18.335 Problem Set 2 Solutions

Problem 1: (14+(10+5) points)

- (a) Trefethen, exercise 15.1. In the following, I abbreviate $\epsilon_{\text{machine}} = \epsilon_m$, and I use the fact (which follows trivially from the definition of continuity) that we can replace any Lipshitz-continuous $g(O(\epsilon))$ with $g(0) + g'(0)O(\epsilon)$. I also assume that fl(x) is deterministic—by a stretch of Trefethen's definitions, it could conceivably be nondeterministic in which case one of the answers changes as noted below, but this seems crazy to me (and doesn't correspond to any real machine). Note also that, at the end of lecture 13, Trefethen points out that the same axioms hold for complex floating-point arithmetic as for real floating-point arithmetic (possibly with ϵ_m increased by a constant factor), so we don't need to do anything special here for \mathbb{C} vs. \mathbb{R} .
 - (i) Backward stable. $x \oplus x = \text{fl}(x) \oplus \text{fl}(x) = [x(1+\epsilon_1) + x(1+\epsilon_1)](1+\epsilon_2) = 2\tilde{x}$ for $|\epsilon_i| \le \epsilon_m$ and $\tilde{x} = x(1+\epsilon_1+\epsilon_2+2\epsilon_1\epsilon_2) = x[1+O(\epsilon_m)]$.
 - (ii) Backward stable. $x \otimes x = \text{fl}(x) \otimes \text{fl}(x) = [x(1+\epsilon_1) \times x(1+\epsilon_1)](1+\epsilon_2) = \tilde{x}^2 \text{ for } |\epsilon_i| \le \epsilon_m$ and $\tilde{x} = x(1+\epsilon_1)\sqrt{1+\epsilon_2} = x[1+O(\epsilon_m)].$
 - (iii) Stable but not backwards stable. $x \oslash x = [f(x)/f(x)](1+\epsilon) = 1+\epsilon$ (not including x=0 or ∞ , which give NaN). This is actually forwards stable, but there is no \tilde{x} such that $\tilde{x}/\tilde{x} \neq 1$ so it is not backwards stable. (Under the stronger assumption of correctly rounded arithmetic, this will give exactly 1, however.)
 - (iv) Backwards stable. $x \ominus x = [f(x) f(x)](1 + \epsilon) = 0$. This is the correct answer for $\tilde{x} = x$. (In the crazy case where fl is not deterministic, then it might give a nonzero answer, in which case it is unstable.)
 - (v) Unstable. It is definitely not backwards stable, because there is no data (and hence no way to choose \tilde{x} to match the output). To be stable, it would have to be forwards stable, but it isn't because the errors decrease more slowly than $O(\epsilon_m)$. More explicitly, $1 \oplus \frac{1}{2} \oplus \frac{1}{6} \oplus \cdots$ summed from left to right will give $((1 + \frac{1}{2})(1 + \epsilon_1) + \frac{1}{6})(1 + \epsilon_2) \cdots = e + \frac{3}{2}\epsilon_1 + \frac{10}{6}\epsilon_2 + \cdots$ dropping terms of $O(\epsilon^2)$, where the coefficients of the ϵ_k factors converge to e. The number of terms is n where n satisfies $n! \approx 1/\epsilon_m$, which is a function that grows very slowly with $1/\epsilon_m$, and hence the error from the additions alone is bounded above by $\approx n\epsilon_m$. The key point is that the errors grow at least as fast as $n\epsilon_m$ (not even counting errors from truncation of the series, approximation of 1/k!, etcetera), which is not $O(\epsilon_m)$ because n grows slowly with decreasing ϵ_m .
 - (vi) Stable. As in (e), it is not backwards stable, so the only thing is to check forwards stability. Again, there will be n terms in the series, where n is a slowly growing function of $1/\epsilon_m$ ($n! \approx 1/\epsilon_m$). However, the summation errors no longer grow as n. From right to left, we are summing $\frac{1}{n!} \oplus \frac{1}{(n-1)!} \oplus \cdots \oplus 1$. But this gives $((\frac{1}{n!} + \frac{1}{(n-1)!})(1 + \epsilon_{n-1}) + \frac{1}{(n-2)!})(1 + \epsilon_{n-2}) \cdots$, and the linear terms in the ϵ_k are then bounded by

$$\left| \sum_{k=1}^{n-1} \epsilon_k \sum_{j=k}^n \frac{1}{j!} \right| \le \epsilon_m \sum_{k=1}^{n-1} \sum_{j=k}^n \frac{1}{j!} = \epsilon_m \left[\frac{n-1}{n!} + \sum_{j=1}^{n-1} \frac{j}{j!} \right] \approx \epsilon_m e = O(\epsilon_m).$$

