence: the Viterbi algorithm

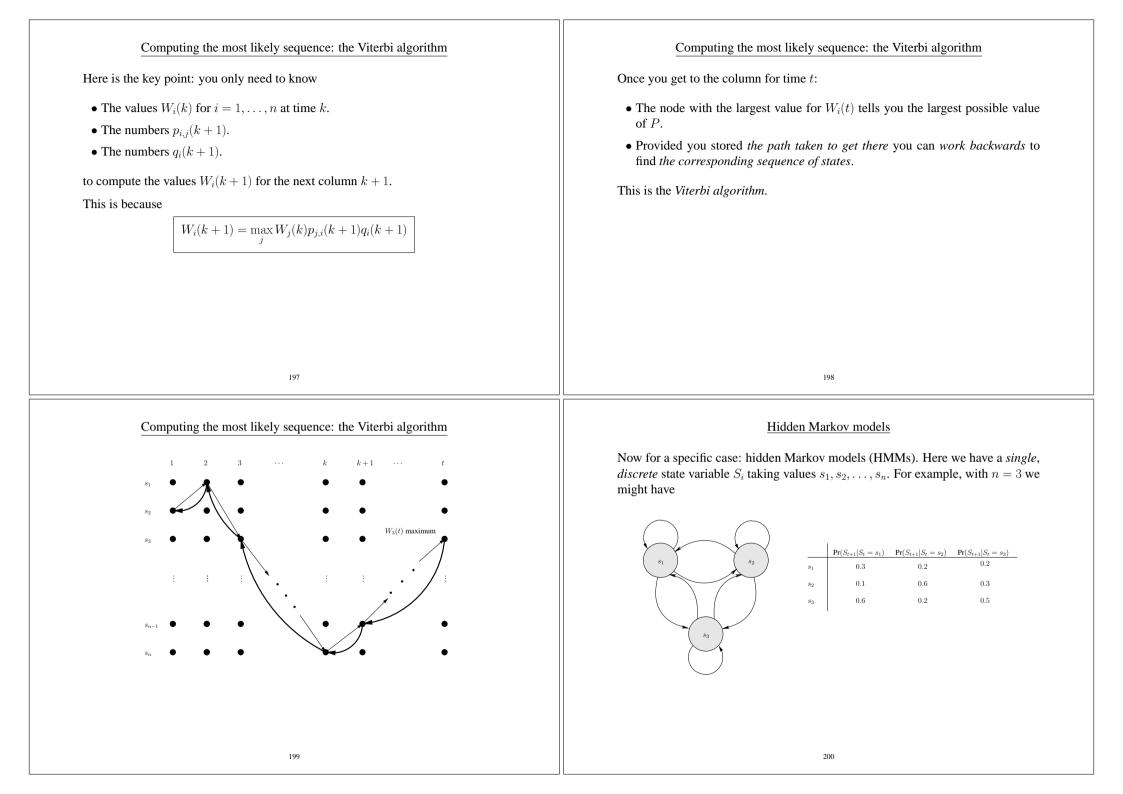
193

- The value of $P = \prod_{t=1}^{T} p_{i,j}(t)q_i(t)$ for any path through the grid is just the product of the corresponding labels that have been added.
- But we don't want to find the maximum by looking at all the possible paths because this would be time-consuming.
- The *Viterbi algorithm* computes the maximum by moving from one column to the next updating as it goes.
- Say you're at column k and for each node m in that column you know the highest value for the product to this point over any possible path. Call this:

$$W_m(k) = \max_{s_{1:k}} \prod_{t=1}^k p_{i,j}(t) q_i(t)$$

Computing the most likely sequence: the Viterbi algorithm





Hidden Markov models

In this simplified case the conditional probabilities $\Pr(S_{t+1}|S_t)$ can be represented using the matrix

$$S_{ij} = \Pr(S_{t+1} = s_j | S_t = s_i)$$

or for the example on the previous slide

$$\mathbf{S} = \begin{pmatrix} 0.3 & 0.1 & 0.6 \\ 0.2 & 0.6 & 0.2 \\ 0.2 & 0.3 & 0.5 \end{pmatrix} \xleftarrow{} \mathbf{Pr}(S|s_1) \\ \xleftarrow{} \mathbf{Pr}(S|s_2) \\ \xleftarrow{} \mathbf{Pr}(S|s_3) \\ = \begin{pmatrix} \mathbf{Pr}(s_1|s_1) & \mathbf{Pr}(s_2|s_1) & \cdots & \mathbf{Pr}(s_n|s_1) \\ \mathbf{Pr}(s_1|s_2) & \mathbf{Pr}(s_2|s_2) & \cdots & \mathbf{Pr}(s_n|s_2) \\ \vdots & \vdots & \ddots & \vdots \\ \mathbf{Pr}(s_1|s_n) & \mathbf{Pr}(s_2|s_n) & \cdots & \mathbf{Pr}(s_n|s_n) \end{pmatrix}$$

To save space, I am abbreviating $Pr(S_{t+1} = s_i | S_t = s_j)$ to $Pr(s_i | s_j)$.

Hidden Markov models

The computations we're making are always conditional on some actual observations $e_{1:T}$.

For each t we can therefore use the sensor model to define a further matrix \mathbf{E}_t :

- \mathbf{E}_t is square and diagonal (all off-diagonal elements are 0).
- The *i*th element of the diagonal is $Pr(e_t|S_t = s_i)$.

So in our present example with 3 states, there will be a matrix

$$\mathbf{E}_t = \begin{pmatrix} \Pr(e_t|s_1) & 0 & 0\\ 0 & \Pr(e_t|s_2) & 0\\ 0 & 0 & \Pr(e_t|s_3) \end{pmatrix}$$
for each $t = 1, \dots, T$.

202

Hidden Markov models

201

In the general case the equation for filtering was

$$\Pr(S_{t+1}|e_{1:t+1}) = c\Pr(e_{t+1}|S_{t+1}) \sum_{s_t} \Pr(S_{t+1}|s_t) \Pr(s_t|e_{1:t})$$

and the message $f_{1:t}$ was introduced as a representation of $Pr(S_t|e_{1:t})$.

In the present case we can define $f_{1:t}$ to be the vector

$$f_{1:t} = \begin{pmatrix} \Pr(s_1|e_{1:t}) \\ \Pr(s_2|e_{1:t}) \\ \vdots \\ \Pr(s_n|e_{1:t}) \end{pmatrix}$$

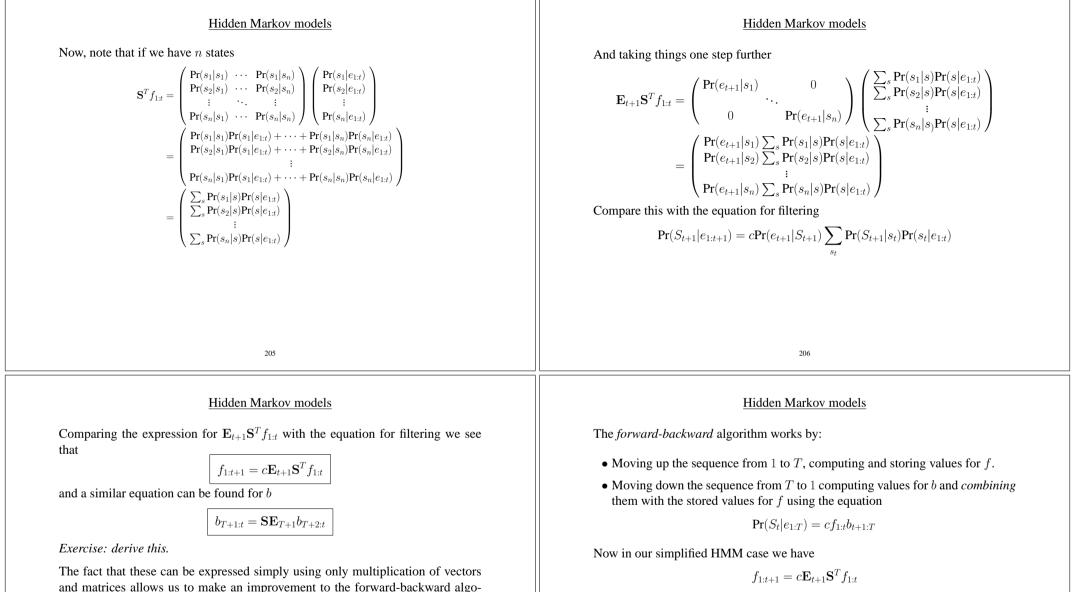
Key point: the filtering equation now reduces to nothing but matrix multiplication.

What does matrix multiplication do?

What does matrix multiplication do? It computes weighted summations:

$$\mathbf{A}b = \begin{pmatrix} a_{1,1} & a_{1,2} & \cdots & a_{1,m} \\ a_{2,1} & a_{2,2} & \cdots & a_{2,m} \\ \vdots & \vdots & \ddots & \vdots \\ a_{n,1} & a_{n,2} & \cdots & a_{n,m} \end{pmatrix} \begin{pmatrix} b_1 \\ b_2 \\ \vdots \\ b_m \end{pmatrix} = \begin{pmatrix} \sum_{i=1}^m a_{1,i}b_i \\ \sum_{i=1}^m a_{2,i}b_i \\ \vdots \\ \sum_{i=1}^m a_{n,i}b_i \end{pmatrix}$$

So the point at the end of the last slide shouldn't come as a big surprise!



or multiplying through by $(\mathbf{E}_{t+1}\mathbf{S}^T)^{-1}$ and re-arranging

$$f_{1:t} = \frac{1}{c} (\mathbf{S}^T)^{-1} (\mathbf{E}_{t+1})^{-1} f_{1:t+1}$$

rithm.

Hidden Markov models Example: 2008, paper 9, question 5 A friend of mine likes to climb on the roofs of Cambridge. To make a good start to the coming week, he climbs on a Sunday with probability 0.98. Being concerned • We know the *final* value for *f*. for his own safety, he is less likely to climb today if he climbed yesterday, so $\Pr(\text{climb today}|\text{climb yesterday}) = 0.4$ • Every observation has non-zero probability in every state. If he did not climb yesterday then he is very unlikely to climb today, so $\Pr(\text{climb today}|\neg\text{climb yesterday}) = 0.1$ We *don't* have to store T different values for f—we just work through, discarding

Unfortunately, he is not a very good climber, and is quite likely to injure himself if he goes climbing, so

Pr(injury|climb today) = 0.8

whereas

 $\Pr(\text{injury}|\neg \text{climb today}) = 0.1$

210

Example: 2008, paper 9, question 5

The update equation is

 $f_{1:t+1} = cE_{t+1}S^T f_{1:t}$

$f_{1:1} = \frac{c}{10,000} \begin{pmatrix} 8 & 0\\ 0 & 1 \end{pmatrix} \begin{pmatrix} 4 & 1\\ 6 & 9 \end{pmatrix} \begin{pmatrix} 98\\ 2 \end{pmatrix} = \begin{pmatrix} 0.83874\\ 0.16126 \end{pmatrix}$

Repeating this twice more using E' rather than E the final time gives

$$f_{1:2} = \begin{pmatrix} 0.81268\\ 0.18732 \end{pmatrix}$$
$$f_{1:3} = \begin{pmatrix} 0.10429\\ 0.89571 \end{pmatrix}$$

212

so the answer is 0.1.

Example: 2008, paper 9, question 5

209

You learn that on Monday and Tuesday evening he obtains an injury, but on Wednesday evening he does not. Use the filtering algorithm to compute the probability that he climbed on Wednesday.

Initially

$$f_{1:0} = \begin{pmatrix} 0.98\\ 0.02 \end{pmatrix}$$
$$S = \begin{pmatrix} 0.4 & 0.6\\ 0.1 & 0.9 \end{pmatrix}$$
$$E = \begin{pmatrix} 0.8 & 0\\ 0 & 0.1 \end{pmatrix}$$
$$E' = \begin{pmatrix} 0.2 & 0\\ 0 & 0.9 \end{pmatrix}$$

so

So as long as:

- \mathbf{S}^T has an inverse.

intermediate values, to obtain the last value and then work backward.

Example: 2008, paper 9, question 5

Over the course of the week, you also learn that he does not obtain an injury on Thursday or Friday. Use the smoothing algorithm to compute the probability that he climbed on Thursday.

