

Large Deviations. A survey

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-

$$P_n(A) = \exp\left[-n \inf_{x \in A} I(x) + o(n)\right]$$

for "nice" sets A .

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- For closed sets $C \subset X$,

$$P_n(C) \leq \exp[-n \inf_{x \in C} I(x) + o(n)]$$

- Equivalently

$$\int \exp[nF(x)]dP_n = \exp[n \sup_x [F(x) - I(x)] + o(n)]$$

for bounded continuous functions $F(x)$ on X .

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for bounded continuous functions $F(x)$ on X .

- The non-negative rate function $I(\cdot)$ is assumed to be lower semicontinuous and with compact level sets

$$K_\ell = \{x : I(x) \leq \ell\}$$

- If X_i are i.i.d. random variables with finite exponential moments

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- $$I(x) = \sup_{\theta} [\theta x - \log M(\theta)]$$

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- More generally $\{X_i\}$ could be i.i.d random variables with values in some X with a common distribution α

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- The dual of $\mathcal{M}(X)$ is $C(X)$ and

$$M(f) = E^\alpha[\exp[\langle \delta_X, f \rangle]] = \int e^{f(x)} d\alpha(x)$$

and for $\beta \in \mathcal{M}(X)$


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$$h(\beta; \alpha) = \int \log b(x) d\beta = \int b(x) \log b(x) d\alpha$$

Contraction principle

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- If P_n on X satisfies an LDP with rate $I(x)$ and $F : X \rightarrow Y$ is a continuous map then $Q_n = P_n F^{-1}$ on Y satisfies an LDP with rate

$$J(y) = \inf_{x:F(x)=y} I(x)$$

We optimize when we project.

- $\bar{X}_n = \int x d\nu_n$. $\nu \rightarrow \int x d\nu$ maps $\mathcal{M} \rightarrow R$.

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- $\bar{X}_n = \int x d\nu_n$. $\nu \rightarrow \int x d\nu$ maps $\mathcal{M} \rightarrow R$.

- $$\inf_{\beta: \int x d\beta(x)=a} h(\beta; \alpha) = I(a)$$

- $$= \sup_{\theta \in R} [a\theta - \log \int e^{\theta x} d\alpha(x)]$$

- The next step is to try and calculate $\psi_P(F)$

$$\lim_{n \rightarrow \infty} \frac{1}{n} \log E^P \left[\exp \left[\sum_{i=1}^n F(x_i, x_{i+1}, \dots, x_{i+k-1}) \right] \right]$$

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- F is a function of k variables (x_1, \dots, x_k)
- P is stationary process with values in X . i.e a shift invariant probability measure on X^∞ , i.e $P \in \mathcal{M}_s(X)$.

- The empirical process which looks at all the finite dimensional distributions

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- Start with (x_1, x_2, \dots, x_n) extend it periodically to get a sequence $\omega \in X^\infty$ and consider the orbital measure

$$R_n(x_1, x_2, \dots, x_n) = \frac{1}{n} \sum_{i=1}^n \delta_{T^{i-1}\omega}$$

- Let P_n be the distribution of $R_n \in \mathcal{M}(\mathcal{M}_s(X))$ It satisfies an LDP with rate function $I_P(Q)$ and

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- The rate function $I_P(Q)$ is universal and is a version of Kolmogorov-Sinai entropy.

- Let $p(dx_1|\omega)$ and $q(dx_1|\omega)$ be the conditional distributions of x_1 given the past $\{x_i : i \leq 0\}$ under P and Q respectively. Then

$$I_P(Q) = E^Q[h(q(\cdot|\omega); p(\cdot|\omega))]$$

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- The problem is that $p(\cdot|\omega)$ is only defined a.e. P and we need to integrate with respect to Q .
- Put assumptions on P so that $p(\cdot|\omega)$ has a nice everywhere defined version. Markov will do it.

- if F is only a function of one variable $F(x_1)$ one can contract

$$I_P(\beta) = \inf_{\substack{Q \in \mathcal{M}_s(X) \\ Qx_1^{-1} = \beta}} I_P(Q)$$

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- Controls the large deviations of $\frac{1}{n} \sum \delta_{x_i}$
- If $\pi(x, dy) = \alpha(dy)$ then we are back in the Sanov case.

- We now turn to a more general problem. Calculate

$$J = \lim_{n \rightarrow \infty} \frac{1}{n} E^P \left[\exp \left[\sum_{i=1}^n a_i F(x_i, x_{i+1}, \dots, x_{i+k-1}) \right] \right]$$

for a given sequence $\{a_i : i \geq 1\}$.

- When will it exist?

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- What will it be?

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 - What will it be?
 - What is it good for?

- We note that when P is a product measure and F is a function of one variable, we need to have the limit

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- This requires the empirical distribution $\frac{1}{n} \sum_{i=1}^n \delta_{a_i}$ to have a limit.

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- More generally we can assume that all the finite dimensional joint distributions $\frac{1}{n} \sum_{i=1}^n \delta_{a_i, a_{i+1}, \dots, a_{i+k-1}}$ have limits.
- The empirical process $R_n(a_1, \dots, a_n)$ has a limit $\nu \in \mathcal{M}_s(X)$
- It looks like a sample from ν .

