

# SKETCH OF THE “ALTERNATING SQP” METHOD FOR FITTING POISSON TOPIC MODELS

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**1. Problem statement.** Given an  $n \times p$  matrix of counts  $X$  with entries  $x_{ij} \geq 0$ , the aim is to fit a Poisson model of the counts,

$$(1) \quad \begin{aligned} p(x) &= \prod_{i=1}^n \prod_{j=1}^p p(x_{ij}) \\ &= \prod_{i=1}^n \prod_{j=1}^p \text{Poisson}(x_{ij}; \lambda_{ij}), \end{aligned}$$

in which the Poisson rates are given by  $\lambda_{ij} = \sum_{k=1}^K l_{ik} f_{jk}$ . The model is determined by a  $p \times K$  matrix  $F$  with entries  $f_{ik} \geq 0$  (the “factors”) and an  $n \times K$  matrix  $L$  with entries  $l_{ik} \geq 0$  (the “loadings”). Fitting  $F$  and  $L$  is equivalent to non-negative matrix factorization with the “beta-divergence” cost function [3]. It can also be used to recover a maximum-likelihood estimate for the latent Dirichlet allocation (LDA) model [1]. So fitting this model is useful for a wide range of applications.

The log-likelihood for the Poisson model is

$$(2) \quad \log p(x | F, L) \propto \sum_{i=1}^n \sum_{j=1}^p x_{ij} \log(\sum_{k=1}^K l_{ik} f_{jk}) - \sum_{i=1}^n \sum_{j=1}^p \sum_{k=1}^K l_{ik} f_{jk},$$

where the constant of proportionality is obtained from factorial terms in the Poisson densities. Our specific aim is to find a  $F$  and  $L$  that maximizes the log-likelihood (2); that is, we would like to solve

$$(3) \quad \begin{aligned} &\text{minimize} && \ell(F, L) \equiv -\log p(x | F, L) \\ &\text{subject to} && F \geq 0, L \geq 0. \end{aligned}$$

In the next section, we derive an efficient approach to doing this.

**2. Block co-ordinate descent strategy.** Our strategy for solving (3) is to alternate between solving for  $F$  with  $L$  fixed, and solving for  $L$  with  $F$  fixed. When solving for  $F$  with  $L$  fixed (and vice versa), the problem naturally decomposes into a collection of much smaller subproblems that are much more tractable to solve. All the subproblems are of the following form:

$$(4) \quad \begin{aligned} &\text{minimize} && \phi(y; B, w) \\ &\text{subject to} && y_k \geq 0 \text{ for all } k = 1, \dots, K, \end{aligned}$$

in which the objective function is defined as

$$(5) \quad \phi(y; B, w) = - \sum_{i=1}^n w_i \log \left( \sum_{k=1}^K b_{ik} y_k \right) + \sum_{i=1}^n \sum_{k=1}^K b_{ik} y_k.$$

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To see the connection between subproblem (4) and the original optimization problem (3), observe that the negative log-likelihood can be recovered as

$$(6) \quad -\log p(x | F, L) = \sum_{i=1}^n \phi(l_i; F, x_i),$$

where  $x_i$  is the  $i$ th row of  $X$  and  $l_i$  is the  $i$ th row of  $L$ . Alternatively, it can be recovered as

$$(7) \quad -\log p(x | F, L) = \sum_{j=1}^p \phi(f_j; L, x_j),$$

in which  $x_j$  is the  $j$ th column of  $X$ , and  $f_j$  is the  $j$ th row of  $F$ . Therefore, when  $F$  is fixed, each row of  $L$  can be separately optimized by solving a problem of the form (4), and when  $L$  is fixed, each row of  $F$  can be separately optimized by solving a problem of the form (4).

Initially this may seem like a sensible strategy, but directly optimizing (4) turns out to be difficult to do for numerical reasons: in practice, the entries of  $B$  can be very large or very small, resulting in solutions  $y$  in which all the entries are either very large or very small. This makes it difficult to devise an algorithm that will work well for all possible input matrices  $B$ .

I propose to solve for  $y$  indirectly by instead solving

$$(8) \quad \begin{aligned} &\text{minimize} && f(t; P, u) \\ &\text{subject to} && t_k \geq 0 \text{ for all } k = 1, \dots, K, \end{aligned}$$

in which the new objective function is

$$(9) \quad f(t; P, u) = -\sum_{i=1}^n u_i \log \left( \sum_{k=1}^K p_{ik} t_k \right) + \sum_{k=1}^K t_k,$$

where I've defined

$$\begin{aligned} u_i &= \frac{w_i}{\sum_{i'=1}^n w_{i'}} \\ p_{ik} &= b_{ik} \times \frac{\sum_{i'=1}^n w_{i'}}{\sum_{i'=1}^n b_{i'k}}. \end{aligned}$$

After finding the solution  $t^*$  to (8), the solution  $y^*$  to (4) is recovered as

$$(10) \quad y_k^* = t_k^* \times \frac{\sum_{i=1}^n w_i}{\sum_{i=1}^n b_{ik}}.$$

The main advantage of solving (8) is that the solution is numerically well behaved; in particular, the entries of the solution  $t^*$  sum to 1 [2]. A straightforward way to appreciate this result is to compare the complementary slackness condition  $\sum_{k=1}^K g_k x_k = 0$  for the original and modified subproblems, where  $g_k$  denotes the partial derivative of the objective (either  $\phi$  or  $f$ ) with respect to the  $k$ th co-ordinate (either  $y_k$  or  $t_k$ ). Further, since problem (8) is obtained by a simple linear transformation of the variables, any iterate that leads to an improvement in the solution to the modified subproblem will produce an improvement in the solution to the original subproblem, and vice versa.

**3. Karush-Kuhn-Tucker conditions.** To assess the quality of a solution, here we derive the first-order Karush-Kuhn-Tucker (KKT) conditions for (3). Written in matrix notation, they are

$$(11) \quad \nabla_F \ell^*(F, L, \Omega, \Gamma) = 0$$

$$(12) \quad \nabla_L \ell^*(F, L, \Omega, \Gamma) = 0$$

$$(13) \quad \Omega \odot F = 0$$

$$(14) \quad \Gamma \odot L = 0.$$

Here, I have introduced matrices of Lagrange multipliers  $\Omega \geq 0, \Gamma \geq 0$  associated with the non-negativity constraints  $F \geq 0, L \geq 0$ ,  $A \odot B$  is the Hadamard (elementwise) product of matrices  $A$  and  $B$ , and  $\ell^*(F, L, \Omega, \Gamma)$  is the Lagrangian function:

$$(15) \quad \phi(F, L, \Omega, \Gamma) = \ell(F, L) - \sum_{i=1}^n \sum_{k=1}^K \gamma_{ik} l_{ik} - \sum_{j=1}^m \sum_{k=1}^K \omega_{jk} f_{jk}.$$

Applying the condition that the Lagrangian vanishes (11, 12), we obtain

$$(16) \quad \Omega = (1 - A)^T L$$

$$(17) \quad \Gamma = (1 - A)F,$$

in which  $A$  is defined to be the  $n \times m$  matrix with entries  $a_{ij} = x_{ij}/\lambda_{ij}$ . Next, applying the complementary conditions (13, 14) results in the following expressions for the residuals of the complementary conditions:

$$(18) \quad r_F = F(1 - A)^T L$$

$$(19) \quad r_L = L(1 - A)F$$

These residuals should vanish near a minimizer of (3).

#### REFERENCES

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