The S, E and E' matrices are the same. The backward message starts as

$$b_{6:5} = \begin{pmatrix} 1\\ 1 \end{pmatrix}$$

and the update equation is

$$b_{t:T} = SE_t b_{t+1:T}$$

Then working backwards

$$b_{5:5} = \frac{1}{100} \begin{pmatrix} 4 & 6 \\ 1 & 9 \end{pmatrix} \begin{pmatrix} 2 & 0 \\ 0 & 9 \end{pmatrix} \begin{pmatrix} 1 \\ 1 \end{pmatrix} = \begin{pmatrix} 0.62 \\ 0.83 \end{pmatrix}$$

Example: 2008, paper 9, question 5

We also need one more forward step, which gives

$$f_{1:4} = \left(\begin{array}{c} 0.03249\\ 0.96751 \end{array}\right)$$

Finally

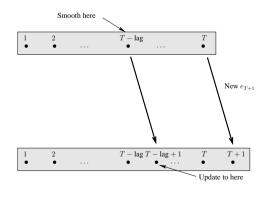
$$cf_{1:4}b_{5:5} = c \left(\begin{array}{c} 0.03249 \times 0.62\\ 0.96751 \times 0.83 \end{array} \right) = \left(\begin{array}{c} 0.02447\\ 0.97553 \end{array} \right)$$

giving the answer 0.02447.

213

Online smoothing

Say we want to smooth at a *fixed number of time steps*. We can also obtain a simple algorithm for updating the result each time a new e_{t+1} appears.



Online smoothing

214

As usual we need to calculate

 $cf_{1:T-lag}b_{T-lag+1:T}$

to smooth at time (T - lag) if we've progressed to time T. So: assume $f_{1:T-lag}$ and $b_{T-lag+1:T}$ are known.

What can we now do when e_{T+1} arrives to obtain $f_{1:T-lag+1}$ and $b_{T-lag+2:T+1}$?

f is easy to update because as usual

$$f_{1:T-\text{lag}+1} = c \mathbf{E}_{T-\text{lag}+1} \mathbf{S}^{T} \underbrace{f_{1:T-\text{lag}}}_{\text{Known}}$$

Online smoothing

b is more tricky.

We know that

$$b_{T-\text{lag}+1:T} = \mathbf{SE}_{T-\text{lag}+1} b_{T-\text{lag}+2:T}$$

and continuing this recursion up to the end of the sequence at T gives

$$b_{T-\text{lag}+1:T} = \prod_{i=T-\text{lag}+1}^{T} \mathbf{SE}_i \times \begin{pmatrix} 1\\1\\\vdots\\1 \end{pmatrix}$$

Define

so

$$b_{T-\text{lag}+1:T} = \boldsymbol{\beta}_{T-\text{lag}+1:T} \times \begin{pmatrix} 1\\ 1\\ \vdots\\ 1 \end{pmatrix}$$

 $oldsymbol{eta}_{a:b} = \prod_{i=a}^{b} \mathbf{SE}_i$

217

Online smoothing

This leads to an easy way to update β

$$\boldsymbol{\beta}_{a+1:b+1} = \mathbf{E}_a^{-1} \mathbf{S}^{-1} \boldsymbol{\beta}_{a:b} \mathbf{S} \mathbf{E}_{b+1}$$

Using this gives the required update for *b*.

Supervised learning II: the Bayesian approach

We now place supervised learning into a probabilistic setting by examining:

- The application of Bayes' theorem to the *supervised learning problem*.
- Priors, the likelihood, and the posterior probability of a hypothesis.
- The *maximum likelihood* and *maximum a posteriori* hypotheses, and some examples.
- Bayesian decision theory: minimising the error rate.
- Application of the approach to *neural networks*, using approximation techniques.

Online smoothing

Now when e_{T+1} arrives we have

$$b_{T-\text{lag}+2:T+1} = \prod_{i=T-\text{lag}+2}^{T+1} \mathbf{SE}_i \times \begin{pmatrix} 1\\1\\\vdots\\1 \end{pmatrix}$$
$$= \boldsymbol{\beta}_{T-\text{lag}+2:T+1} \times \begin{pmatrix} 1\\1\\\vdots\\1 \end{pmatrix}$$
$$= \mathbf{E}_{T-\text{lag}+1}^{-1} \mathbf{S}^{-1} \boldsymbol{\beta}_{T-\text{lag}+1:T} \mathbf{SE}_{T+1} \times \begin{pmatrix} 1\\1\\\vdots\\1 \end{pmatrix}$$

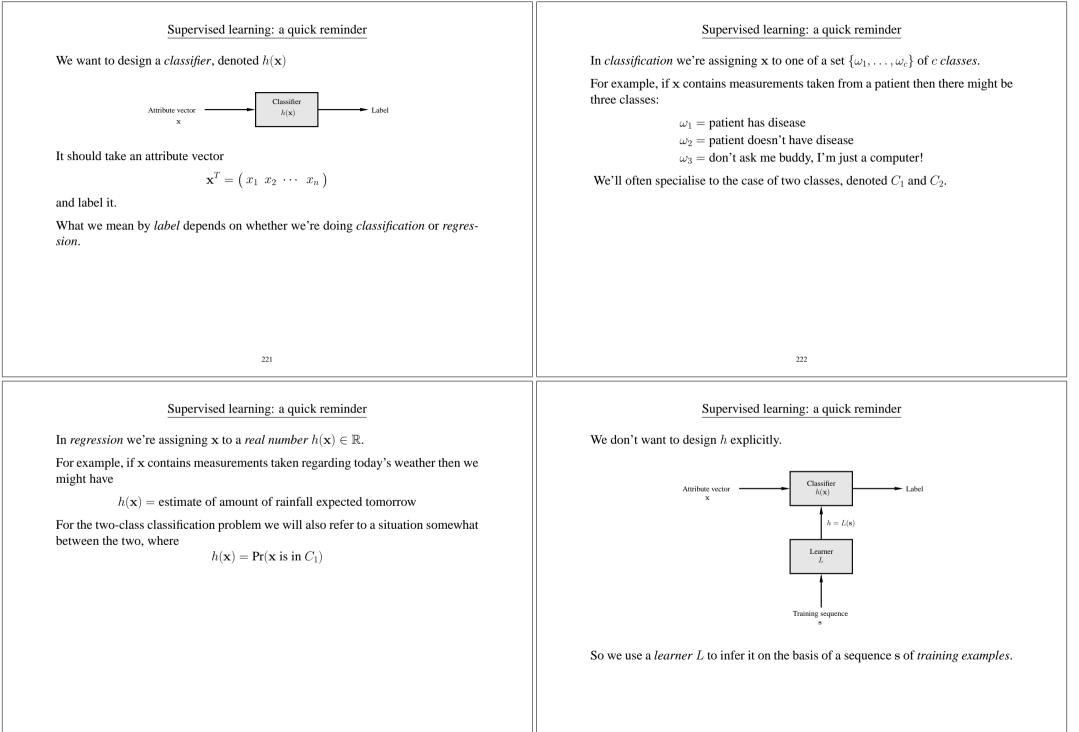
Reading

218

There is some relevant material to be found in *Russell and Norvig* chapters 18 to 20 although the intersection between that material and what I will cover is small.

Almost all of what I cover can be found in:

- Machine Learning. Tom Mitchell, McGraw Hill 1997, chapter 6.
- *Pattern Recognition and Machine Learning*. Christopher M. Bishop, Springer, 2006.



Supervised learning: a quick reminder

The training sequence s is a sequence of m labelled examples.

$$\mathbf{s} = \begin{pmatrix} (\mathbf{x}_1, y_1) \\ (\mathbf{x}_2, y_2) \\ \vdots \\ (\mathbf{x}_m, y_m) \end{pmatrix}$$

That is, examples of attribute vectors x with their correct label attached.

So a learner only gets to see the labels for a—most probably small—subset of the possible inputs x.

Regardless, we aim that the hypothesis h = L(s) will usually be successful at predicting the label of an input it hasn't seen before.

This ability is called *generalization*.

225

Supervised learning: a quick reminder

In AI I you saw the *backpropagation algorithm* for training multilayer perceptrons, in the case of *regression*.

This worked by minimising a function of the weights representing the *error* currently being made:

$$E(\mathbf{w}) = \frac{1}{2} \sum_{i=1}^{m} \left(f(\mathbf{w}; \mathbf{x}_i) - y_i \right)^2$$

The summation here is over the training examples. The expression in the summation grows as f's prediction for \mathbf{x}_i diverges from the known label y_i .

Backpropagation tries to find a w that minimises $E(\mathbf{w})$ by performing gradient descent

$$\mathbf{w}_{t+1} = \mathbf{w}_t - \alpha \left. \frac{\partial E(\mathbf{w})}{\partial \mathbf{w}} \right|_{\mathbf{w}_t}$$

Supervised learning: a quick reminder

There is generally a set ${\mathcal H}$ of hypotheses from which L is allowed to select h

$$L(\mathbf{s}) = h \in \mathcal{H}$$

 \mathcal{H} is called the *hypothesis space*.

The learner can output a hypothesis explicitly or—as in the case of a multilayer perceptron—it can output a vector

$$\mathbf{w} = \left(egin{array}{ccccc} w_1 & w_2 & \cdots & w_W \end{array}
ight)$$

of *weights* which in turn specify h

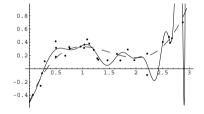
 $h(\mathbf{x}) = f(\mathbf{w}; \mathbf{x})$

where $\mathbf{w} = L(\mathbf{s})$.

226

Difficulties with classical neural networks

There are some well-known difficulties associated with neural network training of this kind.



BEWARE!!!

Sources of uncertainty

So we have to be careful. But let's press on with this approach for a little while longer...

The model used above suggests two sources of uncertainty that we might treat with probabilities.

- Let's *assume* we've selected an \mathcal{H} to use, *and it's the same one nature is using*.
- We don't know how nature chooses h' from \mathcal{H} . We therefore model our uncertainty by introducing the *prior* distribution Pr(h) on \mathcal{H} .
- There is noise on the training examples.

0.

It's worth emphasising at this point that in modelling noise on the training examples we'll only consider noise on the labels. The input vectors \mathbf{x} are not modelled using a probability distribution.

229

The likelihood

We model our uncertainty in the training examples by specifying a *likelihood*:

 $\Pr(Y|h, \mathbf{x})$

<u>Translation</u>: the probability of seeing a given label Y, when the input vector is \mathbf{x} and the underlying hypothesis is h.

Example: two-class classification. A common likelihood is

$$\Pr(Y = C_1 | h, \mathbf{x}) = \sigma(h(\mathbf{x}))$$

where

$$\sigma(z) = \frac{1}{1 + \exp(-z)}$$

(*Note*: strictly speaking x should not appear in these probabilities because it's not a random variable. It is included for clarity.)

230

The likelihood

So: if we're given a training sequence, *what is the probability that it was generated using some h*?

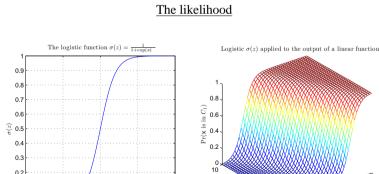
For an example (x, y), y can be C_1 or C_2 . It's helpful here to rename the classes as just 1 and 0 respectively because this leads to a nice simple expression. Now

$$\Pr(Y|h, \mathbf{x}) = \begin{cases} \sigma(h(\mathbf{x})) & \text{if } Y = 1\\ 1 - \sigma(h(\mathbf{x})) & \text{if } Y = 0 \end{cases}$$

Consequently when y has a known value we can write

$$\Pr(y|h, \mathbf{x}) = \left[\sigma(h(\mathbf{x}))\right]^y \left[1 - \sigma(h(\mathbf{x}))\right]^{(1-y)}$$

If we assume that the examples are independent then the probability of seeing the labels in a training sequence s is straightforward.