- If $\{x_i\}$ is a Markov process with positive transition probabilities on a finite set X , then for every $\{a_i\}$ such that $R_n(a_1, \dots, a_n) \rightarrow \nu$ the limit $J = J(\nu)$ exists and is a continuous linear function of ν

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- Sub-additive argument will do it.
- What is it?

- Let \mathcal{M}_ν be the set of stationary process Q with values in $R \times X$ such that the marginal on R^∞ is ν .

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- Let $Q_0 = \nu \times P$, i.e $\{a_i\}$ has ν for its distribution and while $\{x_i\}$ is distributed according to P , the two components are independent.

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- $q_0(dx_1|\omega) = q_0(dx_1|\omega_1), \nu(da_1|\omega) = \nu(da_1|\omega_2)$
- Then $J(\nu)$ is equal to

$$\sup_{Q \in \mathcal{M}_\nu} [E^Q[a_1 F(x_1, \dots, x_k)] - H(Q; Q_0)]$$

- Define $H_n(a_1, \dots, a_n)$

$$\log E^P \left[\exp \left[\sum_{i=1}^n a_i F(x_i, x_{i+1}, \dots, x_{i+k-1}) \right] \right]$$

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$$\begin{aligned} & |H_{n+m}(a_1, \dots, a_{n+m}) - H_n(a_1, \dots, a_n) \\ & \quad - H_m(a_1, \dots, a_m)| \leq C \end{aligned}$$

uniformly in n, m and $\{a_i\}$.

$$\left| H_{nk}(a_1, \dots, a_{nk}) - \sum_{i=1}^k H_n(a_{(i-1)k+1}, \dots, a_{ik}) \right| \leq Cn$$

$$\begin{aligned}
 & \left| H_{nk}(a_1, \dots, a_{nk}) - \sum_{i=1}^k H_n(a_{(i-1)k+1}, \dots, a_{ik}) \right| \\
 & \leq Cn
 \end{aligned}$$

- Partitioning of a block of size kn into blocks of size k allows some freedom as to the location and one can average over this collection

$$\limsup_{n \rightarrow \infty} \left| \frac{1}{n} H_n(a_1, \dots, a_n) - \int \frac{1}{k} H_k(a_1, \dots, a_k) dR_n \right| \leq \frac{C}{k}$$

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- Then for every ν typical $\{x_i\}$ the limit

$$\lim_{n \rightarrow \infty} \frac{1}{n} A_n(x_1, x_2, \dots, x_n) = \mathcal{A}(\nu)$$

exists and is a continuous linear functional on $\mathcal{M}_s(X)$

$$\begin{aligned}\mathcal{H}(\nu) &= \lim_{n \rightarrow \infty} \frac{1}{n} H_n(a_1, a_2, \dots, a_n) \\ &= \lim_{k \rightarrow \infty} \int H_k(a_1, a_2, \dots, a_k) d\nu\end{aligned}$$

exists and depends (linearly and continuously) on ν .

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- If $H_n((a_1, x_1), \dots, (a_n, x_n))$ is an almost additive sequence on $\{(A \times X)^n\}$ with $\frac{H_n}{n} \rightarrow \mathcal{H}(\cdot)$,

- A slight variant is the following theorem.
- If $H_n((a_1, x_1), \dots, (a_n, x_n))$ is an almost additive sequence on $\{(A \times X)^n\}$ with $\frac{H_n}{n} \rightarrow \mathcal{H}(\cdot)$,
- $K_n(a_1, \dots, a_n)$ defined as

$$\log E^P[\exp[H_n((a_1, x_1), \dots, (a_n, x_n))]]$$

is almost additive.

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$$\mathcal{K}(\nu) = \sup_{Q \in \mathcal{M}_\nu} \left[\mathcal{H}(Q) - H_{\nu \times P}(Q) \right]$$

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- To begin with consider the following simpler problem. $\{x_i\}$ are i.i.d and U, V are bounded. Let n be large say $2^k n$ for large k .

- We split the sum

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$$S_n^1 = \sum_{j=1}^n U(x_j) V(x_{2j})$$

- For $i \geq 2$

$$S_n^i = \sum_{j=2^{i-1}n+1}^{2^i n} U(x_j) V(x_{2j})$$

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- S_n^i is measurable w.r.t $\mathcal{F}_{2^{i+1}n}$

$$\begin{aligned} E[e^{S_{2^k n}} | \mathcal{F}_{2^k n}] &= e^{S_{2^{k-1}n}} E[e^{S_n^k} | \mathcal{F}_{2^k n}] \\ &= e^{S_{2^{k-1}n}} e^{\sum_{j=2^{k-1}n+1}^{2^k n} \psi_1(x_j)} \end{aligned}$$

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- where $\psi_1(x) = \log E^y[e^{U(x) V(y)}]$. The odd j 's do not appear again and hence $2^{k-2}n$ terms separate out.