The key point is that the coefficients of the ϵ_k coefficients grow smaller and smaller with k, rather than approaching e as for left-to-right summation, and the sum of the coefficients converges. The truncation error is of $O(\epsilon_m)$, and we assume 1/k! can also be calculated to within $O(\epsilon_m)$, e.g. via Stirling's approximation for large k, so the overall error is $O(\epsilon_m)$ and the algorithm is forwards stable.

- (vii) Forwards stable. Not backwards stable since no data, but what about forwards stability? Supposing $\sin(x)$ is computed in a stable manner, then $\widehat{\sin}(x) = \sin(x + \delta) \cdot [1 + O(\epsilon_m)]$ for $|\delta| = |x|O(\epsilon_m)$. It follows that, in the vicinity of $x = \pi$, the $\widehat{\sin}$ function can only change sign within $|\delta| = \pi O(\epsilon_m)$ of $x = \pi$. Hence, checking for $\widehat{\sin}(x) \otimes \widehat{\sin}(x') \leq 0$, where x' is the floating-point successor to x (nextfloat(x) in Julia) yields $\pi[1 + O(\epsilon_m)]$, a forwards-stable result.
- (b) Trefethen, exercise 16.1. Note that we are free to switch norms as needed, by norm equivalence. Notation: the floating-point algorithm for computing f(A) = QA will be denoted $\tilde{f}(A) = \overline{QA}$; I will assume that we simply use the obvious three-loop algorithm, i.e. computing the row-column dot products with in-order ("recursive") summation, allowing us to re-use the summation error analysis from pset 1.
 - (i) We will proceed by induction on k: first, we will prove the base case, that multiplying A by a *single* Q is backwards stable, and then we will do the inductive step (assume it is true for k, prove it for k+1).

First, the base case: we need to find a δA with $\|\delta A\| = \|A\| O(\epsilon_{\text{machine}})$ such that $\widetilde{QA} = Q(A + \delta A)$. Since $\|\delta A\| = \|Q^*\widetilde{QA} - A\| = \|Q(Q^*\widetilde{QA} - A)\| = \|\widetilde{QA} - QA\|$ in the L_2 norm, however, this is equivalent to showing $\|\widetilde{QA} - QA\| = \|A\| O(\epsilon_{\text{machine}})$; that is, we can look at the forwards error, which is a bit easier. It is sufficient to look at the error in the ij-th element of QA, i.e. the error in computing $\sum_k q_{ik}a_{kj}$. Assuming we do this sum by a straightforward loop, the analysis is exactly the same as in problem 2, except that there is an additional $(1+\epsilon)$ factor in each term for the error in the product $q_{ik}a_{kj}$ [or $(1+2\epsilon)$ if we include the rounding of q_{ik} to $\tilde{q}_{ik}=\mathrm{fl}(q_{ik})$]. Hence, the error in the ij-th element is bounded by $mO(\epsilon_{\mathrm{machine}})\sum_k |q_{ik}a_{kj}|$, and (using the unitarity of Q, which implies that $|q_{ik}| \leq 1$) this in turn is bounded by $mO(\epsilon_{\mathrm{machine}})\sum_k |a_{kj}| \leq mO(\epsilon_{\mathrm{machine}})\|A\|$ (since $\sum_{kj} |a_{kj}|$ is just an L_1 Frobenius norm of A, which is within a constant factor of any other norm). Summing m^2 of these errors in the individual elements of QA, again using norm equivalence, we obtain $\|\widetilde{QA} - QA\| = O(\sum_{ij} |(\widetilde{QA} - QA)_{ij}|) = m^3 O(\epsilon_{\mathrm{machine}})\|A\|$. Thus, we have proved backwards stability for multiplying by one unitary matrix (with a overly pessimistic m^3 coefficient, but that doesn't matter here).