231

-10 -10

The likelihood

The likelihood

Collecting the inputs and outputs in s together into separate matrices, so

$$\mathbf{y}^T = \left(\ y_1 \ \ y_2 \ \cdots \ \ y_m \
ight)$$

and

 $\mathbf{X} = \begin{pmatrix} \mathbf{x}_1 \ \mathbf{x}_2 \ \cdots \ \mathbf{x}_m \end{pmatrix}$ we have the *likelihood of the training sequence*

$$\begin{aligned} \Pr(\mathbf{y}|h, \mathbf{X}) &= \prod_{i=1}^{m} \Pr(y_i|h, \mathbf{x}_i) \\ &= \prod_{i=1}^{m} \left[\sigma(h(\mathbf{x}_i))\right]^{y_i} \left[1 - \sigma(h(\mathbf{x}_i))\right]^{(1-y_i)} \end{aligned}$$

<u>Another example</u>: regression. A common likelihood in the regression case works by assuming that examples are corrupted by Gaussian noise with mean 0 and some specified variance σ^2

$$y = h(\mathbf{x}) + \epsilon$$
, where $\epsilon \sim \mathcal{N}(0, \sigma^2)$

As usual, the density for $\mathcal{N}(\mu,\sigma^2)$ is

$$p(Z) = \frac{1}{\sqrt{2\pi\sigma^2}} \exp\left(-\frac{(z-\mu)^2}{2\sigma^2}\right)$$

by adding $h(\mathbf{x})$ to ϵ we just shift its mean, so

$$p(y|h, \mathbf{x}) = \frac{1}{\sqrt{2\pi\sigma^2}} \exp\left(-\frac{(y - h(\mathbf{x}))^2}{2\sigma^2}\right)$$

234

The likelihood

233

Consequently if the examples are independent then the likelihood of a training sequence s is

 $p(\mathbf{y}|h, \mathbf{X}) = \prod_{i=1}^{m} p(y_i|h, \mathbf{x}_i)$ $= \prod_{i=1}^{m} \frac{1}{\sqrt{2\pi\sigma^2}} \exp\left(-\frac{(y_i - h(\mathbf{x}_i))^2}{2\sigma^2}\right)$ $= \frac{1}{(2\pi\sigma^2)^{m/2}} \exp\left(-\frac{1}{2\sigma^2} \sum_{i=1}^{m} (y_i - h(\mathbf{x}_i))^2\right)$

where we've used the fact that

$$\exp(a)\exp(b) = \exp(a+b)$$

Bayes' theorem appears once more ...

Right: we've take care of the uncertainty by introducing the *prior* p(h) and the *likelihood of the training sequence* $p(\mathbf{y}|h, \mathbf{X})$.

By this point you hopefully want to apply Bayes' theorem and write

$$p(h|\mathbf{y}) = \frac{p(\mathbf{y}|h)p(h)}{p(\mathbf{y})}$$

where

$$p(\mathbf{y}) = \sum_{h \in \mathcal{H}} p(h, \mathbf{y}) = \sum_{h \in \mathcal{H}} p(\mathbf{y}|h) p(h)$$

and to simplify the expression we have now dropped the mention of X as the inputs are fixed. p(h|y) is called the *posterior distribution*.

The denominator $Z = p(\mathbf{y})$ is called the *evidence* and leads on to fascinating issues of its own. Unfortunately we won't have time to explore them.

Bayes' theorem appears once more ...

The boxed equation on the last slide has a very simple interpretation: *what's the probability that this specific h was used to generate the training sequence I've been given?*

Two natural learning algorithms now present themselves:

1. The maximum likelihood hypothesis

$$h_{\mathrm{ML}} = \operatorname*{argmax}_{h \in \mathcal{H}} p(\mathbf{y}|h)$$

2. The maximum a posteriori hypothesis

$$h_{\text{MAP}} = \underset{h \in \mathcal{H}}{\operatorname{argmax}} p(h|\mathbf{y})$$
$$= \underset{h \in \mathcal{H}}{\operatorname{argmax}} p(\mathbf{y}|h) p(h)$$

Obviously h_{ML} corresponds to the case where the prior p(h) is uniform.

237

Example: maximum likelihood learning

The proposition holds because:

$$\begin{split} h_{\text{ML}} &= \operatorname*{argmax}_{h \in \mathcal{H}} p(\mathbf{y}|h) \\ &= \operatorname*{argmax}_{h \in \mathcal{H}} \log p(\mathbf{y}|h) \\ &= \operatorname*{argmax}_{h \in \mathcal{H}} \log \left[\frac{1}{(2\pi\sigma^2)^{m/2}} \exp\left(-\frac{1}{2\sigma^2} \sum_{i=1}^m (y_i - h(\mathbf{x}_i))^2 \right) \right] \\ &= \operatorname*{argmax}_{h \in \mathcal{H}} \log \left[\frac{1}{(2\pi\sigma^2)^{m/2}} \right] - \frac{1}{2\sigma^2} \sum_{i=1}^m (y_i - h(\mathbf{x}_i))^2 \\ &= \operatorname*{argmax}_{h \in \mathcal{H}} - \frac{1}{2\sigma^2} \sum_{i=1}^m (y_i - h(\mathbf{x}_i))^2 \\ &= \operatorname*{argmax}_{h \in \mathcal{H}} \sum_{i=1}^m (y_i - h(\mathbf{x}_i))^2 \end{split}$$

Example: maximum likelihood learning

We derived an exact expression for the likelihood in the regression case above:

$$p(\mathbf{y}|h) = \frac{1}{(2\pi\sigma^2)^{m/2}} \exp\left(-\frac{1}{2\sigma^2} \sum_{i=1}^m (y_i - h(\mathbf{x}_i))^2\right)$$

<u>Proposition</u>: under the assumptions used, *any* learning algorithm that works by minimising the sum of squared errors on s finds $h_{\rm ML}$.

This is clearly of interest: the notable example is the *backpropagation algorithm*.

We now prove the proposition...

238

Example: maximum likelihood learning

Note:

- If the distribution of the noise is not Gaussian a different result is obtained.
- The use of log above to simplify a maximisation problem is a standard trick.
- The Gaussian assumption is sometimes, but not always a good choice. (*Beware the Central Limit Theorem!*).

The next step...

We have so far concentrated throughout our coverage of machine learning on choosing a *single hypothesis*.

Are we asking the right question though?

Ultimately, we want to generalise.

That means being presented with a new x and asking the question: what is the most probable classification of x?

Is it reasonable to expect a single hypothesis to provide the optimal answer?

We need to look at what the optimal solution to this kind of problem might be...

Bayesian decision theory

What is the *optimal* approach to this problem?

Put another way: how should we make decisions in such a way that the outcome obtained is, on average, the best possible? Say we have:

- Attribute vectors $\mathbf{x} \in \mathbb{R}^d$.
- A set of *classes* $\{\omega_1, \ldots, \omega_c\}$.
- Several possible *actions* $\{\alpha_1, \ldots, \alpha_a\}$.

The actions can be thought of as saying "assign the vector to class 1" and so on. There is also a loss $\lambda(\alpha_i, \omega_j)$ associated with taking action α_i when the class is ω_j . The loss will sometimes be abbreviated to $\lambda(\alpha_i, \omega_j) = \lambda_{ij}$.

241

Bayesian decision theory

Say we can also *model* the world as follows:

- Classes have probabilities $Pr(\omega)$ of occurring.
- The probability of seeing x when the class is ω has density $p(\mathbf{x}|\omega)$.

Think of nature choosing classes at random (although not revealing them) and showing us a vector selected at random using $p(\mathbf{x}|\omega)$.

As usual Bayes rule tells us that

$$\Pr(\omega|\mathbf{x}) = \frac{p(\mathbf{x}|\omega)\Pr(\omega)}{p(\mathbf{x})}$$

and now the denominator is

$$p(\mathbf{x}) = \sum_{i=1}^{c} p(\mathbf{x}|\omega_i) \Pr(\omega_i)$$

242

Bayesian decision theory

Say nature shows us x and we take action α_i .

If we *always* take action α_i when we see x then the *average* loss on seeing x is

$$R(\alpha_i | \mathbf{x}) = \mathbb{E}_{\omega \sim p(\omega | \mathbf{x})} [\lambda_{ij} | \mathbf{x}] = \sum_{j=1}^{c} \lambda(\alpha_i, \omega_j) \Pr(\omega_j | \mathbf{x}).$$

The quantity $R(\alpha_i | \mathbf{x})$ is called the *conditional risk*. Note that this particular \mathbf{x} is *fixed*.

Bayesian decision theory

Now say we have a *decision rule* $\alpha : \mathbb{R}^d \to {\alpha_1, \ldots, \alpha_a}$ telling us what action to take on seeing *any* $\mathbf{x} \in \mathbb{R}^d$.

The average loss, or risk, is

$$R = \mathbb{E}_{(\mathbf{x},\omega) \sim p(\mathbf{x},\omega)} [\lambda(\alpha(\mathbf{x}),\omega)]$$

= $\mathbb{E}_{\mathbf{x} \sim p(\mathbf{x})} [\mathbb{E}_{\omega \sim \Pr(\omega|\mathbf{x})} [\lambda(\alpha(\mathbf{x}),\omega)|\mathbf{x}]]$
= $\mathbb{E}_{\mathbf{x} \sim p(\mathbf{x})} [R(\alpha(\mathbf{x})|\mathbf{x})]$ (11)
= $\int R(\alpha(\mathbf{x})|\mathbf{x})p(\mathbf{x})d\mathbf{x}$

where we have used the standard result from probability theory that

$$\mathbb{E}\left[\mathbb{E}\left[X|Y\right]\right] = \mathbb{E}\left[X\right].$$

(See the supplementary notes for a proof.)

245

Example: minimum error rate classification

In supervised learning our aim is often to work in such a way that we *minimise* the probability of error.

What loss should we consider in these circumstances? From basic probability theory $\Pr(A) = \mathbb{E}\left[I(A)\right]$

where

$$I(A) = \begin{cases} 1 & \text{if } A \text{ happens} \\ 0 & \text{otherwise} \end{cases}$$

(See the supplementary notes for a proof.)

Bayesian decision theory

Clearly the risk is minimised for the decision rule defined as follows:

 α outputs the action α_i that minimises $R(\alpha_i | \mathbf{x})$, for all $\mathbf{x} \in \mathbb{R}^d$.

The provides us with the minimum possible risk, or *Bayes risk* R^* . The rule specified is called the *Bayes decision rule*.

246

Example: minimum error rate classification

So if we are addressing a supervised learning problem with c classes $\{\omega_1, \ldots, \omega_c\}$ and we interpret action α_i as meaning 'the input is in class ω_i ', then a loss

$$\lambda_{ij} = \begin{cases} 1 & \text{if } i \neq j \\ 0 & \text{otherwise} \end{cases}$$

means that the risk R is

 $R = \mathbb{E}[\lambda] = \Pr(\alpha(\mathbf{x}) \text{ is in error})$

and the Bayes decision rule minimises the probability of error.

Example: minimum error rate classification

Bayesian learning II

Now, what is the Bayes decision rule?

$$R(\alpha_i | \mathbf{x}) = \sum_{j=1}^{c} \lambda(\alpha_i, \omega_j) \Pr(\omega_j | \mathbf{x})$$
$$= \sum_{i \neq j} \Pr(\omega_j | \mathbf{x})$$
$$= 1 - \Pr(\omega_i | \mathbf{x})$$

so $\alpha(\mathbf{x})$ should be the class that maximises $Pr(\omega_i | \mathbf{x})$.

THE IMPORTANT SUMMARY: Given a new x to classify, choosing the class that maximises $Pr(\omega_i | \mathbf{x})$ is the best strategy if your aim is to obtain the minimum error rate!

Bayes decision theory tells us that in this context we should consider the quantity $Pr(\omega_i | \mathbf{s}, \mathbf{x})$ where the involvement of the training sequence has been made explicit.

$$\begin{split} \mathbf{Pr}(\omega_i | \mathbf{s}, \mathbf{x}) \; &=\; \sum_{h \in \mathcal{H}} \mathbf{Pr}(\omega_i, h | \mathbf{s}, \mathbf{x}) \\ &=\; \sum_{h \in \mathcal{H}} \mathbf{Pr}(\omega_i | h, \mathbf{s}, \mathbf{x}) \mathbf{Pr}(h | \mathbf{s}, \mathbf{x}) \\ &=\; \sum_{h \in \mathcal{H}} \mathbf{Pr}(\omega_i | h, \mathbf{x}) \mathbf{Pr}(h | \mathbf{s}). \end{split}$$

Here we have re-introduced \mathcal{H} using marginalisation. In moving from line 2 to line 3 we are assuming some independence properties.

249

Bayesian learning II

So our classification should be

$$\omega = \operatorname*{argmax}_{\omega \in \{\omega_1, \dots, \omega_c\}} \sum_{h \in \mathcal{H}} \Pr(\omega|h, \mathbf{x}) \Pr(h|\mathbf{s})$$

If \mathcal{H} is infinite the sum becomes an integral. So for example for a neural network

$$\omega = \operatorname*{argmax}_{\omega \in \{\omega_1, ..., \omega_c\}} \int_{\mathbb{R}^W} \Pr(\omega | \mathbf{w}, \mathbf{x}) \Pr(\mathbf{w} | \mathbf{s}) \, d\mathbf{w}$$

where W is the number of weights in w.

250

Bayesian learning II

Why might this make any difference? (Aside from the fact that we now know it's optimal!)

Example 1: Say $|\mathcal{H}| = 3$ and $h(\mathbf{x}) = \Pr(\mathbf{x} \text{ is in class } C_1)$ for a 2 class problem.

$$\mathbf{Pr}(h_1|\mathbf{s}) = 0.4$$
$$\mathbf{Pr}(h_2|\mathbf{s}) = \mathbf{Pr}(h_3|\mathbf{s}) = 0.3$$

Now, say we have an x for which

$$h_1(\mathbf{x}) = 1$$
$$h_2(\mathbf{x}) = h_3(\mathbf{x}) = 0$$

so h_{MAP} says that x is in class C_1 .