- We can remove $2^{k-2}n$ factors of

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- Now condition with respect to $\mathcal{F}_{2^{k-1}n}$. We get

$$e^{S_{2^{k-2}n}} E\left[e^{\sum_{j=2^{k-2}n+1}^{2^{k-1}n} [U(x_j)V(x_{2j}) + \psi_1(x_{2j})]} \mid \mathcal{F}_{2^{k-1}n}\right]$$

- If we define

$$\psi_2(x) = \log E[e^{U(x)V(y)+\psi_1(y)}]$$

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- This reduces to

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- Again the odd ones stand alone. Can factor $2^{k-3}n$ factors of

$$E[e^{\psi_2(x)}] = E[e^{U(x)V(y)+\psi_1(y)}] = e^{c_2}$$

- Left with

$$e^{S_{2^{k-2}n}} e^{\sum_{j=2^{k-3}n+1}^{2^{k-2}n} \psi_2(x_{2j})}$$

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- Continuing recursively

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- Continuing recursively

$$\psi_k(x) = \log E[e^{U(x)V(y)+\psi_{k-1}(y)}]$$



$$E[e^{\psi_k(x)}] = E[e^{U(x)V(y)+\psi_{k-1}(y)}] = e^{c_k}$$

We have some terms left over from $1 \leq j \leq n$. But for k large they can be ignored. Hence the limit equals

$$\lim_{n \rightarrow \infty} \frac{1}{n} \log E \left[\exp \left[\sum_{i=1}^n U(x_i) V(x_{2i}) \right] \right]$$
$$= \sum_{j=1}^{\infty} \frac{c_j}{2^{j+1}}$$

$$\lim_{n \rightarrow \infty} \frac{1}{n} \log E \left[\exp \left[\sum_{i=1}^n U(x_i) V(x_{2i}) \right] \right] \\ = \sum_{j=1}^{\infty} \frac{c_j}{2^{j+1}}$$

- Now we turn to the more general case where $\{x_i\}$ is a finite state space Markov chain that is mixing.

■ Let

$$f_n = \frac{1}{n} \log E^P \left[\exp \left[\sum_{i=1}^n U(x_i) V(x_{2j}) \right] \right]$$

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- Replace n by $2^k n$. We can pretend x_{2j} for $j \geq 2^{k-1}n + 1$ are independent of what happened before i.e. x_j for $j \leq 2^k n$.

$$\begin{aligned}
f_{2^k n} &= \frac{1}{2^k n} \log E^{P \times P} \left[\exp \left[\sum_{i=1}^{2^{k-1} n} U(x_i) V(x_{2j}) \right. \right. \\
&\quad \left. \left. + \sum_{i=2^{k-1} n+1}^{2^k n} U(x_i) V(y_{2j}) \right] \right] \\
&= \frac{1}{2^k n} \log E^P \left[\exp \left[\sum_{i=1}^{2^{k-1} n} U(x_i) V(x_{2j}) \right. \right. \\
&\quad \left. \left. + H_{2^{k-1} n} (U(x_{2^{k-1} n+1}), \dots, U(x_{2^k n})) \right] \right]
\end{aligned}$$

- Now we peel off $j \leq 2^{k-2}n$ and pretend the rest is independent.

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$$\begin{aligned}
 f_{2^k n} &\simeq \frac{1}{2^k n} \log E^{P \times P} \left[\exp \left[\sum_{i=1}^{2^{k-2} n} U(x_i) V(x_{2j}) \right. \right. \\
 &+ \left. \left. \sum_{i=2^{k-2} n+1}^{2^{k-1} n} U(x_i) V(y_{2j}) \right] \right] \\
 &+ H_{2^{k-1} n} \left(U(y_{2^{k-1} n+1}), \dots, U(y_{2^k n}) \right) \left. \right]
 \end{aligned}$$

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$$H_{2^{k-r-1}n}^{r+1}(a_1, \dots, a_{2^{k-r-1}n})$$

$$= \log E^P \left[\exp \left[\sum_{i=2^{k-r-1}+1}^{2^{k-r}} a_i V(x_{2i}) \right. \right.$$

$$\left. \left. + H_{2^{k-r}n}^r(U(x_{2^{k-r}n+1}), \dots, U(x_{2^{k-r+1}n})) \right] \right]$$

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$$\begin{aligned}
 & H_{2^{k-r-1}n}^{r+1}(a_1, \dots, a_{2^{k-r-1}n}) \\
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 & \quad \left. \left. + H_{2^{k-r}n}^r(U(x_{2^{k-r}n+1}), \dots, U(x_{2^{k-r+1}n})) \right] \right]
 \end{aligned}$$

- This leads to a map

$$\mathcal{H}_r(\cdot) \rightarrow \mathcal{H}_{r+1}(\cdot)$$

- Starting from $\mathcal{H}_0(\cdot) \equiv 0$. After many iterations $\mathcal{H}_r(\nu)$ will be nearly a constant. Or we drop the first n terms and just calculate

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- Remark. One can extend this to the calculation of

$$\lim_{n \rightarrow \infty} \frac{1}{n} \log E^P [\exp[\sum_{i=1}^n f(x_i, x_{2i}, \dots, x_{ki})]]$$

Last Slide

THE END