Now, we will show by induction that multiplying by k unitary matrices is backwards stable. Suppose we have proved it for k, and want to prove for k+1. That, consider $QQ_k\cdots Q_1A$. By assumption, $Q_k\cdots Q_1A$ is backwards stable, and hence $\tilde{B}=Q_k\cdots Q_1A=Q_k\cdots Q_1(A+\delta A_k)$ for some $\|\delta A_k\|=O(\epsilon_{\text{machine}})\|A\|$. Also, from above, $Q\tilde{B}=Q(\tilde{B}+\delta \tilde{B})$ for some $\|\delta \tilde{B}\|=O(\epsilon_{\text{machine}})\|\tilde{B}\|=\|Q_k\cdots Q_1(A+\delta A_k)\|O(\epsilon_{\text{machine}})=\|A+\delta A_k\|O(\epsilon_{\text{machine}})\leq \|A\|O(\epsilon_{\text{machine}})+\|\delta A_k\|O(\epsilon_{\text{machine}})=\|A\|O(\epsilon_{\text{machine}})$. Hence, $QQ_k\cdots Q_1A=Q\tilde{B}=Q[Q_k\cdots Q_1(A+\delta A_k)+\delta \tilde{B}]=QQ_k\cdots Q_1(A+\delta A)$ where $\delta A=\delta A_k+[Q_1^*\cdots Q_k^*]\delta \tilde{B}$ and $\|\delta A\|\leq \|\delta A_k\|+\|\delta \tilde{B}\|=O(\epsilon_{\text{machine}})\|A\|$. Q.E.D.

(ii) Consider XA, where X is some rank-1 matrix xy^* and A has rank > 1. The product XA has rank 1 in exact arithmetic, but after floating-point errors it is unlikely that \widetilde{XA} will be exactly rank 1. Hence it is not backwards stable, because $X\tilde{A}$ will be rank 1 regardless of \tilde{A} , and thus is $\neq \widetilde{XA}$. (See also example 15.2 in the text.)

Problem 2: (10+10 points)

(a) Denote the rows of A by a_1^T, \ldots, a_m^T . Consider the unit ball in the L_{∞} norm, the set $\{x \in \mathbb{C}^n : \|x\|_{\infty} \leq 1\}$. Any vector Ax in the image of this set satisfies:

$$||Ax||_{\infty} = \max_{j \in 1:m} |a_j^T x| = \max_{j \in 1:m} \left| \sum_{k \in 1:n} a_{j,k} x_k \right| \le \max_{j \in 1:m} \sum_{k \in 1:n} |a_{j,k}| = \max_{j \in 1:m} ||a_j||,$$

since $|x_k| \le 1$ in the L_{∞} unit ball. Furthermore, this bound is achieved when $x_k = \text{sign}(a_{j,k})$ where $j = \text{argmax}_j ||a_j||$. Hence $||A||_{\infty} = \text{max}_j ||a_j||$, corresponding to (3.10). Q.E.D.

If we look in the Julia source code, we find that this norm is computed by summing the absolute values of each row of A and then takes the maximum, exactly as in (3.10).

(b) To obtain $\mu \times \nu$ submatrix B of the $m \times n$ matrix A by selecting a subset of the rows and columns of A, we simply multiply A on the left and right by $\mu \times m$ and $n \times \nu$ matrices as follows:

$$B = \begin{pmatrix} 1 & & & \\ & & 1 & \\ & & \ddots & \end{pmatrix} A \begin{pmatrix} 1 & & \\ & 1 & \\ & & \ddots & \end{pmatrix}$$

where there are 1's in the columns/rows to be selected. More precisely, if we want a subset \mathcal{R} of the rows of A and a subset \mathcal{C} of the columns of A, then we compute $B = D_{\mathcal{R}}AD_{\mathcal{C}}^T$, where the "deletion matrix" for an ordered set \mathcal{S} of indices is given by $(D_{\mathcal{S}})_{ij} = 1$ if j equals the i-th element of \mathcal{S} and $(D_{\mathcal{S}})_{ij} = 0$ otherwise; $D_{\mathcal{R}}$ is $\mu \times m$ and $D_{\mathcal{C}}$ is $\nu \times n$.