Bayesian learning II A more in-depth example However, Let's take this a step further and work through something a little more complex in detail. For a two-class classification problem with h(x) denoting $Pr(C_1|h, x)$ and $Pr(class \ 1|s, x) = 1 \times 0.4 + 0 \times 0.3 + 0 \times 0.3$ $x \in \mathbb{R}$: = 0.4Hypotheses: We have three hypotheses $Pr(class \ 2|s, x) = 0 \times 0.4 + 1 \times 0.3 + 1 \times 0.3$ = 0.6 $h_1(x) = \exp(-(x-1)^2)$ $h_2(x) = \exp(-(2x - 2)^2)$ so class C_2 is the more probable! $h_3(x) = \exp(-(1/10)(x-3)^2)$ In this case the Bayes optimal approach in fact leads to a different answer. Prior: The prior is $Pr(h_1) = 0.1$, $Pr(h_2) = 0.05$ and $Pr(h_3) = 0.85$. 253 254 A more in-depth example A more in-depth example We see the examples $(0.5, C_1), (0.9, C_1), (3.1, C_2)$ and $(3.4, C_1)$. Now let's classify the point x' = 2.5. Likelihood: For the individual hypotheses the likelihoods are given by We need $Pr(C_1|\mathbf{s}, x') = Pr(C_1|h_1)Pr(h_1|\mathbf{s}) + Pr(C_1|h_2)Pr(h_2|\mathbf{s}) + Pr(C_1|h_3)Pr(h_3|\mathbf{s})$ $\Pr(\mathbf{s}|h) = h(x_1)h(x_2)[1 - h(x_3)]h(x_4)$ = 0.6250705317Which in this case tells us So: it's most likely to be in class C_1 , but not with great certainty. $\Pr(\mathbf{s}|h_1) = 0.0024001365$ $\Pr(\mathbf{s}|h_2) = 0.0031069836$ $\Pr(\mathbf{s}|h_3) = 0.0003387476$ Posterior: Multiplying by the priors and normalising gives $\Pr(h_1|\mathbf{s}) = 0.3512575000$ $\Pr(h_2|\mathbf{s}) = 0.2273519164$ $\Pr(h_3|\mathbf{s}) = 0.4213905836$ 255 256

The Bayesian approach to neural networks

Let's now see how this can be applied to *neural networks*. We have:

- A neural network computing a function $f(\mathbf{w}; \mathbf{x})$.
- A training sequence $\mathbf{s} = ((\mathbf{x}_1, y_1), \dots, (\mathbf{x}_m, y_m))$, split into

$$\mathbf{y} = (y_1 \ y_2 \ \cdots \ y_m)$$

and

$$\mathbf{X} = \left(\mathbf{x}_1 \ \mathbf{x}_2 \ \cdots \ \mathbf{x}_m \right)$$

The *prior distribution* $p(\mathbf{w})$ is now on the weight vectors and Bayes' theorem tells us that

$$p(\mathbf{w}|\mathbf{s}) = p(\mathbf{w}|\mathbf{X}, \mathbf{y}) = \frac{p(\mathbf{y}|\mathbf{w}, \mathbf{X})p(\mathbf{w}|\mathbf{X})}{p(\mathbf{y}|\mathbf{X})}$$

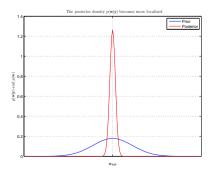
Nothing new so far...

257

The Bayesian approach to neural networks

 $p(\mathbf{w})$ is typically a *broad distribution* to reflect the fact that in the absence of any data we have little idea of what \mathbf{w} might be.

When we see some data the above equation tells us how to obtain $p(\mathbf{w}|\mathbf{y})$. This will typically be *more localised*.



To put this into practice we need expressions for $p(\mathbf{w})$ and $p(\mathbf{y}|\mathbf{w})$.

The Bayesian approach to neural networks

As usual, we don't consider uncertainty in ${\bf x}$ and so ${\bf X}$ will be omitted. Consequently

 $p(\mathbf{w}|\mathbf{y}) = \frac{p(\mathbf{y}|\mathbf{w})p(\mathbf{w})}{p(\mathbf{y})}$

where

$$p(\mathbf{y}) = \int_{\mathbb{R}^W} p(\mathbf{y}|\mathbf{w}) p(\mathbf{w}) d\mathbf{w}$$

 $p(\mathbf{y}|\mathbf{w})$ is a model of the noise corrupting the labels and as previously is the *like-lihood function*.

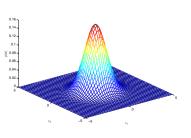
258

Reminder: the general Gaussian density

<u>Reminder</u>: we're going to be making a lot of use of the general *Gaussian density* $\mathcal{N}(\mu, \Sigma)$ in *d* dimensions

$$p(\mathbf{z}) = (2\pi)^{-d/2} |\mathbf{\Sigma}|^{-1/2} \exp\left[-\frac{1}{2}\left((\mathbf{z} - \boldsymbol{\mu})^T \mathbf{\Sigma}^{-1} (\mathbf{z} - \boldsymbol{\mu})\right)\right]$$

where μ is the *mean vector* and Σ is the *covariance matrix*.



Gaussian density, $\mu = [0 \ 0]$, $\Sigma = I$

The Gaussian prior

A common choice for $p(\mathbf{w})$ is the *Gaussian prior* with zero mean and

 $\boldsymbol{\Sigma} = \sigma^2 \mathbf{I}$

 \mathbf{so}

$$p(\mathbf{w}) = (2\pi)^{-W/2} \sigma^{-W} \exp\left[-\frac{\mathbf{w}^T \mathbf{w}}{2\sigma^2}\right]$$

Note that σ controls the distribution of other parameters.

• Such parameters are called *hyperparameters*.

• Assume for now that they are both fixed and known.

Hyperparameters can be learnt using s through the application of more advanced techniques.

261

The Gaussian noise model for regression

We've already seen that for a regression problem with zero mean Gaussian noise having variance σ_n^2

$$y_i = f(\mathbf{x}_i) + \epsilon_i$$
$$p(\epsilon_i) = \frac{1}{\sqrt{2\pi\sigma_n^2}} \exp\left(-\frac{\epsilon_i^2}{2\sigma_n^2}\right)$$

where f corresponds to some unknown network, the likelihood function is

$$p(\mathbf{y}|\mathbf{w}) = \frac{1}{(2\pi\sigma_n^2)^{m/2}} \exp\left(-\frac{1}{2\sigma_n^2} \sum_{i=1}^m (y_i - f(\mathbf{w}; \mathbf{x}_i))^2\right)$$

Note that there are now two variances: σ^2 for the prior and σ_n^2 for the noise.

The Bayesian approach to neural networks

Physicists like to express quantities such as $p(\mathbf{w})$ in terms of a measure of "*energy*". The expression is therefore usually re-written as

$$p(\mathbf{w}) = \frac{1}{Z_W(\alpha)} \exp\left(-\frac{\alpha}{2} ||\mathbf{w}||^2\right)$$

where

$$E_W(\mathbf{w}) = \frac{1}{2} ||\mathbf{w}||^2$$
$$Z_W(\alpha) = \left(\frac{2\pi}{\alpha}\right)^{d/2}$$
$$\alpha = \frac{1}{\sigma^2}$$

This is simply a re-arranged version of the more usual equation.

262

The Bayesian approach to neural networks

This expression can also be rewritten in physicist-friendly form

$$\left| \begin{array}{c} p(\mathbf{y}|\mathbf{w}) = \frac{1}{Z_{\mathbf{y}}(\beta)} \exp\left(-\beta E_{\mathbf{y}}(\mathbf{w})\right) \end{array} \right.$$

where

$$\begin{split} \boldsymbol{\beta} &= \frac{1}{\sigma_n^2} \\ Z_{\mathbf{y}}(\boldsymbol{\beta}) &= \left(\frac{2\pi}{\beta}\right)^{m/2} \\ E_{\mathbf{y}}(\mathbf{w}) &= \frac{1}{2} \sum_{i=1}^m (y_i - f(\mathbf{w}; \mathbf{x}_i))^2 \end{split}$$

Here, β is a second *hyperparameter*. Again, we assume it is fixed and known, although it can be learnt using s using more advanced techniques.

The Bayesian approach to neural networks

Combining the two boxed equations gives

$$p(\mathbf{w}|\mathbf{y}) = \frac{1}{Z_S(\alpha, \beta)} \exp(-S(\mathbf{w}))$$

where

$$S(\mathbf{w}) = \alpha E_W(\mathbf{w}) + \beta E_{\mathbf{y}}(\mathbf{w})$$

The quantity

$$Z_S(\alpha,\beta) = \int_{\mathbb{R}^W} \exp(-S(\mathbf{w})) d\mathbf{w}$$

265

Example II: two-class classification in two dimensions

x 10

x 10⁻⁴

0.5

1.5

Prior density n(w

-10 -10

w2 -10 -10

Posterior density $p(\mathbf{w}|\mathbf{y})$

 w_1

 w_1

normalises the density. Recall that this is called the *evidence*.

Examples

0 1 x1

Likelihood $p(\mathbf{y}|\mathbf{w})$

-10 -10

§ (

-3L -3

0.06

0.02

-2 -1

 w_2

0

0

2 3

Example I: gradient descent revisited...

To find h_{MAP} (in this scenario by finding w_{MAP}) we therefore maximise

$$p(\mathbf{w}|\mathbf{y}) = \frac{1}{Z_S(\alpha,\beta)} \exp(-(\alpha E_W(\mathbf{w}) + \beta E_\mathbf{y}(\mathbf{w})))$$

or equivalently find

$$\mathbf{w}_{\text{MAP}} = \underset{\mathbf{w}}{\operatorname{argmin}} \frac{\alpha}{2} ||\mathbf{w}||^2 + \frac{\beta}{2} \sum_{i=1}^{m} (y_i - f(\mathbf{w}; \mathbf{x}_i))^2$$

This algorithm has also been used a lot in the neural network literature and is called the *weight decay* technique.

266

The Bayesian approach to neural networks

What happens as the number m of examples increases?

- The first term *corresponding to the prior* remains fixed.
- The second term *corresponding to the likelihood* increases.

So for small training sequences the prior dominates, but for large ones $h_{\rm ML}$ is a good approximation to $h_{\rm MAP}$.

The Bayesian approach to neural networks

Where have we got to ...? We have obtained

$$p(\mathbf{w}|\mathbf{y}) = \frac{1}{Z_S(\alpha,\beta)} \exp(-(\alpha E_W(\mathbf{w}) + \beta E_\mathbf{y}(\mathbf{w})))$$
$$Z_S(\alpha,\beta) = \int_{\mathbb{R}^W} \exp(-(\alpha E_W(\mathbf{w}) + \beta E_\mathbf{y}(\mathbf{w}))) d\mathbf{w}$$

Translating the expression for the *Bayes optimal* solution given earlier into the current scenario, we need to compute

$$p(Y|\mathbf{y}, \mathbf{x}) = \int_{\mathbb{R}^W} p(y|\mathbf{w}, \mathbf{x}) p(\mathbf{w}|\mathbf{y}) \, d\mathbf{w}$$

Easy huh? Unfortunately not...

The Bayesian approach to neural networks

In order to make further progress it's necessary to perform integrals of the general form

$$\sum_{\mathbf{w} \in W} F(\mathbf{w}) p(\mathbf{w} | \mathbf{y}) d\mathbf{w}$$

for various functions F and this is generally not possible.

There are two ways to get around this:

1. We can use an *approximate form* for $p(\mathbf{w}|\mathbf{y})$.

2. We can use *Monte Carlo* methods.

269

Method 1: approximation to $p(\mathbf{w}|\mathbf{y})$

The first approach introduces a *Gaussian approximation* to $p(\mathbf{w}|\mathbf{y})$ by using a *Taylor expansion* of

 $S(\mathbf{w}) = \alpha E_W(\mathbf{w}) + \beta E_\mathbf{y}(\mathbf{w})$

at \mathbf{w}_{MAP} .

This allows us to use a *standard integral*.

The result will be *approximate* but we hope it's good!

Let's recall how Taylor series work ...