From Trefethen, chapter 3, we have $||B||_p \leq ||A||_p ||D_{\mathcal{R}}||_p ||D_{\mathcal{C}}||_p$. So, we merely need to show $||D_{\mathcal{S}}||_p \leq 1$ and the result follows. But this is trivial: $||D_{\mathcal{S}}x||_p = \left[\sum_{i \in \mathcal{S}} |x_i|^p\right]^{1/p} \leq \left[\sum_i |x_i|^p\right]^{1/p} = ||x||_p$, so $||D_{\mathcal{S}}||_p \leq 1$ and we obtain $||B||_p \leq ||A||_p$.

In Julia, we construct a random 10×7 A by A=randn(10,7), and an arbitrary 3×4 subset of this matrix by B = A[[1,3,4],[2,3,5,6]]. Then norm(B) <= norm(A) (the p=2 norm) returns true. As a more careful test, we can also try computing thousands of such random matrices and check that the maximum of norm(B)/norm(A) is < 1; a one-liner to do this in Julia is maximum(Float64[let A=randn(10,7); norm(A[1:3,1:4])/norm(A); end for i=1:10000]), which returns roughly 0.92. However, a quick check with a single matrix is acceptable here—such numerical "spot checks" are extremely useful to catch gross errors, but of course they aren't a substitute for proof, only a supplement (or sometimes a suggestive guide, if the numerical results precede the proof).

Problem 3: (21 points)

To get the condition number of g(A) = Ax, we first need to to get the Jacobian. The most natural definition would have entries $\frac{\partial g_i}{\partial A_{jk}}$, but this would result in a 3-index object (a "3d matrix" or "rank 3 tensor"), which would require us to involve techniques from multilinear algebra. Instead, we can use ordinary linear algebra with an ordinary Jacobian matrix if we treat the input A as a "1d" vector a of length mn:

$$a = \begin{pmatrix} (A_{1,:})^T \\ (A_{2,:})^T \\ \vdots \\ (A_{m,:})^T \end{pmatrix},$$

i.e. a consists of the rows of A (transposed to column vectors), one after the other, in sequence (i.e., row-major storage of A). (We could also stack the columns; this isomorphism is called the "vectorization" of A.) There are m outputs f_i of Ax, each one of which dots one row of A with x. Hence, in terms of a, the $m \times (mn)$ Jacobian matrix looks like

$$J = \left(\begin{array}{ccc} x^T & & & \\ & x^T & & \\ & & \ddots & \\ & & & x^T \end{array}\right).$$

Since this is block-diagonal, it is easy to figure out $\sup_{z\neq 0} \frac{\|Jz\|}{\|z\|}$, as long as we pick a norm of the inputs A that allows us to easily compute the corresponding norm of the "vector" z, and so it is convenient to choose the Frobenius norm so that $\|A\|_F = \|z\|_2$. To maximize $\frac{\|Jz\|}{\|z\|}$, let's write z as

$$z = \left(\begin{array}{c} x_1 \\ x_2 \\ \vdots \\ x_m \end{array}\right)$$

in terms of vectors $x_k \in \mathbb{C}^n$. Then

$$||J||_2 = \frac{||Jz||_2}{||z||_2} = \sqrt{\frac{|x^Tx_1|^2 + |x^Tx_2|^2 + \dots + |x^Tx_m|^2}{x_1^*x_1 + x_2^*x_2 + \dots + x_m^*x_m}},$$

which is clearly maximized when $x_k = \alpha_k \bar{x}$ (to maximize the dot products $x^T x_k \to |\alpha_k|^2 |x||_2^2$ over all vectors x_k of a given length) for some scalar $\alpha_k \in \mathbb{C}$, giving $||J||_2 = ||x||_2 \sqrt{\sum |\alpha_k|^2} = ||x||_2$. Hence, the condition number is $\kappa(A) = \frac{||J||}{||Ax||/||A||} = \frac{||x||_2 ||A||_F}{||Ax||_2}$, which is almost exactly the same the condition number for f(x) = Ax, except that we substitute $||A||_F$ for $||A||_2$. Due to the equivalence of norms, however, this means that the condition numbers differ only by at most a constant factor independent of A or x.