Reminder: Taylor expansion

270

In one dimension the Taylor expansion about a point $x_0 \in \mathbb{R}$ for a function $f : \mathbb{R} \to \mathbb{R}$ is

$$f(x) \approx f(x_0) + \frac{1}{1!}(x - x_0)f'(x_0) + \frac{1}{2!}(x - x_0)^2 f''(x_0) + \dots + \frac{1}{k!}(x - x_0)^k f^k(x_0)$$

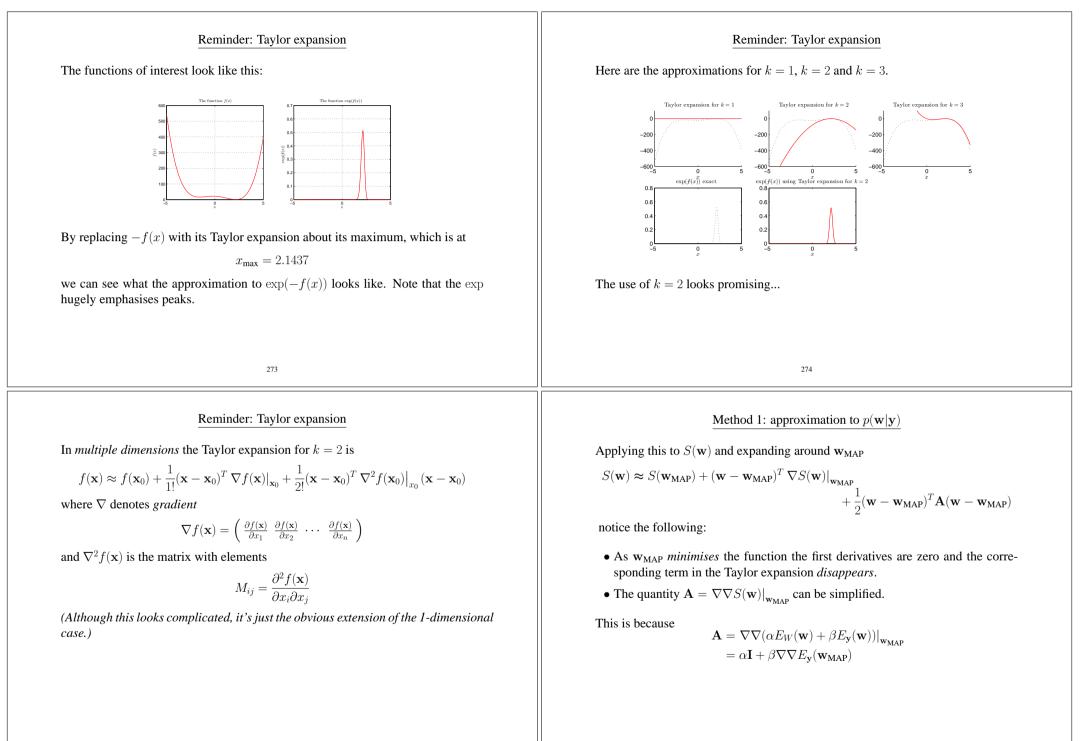
What does this look like for the kinds of function we're interested in? We can try to approximate

 $\exp\left(-f(x)\right)$

where

$$f(x) = x^4 - \frac{1}{2}x^3 - 7x^2 - \frac{5}{2}x + 22$$

This has a form similar to $S(\mathbf{w})$, but in one dimension.



Method 1: approximation to $p(\mathbf{w}|\mathbf{y})$

Defining

 $\Delta \mathbf{w} = \mathbf{w} - \mathbf{w}_{MAP}$

we now have

$$S(\mathbf{w}) \approx S(\mathbf{w}_{\text{MAP}}) + \frac{1}{2} \Delta \mathbf{w}^T \mathbf{A} \Delta \mathbf{w}$$

The vector \mathbf{w}_{MAP} can be obtained using any standard optimisation method (such as *backpropagation*).

The quantity $\nabla \nabla E_{\mathbf{y}}(\mathbf{w})$ can be evaluated using an extended form of backpropagation.

A useful integral

Dropping *for this slide only* the special meanings usually given to vectors **x** and **y**, here is a useful standard integral:

If $\mathbf{A} \in \mathbb{R}^{n \times n}$ is symmetric then for $\mathbf{b} \in \mathbb{R}^n$ and $c \in \mathbb{R}$

$$\int_{\mathbb{R}^n} \exp\left(-\frac{1}{2}\left(\mathbf{x}^T \mathbf{A} \mathbf{x} + \mathbf{x}^T \mathbf{b} + c\right)\right) d\mathbf{x}$$
$$= (2\pi)^{n/2} |\mathbf{A}|^{-1/2} \exp\left(-\frac{1}{2}\left(c - \frac{\mathbf{b}^T \mathbf{A}^{-1} \mathbf{b}}{4}\right)\right)$$

At the beginning of the course, two exercises were set involving the evaluation of this integral.

To make this easy to refer to, let's call it the BIG INTEGRAL.

278

Method 1: approximation to $p(\mathbf{w}|\mathbf{y})$

277

We now have

$$p(\mathbf{w}|\mathbf{y}) \approx \frac{1}{Z(\alpha,\beta)} \exp\left(-S(\mathbf{w}_{\text{MAP}}) - \frac{1}{2}\Delta \mathbf{w}^T \mathbf{A} \Delta \mathbf{w}\right)$$

where $\Delta \mathbf{w} = \mathbf{w} - \mathbf{w}_{MAP}$ and using the *BIG INTEGRAL*

$$Z(\alpha,\beta) = (2\pi)^{W/2} |\mathbf{A}|^{-1/2} \exp(-S(\mathbf{w}_{\mathrm{MAP}}))$$

Our earlier discussion tells us that given a new input x we should calculate

$$p(Y|\mathbf{y},\mathbf{x}) = \int_{\mathbb{R}^W} p(y|\mathbf{w},\mathbf{x}) p(\mathbf{w}|\mathbf{y}) d\mathbf{w}$$

 $p(y|\mathbf{w}, \mathbf{x})$ is just the *likelihood* so...

Method 1: approximation to $p(\mathbf{w}|\mathbf{y})$

The likelihood we're using is

$$p(y|\mathbf{w}, \mathbf{x}) = \frac{1}{\sqrt{2\pi\sigma^2}} \exp\left(-\frac{(y - f(\mathbf{w}; \mathbf{x}))^2}{2\sigma^2}\right)$$
$$\propto \exp\left(-\frac{\beta}{2}(y - f(\mathbf{w}; \mathbf{x}))^2\right)$$

and plugging it into the integral gives

$$p(y|\mathbf{x}, \mathbf{y}) \propto \int_{\mathbb{R}^W} \exp\left(-\frac{\beta}{2}(y - f(\mathbf{w}; \mathbf{x}))^2\right) \exp\left(-\frac{1}{2}\Delta\mathbf{w}^T \mathbf{A}\Delta\mathbf{w}\right) d\mathbf{w}$$

which has no solution!

We need another approximation...

Method 1: approximation to $p(\mathbf{w}|\mathbf{y})$

If we *assume* that $p(\mathbf{w}|\mathbf{y})$ is narrow (this depends on A) then we can introduce a *linear approximation* of $f(\mathbf{w}; \mathbf{x})$ at \mathbf{w}_{MAP} :

$$f(\mathbf{w}; \mathbf{x}) \approx f(\mathbf{w}_{\text{MAP}}; \mathbf{x}) + \mathbf{g}^T \Delta \mathbf{w}$$

where $\mathbf{g} = \nabla f(\mathbf{w}; \mathbf{x})|_{\mathbf{w}_{\mathsf{MAP}}}$.

That is, the MAP solution.

Or *interpreted* graphically...

By linear approximation we just mean the Taylor expansion for k = 1.

This leads to

$$p(Y|\mathbf{y}, \mathbf{x}) \propto \int_{\mathbb{R}^W} \exp\left(-\frac{\beta}{2} \left(y - f(\mathbf{w}_{\text{MAP}}; \mathbf{x}) - \mathbf{g}^T \Delta \mathbf{w}\right)^2 - \frac{1}{2} \Delta \mathbf{w}^T \mathbf{A} \Delta \mathbf{w}\right) d\mathbf{w}$$

and this integral can be evaluated using the *BIG INTEGRAL* to give *THE AN-SWER*...

281

Method 1: approximation to $p(\mathbf{w}|\mathbf{y})$

This is a *Gaussian density*, so we can now see that $p(Y|\mathbf{y}, \mathbf{x})$ peaks at $f(\mathbf{w}_{MAP}; \mathbf{x})$.

• The second term of σ_u^2 is $\mathbf{g}^T \mathbf{A}^{-1} \mathbf{g}$ and corresponds to the width of $p(\mathbf{w}|\mathbf{y})$.

The variance σ_y^2 can be interpreted as a measure of *certainty*.

• The first term of σ_y^2 is $1/\beta$ and corresponds to the noise.

Method 1: approximation to $p(\mathbf{w}|\mathbf{y})$

Finally

$$p(Y|\mathbf{y}, \mathbf{x}) = \frac{1}{\sqrt{2\pi\sigma_y^2}} \exp\left(-\frac{(y - f(\mathbf{w}_{\text{MAP}}; \mathbf{x}))^2}{2\sigma_y^2}\right)$$

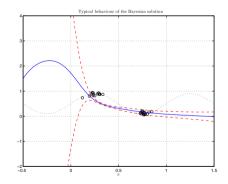
where

$$\sigma_y^2 = \frac{1}{\beta} + \mathbf{g}^T \mathbf{A}^{-1} \mathbf{g}.$$

Hooray! But what does it mean?

282

Method 1: approximation to $p(\mathbf{w}|\mathbf{y})$



Method II: Markov chain Monte Carlo (MCMC) methods

MCMC methods

The second solution to the problem of performing integrals

$$I = \int F(\mathbf{w}) p(\mathbf{w}|\mathbf{y}) d\mathbf{w}$$

is to use *Monte Carlo* methods. The basic approach is to make the approximation

$$I \approx \frac{1}{N} \sum_{i=1}^{N} F(\mathbf{w}_i)$$

where the \mathbf{w}_i have distribution $p(\mathbf{w}|\mathbf{y})$. Unfortunately, generating \mathbf{w}_i with a *given distribution* can be non-trivial.

285

Approximate inference for Bayesian networks

MCMC methods also provide a method for performing *approximate inference* in *Bayesian networks*.

Say a system can be in a state s and moves from state to state in discrete time steps according to a probabilistic transition

$$\Pr(\mathbf{s} \to \mathbf{s}')$$

Let $\pi_t(\mathbf{s})$ be the probability distribution for the state after t steps, so

$$\pi_{t+1}(\mathbf{s}') = \sum_{\mathbf{s}} \Pr(\mathbf{s} \to \mathbf{s}') \pi_t(\mathbf{s})$$

If at some point we obtain $\pi_{t+1}(\mathbf{s}) = \pi_t(\mathbf{s})$ for all \mathbf{s} then we have reached a *stationary distribution* π . In this case

$$\forall \mathbf{s'} \pi(\mathbf{s'}) = \sum_{\mathbf{s}} \Pr(\mathbf{s} \to \mathbf{s'}) \pi(\mathbf{s})$$

There is exactly one stationary distribution for a given $Pr(s \rightarrow s')$ provided the latter obeys some simple conditions.

$\mathbf{w}_{i+1} = \mathbf{w}_i + \epsilon$

where ϵ is zero mean spherical Gaussian and has small variance. Obviously the sequence \mathbf{w}_i does not have the required distribution. However, we can use the *Metropolis algorithm*, which does *not* accept all the steps in the random walk:

1. If $p(\mathbf{w}_{i+1}|\mathbf{y}) > p(\mathbf{w}_i|\mathbf{y})$ then accept the step.

A simple technique is to introduce a random walk, so

2. Else accept the step with probability $\frac{p(\mathbf{w}_{i+1}|\mathbf{y})}{p(\mathbf{w}_i|\mathbf{y})}$.

In practice, the Metropolis algorithm has several shortcomings, and a great deal of research exists on improved methods, see:

R. Neal, "Probabilistic inference using Markov chain Monte Carlo methods," University of Toronto, Department of Computer Science Technical Report CRG-TR-93-1, 1993.

286

Approximate inference for Bayesian networks

The condition of detailed balance

$$\mathbf{A}\mathbf{s}, \mathbf{s}'\pi(\mathbf{s})\operatorname{Pr}(\mathbf{s}\to\mathbf{s}') = \pi(\mathbf{s}')\operatorname{Pr}(\mathbf{s}'\to\mathbf{s}')$$

is sufficient to provide a π that is a stationary distribution. To see this simply sum:

$$\sum_{\mathbf{s}} \pi(\mathbf{s}) \Pr(\mathbf{s} \to \mathbf{s}') = \sum_{\mathbf{s}} \pi(\mathbf{s}') \Pr(\mathbf{s}' \to \mathbf{s})$$
$$= \pi(\mathbf{s}') \underbrace{\sum_{\mathbf{s}} \Pr(\mathbf{s}' \to \mathbf{s})}_{=1}$$
$$= \pi(\mathbf{s}')$$

If all this is looking a little familiar, it's because we now have an excellent application for the material in *Mathematical Methods for Computer Science*. That course used the alternative term *local balance*.

Approximate inference for Bayesian networks Approximate inference for Bayesian networks The evidence is fixed. Let the *state* of our system be a specific set of values for Recalling once again the basic equation for performing probabilistic inference the query variable and the unobserved variables $\Pr(Q|e) = \frac{1}{Z} \Pr(Q \wedge e) = \frac{1}{Z} \sum \Pr(Q, u, e)$ $\mathbf{s} = (q, u_1, u_2, \dots, u_n) = (s_1, s_2, \dots, s_{n+1})$ where and define $\bar{\mathbf{s}}_i$ to be the state vector with s_i removed $\overline{\mathbf{s}}_i = (s_1, \ldots, s_{i-1}, s_{i+1}, \ldots, s_{n+1})$ $\bullet Q$ is the query variable. To move from s to s' we replace one of its elements, say s_i , with a new value s'_i • *e* is the evidence. sampled according to • *u* are the unobserved variables. $s_i' \sim \Pr(S_i | \overline{\mathbf{s}}_i, e)$ • 1/Z normalises the distribution. This has detailed balance, and has Pr(Q, U|e) as its stationary distribution. We are going to consider obtaining samples from the distribution Pr(Q, U|e). 290 289 Approximate inference for Bayesian networks Approximate inference for Bayesian networks To see that Pr(Q, U|e) is the stationary distribution So: $\pi(\mathbf{s})\operatorname{Pr}(\mathbf{s}\to\mathbf{s}') = \operatorname{Pr}(\mathbf{s}|e)\operatorname{Pr}(s'_i|\overline{\mathbf{s}}_i,e)$ • We successively sample the query variable and the unobserved variables, con- $= \Pr(s_i, \overline{\mathbf{s}}_i | e) \Pr(s'_i | \overline{\mathbf{s}}_i, e)$ ditional on their parents, children and children's parents. $= \Pr(s_i | \overline{\mathbf{s}}_i, e) \Pr(\overline{\mathbf{s}}_i | e) \Pr(s'_i | \overline{\mathbf{s}}_i, e)$ • This gives us a sequence s_1, s_2, \ldots which has been sampled according to Pr(Q, U|e). $= \Pr(s_i | \overline{\mathbf{s}}_i, e) \Pr(s'_i, \overline{\mathbf{s}}_i | e)$ $= \Pr(\mathbf{s}' \to \mathbf{s}) \pi(\mathbf{s}')$ Finally, note that as $\Pr(Q|e) = \sum \Pr(Q, u|e)$ As a further simplification, sampling from $\Pr(S_i | \bar{\mathbf{s}}_i, e)$ is equivalent to sampling S_i conditional on its parents, children and children's parents. we can just ignore the values obtained for the unobserved variables. This gives

us q_1, q_2, \ldots with

 $q_i \sim \Pr(Q|e)$

Approximate inference for Bayesian networks

To see that the final step works, consider what happens when we estimate the expected value of some function of Q.

$$\begin{split} \mathbb{E}[f(Q)] &= \sum_{q} f(q) \mathbf{Pr}(q|e) \\ &= \sum_{q} f(q) \sum_{u} \mathbf{Pr}(q, u|e) \\ &= \sum_{q} \sum_{u} f(q) \mathbf{Pr}(q, u|e) \end{split}$$

so sampling using $\Pr(q,u|e)$ and ignoring the values for u obtained works exactly as required.

A (very) brief introduction into how to learn hyperparameters

So far in our coverage of the Bayesian approach to neural networks, the hyperparameters α and β were assumed to be known and fixed. • But this is not a good assumption because... • ... α corresponds to the width of the prior and β to the noise variance. • So we really want to learn these from the data as well. • How can this be done? We now take a look at one of several ways of addressing this problem. 294 The Bayesian approach to neural networks The prior and likelihood depend on α and β respectively so we now make this clear and write $p(\mathbf{w}|\mathbf{y}, \alpha, \beta) = \frac{p(\mathbf{y}|\mathbf{w}, \beta)p(\mathbf{w}|\alpha)}{p(\mathbf{y}|\alpha, \beta)}$ (Don't worry about recalling the actual expressions for the prior and likelihood just yet, they appear in a few slides time.) In the earlier slides we found that the Bayes classifier should in fact compute

$$p(Y|\mathbf{y}, \mathbf{x}, \alpha, \beta) = \int_{\mathbb{R}^W} p(y|\mathbf{w}, \mathbf{x}, \beta) p(\mathbf{w}|\mathbf{y}, \alpha, \beta) \, d\mathbf{w}$$

and we found an approximation to this integral. (Again, the necessary parts of the result are repeated later.)

The Bayesian approach to neural networks

Earlier we looked at the Bayesian approach to *neural networks* using the following notation. We have:

- A neural network computing a function $f(\mathbf{w}; \mathbf{x})$.
- A training sequence $\mathbf{s} = ((\mathbf{x}_1, y_1), \dots, (\mathbf{x}_m, y_m))$, split into

$$\mathbf{y} = (y_1 \ y_2 \ \cdots \ y_m)$$

and

$$\mathbf{X} = (\mathbf{x}_1 \ \mathbf{x}_2 \ \cdots \ \mathbf{x}_m)$$

The *prior distribution* $p(\mathbf{w})$ is now on the weight vectors and Bayes' theorem tells us that

$$p(\mathbf{w}|\mathbf{y}) = \frac{p(\mathbf{y}|\mathbf{w})p(\mathbf{w})}{p(\mathbf{y})}$$

In addition we have a Gaussian prior and a likelihood assuming Gaussian noise.

Hierarchical Bayes and the evidence

Let's write down directly something that might be useful to know:

$$p(\alpha, \beta | \mathbf{y}) = \frac{p(\mathbf{y} | \alpha, \beta) p(\alpha, \beta)}{p(\mathbf{y})}$$

If we know $p(\alpha, \beta | \mathbf{y})$ then a straightforward approach is to use the values for α and β that maximise it.

Here is a standard trick: assume that the prior $p(\alpha,\beta)$ is flat, so that we can just maximise

$$p(\mathbf{y}|\alpha,\beta)$$

This is called *type II maximum likelihood* and is one common way of doing the job.

As usual there are other ways of handling α and $\beta,$ some of which are regarded as more "correct".

297

An expression for the evidence

We have already *derived everything necessary* to write an *explicit equation for the evidence* for the case of regression that we've been following.

First, as we know about a lot of expressions involving w we can introduce it by the standard trick of *marginalising*:

$$p(\mathbf{y}|\alpha,\beta) = \int p(\mathbf{y},\mathbf{w}|\alpha,\beta)d\mathbf{w}$$
$$= \int p(\mathbf{y}|\mathbf{w},\alpha,\beta)p(\mathbf{w}|\alpha,\beta)d\mathbf{w}$$
$$= \int p(\mathbf{y}|\mathbf{w},\beta)p(\mathbf{w}|\alpha)d\mathbf{w}$$

where we've made the obvious independence simplifications.

The two densities in this integral *are just the likelihood and prior we've already studied*.

We've just conditioned on α and β , which previously were constants but are now being treated as random variables.

Hierarchical Bayes and the evidence

The quantity

 $p(\mathbf{y}|\alpha,\beta)$

is called the *evidence*.

When we re-wrote our earlier equation for the posterior density of the weights, making α and β explicit, we found

$$p(\mathbf{w}|\mathbf{y}, \alpha, \beta) = \frac{p(\mathbf{y}|\mathbf{w}, \alpha, \beta)p(\mathbf{w}|\alpha, \beta)}{p(\mathbf{y}|\alpha, \beta)}$$

So the evidence is the denominator in this equation.

This is the *common pattern* and leads to the idea of *hierarchical Bayes*: the *evidence for the hyperparameters* at one level is the *denominator in the relevant application of Bayes theorem*.

298

An expression for the evidence

Here are the actual expression for the prior and likelihood.

The prior is

$$p(\mathbf{w}|\alpha) = \frac{1}{Z_W(\alpha)} \exp\left(-\alpha E_W(\mathbf{w})\right)$$

117/0

where

$$Z_W(\alpha) = \left(\frac{2\pi}{\alpha}\right)^{w/2}$$
 and $E_W(\mathbf{w}) = \frac{1}{2}||\mathbf{w}||^2$

and the likelihood is

$$p(\mathbf{y}|\mathbf{w},\beta) = \frac{1}{Z_{\mathbf{y}}(\beta)} \exp\left(-\beta E_{\mathbf{y}}(\mathbf{w})\right)$$

where

$$Z_{\mathbf{y}}(\beta) = \left(\frac{2\pi}{\beta}\right)^{m/2} \text{ and } E_{\mathbf{y}}(\mathbf{w}) = \frac{1}{2}\sum_{i=1}^{m} (y_i - h(\mathbf{w}; \mathbf{x}_i))^2$$

Both of these equations have been copied directly from earlier slides: *there is nothing to add*.

An expression for the evidence

That gives us

$$p(\mathbf{y}|\alpha,\beta) = \left(\frac{2\pi}{\alpha}\right)^{-W/2} \left(\frac{2\pi}{\beta}\right)^{-m/2} \int \exp\left(-S(\mathbf{w})\right) d\mathbf{w}$$

where

$$S(\mathbf{w}) = \alpha E_W(\mathbf{w}) + \beta E_\mathbf{y}(\mathbf{w})$$

This is exactly the integral we first derived an approximation for.

Specifically

$$\int \exp\left(-S(\mathbf{w})\right) d\mathbf{w} \simeq (2\pi)^{W/2} |\mathbf{A}|^{-1/2} \exp\left(-S(\mathbf{w}_{\mathbf{MAP}})\right)$$

where

$$\mathbf{A} = \alpha \mathbf{I} + \beta \nabla \nabla E_{\mathbf{y}}(\mathbf{w}_{\text{MAP}})$$

and \mathbf{w}_{MAP} is the maximum a posteriori solution.

301

Maximising the evidence

We want to maximise this, so let's differentiate it with respect to α and β .

For α

$$\frac{\partial \log p(\mathbf{y}|\alpha,\beta)}{\partial \alpha} = \frac{W}{2\alpha} - E_W(\mathbf{w}_{\text{MAP}}) - \frac{1}{2} \frac{\partial \log |\mathbf{A}|}{\partial \alpha}$$

How do we handle the final term? This is straightforward if we can compute the *eigenvalues* of **A**.

Recall that the *n* eigenvalues λ_i and *n* eigenvectors \mathbf{v}_i of an $n \times n$ matrix M are defined such that

$$\mathbf{M}\mathbf{v}_i = \lambda_i \mathbf{v}_i \text{ for } i = 1, \dots, n$$

and the eigenvectors are orthonormal

$$\mathbf{v}_i^T \mathbf{v}_j = \begin{cases} 1 & \text{if } i = j \\ 0 & \text{otherwise} \end{cases}$$

One standard result is that *the determinant of a matrix is the product of its eigenvalues*.

$$|\mathbf{M}| = \prod_{i=1}^n \lambda_i$$

An expression for the evidence

Putting all that together we get an *expression for the logarithm of the evidence*:

$$\begin{split} \log p(\mathbf{y}|\alpha,\beta) \simeq & \frac{W}{2} \log \alpha - \frac{m}{2} \log 2\pi + \frac{m}{2} \log \beta \\ & -\frac{1}{2} \log |\mathbf{A}| \\ & -\alpha E_W(\mathbf{w}_{\mathrm{MAP}}) - \beta E_{\mathbf{y}}(\mathbf{w}_{\mathrm{MAP}}) \end{split}$$

Again, we're using the fact that we want to *maximise the evidence* and this is equivalent to *maximising its logarithm* which turns a product into a more friendly sum.

302

Maximising the evidence

We have

$$\mathbf{A} = \alpha \mathbf{I} + \beta \nabla \nabla E_{\mathbf{y}}(\mathbf{w}_{\mathbf{MAP}})$$

Say the eigenvalues of $\beta \nabla \nabla E_{\mathbf{y}}(\mathbf{w}_{\text{MAP}})$ are λ_i . (These can be computed using standard numerical algorithms.)

Then the eigenvalues of **A** are $\alpha + \lambda_i$ and

$$\frac{\partial \log |\mathbf{A}|}{\partial \alpha} = \frac{\partial}{\partial \alpha} \left(\log \prod_{i=1}^{W} (\alpha + \lambda_i) \right)$$
$$= \frac{\partial}{\partial \alpha} \left(\sum_{i=1}^{W} \log(\alpha + \lambda_i) \right)$$
$$= \sum_{i=1}^{W} \frac{1}{\alpha + \lambda_i} \frac{\partial(\alpha + \lambda_i)}{\partial \alpha}$$

This remains tricky because the eigenvalues might be functions of α .