(There are other possible choices of norm where this computation is reasonable to carry out, but in any case you should get something similar up to a constant factor.)

Problem 4: (10+10+10 points)

(a) Trefethen, probem 4.5. It is sufficient to show that the reduced SVD $A\hat{V} = \hat{U}\hat{\Sigma}$ is real, since the remaining columns of U and V are formed as a basis for the orthogonal complement of the columns of \hat{U} and \hat{V} , and if the latter are real then their complement is obviously also real. Furthermore, it is sufficient to show that \hat{U} can be chosen real, since (from class) $A^*u_i/\sigma_i = v_i$ for each column u_i of \hat{U} and v_i of \hat{U} , and A^* is real. The columns u_i are eigenvectors of $A^*A = B$, which is a real-symmetric matrix, i.e. $Bu_i = \sigma_i^2 u_i$. Suppose that the u_i are not real. Then the real and imaginary parts of u_i are themselves eigenvectors (if they are nonzero) with eigenvalue σ_i^2 (proof: take the real and imaginary parts of $Bu_i = \sigma_i^2 u_i$, since B and σ_i^2 are real). Hence, taking either the real or imaginary parts of the complex u_i (whichever is nonzero) and normalizing them to unit length, we obtain a new purely real \hat{U} . Q.E.D.¹

¹There is a slight wrinkle if there are repeated eigenvalues, e.g. $\sigma_1 = \sigma_2$, because the real or imaginary parts of u_1 and u_2 might not be orthogonal. However, taken together, the real and imaginary parts of any multiple eigenvalues must span the same space, and hence we can find a real orthonormal basis with Gram-Schmidt or whatever.

- (b) Trefethen, problem 5.2. We just need to show that, for any $A \in \mathbb{C}^{m \times n}$ with rank < n and for any $\epsilon > 0$, we can find sequence of full-rank matrices B that eventually satisfies $||A B||_2 < \epsilon$. Form the SVD $A = U\Sigma V^*$ with singular values $\sigma_1, \ldots, \sigma_r$ where r < n is the rank of A. Let $B = U\widetilde{\Sigma}V^*$ where $\widetilde{\Sigma}$ is the same as Σ except that it has n r additional nonzero singular values $\sigma_{k>r} = \epsilon/2$. From equation 5.4 in the book, $||B A||_2 = \sigma_{r+1} = \epsilon/2 < \epsilon$, noting that $A = B_r$ in the notation of the book. We can then make a sequence of such matrices e.g. by letting $\epsilon = \sigma_r 2^{-k}$ for $k = 1, 2, \ldots$
- (c) Trefethen, problem 5.4. From $A = U\Sigma V^*$, recall that $AV = U\Sigma$ and $A^*U = V\Sigma$. Therefore,

$$\left(\begin{array}{cc} A^* \\ A \end{array}\right) \left(\begin{array}{c} V \\ \pm U \end{array}\right) = \left(\begin{array}{c} \pm A^* U \\ A V \end{array}\right) = \pm \left(\begin{array}{c} V \Sigma \\ \pm U \Sigma \end{array}\right) = \pm \left(\begin{array}{c} V \\ \pm U \end{array}\right) \Sigma$$

and hence $(v_i; \pm u_i)$ is an eigenvector of $\begin{pmatrix} A^* \\ A \end{pmatrix}$ with eigenvalue $\pm \sigma_i$. Noting that these vectors $(v_i; \pm u_i)$ are orthogonal by construction and only need to be divided by $\sqrt{2}$ to be normalized, we immediately obtain the diagonalization

$$\left(\begin{array}{c} A^* \\ A \end{array}\right) = Q \left(\begin{array}{cc} +\Sigma \\ -\Sigma \end{array}\right) Q^*$$

for

$$Q = \left(\begin{array}{cc} V & V \\ +U & -U \end{array} \right) / \sqrt{2}.$$