303

Maximising the evidence

To make further progress, assume (sometimes correct, sometimes not!) that the λ_i do not depend on α .

In that case

$$\frac{\partial \log |\mathbf{A}|}{\partial \alpha} = \sum_{i=1}^{W} \frac{1}{\alpha + \lambda_i}$$
$$= \operatorname{Trace}(\mathbf{A}^{-1})$$

because M^{-1} has eigenvalues $1/\lambda_i$ and the trace of a matrix is equal to the sum of its eigenvalues.

Finally, equating the derivative to zero gives:

$$\frac{W}{2\alpha} - E_W(\mathbf{w}_{\text{MAP}}) - \frac{1}{2}\text{Trace}(\mathbf{A}^{-1}) = 0$$

or

$$\alpha = \frac{1}{2E_W(\mathbf{w}_{\text{MAP}})} \left(W - \sum_{i=1}^{W} \frac{\alpha}{\alpha + \lambda_i} \right)$$

which can be used to update the value for α .

305

Maximising the evidence

As the λ_i are the eigenvalues of $\beta \nabla \nabla E_y(\mathbf{w}_{MAP})$ we have

$$\frac{\partial \lambda_i}{\partial \beta} = \frac{\lambda_i}{\beta}$$

(can you see why?) so

$$\frac{\partial \log |\mathbf{A}|}{\partial \beta} = \frac{1}{\beta} \sum_{i=1}^{W} \frac{\lambda_i}{\alpha + \lambda}$$

Equating the derivative to zero gives

$$\beta = \frac{1}{2E_{\mathbf{y}}(\mathbf{w}_{\text{MAP}})} \left(m - \sum_{i=1}^{W} \frac{\lambda_i}{\alpha + \lambda_i} \right)$$

which can be used to update the value for β .

Maximising the evidence

We can now repeat the process to obtain an update for β :

$$\frac{\partial \log p(\mathbf{y}|\alpha,\beta)}{\partial \beta} = \frac{m}{2\beta} - E_{\mathbf{y}}(\mathbf{w}_{\text{MAP}}) - \frac{1}{2} \frac{\partial \log |\mathbf{A}|}{\partial \beta}$$

In this case

$$\frac{\partial \log |\mathbf{A}|}{\partial \beta} = \frac{\partial}{\partial \beta} \left(\sum_{i=1}^{W} \log(\alpha + \lambda_i) \right)$$
$$= \sum_{i=1}^{W} \frac{1}{\alpha + \lambda_i} \frac{\partial}{\partial \beta} (\alpha + \lambda_i)$$
$$= \sum_{i=1}^{W} \frac{1}{\alpha + \lambda_i} \frac{\partial \lambda_i}{\partial \beta}$$

and again we have a *potentially tricky derivative*.

306

Maximising the evidence

Here's why the derivative works.

Say

$$\mathbf{M} = \nabla \nabla E_{\mathbf{y}}(\mathbf{w}_{\mathrm{MAP}})$$

so we're interested in $\partial \lambda_i / \partial \beta$ when the λ_i are the eigenvalues of $\beta \mathbf{M}$. Thus

$$(\beta \mathbf{M})\mathbf{v}_i = \lambda_i \mathbf{v}$$

and using the fact that the eigenvectors are orthonormal

$$\beta \mathbf{v}_i^T \mathbf{M} \mathbf{v}_i = \lambda_i \mathbf{v}_i^T \mathbf{v}_i = \lambda_i.$$

 $\mathbf{v}_i^T \mathbf{M} \mathbf{v}_i = \frac{\lambda_i}{\beta}$ $\frac{\partial \lambda_i}{\partial \beta} = \mathbf{v}_i^T \mathbf{M} \mathbf{v}_i = \frac{\lambda_i}{\beta}.$

308

So

and

Maximising the evidence

Summary:

Define

 $heta_t = \sum_{i=1}^W rac{\lambda_i}{lpha_t + \lambda_i}$

where the subscript denotes the fact that we're using the following equations to periodically update our estimates of α and β .

Collecting the two update equations together we have

$$\alpha_{t+1} = \frac{\theta_t}{2E_W(\mathbf{w}_{\mathrm{MAP}})}$$

and

$$\beta_{t+1} = \frac{m - \theta_t}{2E_{\mathbf{y}}(\mathbf{w}_{\mathrm{MAP}})}$$

An alternative: integrate the hyperparameters out

309

While choosing α and β by maximising the evidence leads to an effective algorithm, it might be argued that a more correct way to deal with these parameters would be to *integrate them out*.

$$p(\mathbf{w}|\mathbf{y}) = \int \int p(\mathbf{w}, \alpha, \beta|\mathbf{y}) d\alpha d\beta.$$

(Recall the *general equation for probabilistic inference* where we integrate out unobserved random variables.)

Re-arranging this we have

$$\begin{split} \int \int p(\mathbf{w}, \alpha, \beta | \mathbf{y}) d\alpha d\beta &= \frac{1}{p(\mathbf{y})} \int \int p(\mathbf{y} | \mathbf{w}, \alpha, \beta) p(\mathbf{w}, \alpha, \beta) d\alpha d\beta \\ &= \frac{1}{p(\mathbf{y})} \int \int p(\mathbf{y} | \mathbf{w}, \alpha, \beta) p(\mathbf{w} | \alpha, \beta) p(\alpha, \beta) d\alpha d\beta \\ &= \frac{1}{p(\mathbf{y})} \int \int p(\mathbf{y} | \mathbf{w}, \beta) p(\mathbf{w} | \alpha) p(\alpha) p(\beta) d\alpha d\beta \end{split}$$

where we're assuming α and β are independent.

Maximising the evidence

This suggests a method for the overall learning process:

- 1. Choose the initial values α_0 and β_0 at random.
- 2. Choose an initial weight vector w according to the prior.
- 3. Use a standard optimisation algorithm to iteratively estimate \mathbf{w}_{MAP} .
- 4. While the optimisation progresses, periodically use the equations above to reestimate α and $\beta.$

Step 4 requires that we compute an eigendecomposition, which might well be time-consuming. If necessary we can make a simplification.

When m >> W it is reasonable to expect that $\theta_t \simeq W$ an so we can use

$$\alpha_{t+1} = \frac{W}{2E_W(\mathbf{w}_{\mathrm{MAP}})}$$

and

$$\beta_{t+1} = \frac{m}{2E_{\mathbf{y}}(\mathbf{w}_{\mathrm{MAP}})}$$

310

An alternative: integrate the hyperparameters out

In order to continue we need to specify priors on α and β .

On this occasion we have a good reason to choose particular priors, as α and β are *scale parameters*.

In general, a scale parameter σ is one that appears in a density of the form

$$p(x|\sigma) = \frac{1}{\sigma} f\left(\frac{x}{\sigma}\right)$$

The standard deviation of a Gaussian density is an example.

What happens to this density if we *scale* x such that x' = cx?

Standard result number 1

We need to recall how to deal with *transformations of continuous random variables*.

Say we have a random variable x with probability density $p_x(x)$.

We then transform x to y = f(x) where f is strictly increasing.

What is the probability density function of y? There is a standard method for computing this. (See NST maths, or the 1A Probability course.)

$$p_y(y) = \frac{p_x(f^{-1}(y))}{f'(f^{-1}(y))}$$

313

Standard result number 2

Returning to the integral of interest

$$\frac{1}{p(\mathbf{y})} \int \int p(\mathbf{y}|\mathbf{w},\beta) p(\mathbf{w}|\alpha) p(\alpha) p(\beta) d\alpha d\beta$$

Taking the integral for α first we have

$$\int p(\mathbf{w}|\alpha)p(\alpha)d\alpha = \int \frac{1}{\alpha Z_W(\alpha)} \exp(-\alpha E_W(\mathbf{w}))d\alpha$$
$$= \int \frac{1}{\alpha} \left(\frac{\alpha}{2\pi}\right)^{W/2} \exp\left(-\frac{\alpha}{2}||\mathbf{w}||^2\right) d\alpha$$

and to evaluate this we use the following *standard result*:

$$\int_0^\infty x^n \exp(-ax) dx = \frac{\Gamma(n+1)}{a^{n+1}}$$

where n > -1 and a > 0. So the integral becomes

$$(2\pi)^{-W/2} \frac{\Gamma(W/2)}{E_W(\mathbf{w})^{W/2}}$$

_ /--- / . .

Applying this when x' = cx we have

$$f(x) = cx$$
$$f^{-1}(x') = \frac{x'}{c}$$
$$f'(x) = c$$

and so

$$p_{x'}(x') = \frac{1}{c\sigma} f\left(\frac{x'}{c\sigma}\right) = \frac{1}{\sigma'} f\left(\frac{x'}{\sigma'}\right)$$

Thus the transformation leaves the density essentially unchanged, and in particular we want the densities $p(\sigma)$ and $p(\sigma')$ to be identical.

It turns out that this forces the choice

$$p(\sigma) = \frac{c'}{\sigma}.$$

This is an *improper prior* and it is conventional to take c' = 1.

314

An alternative: integrate the hyperparameters out

Repeating the process for β and using the same standard result we have

$$\int p(\mathbf{y}|\mathbf{w},\beta)p(\beta)d\beta = \int \frac{1}{\beta} \left(\frac{\beta}{2\pi}\right)^{m/2} \exp(-\beta E_{\mathbf{y}}(\mathbf{w}))d\beta$$
$$= (2\pi)^{-m/2} \frac{\Gamma(m/2)}{E_{\mathbf{y}}(\mathbf{w})^{m/2}}$$

Combining the two expression we obtain

$$\log p(\mathbf{w}|\mathbf{y}) = -\log\left(\frac{1}{p(\mathbf{y})}(2\pi)^{-W/2}\frac{\Gamma(W/2)}{E_W(\mathbf{w})^{W/2}}(2\pi)^{-m/2}\frac{\Gamma(m/2)}{E_{\mathbf{y}}(\mathbf{w})^{m/2}}\right)$$
$$= \frac{W}{2}\log E_W(\mathbf{w}) + \frac{m}{2}\log E_{\mathbf{y}}(\mathbf{w}) + \text{constant}$$

and we want to minimise this so we need

$$\frac{W}{2} \frac{1}{E_W(\mathbf{w})} \frac{\partial E_W(\mathbf{w})}{\partial \mathbf{w}} + \frac{m}{2} \frac{1}{E_{\mathbf{y}}(\mathbf{w})} \frac{\partial E_{\mathbf{y}}(\mathbf{w})}{\partial \mathbf{w}} = 0$$

An alternative: integrate the hyperparameters out

The actual value for the evidence is

$$-\log p(\mathbf{w}|\mathbf{y}) = -\log\left(\frac{1}{p(\mathbf{y})}\frac{1}{Z_{\mathbf{y}}(\alpha,\beta)}\exp(-(\alpha E_W(\mathbf{w}) + \beta E_{\mathbf{y}}(\mathbf{w})))\right)$$
$$= \alpha E_W(\mathbf{w}) + \beta E_{\mathbf{y}}(\mathbf{w}) + \text{constant}$$

and we want to minimise this so we need

$$\alpha \frac{\partial E_W(\mathbf{w})}{\partial \mathbf{w}} + \beta \frac{\partial E_{\mathbf{y}}(\mathbf{w})}{\partial \mathbf{w}} = 0$$

This should make us *VERY VERY HAPPY* because if we equate the two boxed equations we get

$$\alpha = \frac{W}{2E_W(\mathbf{w})}$$

and

$$\beta = \frac{m}{2E_{\mathbf{v}}(\mathbf{w})}$$

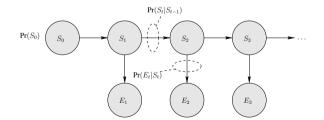
and so the result for *integrating out the hyperparameters* agrees with the result for *optimising the evidence*.

Reinforcement Learning

317

Reinforcement learning and HMMs

Hidden Markov Models (HMMs) are appropriate when our agent models the world as follows



and only wants to infer information about the *state* of the world on the basis of observing the available *evidence*.

This might be criticised as un-necessarily restricted, although it is very effective for the right kind of problem.

We now examine:

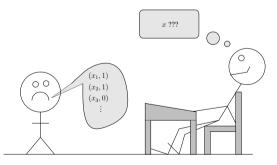
- Some potential shortcomings of hidden Markov models, and of supervised learning.
- An extension know as the Markov Decision Process (MDP).
- The way in which we might *learn from rewards* gained as a result of *acting within an environment*.
- Specific, simple algorithms for performing such learning, and their convergence properties.

Reading: Russell and Norvig, chapter 21. Mitchell chapter 13.

318

Reinforcement learning and supervised learning

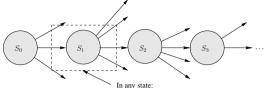
Supervised learners learn from *specifically labelled chunks of information*:



This might also be criticised as un-necessarily restricted: there are other ways to learn.

Reinforcement learning: the basic case

We now begin to model the world in a more realistic way as follows:



Perform an action a to move to a new state. (There may be many possibilities.) Receive a reward r depending on the start state and action.

The agent can *perform actions* in order to *change the world's state*.

If the agent performs an action in a particular state, then it *gains a corresponding reward*.

321

Deterministic Markov Decision Processes

From the point of view of the agent, there is a matter of considerable importance:

The agent does not have access to the functions \mathcal{S} and \mathcal{R}

It therefore has to *learn* a *policy*, which is a function

 $p:S\to A$

such that $p(\boldsymbol{s})$ provides the action \boldsymbol{a} that should be executed in state $\boldsymbol{s}.$

What might the agent use as its criterion for learning a policy?

Deterministic Markov Decision Processes

Formally, we have a set of states

$$S = \{s_1, s_2, \ldots, s_n\}$$

and in each state we can perform one of a set of actions

$$A = \{a_1, a_2, \ldots, a_m\}.$$

We also have a function

$$\mathcal{S}:S\times A\to S$$

such that $\mathcal{S}(s,a)$ is the new state resulting from performing action a in state s, and a function

 $\mathcal{R}: S \times A \to \mathbb{R}$

such that $\mathcal{R}(s, a)$ is the *reward* obtained by executing action a in state s.

322

Measuring the quality of a policy

Say we start in a state at time t, denoted s_t , and we follow a policy p. At each future step in time we get a reward. Denote the rewards r_t, r_{t+1}, \ldots and so on.

A common measure of the quality of a policy \boldsymbol{p} is the discounted cumulative reward

$$V^{p}(s_{t}) = \sum_{i=0}^{\infty} \epsilon^{i} r_{t+i}$$
$$= r_{t} + \epsilon r_{t+1} + \epsilon^{2} r_{t+2} + \cdot$$

where $0 \le \epsilon \le 1$ is a constant, which defines a trade-off for how much we value immediate rewards against future rewards.

The intuition for this measure is that, on the whole, we should like our agent to prefer rewards gained quickly.

Measuring the quality of a policy	Two important issues
Other common measures are the <i>average reward</i> $\lim_{T\to\infty} \frac{1}{T} \sum_{i=0}^{T} r_{t+i}$ and the <i>finite horizon reward</i> $\sum_{i=0}^{T} r_{t+i}$ In these notes we will only address the discounted cumulative reward.	 Note that in this kind of problem we need to address two particularly relevant issues: The <i>temporal credit assignment</i> problem: that is, how do we decide which specific actions are important in obtaining a reward? The <i>exploration/exploitation</i> problem. How do we decide between <i>exploiting</i> the knowledge we already have, and <i>exploring</i> the environment in order to possibly obtain new (and more useful) knowledge? We will see later how to deal with these.
325	326
$\frac{\text{The optimal policy}}{\text{Ultimately, our learner's aim is to learn the optimal policy} p_{\text{opt}} = \underset{p}{\operatorname{argmax}} V^p(s)$ for all s. We will denote the optimal discounted cumulative reward as $V_{\text{opt}}(s) = V^{p_{\text{opt}}}(s).$ How might we go about learning the optimal policy?	$\label{eq:linear} \underbrace{ \text{Learning the optimal policy} } \\ \text{The only information we have during learning is the individual rewards obtained from the environment.} \\ \text{We could try to learn } V_{\text{opt}}(s) \text{ directly, so that states can be compared:} \\ \text{Consider } s \text{ as better than } s' \text{ if } V_{\text{opt}}(s) > V_{\text{opt}}(s'). \\ \text{However we actually want to compare } actions, not states. Learning } V_{\text{opt}}(s) \text{ might help as} \\ p_{\text{opt}}(s) = \operatorname*{argmax}_{a} \left[\mathcal{R}(s, a) + \epsilon V_{\text{opt}}(\mathcal{S}(s, a)) \right] \\ \text{but } only \text{ if } we \text{ know } \mathcal{S} \text{ and } \mathcal{R}. \\ \text{As we are interested in the case where these functions are not known, we need something slightly different.} \\ \end{aligned}$

The Q function

The trick is to define the following function:

$$\mathcal{Q}(s,a) = \mathcal{R}(s,a) + \epsilon V_{\text{opt}}(\mathcal{S}(s,a))$$

This function specifies the discounted cumulative reward obtained if you do action *a* in state *s* and then follow the optimal policy.

As

$$p_{\text{opt}}(s) = \operatorname{argmax} \mathcal{Q}(s, a)$$

then provided one can learn Q it is not necessary to have knowledge of S and R to obtain the optimal policy.

The Q function

Note also that

$$V_{\text{opt}}(s) = \max_{\alpha} \mathcal{Q}(s, \alpha)$$

and so

$$\mathcal{Q}(s,a) = \mathcal{R}(s,a) + \epsilon \max_{\alpha} \mathcal{Q}(\mathcal{S}(s,a),\alpha)$$

which suggests a simple learning algorithm.

Let Q' be our learner's estimate of what the exact Q function is.

That is, in the current scenario Q' is a table containing the estimated values of Q(s, a) for all pairs (s, a).

329

$\mathcal{Q}\text{-learning}$

Start with all entries in Q' set to 0. (In fact we will see in a moment that random entries will do.)

Repeat the following:

- 1. Look at the current state *s* and choose an action *a*. (We will see how to do this in a moment.)
- 2. Do the action a and obtain some reward $\mathcal{R}(s, a)$.
- 3. Observe the new state S(s, a).
- 4. Perform the update

$$Q'(s,a) = \mathcal{R}(s,a) + \epsilon \max_{\alpha} Q'(\mathcal{S}(s,a),\alpha)$$

Note that this can be done in *episodes*. For example, in learning to play games, we can play multiple games, each being a single episode.

Convergence of *Q*-learning

330

This looks as though it might converge!

Note that, if the rewards are at least 0 and we initialise Q' to 0 then,

$$\forall n, s, a \ Q'_{n+1}(s, a) \ge Q'_n(s, a)$$

and

 $\forall n,s,a \; Q(s,a) \geq Q_n'(s,a) \geq 0$

However, we need to be a bit more rigorous than this...

Convergence of Q-learning

Convergence of Q-learning

If:

- 1. The agent is operating in an environment that is a deterministic MDP.
- 2. Rewards are bounded in the sense that there is a constant $\delta > 0$ such that

 $\forall s, a \ |\mathcal{R}(s, a)| < \delta$

3. All possible pairs s and a are visited infinitely often.

Then the \mathcal{Q} -learning algorithm converges, in the sense that

$$\forall a, s \ Q'_n(s, a) \to \mathcal{Q}(s, a)$$

as $n \to \infty$.

333

Convergence of Q-learning

We have,

$$\begin{aligned} |Q'_{n+1}(s,a) - \mathcal{Q}(s,a)| \\ &= |(\mathcal{R}(s,a) + \epsilon \max_{\alpha} Q'_n(\mathcal{S}(s,a),\alpha)) - (\mathcal{R}(s,a) + \epsilon \max_{\alpha} \mathcal{Q}(\mathcal{S}(s,a),\alpha))| \\ &= \epsilon |\max_{\alpha} Q'_n(\mathcal{S}(s,a),\alpha) - \max_{\alpha} \mathcal{Q}(\mathcal{S}(s,a),\alpha)| \\ &\leq \epsilon \max_{\alpha} |Q'_n(\mathcal{S}(s,a),\alpha) - \mathcal{Q}(\mathcal{S}(s,a),\alpha)| \\ &\leq \epsilon \max_{s,a} |Q'_n(s,a) - \mathcal{Q}(s,a)| \\ &= \epsilon \xi(n). \end{aligned}$$

Convergence as described follows.

Using condition 3, take two stretches of time in which all s and a pairs occur:

Define

$$\xi(n) = \max_{s,a} |Q'_n(s,a) - \mathcal{Q}(s,a)|$$

the maximum error in Q' at n.

This is straightforward to demonstrate.

What happens when $Q'_n(s, a)$ is updated to $Q'_{n+1}(s, a)$?

334

Choosing actions to perform

We have not yet answered the question of how to choose actions to perform during learning.

One approach is to choose actions based on our current estimate Q'. For instance

action chosen in current state $s = \operatorname{argmax} Q'(s, a)$.

However we have already noted the trade-off between exploration and exploitation. It makes more sense to:

- Explore during the early stages of training.
- *Exploit* during the later stages of training.

This seems particularly important in the light of condition 3 of the convergence proof.

<u>Choosing actions to perform</u> One way in which to choose actions that incorporates these requirements is to introduce a constant λ and choose actions <i>probabilistically</i> according to $Pr(action \ a state \ s) = \frac{\lambda^{Q'(s,a)}}{\sum_a \lambda^{Q'(s,a)}}$ Note that: • If λ is <i>small</i> this promotes <i>exploration</i> . • If λ is <i>large</i> this promotes <i>exploitation</i> . We can vary λ as training progresses.	Improving the training process There are two simple ways in which the process can be improved: If training is episodic, we can store the rewards obtained during an episode and update <i>backwards</i> at the end. This allows better updating at the expense of requiring more memory. We can remember information about rewards and occasionally <i>re-use</i> it by re-training.
337	338
Nondeterministic MDPs	\underline{Q} -learning for nondeterministic MDPs
The Q-learning algorithm generalises easily to a more realistic situation, where the outcomes of actions are <i>probabilistic</i> . Instead of the functions S and R we have <i>probability distributions</i> Pr(new state current state, action) and Pr(reward current state, action). and we now use $S(s, a)$ and $\mathcal{R}(s, a)$ to denote the corresponding random variables. We now have $V^p = \mathbb{E}\left(\sum_{i=0}^{\infty} \epsilon^i r_{t+i}\right)$ and the best policy p_{opt} maximises V^p .	We now have $\begin{aligned} \mathcal{Q}(s,a) &= \mathbb{E}(\mathcal{R}(s,a)) + \epsilon \sum_{\sigma} \Pr(\sigma s,a) V^{\text{opt}}(\sigma) \\ &= \mathbb{E}(\mathcal{R}(s,a)) + \epsilon \sum_{\sigma} \Pr(\sigma s,a) \max_{\alpha} \mathcal{Q}(\sigma,\alpha) \end{aligned}$ and the rule for learning becomes $\begin{aligned} \mathcal{Q}'_{n+1} &= (1 - \theta_{n+1}) \mathcal{Q}'_n(s,a) + \theta_{n+1} \left[\mathcal{R}(s,a) + \max_{\alpha} \mathcal{Q}'_n(\mathcal{S}(s,a),\alpha) \right] \end{aligned}$ with $\begin{aligned} \theta_{n+1} &= \frac{1}{1 + v_{n+1}(s,a)} \end{aligned}$ where $v_{n+1}(s,a)$ is the number of times the pair s and a has been visited so far.

Convergence of Q-learning for nondeterministic MDPs

The agent is operating in an environment that is a nondeterministic MDP.
 Rewards are bounded in the sense that there is a constant δ > 0 such that

3. All possible pairs *s* and *a* are visited infinitely often.

4. $n_i(s, a)$ is the *i*th time that we do action *a* in state *s*.

 $\forall s, a \ |\mathcal{R}(s, a)| < \delta$

Convergence of Q-learning for nondeterministic MDPs

...we have

$$0 \leq \theta_n < 1$$
$$\sum_{i=1}^{\infty} \theta_{n_i(s,a)} = \infty$$
$$\sum_{i=1}^{\infty} \theta_{n_i(s,a)}^2 < \infty$$

then with probability 1 the Q-learning algorithm converges, in the sense that

 $\forall a, s \ Q'_n(s, a) \to \mathcal{Q}(s, a)$

342

as $n \to \infty$.

Alternative representation for the Q' table

341

But there's always a catch...

If:

and also ...

We have to store the table for Q':

- Even for quite straightforward problems it is HUGE!!! certainly big enough that it can't be stored.
- A standard approach to this problem is, for example, to represent it as a *neural network*.
- One way might be to make s and a the inputs to the network and train it to produce Q'(s, a) as its output.

This, of course, introduces its own problems, although it has been used very successfully in practice.

It might be covered in *Artificial Intelligence III*, which unfortunately does not yet exist.

343